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CHAPTER 1

Introduction

I can illustrate the .... approach with the ... image of a nut to be opened. The first analogy that came to my mind is of immersing the nut in some softening liquid, and why not simply water? From time to time you rub so the liquid penetrates better, and otherwise you let time pass. The shell becomes more flexible through weeks and months — when the time is ripe, hand pressure is enough, the shell opens like a perfectly ripened avocado!

A different image came to me a few weeks ago. The unknown thing to be known appeared to me as some stretch of earth or hard marl, resisting penetration ... the sea advances insensibly in silence, nothing seems to happen, nothing moves, the water is so far off you hardly hear it ... yet finally it surrounds the resistant substance.

— Alexandre Grothendieck, Récoltes et Semailles p. 552-3, translation by Colin McLarty

1.1 Goals

These are an updated version of notes accompanying a hard year-long class taught at Stanford in 2009-2010. I am currently editing them and adding a few more sections, and I hope a reasonably complete (if somewhat rough) version over the 2010-11 academic year at the site http://math216.wordpress.com/.

In any class, choices must be made as to what the course is about, and who it is for — there is a finite amount of time, and any addition of material or explanation or philosophy requires a corresponding subtraction. So these notes are highly inappropriate for most people and most classes. Here are my goals. (I do not claim that these goals are achieved; but they motivate the choices made.)

These notes currently have a very particular audience in mind: Stanford Ph.D. students, postdocs and faculty in a variety of fields, who may want to use algebraic geometry in a sophisticated way. This includes algebraic and arithmetic geometers, but also topologists, number theorists, symplectic geometers, and others.

The notes deal purely with the algebraic side of the subject, and completely neglect analytic aspects.

They assume little prior background (see §1.2), and indeed most students have little prior background. Readers with less background will necessarily have to work harder. It would be great if the reader had seen varieties before, but many students haven’t, and the course does not assume it — and similarly for category theory, homological algebra, more advanced commutative algebra, differential geometry, .... Surprisingly often, what we need can be developed quickly from scratch. The cost is that the course is much denser; the benefit is that more people can follow it; they don’t reach a point where they get thrown. (On the other hand,
people who already have some familiarity with algebraic geometry, but want to understand the foundations more completely should not be bored, and will focus on more subtle issues.)

The notes seek to cover everything that one should see in a first course in the subject, including theorems, proofs, and examples.

They seek to be complete, and not leave important results as black boxes pulled from other references.

There are lots of exercises. I have found that unless I have some problems I can think through, ideas don’t get fixed in my mind. Some are trivial — that’s okay, and even desirable. As few necessary ones as possible should be hard, but the reader should have the background to deal with them — they are not just an excuse to push material out of the text.

There are optional starred (⋆) sections of topics worth knowing on a second or third (but not first) reading. You should not read double-starred sections (★★) unless you really really want to, but you should be aware of their existence.

The notes are intended to be readable, although certainly not easy reading.

In short, after a year of hard work, students should have a broad familiarity with the foundations of the subject, and be ready to attend seminars, and learn more advanced material. They should not just have a vague intuitive understanding of the ideas of the subject; they should know interesting examples, know why they are interesting, and be able to prove interesting facts about them.

I have greatly enjoyed thinking through these notes, and teaching the corresponding classes, in a way I did not expect. I have had the chance to think through the structure of algebraic geometry from scratch, not blindly accepting the choices made by others. (Why do we need this notion? Aha, this forces us to consider this other notion earlier, and now I see why this third notion is so relevant...) I have repeatedly realized that ideas developed in Paris in the 1960’s are simpler than I initially believed, once they are suitably digested.

1.1.1. Implications. We will work with as much generality as we need for most readers, and no more. In particular, we try to have hypotheses that are as general as possible without making proofs harder. The right hypotheses can make a proof easier, not harder, because one can remember how they get used. As an inflammatory example, the notion of quasiseparated comes up early and often. The cost is that one extra word has to be remembered, on top of an overwhelming number of other words. But once that is done, it is not hard to remember that essentially every scheme anyone cares about is quasiseparated. Furthermore, whenever the hypotheses “quasicompact and quasiseparated” turn up, the reader will likely immediately see a key idea of the proof.

Similarly, there is no need to work over an algebraically closed field, or even a field. Geometers needn’t be afraid of arithmetic examples or of algebraic examples; a central insight of algebraic geometry is that the same formalism applies without change.

1.1.2. Costs. Choosing these priorities requires that others be shortchanged, and it is best to be up front about these. Because of our goal is to be comprehensive, and to understand everything one should know after a first course, it will necessarily take longer to get to interesting sample applications. You may be misled
into thinking that one has to work this hard to get to these applications — it is not true!

1.2 Background and conventions

All rings are assumed to be commutative unless explicitly stated otherwise. All rings are assumed to contain a unit, denoted 1. Maps of rings must send 1 to 1. We don’t require that 0 ≠ 1; in other words, the “0-ring” (with one element) is a ring. (There is a ring map from any ring to the 0-ring; the 0-ring only maps to itself. The 0-ring is the final object in the category of rings.) The definition of “integral domain” includes 1 ≠ 0, so the 0-ring is not an integral domain. We accept the axiom of choice. In particular, any proper ideal in a ring is contained in a maximal ideal. (The axiom of choice also arises in the argument that the category of A-modules has enough injectives, see Exercise [23.2.E].)

The reader should be familiar with some basic notions in commutative ring theory, in particular the notion of ideals (including prime and maximal ideals) and localization. For example, the reader should be able to show that if S is a multiplicative set of a ring A (which we assume to contain 1), then the primes of $S^{-1}A$ are in natural bijection with those primes of A not meeting S ([4.2.6]). Tensor products and exact sequences of A-modules will be important. We will use the notation $(A, m)$ or $(A, m, k)$ for local rings — A is the ring, m its maximal ideal, and $k = A/m$ its residue field. We will use (in Proposition [13.7.2]) the structure theorem for finitely generated modules over a principal ideal domain A: any such module can be written as the direct sum of principal modules $A/(a)$.

1.2.1. Caution about foundational issues. We will not concern ourselves with subtle foundational issues (set-theoretic issues involving universes, etc.). It is true that some people should be careful about these issues. (If you are one of these rare people, a good start is [KS §1.1].)

1.2.2. Further background. It may be helpful to have books on other subjects handy that you can dip into for specific facts, rather than reading them in advance. In commutative algebra, Eisenbud [E] is good for this. Other popular choices are Atiyah-Macdonald [AM] and Matsumura [M-CRT]. For homological algebra, Weibel [W] is simultaneously detailed and readable.

Background from other parts of mathematics (topology, geometry, complex analysis) will of course be helpful for developing intuition.

Finally, it may help to keep the following quote in mind.

[Algebraic geometry] seems to have acquired the reputation of being esoteric, exclusive, and very abstract, with adherents who are secretly plotting to take over all the rest of mathematics! In one respect this last point is accurate ...

— David Mumford, 1975 [M-Red2, p. 227]
Part I

Preliminaries
CHAPTER 2

Some category theory

That which does not kill me, makes me stronger. — Nietzsche

2.1 Motivation

Before we get to any interesting geometry, we need to develop a language to discuss things cleanly and effectively. This is best done in the language of categories. There is not much to know about categories to get started; it is just a very useful language. Like all mathematical languages, category theory comes with an embedded logic, which allows us to abstract intuitions in settings we know well to far more general situations.

Our motivation is as follows. We will be creating some new mathematical objects (such as schemes, and certain kinds of sheaves), and we expect them to act like objects we have seen before. We could try to nail down precisely what we mean by “act like”, and what minimal set of things we have to check in order to verify that they act the way we expect. Fortunately, we don’t have to — other people have done this before us, by defining key notions, such as abelian categories, which behave like modules over a ring.

Our general approach will be as follows. I will try to tell what you need to know, and no more. (This I promise: if I use the word “topoi”, you can shoot me.) I will begin by telling you things you already know, and describing what is essential about the examples, in a way that we can abstract a more general definition. We will then see this definition in less familiar settings, and get comfortable with using it to solve problems and prove theorems.

For example, we will define the notion of product of schemes. We could just give a definition of product, but then you should want to know why this precise definition deserves the name of “product”. As a motivation, we revisit the notion of product in a situation we know well: (the category of) sets. One way to define the product of sets $U$ and $V$ is as the set of ordered pairs $\{(u, v) : u \in U, v \in V\}$. But someone from a different mathematical culture might reasonably define it as the set of symbols $\{uv : u \in U, v \in V\}$. These notions are “obviously the same”. Better: there is “an obvious bijection between the two”.

This can be made precise by giving a better definition of product, in terms of a universal property. Given two sets $M$ and $N$, a product is a set $P$, along with maps $\mu : P \to M$ and $\nu : P \to N$, such that for any set $P'$ with maps $\mu' : P' \to M$ and $\nu' : P' \to N$, there is a unique map $\delta : P' \to P$ such that $\mu \circ \delta = \mu'$ and $\nu \circ \delta = \nu'$.
\( \nu' : P' \to N \), these maps must factor \textit{uniquely} through \( P \):

\[(2.1.0.1)\]

\[
\begin{array}{c}
P' \\
\mu' \\
\downarrow \\
\mu \\
\downarrow \\
M \\
\nu \\
\downarrow \\
P \\
\mu \\
\downarrow \\
N \\
\end{array}
\]

(The symbol \( \exists \) means “there exists”, and the symbol \( ! \) here means “unique”.) Thus a \textbf{product} is a diagram

\[
\begin{array}{c}
P \\
\mu \\
\downarrow \\
M \\
\nu \\
\downarrow \\
N \\
\end{array}
\]

and not just a set \( P \), although the maps \( \mu \) and \( \nu \) are often left implicit.

This definition agrees with the traditional definition, with one twist: there isn’t just a single product; but any two products come with a \textit{unique} isomorphism between them. In other words, the product is unique up to \textit{unique} isomorphism. Here is why: if you have a product

\[
\begin{array}{c}
P_1 \\
\mu_1 \\
\downarrow \\
M \\
\nu_1 \\
\downarrow \\
N \\
\end{array}
\]

and I have a product

\[
\begin{array}{c}
P_2 \\
\mu_2 \\
\downarrow \\
M \\
\nu_2 \\
\downarrow \\
N \\
\end{array}
\]

then by the universal property of my product (letting \( (P_2, \mu_2, \nu_2) \) play the role of \( (P, \mu, \nu) \), and \( (P_1, \mu_1, \nu_1) \) play the role of \( (P', \mu', \nu') \) in \((2.1.0.1)\)), there is a unique map \( f : P_1 \to P_2 \) making the appropriate diagram commute (i.e. \( \mu_1 = \mu_2 \circ f \) and \( \nu_1 = \nu_2 \circ f \)). Similarly by the universal property of your product, there is a unique map \( g : P_2 \to P_1 \) making the appropriate diagram commute. Now consider the universal property of my product, this time letting \( (P_2, \mu_2, \nu_2) \) play the role of both \( (P, \mu, \nu) \) and \( (P', \mu', \nu') \) in \((2.1.0.1)\). There is a unique map \( h : P_2 \to P_2 \) such that

\[
\begin{array}{c}
P_2 \\
\mu_2 \\
\downarrow \\
M \\
\nu_2 \\
\downarrow \\
N \\
\end{array}
\]

commutes. However, I can name two such maps: the identity map \( \text{id}_{P_2} \), and \( g \circ f \). Thus \( g \circ f = \text{id}_{P_2} \). Similarly, \( f \circ g = \text{id}_{P_1} \). Thus the maps \( f \) and \( g \) arising from
the universal property are bijections. In short, there is a unique bijection between $P_1$ and $P_2$ preserving the “product structure” (the maps to $M$ and $N$). This gives us the right to name any such product $M \times N$, since any two such products are uniquely identified.

This definition has the advantage that it works in many circumstances, and once we define categories, we will soon see that the above argument applies verbatim in any category to show that products, if they exist, are unique up to unique isomorphism. Even if you haven’t seen the definition of category before, you can verify that this agrees with your notion of product in some category that you have seen before (such as the category of vector spaces, where the maps are taken to be linear maps; or the category of smooth manifolds, where the maps are taken to be smooth maps (submersions)).

This is handy even in cases that you understand. For example, one way of defining the product of two manifolds $M$ and $N$ is to cut them both up into charts, then take products of charts, then glue them together. But if I cut up the manifolds in one way, and you cut them up in another, how do we know our resulting manifolds are the “same”? We could wave our hands, or make an annoying argument about refining covers, but instead, we should just show that they are “categorical products” and hence canonically the “same” (i.e. isomorphic). We will formalize this argument in §2.3.

Another set of notions we will abstract are categories that “behave like modules”. We will want to define kernels and cokernels for new notions, and we should make sure that these notions behave the way we expect them to. This leads us to the definition of abelian categories, first defined by Grothendieck in his Tōhoku paper [Gr].

In this chapter, we will give an informal introduction to these and related notions, in the hope of giving just enough familiarity to comfortably use them in practice.

## 2.2 Categories and functors

We begin with an informal definition of categories and functors.

### 2.2.1. Categories.

A category consists of a collection of objects, and for each pair of objects, a set of maps, or morphisms (or arrows), between them. The collection of objects of a category $C$ are often denoted obj$(C)$, but we will usually denote the collection also by $C$. If $A, B \in C$, then the morphisms from $A$ to $B$ are denoted Mor$(A, B)$. A morphism is often written $f : A \to B$, and $A$ is said to be the source of $f$, and $B$ the target of $f$. (Of course, Mor$(A, B)$ is taken to be disjoint from Mor$(A', B')$ unless $A = A'$ and $B = B'$.)

Morphisms compose as expected: there is a composition Mor$(A, B) \times$ Mor$(B, C) \to$ Mor$(A, C)$, and if $f \in$ Mor$(A, B)$ and $g \in$ Mor$(B, C)$, then their composition is denoted $g \circ f$. Composition is associative: if $f \in$ Mor$(A, B)$, $g \in$ Mor$(B, C)$, and $h \in$ Mor$(C, D)$, then $h \circ (g \circ f) = (h \circ g) \circ f$. For each object $A \in C$, there is always an identity morphism id$_A : A \to A$, such that when you (left- or right-)compose a morphism with the identity, you get the same morphism. More precisely, if
f : A → B is a morphism, then f ◦ id_A = f = id_B ◦ f. (If you wish, you may check that “identity morphisms are unique”: there is only one morphism deserving the name id_A.)

If we have a category, then we have a notion of isomorphism between two objects (a morphism f : A → B such that there exists some — necessarily unique — morphism g : B → A, where f ◦ g and g ◦ f are the identity on B and A respectively), and a notion of automorphism of an object (an isomorphism of the object with itself).

2.2.2. Example. The prototypical example to keep in mind is the category of sets, denoted Sets. The objects are sets, and the morphisms are maps of sets. (Because Russell’s paradox shows that there is no set of all sets, we did not say earlier that there is a set of all objects. But as stated in §1.2 we are deliberately omitting all set-theoretic issues.)

2.2.3. Example. Another good example is the category Vec_k of vector spaces over a given field k. The objects are k-vector spaces, and the morphisms are linear transformations. (What are the isomorphisms?)

2.2.A. UNIMPORTANT EXERCISE. A category in which each morphism is an isomorphism is called a groupoid. (This notion is not important in these notes. The point of this exercise is to give you some practice with categories, by relating them to an object you know well.)

(a) A perverse definition of a group is: a groupoid with one object. Make sense of this.

(b) Describe a groupoid that is not a group.

2.2.B. EXERCISE. If A is an object in a category C, show that the invertible elements of Mor(A, A) form a group (called the automorphism group of A, denoted Aut(A)). What are the automorphism groups of the objects in Examples 2.2.2 and 2.2.3? Show that two isomorphic objects have isomorphic automorphism groups. (For readers with a topological background: if X is a topological space, then the fundamental groupoid is the category where the objects are points of X, and the morphisms x → y are paths from x to y, up to homotopy. Then the automorphism group of x_0 is the (pointed) fundamental group π_1(X, x_0). In the case where X is connected, and π_1(X) is not abelian, this illustrates the fact that for a connected groupoid — whose definition you can guess — the automorphism groups of the objects are all isomorphic, but not canonically isomorphic.)

2.2.4. Example: abelian groups. The abelian groups, along with group homomorphisms, form a category Ab.

2.2.5. Important example: modules over a ring. If A is a ring, then the A-modules form a category Mod_A. (This category has additional structure; it will be the prototypical example of an abelian category, see §2.6) Taking A = k, we obtain Example 2.2.3; taking A = Z, we obtain Example 2.2.4.

2.2.6. Example: rings. There is a category Rings, where the objects are rings, and the morphisms are morphisms of rings (which send 1 to 1 by our conventions, §1.2).
2.2.7. Example: topological spaces. The topological spaces, along with continuous maps, form a category $\text{Top}$. The isomorphisms are homeomorphisms.

In all of the above examples, the objects of the categories were in obvious ways sets with additional structure. This needn’t be the case, as the next example shows.

2.2.8. Example: partially ordered sets. A partially ordered set, or poset, is a set $S$ along with a binary relation $\geq$ on $S$ satisfying:

(i) $x \geq x$ (reflexivity),
(ii) $x \geq y$ and $y \geq z$ imply $x \geq z$ (transitivity), and
(iii) if $x \geq y$ and $y \geq x$ then $x = y$.

A partially ordered set $(S, \geq)$ can be interpreted as a category whose objects are the elements of $S$, and with a single morphism from $x$ to $y$ if and only if $x \geq y$ (and no morphism otherwise).

A trivial example is $(S, \geq)$ where $x \geq y$ if and only if $x = y$. Another example is

(2.2.8.1)

Here there are three objects. The identity morphisms are omitted for convenience, and the two non-identity morphisms are depicted. A third example is

(2.2.8.2)

Here the “obvious” morphisms are again omitted: the identity morphisms, and the morphism from the upper left to the lower right. Similarly,

\[
\cdots \longrightarrow \bullet \longrightarrow \bullet
\]

depicts a partially ordered set, where again, only the “generating morphisms” are depicted.

2.2.9. Example: the category of subsets of a set, and the category of open sets in a topological space. If $X$ is a set, then the subsets form a partially ordered set, where the order is given by inclusion. Similarly, if $X$ is a topological space, then the open sets form a partially ordered set, where the order is given by inclusion.

2.2.10. Example. A subcategory $\mathcal{A}$ of a category $\mathcal{B}$ has as its objects some of the objects of $\mathcal{B}$, and some of the morphisms, such that the morphisms of $\mathcal{A}$ include the identity morphisms of the objects of $\mathcal{A}$, and are closed under composition. (For example, (2.2.8.1) is in an obvious way a subcategory of (2.2.8.2).)

2.2.11. Functors.

A covariant functor $F$ from a category $\mathcal{A}$ to a category $\mathcal{B}$, denoted $F : \mathcal{A} \to \mathcal{B}$, is the following data. It is a map of objects $F : \text{obj}(\mathcal{A}) \to \text{obj}(\mathcal{B})$, and for each $A_1, A_2 \in \mathcal{A}$, and morphism $m : A_1 \to A_2$, a morphism $F(m) : F(A_1) \to F(A_2)$ in $\mathcal{B}$. We require that $F$ preserves identity morphisms (for $A \in \mathcal{A}$, $F(\text{id}_A) = \text{id}_{F(A)}$), and that
F preserves composition \((F(m_1 \circ m_2) = F(m_1) \circ F(m_2))\). (You may wish to verify that covariant functors send isomorphisms to isomorphisms.)

If \(F : \mathcal{A} \to \mathcal{B}\) and \(G : \mathcal{B} \to \mathcal{C}\) are covariant functors, then we define a functor \(G \circ F : \mathcal{A} \to \mathcal{C}\) in the obvious way. Composition of functors is associative in an evident sense.

**2.2.12. Example: a forgetful functor.** Consider the functor from the category of vector spaces (over a field \(k\)) \(\text{Vec}_k\) to \(\text{Sets}\), that associates to each vector space its underlying set. The functor sends a linear transformation to its underlying map of sets. This is an example of a **forgetful functor**, where some additional structure is forgotten. Another example of a forgetful functor is \(\text{Mod}_A \to \text{Ab}\) from \(A\)-modules to abelian groups, remembering only the abelian group structure of the \(A\)-module.

**2.2.13. Topological examples.** Examples of covariant functors include the fundamental group functor \(\pi_1\), which sends a topological space \(X\) with choice of a point \(x_0 \in X\) to a group \(\pi_1(X, x_0)\) (what are the objects and morphisms of the source category?), and the \(i\)th homology functor \(\text{Top} \to \text{Ab}\), which sends a topological space \(X\) to its \(i\)th homology group \(H_i(X, \mathbb{Z})\). The covariance corresponds to the fact that a (continuous) morphism of pointed topological spaces \(f : X \to Y\) with \(f(x_0) = y_0\) induces a map of fundamental groups \(\pi_1(X, x_0) \to \pi_1(Y, y_0)\), and similarly for homology groups.

**2.2.14. Example.** Suppose \(A\) is an object in a category \(\mathcal{C}\). Then there is a functor \(h^A : \mathcal{C} \to \text{Sets}\) sending \(B \in \mathcal{C}\) to \(\text{Mor}(A, B)\), and sending \(f : B_1 \to B_2\) to \(\text{Mor}(A, B_1) \to \text{Mor}(A, B_2)\) described by

\[
[g : A \to B_1] \mapsto [f \circ g : A \to B_1 \to B_2].
\]

This seemingly silly functor ends up surprisingly being an important concept, and will come up repeatedly for us. (Warning only for experts: this is strictly speaking a lie: why should \(\text{Mor}(A, B)\) be a set? But as stated in Caution \[1.2.1\] we will deliberately ignore these foundational issues, and we will in general pass them by without comment. Feel free to patch the problem on your time, perhaps by working in a **small category**, defined in \[2.4.1\] But don’t be distracted from our larger goal.)

**2.2.15. Full and faithful functors.** A covariant functor \(F : \mathcal{A} \to \mathcal{B}\) is **faithful** if for all \(A, A' \in \mathcal{A}\), the map \(\text{Mor}_\mathcal{A}(A, A') \to \text{Mor}_\mathcal{B}(F(A), F(A'))\) is injective, and **full** if it is surjective. A functor that is full and faithful is **fully faithful**. A subcategory \(i : \mathcal{A} \to \mathcal{B}\) is a **full subcategory** if \(i\) is full. Thus a subcategory \(\mathcal{A}'\) of \(\mathcal{A}\) is full if and only if for all \(A, B \in \text{obj}(\mathcal{A}')\), \(\text{Mor}_{\mathcal{A}'}(A, B) = \text{Mor}_{\mathcal{A}}(A, B)\).

**2.2.16. Definition.** A **contravariant functor** is defined in the same way as a covariant functor, except the arrows switch directions: in the above language, \(F(A_1 \to A_2)\) is now an arrow from \(F(A_2)\) to \(F(A_1)\). (Thus \(F(m_2 \circ m_1) = F(m_1) \circ F(m_2)\), not \(F(m_2) \circ F(m_1)\).)

It is wise to always state whether a functor is covariant or contravariant. If it is not stated, the functor is often assumed to be covariant.

(Sometimes people describe a contravariant functor \(\mathcal{C} \to \mathcal{D}\) as a covariant functor \(\mathcal{C}^{\text{opp}} \to \mathcal{D}\), where \(\mathcal{C}^{\text{opp}}\) is the same category as \(\mathcal{C}\) except that the arrows go in the opposite direction. Here \(\mathcal{C}^{\text{opp}}\) is said to be the **opposite category** to \(\mathcal{C}\).)
2.2.17. Linear algebra example. If Vec\(_k\) is the category of \(k\)-vector spaces (introduced in Example 2.2.12), then taking duals gives a contravariant functor \(\cdot ^\vee : \text{Vec}_k \to \text{Vec}_k\). Indeed, to each linear transformation \(f : V \to W\), we have a dual transformation \(f^\vee : W^\vee \to V^\vee\), and \((f \circ g)^\vee = g^\vee \circ f^\vee\).

2.2.18. Topological example (cf. Example 2.2.13) for those who have seen cohomology. The \(i\)th cohomology functor \(H^i(\cdot, \mathbb{Z}) : \text{Top} \to \text{Ab}\) is a contravariant functor.

2.2.19. Example. There is a contravariant functor \(\text{Top} \to \text{Rings}\) taking a topological space \(X\) to the real-valued continuous functions on \(X\). A morphism of topological spaces \(X \to Y\) (a continuous map) induces the pullback map from functions on \(Y\) to maps on \(X\).

2.2.20. Example (cf. Example 2.2.14). Suppose \(A\) is an object of a category \(\mathcal{C}\). Then there is a contravariant functor \(\text{h}_A : \mathcal{C} \to \text{Sets}\) sending \(B \in \mathcal{C}\) to \(\text{Mor}(B, A)\), and sending the morphism \(f : B_1 \to B_2\) to the morphism \(\text{Mor}(B_2, A) \to \text{Mor}(B_1, A)\) via

\[
[g : B_2 \to A] \mapsto [g \circ f : B_1 \to B_2 \to A].
\]

This example initially looks weird and different, but Examples 2.2.17 and 2.2.19 may be interpreted as special cases; do you see how? What is \(A\) in each case?

2.2.21. * Natural transformations (and natural isomorphisms) of functors, and equivalences of categories.

(This notion won’t come up in an essential way until at least Chapter 7, so you shouldn’t read this section until then.) Suppose \(F\) and \(G\) are two functors from \(\mathcal{A}\) to \(\mathcal{B}\). A natural transformation of functors \(F \to G\) is the data of a morphism \(m_a : F(a) \to G(a)\) for each \(a \in \mathcal{A}\) such that for each \(f : a \to a'\) in \(\mathcal{A}\), the diagram

\[
\begin{array}{ccc}
F(a) & \xrightarrow{F(f)} & F(a') \\
m_a & & m_{a'} \\
G(a) & \xrightarrow{G(f)} & G(a')
\end{array}
\]

commutes. A natural isomorphism of functors is a natural transformation such that each \(m_a\) is an isomorphism. The data of functors \(F : \mathcal{A} \to \mathcal{B}\) and \(F' : \mathcal{B} \to \mathcal{A}\) such that \(F \circ F'\) is naturally isomorphic to the identity functor \(I_{\mathcal{B}}\) on \(\mathcal{B}\) and \(F' \circ F\) is naturally isomorphic to \(I_{\mathcal{A}}\) is said to be an equivalence of categories. “Equivalence of categories” is an equivalence relation on categories. The right meaning of when two categories are “essentially the same” is not isomorphism (a functor giving bijections of objects and morphisms) but an equivalence. Exercises 2.2.C and 2.2.D might give you some vague sense of this. Later exercises (for example, that “rings” and “affine schemes” are essentially the same, once arrows are reversed, Exercise 7.3.D) may help too.

Two examples might make this strange concept more comprehensible. The double dual of a finite-dimensional vector space \(V\) is not \(V\), but we learn early to say that it is canonically isomorphic to \(V\). We can make that precise as follows. Let \(\text{f.d.Vec}_k\) be the category of finite-dimensional vector spaces over \(k\). Note that this category contains oodles of vector spaces of each dimension.
2.2.C. **Exercise.** Let $\cdot^\vee \vee : \text{f.d.} \text{Vec}_k \to \text{f.d.} \text{Vec}_k$ be the double dual functor from the category of finite-dimensional vector spaces over $k$ to itself. Show that $\cdot^\vee \vee$ is naturally isomorphic to the identity functor on $\text{f.d.} \text{Vec}_k$. (Without the finite-dimensional hypothesis, we only get a natural transformation of functors from $\text{id}$ to $\cdot^\vee \vee$.)

Let $\mathcal{V}$ be the category whose objects are $k^n$ for each $n$ (there is one vector space for each $n$), and whose morphisms are linear transformations. This latter space can be thought of as vector spaces with bases, and the morphisms are honest matrices. There is an obvious functor $\mathcal{V} \to \text{f.d.} \text{Vec}_k$, as each $k^n$ is a finite-dimensional vector space.

2.2.D. **Exercise.** Show that $\mathcal{V} \to \text{f.d.} \text{Vec}_k$ gives an equivalence of categories, by describing an “inverse” functor. (Recall that we are being cavalier about set-theoretic assumption, see Caution 1.2.1, so feel free to simultaneously choose bases for each vector space in $\text{f.d.} \text{Vec}_k$. To make this precise, you will need to use Godel-Bernays set theory or else replace $\text{f.d.} \text{Vec}_k$ with a very similar small category, but we won’t worry about this.)

2.2.22. **Aside for experts.** Your argument for Exercise 2.2.D will show that (modulo set-theoretic issues) this definition of equivalence of categories is the same as another one commonly given: a covariant functor $F : A \to B$ is an equivalence of categories if it is fully faithful and every object of $B$ is isomorphic to an object of the form $F(a)$ ($F$ is essentially surjective). One can show that such a functor has a quasi-inverse, i.e., that there is a functor $G : B \to A$, which is also an equivalence, and for which there exist natural isomorphisms $G(F(a)) \cong A$ and $F(G(b)) \cong B$. 

2.3 Universal properties determine an object up to unique isomorphism

Given some category that we come up with, we often will have ways of producing new objects from old. In good circumstances, such a definition can be made using the notion of a universal property. Informally, we wish that there were an object with some property. We first show that if it exists, then it is essentially unique, or more precisely, is unique up to unique isomorphism. Then we go about constructing an example of such an object to show existence.

Explicit constructions are sometimes easier to work with than universal properties, but with a little practice, universal properties are useful in proving things quickly and slickly. Indeed, when learning the subject, people often find explicit constructions more appealing, and use them more often in proofs, but as they become more experienced, they find universal property arguments more elegant and insightful.

2.3.1. **Products were defined by universal property.** We have seen one important example of a universal property argument already in §2.1 products. You should go back and verify that our discussion there gives a notion of product in any category, and shows that products, if they exist, are unique up to unique isomorphism.
2.3.2. Initial, final, and zero objects. Here are some simple but useful concepts that will give you practice with universal property arguments. An object of a category $\mathcal{C}$ is an initial object if it has precisely one map to every object. It is a final object if it has precisely one map from every object. It is a zero object if it is both an initial object and a final object.

2.3.A. Exercise. Show that any two initial objects are uniquely isomorphic. Show that any two final objects are uniquely isomorphic.

In other words, if an initial object exists, it is unique up to unique isomorphism, and similarly for final objects. This (partially) justifies the phrase “the initial object” rather than “an initial object”, and similarly for “the final object” and “the zero object”.

2.3.B. Exercise. What are the initial and final objects in $\text{Sets}$, $\text{Rings}$, and $\text{Top}$ (if they exist)? How about in the two examples of $\{2,2.9\}$?

2.3.3. Localization of rings and modules. Another important example of a definition by universal property is the notion of localization of a ring. We first review a constructive definition, and then reinterpret the notion in terms of universal property. A multiplicative subset $S$ of a ring $A$ is a subset closed under multiplication containing 1. We define a ring $S^{-1}A$. The elements of $S^{-1}A$ are of the form $a/s$ where $a \in A$ and $s \in S$, and where $a_1/s_1 = a_2/s_2$ if (and only if) for some $s \in S$, $s(s_2a_1 - s_1a_2) = 0$. (This implies that $S^{-1}A$ is the 0-ring if $0 \in S$.) We define $(a_1/s_1) \times (a_2/s_2) = (a_1a_2)/(s_1s_2)$, and $(a_1/s_1) + (a_2/s_2) = (s_2a_1 + s_1a_2)/(s_1s_2)$.

We have a canonical ring map $A \to S^{-1}A$ given by $a \mapsto a/1$.

There are two particularly important flavors of multiplicative subsets. The first is $\{1, f, f^2, \ldots \}$, where $f \in A$. This localization is denoted $A_f$. The second is $A - p$, where $p$ is a prime ideal. This localization $S^{-1}A$ is denoted $A_p$. (Notational warning: If $p$ is a prime ideal, then $A_p$ means you’re allowed to divide by elements not in $p$. However, if $f \in A$, $A_f$ means you’re allowed to divide by $f$. This can be confusing. For example, if $(f)$ is a prime ideal, then $A_f \neq A_{(f)}$.)

Warning: sometimes localization is first introduced in the special case where $A$ is an integral domain and $0 \notin S$. In that case, $A \cong S^{-1}A$, but this isn’t always true, as shown by the following exercise. (But we will see that noninjective localizations needn’t be pathological, and we can sometimes understand them geometrically, see Exercise 2.2.1.)

2.3.C. Exercise. Show that $A \to S^{-1}A$ is injective if and only if $S$ contains no zero-divisors. ($A$ zero-divisor of a ring $A$ is an element $a$ such that there is a non-zero element $b$ with $ab = 0$. The other elements of $A$ are called non-zero-divisors. For example, a unit is never a zero-divisor. Counter-intuitively, 0 is a zero-divisor in a ring $A$ if and only if $A$ is not the 0-ring.)

If $A$ is an integral domain and $S = A \setminus \{0\}$, then $S^{-1}A$ is called the fraction field of $A$, which we denote $K(A)$. The previous exercise shows that $A$ is a subring of its fraction field $K(A)$. We now return to the case where $A$ is a general (commutative) ring.

2.3.D. Exercise. Verify that $A \to S^{-1}A$ satisfies the following universal property: $S^{-1}A$ is initial among $A$-algebras $B$ where every element of $S$ is sent to a unit in
B. (Recall: the data of “an $A$-algebra $B$” and “a ring map $A \to B$” the same.)

Translation: any map $A \to B$ where every element of $S$ is sent to a unit must factor uniquely through $A \to S^{-1}A$.

In fact, it is cleaner to define $A \to S^{-1}A$ by the universal property, and to show that it exists, and to use the universal property to check various properties $S^{-1}A$ has. Let’s get some practice with this by defining localizations of modules by universal property. Suppose $M$ is an $A$-module. We define the $A$-module map $\phi : M \to S^{-1}M$ as being initial among $A$-module maps $M \to N$ such that elements of $S$ are invertible in $N$ ($s \times : N \to N$ is an isomorphism for all $s \in S$). More precisely, any such map $\alpha : M \to N$ factors uniquely through $\phi$:

\[
\begin{array}{ccc}
M & \longrightarrow & S^{-1}M \\
\alpha & \downarrow & \exists! \\
N & \end{array}
\]

(Translation: $M \to S^{-1}M$ is universal (initial) among $A$-module maps from $M$ to modules that are actually $S^{-1}A$-modules. Can you make this precise by defining clearly the objects and morphisms in this category?)

Notice: (i) this determines $\phi : M \to S^{-1}M$ up to unique isomorphism (you should think through what this means); (ii) we are defining not only $S^{-1}M$, but also the map $\phi$ at the same time; and (iii) essentially by definition the $A$-module structure on $S^{-1}M$ extends to an $S^{-1}A$-module structure.

**2.3.E. Exercise.** Show that $\phi : M \to S^{-1}M$ exists, by constructing something satisfying the universal property. Hint: define elements of $S^{-1}M$ to be of the form $m/s$ where $m \in M$ and $s \in S$, and $m_1/s_1 = m_2/s_2$ if and only if for some $s \in S$, $s(s_1m_1 - s_1m_2) = 0$. Define the additive structure by $(m_1/s_1) + (m_2/s_2) = (s_2m_1 + s_1m_2)/(s_1s_2)$, and the $S^{-1}A$-module structure (and hence the $A$-module structure) is given by $(a_1/s_1) \cdot (m_2/s_2) = (a_1m_2)/(s_1s_2)$.

**2.3.F. Exercise.** Show that localization commutes with finite products. In other words, if $M_1, \ldots, M_n$ are $A$-modules, describe an isomorphism $S^{-1}(M_1 \times \cdots \times M_n) \to S^{-1}M_1 \times \cdots \times S^{-1}M_n$. Show that localization does not necessarily commute with infinite products. (Hint: $\{1, 1/2, 1/3, 1/4, \ldots\} \in \mathbb{Q} \times \mathbb{Q} \times \cdots$)

**2.3.4. Tensor products.** Another important example of a universal property construction is the notion of a tensor product of $A$-modules

\[
\otimes_A : \text{obj}(\text{Mod}_A) \times \text{obj}(\text{Mod}_A) \longrightarrow \text{obj}(\text{Mod}_A)
\]

\[
(M, N) \longmapsto M \otimes_A N
\]

The subscript $A$ is often suppressed when it is clear from context. The tensor product is often defined as follows. Suppose you have two $A$-modules $M$ and $N$. Then elements of the tensor product $M \otimes_A N$ are finite $A$-linear combinations of symbols $m \otimes n$ ($m \in M, n \in N$), subject to relations $(m_1 + m_2) \otimes n = m_1 \otimes n + m_2 \otimes n, m \otimes (n_1 + n_2) = m \otimes n_1 + m \otimes n_2, a(m \otimes n) = (am) \otimes n = m \otimes (an)$ (where $a \in A, m_1, m_2 \in M, n_1, n_2 \in N$). More formally, $M \otimes_A N$ is the free $A$-module
generated by \( M \times N \), quotiented by the submodule generated by \((m_1 + m_2) \otimes n - m_1 \otimes n - n_2 \otimes n, m \otimes (n_1 + n_2) - m \otimes n_1 - m \otimes n_2, a(m \otimes n) - (am) \otimes n \), and \( a(m \otimes n) - m \otimes (an) \) for \( a \in A, m, m_1, m_2 \in M, n, n_1, n_2 \in N \).

If \( A \) is a field \( k \), we recover the tensor product of vector spaces.

2.3.G. Exercise (if you haven’t seen tensor products before). Show that \( \mathbb{Z}/(10) \otimes_{\mathbb{Z}} \mathbb{Z}/(12) \cong \mathbb{Z}/(2) \). (This exercise is intended to give some hands-on practice with tensor products.)

2.3.H. Important Exercise: right-exactness of \( \cdot \otimes_A N \). Show that \( \cdot \otimes_A N \) gives a covariant functor \( \text{Mod}_A \to \text{Mod}_A \). Show that \( \cdot \otimes_A N \) is a right-exact functor, i.e. if

\[ M' \to M \to M'' \to \mathbb{0} \]

is an exact sequence of \( A \)-modules (which means \( f : M \to M'' \) is surjective, and \( M' \) surjects onto the kernel of \( f \); see (2.6) then the induced sequence

\[ M' \otimes_A N \to M \otimes_A N \to M'' \otimes_A N \to \mathbb{0} \]

is also exact. This exercise is repeated in Exercise 2.6.F, but you may get a lot out of doing it now: (You will be reminded of the definition of right-exactness in §2.6.4).

The constructive definition \( \otimes \) is a weird definition, and really the “wrong” definition. To motivate a better one: notice that there is a natural \( A \)-bilinear map \( M \times N \to M \otimes_A N \). (If \( M, N, P \in \text{Mod}_A \), a map \( f : M \times N \to P \) is \( A \)-bilinear if \( f(m_1 + m_2, n) = f(m_1, n) + f(m_2, n), f(m, n_1 + n_2) = f(m, n_1) + f(m, n_2), \) and \( f(am, n) = f(m, an) = af(m, n) \).) Any \( A \)-bilinear map \( M \times N \to C \) factors through the tensor product uniquely: \( M \times N \to M \otimes_A N \to C \). (Think this through!)

We can take this as the definition of the tensor product as follows. It is a \( A \)-module \( T \) along with an \( A \)-bilinear map \( t : M \times N \to T \), such that given any \( A \)-bilinear map \( t' : M \times N \to T' \), there is a unique \( A \)-linear map \( f : T \to T' \) such that \( t' = f \circ t \).

2.3.I. Exercise. Show that \( (T, t : M \times N \to T) \) is unique up to unique isomorphism. Hint: first figure out what “unique up to unique isomorphism” means for such pairs. Then follow the analogous argument for the product.

In short: given \( M \) and \( N \), there is an \( A \)-bilinear map \( t : M \times N \to M \otimes_A N \), unique up to unique isomorphism, defined by the following universal property: for any \( A \)-bilinear map \( t' : M \times N \to T' \) there is a unique \( A \)-linear map \( f : M \otimes_A N \to T' \) such that \( t' = f \circ t \).

As with all universal property arguments, this argument shows uniqueness assuming existence. To show existence, we need an explicit construction.

2.3.J. Exercise. Show that the construction of 2.3.I satisfies the universal property of tensor product.
The two exercises below are some useful facts about tensor products with which you should be familiar.

**2.3.K. Important Exercise.** (a) If $M$ is an $A$-module and $A \to B$ is a morphism of rings, show that $B \otimes_A M$ naturally has the structure of a $B$-module. Show that this describes a functor $\text{Mod}_A \to \text{Mod}_B$.

(b) If further $A \to C$ is a morphism of rings, show that $B \otimes_A C$ has the structure of a ring. Hint: multiplication will be given by $(b_1 \otimes c_1)(b_2 \otimes c_2) = (b_1 b_2) \otimes (c_1 c_2)$. (Exercise 2.3.T will interpret this construction as a coproduct.)

**2.3.L. Important Exercise.** If $S$ is a multiplicative subset of $A$ and $M$ is an $A$-module, describe a natural isomorphism $(S^{-1}A) \otimes_A M \cong S^{-1}M$ (as $S^{-1}A$-modules and as $A$-modules).

**2.3.5. Important Example: Fibered products.** (This notion will be essential later.) Suppose we have morphisms $f : X \to Z$ and $g : Y \to Z$ (in any category). Then the fibered product is an object $X \times_Z Y$ along with morphisms $\pi_X : X \times_Z Y \to X$ and $\pi_Y : X \times_Z Y \to Y$, where the two compositions $f \circ \pi_X, g \circ \pi_Y : X \times_Z Y \to Z$ agree, such that given any object $W$ with maps to $X$ and $Y$ (whose compositions to $Z$ agree), these maps factor through some unique $W \to X \times_Z Y$:

```
        W
       /\    \\
  X \times_Z Y ---\ π_Y
    /   \        ↓ g
   /     \    X  ↓ f
     \   /    ↓ Z
      \ \ /   
       \ π_X
```

(Warning: the definition of the fibered product depends on $f$ and $g$, even though they are omitted from the notation $X \times_Z Y$.)

By the usual universal property argument, if it exists, it is unique up to unique isomorphism. (You should think this through until it is clear to you.) Thus the use of the phrase “the fibered product” (rather than “a fibered product”) is reasonable, and we should reasonably be allowed to give it the name $X \times_Z Y$. We know what maps to it are: they are precisely maps to $X$ and maps to $Y$ (whose compositions to $Z$ agree). Depending on your religion, the diagram

```
        X \times_Z Y
       /\    \\
  X \times_Z Y ---\ π_Y
    /   \        ↓ g
   /     \    X  ↓ f
     \   /    ↓ Z
      \ \ /   
       \ π_X
```

is called a fibered/pullback/Cartesian diagram/square (six possibilities).

The right way to interpret the notion of fibered product is first to think about what it means in the category of sets.

**2.3.M. Exercise.** Show that in $\text{Sets}$,

$$X \times_Z Y = \{ (x \in X, y \in Y) : f(x) = g(y) \}.$$
More precisely, show that the right side, equipped with its evident maps to $X$ and $Y$, satisfies the universal property of the fibered product. (This will help you build intuition for fibered products.)

2.3.N. Exercise. If $X$ is a topological space, show that fibered products always exist in the category of open sets of $X$, by describing what a fibered product is. (Hint: it has a one-word description.)

2.3.O. Exercise. If $Z$ is the final object in a category $C$, and $X, Y \in C$, show that “$X \times_Z Y = X \times Y$”: “the” fibered product over $Z$ is uniquely isomorphic to “the” product. (This is an exercise about unwinding the definition.)

2.3.P. Useful Exercise: Towers of fiber diagrams are fiber diagrams. If the two squares in the following commutative diagram are fiber diagrams, show that the “outside rectangle” (involving $U, V, Y,$ and $Z$) is also a fiber diagram.

\[
\begin{array}{ccc}
U & \rightarrow & V \\
\downarrow & & \downarrow \\
W & \rightarrow & X \\
\downarrow & & \downarrow \\
Y & \rightarrow & Z
\end{array}
\]

2.3.Q. Exercise. Given $X \rightarrow Y \rightarrow Z$, show that there is a natural morphism $X \times_Y X \rightarrow X \times_Z X$, assuming that both fibered products exist. (This is trivial once you figure out what it is saying. The point of this exercise is to see why it is trivial.)

2.3.R. Useful exercise: The magic diagram. Suppose we are given morphisms $X_1, X_2 \rightarrow Y$ and $Y \rightarrow Z$. Describe the natural morphism $X_1 \times_Y X_2 \rightarrow X_1 \times_Z X_2$. Show that the following diagram is a fibered square.

\[
\begin{array}{ccc}
X_1 \times_Y X_2 & \rightarrow & X_1 \times_Z X_2 \\
\downarrow & & \downarrow \\
Y & \rightarrow & Y \times_Z Y
\end{array}
\]

This diagram is surprisingly incredibly useful — so useful that we will call it the magic diagram.

2.3.6. Coproducts. Define coproduct in a category by reversing all the arrows in the definition of product. Define fibered coproduct in a category by reversing all the arrows in the definition of fibered product.

2.3.S. Exercise. Show that coproduct for $Sets$ is disjoint union. (This is why we use the notation $\coprod$ for disjoint union.)

2.3.T. Exercise. Suppose $A \rightarrow B, C$ are two ring morphisms, so in particular $B$ and $C$ are $A$-modules. Recall (Exercise 2.3.K) that $B \otimes_A C$ has a ring structure. Show that there is a natural morphism $B \rightarrow B \otimes_A C$ given by $b \mapsto b \otimes 1$. (This is not necessarily an inclusion, see Exercise 2.3.C.) Similarly, there is a natural morphism...
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C → B ⊗_A C. Show that this gives a fibered coproduct on rings, i.e. that

\[
\begin{array}{c}
B \otimes_A C \\
\uparrow \\
B \\
\end{array}
\begin{array}{c}
\downarrow \\
A \\
\end{array}
\]

satisfies the universal property of fibered coproduct.

2.3.7. Monomorphisms and epimorphisms.

2.3.8. Definition. A morphism \( f : X \rightarrow Y \) is a monomorphism if any two morphisms \( g_1, g_2 : Z \rightarrow X \) such that \( f \circ g_1 = f \circ g_2 \) must satisfy \( g_1 = g_2 \). In other words, for any other object \( Z \), the natural map \( \text{Hom}(Z, X) \rightarrow \text{Hom}(Z, Y) \) is an injection. This a generalization of an injection of sets. In other words, there is at most one way of filling in the dotted arrow so that the following diagram commutes.

\[
\begin{array}{c}
Z \\
\downarrow \downarrow \\
X \\
\downarrow \nearrow \\
Y.
\end{array}
\]

Intuitively, it is the categorical version of an injective map, and indeed this notion generalizes the familiar notion of injective maps of sets. (The reason we don’t use the word “injective” is that in some contexts, “injective” will have an intuitive meaning which may not agree with “monomorphism”. This is also the case with “epimorphism” vs. “surjective”.)

2.3.U. Exercise. Show that the composition of two monomorphisms is a monomorphism.

2.3.V. Exercise. Prove that a morphism \( X \rightarrow Y \) is a monomorphism if and only if the induced morphism \( X \rightarrow X \times_Y X \) is an isomorphism. We may then take this as the definition of monomorphism. (Monomorphisms aren’t central to future discussions, although they will come up again. This exercise is just good practice.)

2.3.W. Exercise. Suppose \( Y \rightarrow Z \) is a monomorphism, and \( X_1, X_2 \rightarrow Y \) are two morphisms. Show that \( X_1 \times_Y X_2 \) and \( X_1 \times_Z X_2 \) are canonically isomorphic. We will use this later when talking about fibered products. (Hint: for any object \( V \), give a natural bijection between maps from \( V \) to the first and maps from \( V \) to the second. It is also possible to use the magic diagram, Exercise 2.3.X.)

The notion of an epimorphism is “dual” to the definition of monomorphism, where all the arrows are reversed. This concept will not be central for us, although it turns up in the definition of an abelian category. Intuitively, it is the categorical version of a surjective map.

2.3.9. Representable functors and Yoneda’s lemma. Much of our discussion about universal properties can be cleanly expressed in terms of representable functors, under the rubric of “Yoneda’s Lemma”. Yoneda’s lemma is an easy fact stated in a complicated way. Informally speaking, you can essentially recover an object in a category by knowing the maps into it. For example, we have seen that the
data of maps to $X \times Y$ are naturally (canonically) the data of maps to $X$ and to $Y$. Indeed, we have now taken this as the definition of $X \times Y$.

Recall Example 2.2.20 Suppose $A$ is an object of category $C$. For any object $C \in C$, we have a set of morphisms $\text{Mor}(C, A)$. If we have a morphism $f : B \to C$, we get a set of morphisms (2.3.9.1) $\text{Mor}(C, A) \to \text{Mor}(B, A)$, by composition: given a map from $C$ to $A$, we get a map from $B$ to $A$ by precomposing with $f : B \to C$. Hence this gives a contravariant functor $h_A : C \to \text{Sets}$. Yoneda’s Lemma states that the functor $h_A$ determines $A$ up to unique isomorphism. More precisely:

**2.3.X. Important Exercise That Everyone Should Do Once in Their Life (Yoneda’s Lemma).** Given two objects $A$ and $A'$ in a category $C$, and bijections (2.3.9.2) $i_C : \text{Mor}(C, A) \to \text{Mor}(C, A')$ that commute with the maps (2.3.9.1). Prove $i_C$ is induced from a unique isomorphism $A \to A'$. (Hint: This sounds hard, but it really is not. This statement is so general that there are really only a couple of things that you could possibly try. For example, if you’re hoping to find an isomorphism $A \to A'$, where will you find it? Well, you are looking for an element $\text{Mor}(A, A')$. So just plug in $C = A$ to (2.3.9.2), and see where the identity goes. You will quickly find the desired morphism; show that it is an isomorphism, then show that it is unique.)

There is an analogous statement with the arrows reversed, where instead of maps into $A$, you think of maps from $A$. The role of the contravariant functor $h_A$ of Example 2.2.20 is played by the covariant functor $h^A$ of Example 2.2.14. Because the proof is the same (with the arrows reversed), you needn’t think it through.

Yoneda’s lemma properly refers to a more general statement. Although it looks more complicated, it is no harder to prove.

**2.3.Y. * Exercise.**
(a) Suppose $A$ and $B$ are objects in a category $C$. Give a bijection between the natural transformations $h^A \to h^B$ of covariant functors $C \to \text{Sets}$ (see Example 2.2.14 for the definition) and the morphisms $B \to A$.
(b) State and prove the corresponding fact for contravariant functors $h_A$ (see Exercise 2.2.20). Remark: A contravariant functor $F$ from $C$ to $\text{Sets}$ is said to be representable if there is a natural isomorphism

$$
i : F \cong h_A.$$

Thus the representing object $A$ is determined up to unique isomorphism by the pair $(F, \ni)$. There is a similar definition for covariant functors. (We will revisit this in §7.6 and this problem will appear again as Exercise 7.6.B)
(c) **Yoneda’s lemma.** Suppose $F$ is a covariant functor $C \to \text{Sets}$, and $A \in C$. Give a bijection between the natural transformations $h^A \to F$ and $F(A)$. (The corresponding fact for contravariant functors is essentially Exercise 10.1.C)

In fancy terms, Yoneda’s lemma states the following. Given a category $C$, we can produce a new category, called the functor category of $C$, where the objects are contravariant functors $C \to \text{Sets}$, and the morphisms are natural transformations
of such functors. We have a functor (which we can usefully call \( h \)) from \( C \) to its functor category, which sends \( A \) to \( h_A \). Yoneda’s Lemma states that this is a fully faithful functor, called the \textit{Yoneda embedding}. (Fully faithful functors were defined in \([2.2.15]\)).

2.4 Limits and colimits

Limits and colimits are two important definitions determined by universal properties. They generalize a number of familiar constructions. I will give the definition first, and then show you why it is familiar. For example, fractions will be motivating examples of colimits (Exercise \( 2.4.B(a) \)), and the \( p \)-adic numbers (Example \( 2.4.3 \)) will be motivating examples of limits.

2.4.1. Limits. We say that a category is a \textbf{small category} if the objects and the morphisms are sets. (This is a technical condition intended only for experts.) Suppose \( I \) is any small category, and \( C \) is any category. Then a functor \( F : I \to C \) (i.e. with an object \( A_i \in C \) for each element \( i \in I \), and appropriate commuting morphisms dictated by \( I \)) is said to be a \textbf{diagram indexed by} \( I \). We call \( I \) an \textbf{index category}. Our index categories will be partially ordered sets (Example \( 2.2.8 \)), in which in particular there is at most one morphism between any two objects. (But other examples are sometimes useful.) For example, if \( \square \) is the category

and \( A \) is a category, then a functor \( \square \to A \) is precisely the data of a commuting square in \( A \).

Then the \textbf{limit} is an object \( \operatorname{lim} \downarrow_{\downarrow I} A_i \) of \( C \) along with morphisms \( f_j : \operatorname{lim} \downarrow_{\downarrow I} A_i \to A_j \) such that if \( m : j \to k \) is a morphism in \( I \), then

\[
\begin{array}{ccc}
\operatorname{lim} \downarrow_{\downarrow I} A_i & \xrightarrow{f_j} & A_j \\
\downarrow f_i & & \downarrow \text{F(m)} \\
A_j & \xrightarrow{g} & A_k
\end{array}
\]

commutes, and this object and maps to each \( A_i \) are universal (final) with respect to this property. More precisely, given any other object \( W \) along with maps \( g_i : W \to A_i \) commuting with the \( \text{F(m)} \) (if \( m : i \to j \) is a morphism in \( I \), then \( g_j = \text{F(m)} \circ g_i \)), then there is a unique map \( g : W \to \operatorname{lim} \downarrow_{\downarrow I} A_i \) so that \( g_i = f_i \circ g \) for all \( i \). (In some cases, the limit is sometimes called the \textbf{inverse limit} or \textbf{projective limit}. We won’t use this language.) By the usual universal property argument, if the limit exists, it is unique up to unique isomorphism.
2.4.2. Examples: products. For example, if $I$ is the partially ordered set
\[ \bullet \to \bullet \]
we obtain the fibered product.

If $I$ is
\[ \bullet \]
we obtain the product.

If $I$ is a set (i.e. the only morphisms are the identity maps), then the limit is called the product of the $A_i$, and is denoted $\prod_i A_i$. The special case where $I$ has two elements is the example of the previous paragraph.

If $I$ has an initial object $e$, then $A_e$ is the limit, and in particular the limit always exists.

2.4.3. Example: the $p$-adic numbers. The $p$-adic numbers, $\mathbb{Z}_p$, are often described informally (and somewhat unnaturally) as being of the form $\mathbb{Z}_p = ? + ? p + ? p^2 + ? p^3 + \cdots$. They are an example of a limit in the category of rings:

\[
\mathbb{Z}_p \to \cdots \to \mathbb{Z}/p^3 \to \mathbb{Z}/p^2 \to \mathbb{Z}/p
\]

Limits do not always exist for any index category $I$. However, you can often easily check that limits exist if the objects of your category can be interpreted as sets with additional structure, and arbitrary products exist (respecting the set-like structure).

2.4.A. Important Exercise. Show that in the category $\text{Sets}$,

\[
\left\{ \{a_i\}_{i \in I} \in \prod_i A_i : F(m)(a_i) = a_j \text{ for all } m \in \text{Mor}_I(i, j) \right\}
\]

along with the obvious projection maps to each $A_i$, is the limit $\lim_I A_i$.

This clearly also works in the category $\text{Mod}_A$ of $A$-modules, and its specializations such as $\text{Vec}_k$ and $\text{Ab}$.

From this point of view, $2 + 3p + 2p^2 + \cdots \in \mathbb{Z}_p$ can be understood as the sequence $(2, 2 + 3p, 2 + 3p + 2p^2, \ldots)$.

2.4.4. Colimits. More immediately relevant for us will be the dual (arrow-reversed version) of the notion of limit (or inverse limit). We just flip all the arrows in that definition, and get the notion of a colimit. Again, if it exists, it is unique up to unique isomorphism. (In some cases, the colimit is sometimes called the direct limit, inductive limit, or injective limit. We won’t use this language. I prefer using limit/colimit in analogy with kernel/cokernel and product/coproduct. This is more than analogy, as kernels and products may be interpreted as limits, and similarly with cokernels and coproducts. Also, I remember that kernels “map to”,
and cokernels are "mapped to", which reminds me that a limit maps to all the objects in the big commutative diagram indexed by $I$; and a colimit has a map from all the objects.)

Even though we have just flipped the arrows, colimits behave quite differently from limits.

### 2.4.5. Example

The group $5^{\sim \infty} \mathbb{Z}$ of rational numbers whose denominators are powers of 5 is a colimit $\lim \cdots \rightarrow 5^{-1}\mathbb{Z} \rightarrow 5^{-2}\mathbb{Z} \rightarrow \cdots$

The colimit over an index set $I$ is called the coproduct, denoted $\bigsqcup I A_i$, and is the dual (arrow-reversed) notion to the product.

### 2.4.B. Exercise

(a) Interpret the statement "$\mathbb{Q} = \lim_{\rightarrow} \frac{1}{n} \mathbb{Z}$". (b) Interpret the union of the some subsets of a given set as a colimit. (Dually, the intersection can be interpreted as a limit.) The objects of the category in question are the subsets of the given set.

Colimits don’t always exist, but there are two useful large classes of examples for which they do.

### 2.4.6. Definition

A nonempty partially ordered set $(S, \geq)$ is filtered (or is said to be a filtered set) if for each $x, y \in S$, there is a $z$ such that $x \geq z$ and $y \geq z$. More generally, a nonempty category $I$ is filtered if:

(i) for each $x, y \in I$, there is a $z \in I$ and arrows $x \to z$ and $y \to z$, and

(ii) for every two arrows $u, v : x \to y$, there is an arrow $w : y \to z$ such that $w \circ u = w \circ v$.

(Other terminologies are also commonly used, such as “directed partially ordered set” and “filtered index category”, respectively.)

### 2.4.C. Exercise

Suppose $I$ is filtered. (We will almost exclusively use the case where $I$ is a filtered set.) Show that any diagram in $\text{Sets}$ indexed by $I$ has the following as a colimit:

$$\left\{ a \in \bigsqcup_{i \in I} A_i \right\} / (a_i \in A_i \sim (f(a_i) \in A_j) \text{ for every } f : A_i \to A_j \text{ in the diagram}.$$

(Hint: Verify that $\sim$ is indeed an equivalence relation, by writing it as $(a_i \in A_i) \sim (a_j \in A_j)$ if there are $f : A_i \to A_k$ and $g : A_j \to A_k$ with $f(a_i) = g(a_j)$.)

This idea applies to many categories whose objects can be interpreted as sets with additional structure (such as abelian groups, $A$-modules, groups, etc.). For example, in Example 2.4.5 each element of the colimit is an element of something upstairs, but you can’t say in advance what it is an element of. For example, $17/125$ is an element of the $5^{-3}\mathbb{Z}$ (or $5^{-4}\mathbb{Z}$, or later ones), but not $5^{-2}\mathbb{Z}$. More generally, in the category of $A$-modules $\text{Mod}_A$, each element $a$ of the colimit $\lim A_i$ can be interpreted as an element of some $a \in A_i$. The element $a \in \lim A_i$ is 0 if there is some $m : i \to j$ such that $F_m(a) = 0$ (i.e. if it becomes 0 “later in the diagram”). Furthermore, two elements interpreted as $a_i \in A_i$ and $a_j \in A_j$ are the same if there are some arrows $m : i \to k$ and $n : j \to k$ such that $F_m(a_i) = F_n(a_j)$, i.e. if they become the same “later in the diagram”. To add two elements interpreted
as \( a_i \in A_i \) and \( a_j \in A_j \), we choose arrows \( m : i \to k \) and \( n : j \to k \), and then interpret their sum as \( F(m)(a_i) + F(n)(a_j) \).

### 2.4.D. Exercise
Verify that the \( A \)-module described above is indeed the colimit.

### 2.4.E. Useful Exercise (Localization as Colimit)
Generalize Exercise 2.4.B(a) to interpret localization of an integral domain as a colimit over a filtered set: suppose \( S \) is a multiplicative set of \( A \) and interpret \( S^{-1}A = \varprojlim_s A_s \) where the limit is over \( s \in S \). (Aside: Can you make some version of this work even if \( A \) isn’t an integral domain, e.g. \( S^{-1}A = \varprojlim_s A_s \)?)

A variant of this construction works without the filtered condition, if you have another means of “connecting elements in different objects of your diagram”. For example:

### 2.4.F. Exercise: Colimits of \( A \)-modules without the filtered condition.
Suppose you are given a diagram of \( A \)-modules indexed by \( I \): \( F : I \to \text{Mod}_A \), where we let \( A_i := F(i) \). Show that the colimit is \( \bigoplus_{i \in I} A_i \) modulo the relations \( a_j - F(m)(a_i) \) for every \( m : i \to j \) in \( I \) (i.e. for every arrow in the diagram).

The following exercise shows that you have to be careful to remember which category you are working in.

### 2.4.G. Unimportant Exercise
Consider the filtered set of abelian groups \( \mathbb{Z}_p / p^{-n} \mathbb{Z}_p \) (here \( p \) is a fixed prime, and \( n \) varies — you should be able to figure out the index set). Show that this system has colimit \( \mathbb{Q}_p / \mathbb{Z}_p \) in the category of abelian groups, and has colimit \( 0 \) in the category of finite abelian groups. Here \( \mathbb{Q}_p \) is the fraction field of \( \mathbb{Z}_p \), which can be interpreted as \( \bigcup p^{-n} \mathbb{Z}_p \).

### 2.4.7. Summary
One useful thing to informally keep in mind is the following. In a category where the objects are “set-like”, an element of a limit can be thought of as an element in each object in the diagram, that are “compatible” (Exercise 2.4.A). And an element of a colimit can be thought of (“has a representative that is”) an element of a single object in the diagram (Exercise 2.4.C). Even though the definitions of limit and colimit are the same, just with arrows reversed, these interpretations are quite different.

### 2.5 Adjoint

We next come to a very useful construction closely related to universal properties. Just as a universal property “essentially” (up to unique isomorphism) determines an object in a category (assuming such an object exists), “adjoints” essentially determine a functor (again, assuming it exists). Two **covariant** functors \( F : A \to B \) and \( G : B \to A \) are **adjoint** if there is a natural bijection for all \( A \in A \) and \( B \in B \)

\[
\tau_{AB} : \text{Mor}_B(F(A), B) \to \text{Mor}_A(A, G(B)).
\]

We say that \( (F, G) \) form an **adjoint pair**, and that \( F \) is **left-adjoint** to \( G \) (and \( G \) is **right-adjoint** to \( F \)). By “natural” we mean the following. For all \( f : A \to A' \) in \( A \),
we require

\[(2.5.0.2) \quad \text{Mor}_G(F(A'), B) \xrightarrow{Ff^*} \text{Mor}_G(F(A), B) \]

\[
\downarrow \tau_{A'B} \quad \downarrow \tau_{AB}
\]

\[
\text{Mor}_A(A', G(B)) \xrightarrow{f^*} \text{Mor}_A(A, G(B))
\]

to commute, and for all \(g : B \to B'\) in \(B\) we want a similar commutative diagram to commute. (Here \(f^*\) is the map induced by \(f : A \to A'\), and \(Ff^*\) is the map induced by \(Ff : F(A) \to F(A')\).)

**2.5.A. Exercise.** Write down what this diagram should be. (Hint: do it by extending diagram \((2.5.0.2)\) above.)

**2.5.B. Exercise.** Show that the map \(\tau_{AB}\) \((2.5.0.1)\) is given as follows. For each \(A\) there is a map \(\eta_A : A \to GF(A)\) so that for any \(g : F(A) \to B\), the corresponding \(f : A \to G(B)\) is given by the composition

\[
A \xrightarrow{\eta_A} GF(A) \xrightarrow{Gg} G(B).
\]

Similarly, there is a map \(\epsilon_B : FG(B) \to B\) for each \(B\) so that for any \(f : A \to G(B)\), the corresponding map \(g : F(A) \to B\) is given by the composition

\[
F(A) \xrightarrow{Ff} FG(B) \xrightarrow{\epsilon_B} B.
\]

Here is an example of an adjoint pair.

**2.5.C. Exercise.** Suppose \(M, N,\) and \(P\) are \(A\)-modules. Describe a bijection \(\text{Hom}_A(M \otimes_A N, P) \leftrightarrow \text{Hom}_A(M, \text{Hom}_A(N, P))\). (Hint: try to use the universal property.)

**2.5.D. Exercise.** Show that \(\otimes_A N\) and \(\text{Hom}_A(N, \cdot)\) are adjoint functors.

**2.5.1. Fancier remarks we won’t use.** You can check that the left adjoint determines the right adjoint up to unique natural isomorphism, and vice versa, by a universal property argument. The maps \(\eta_A\) and \(\epsilon_B\) of Exercise \((2.5.0.1)\) are called the unit and counit of the adjunction. This leads to a different characterization of adjunction. Suppose functors \(F : A \to B\) and \(G : B \to A\) are given, along with natural transformations \(\epsilon : FG \to \text{id}\) and \(\eta : \text{id} \to GF\) with the property that \(G\epsilon \circ \eta G = \text{id}_G\) for each \(B \in B\), the composition of \(\eta_{G(B)} : G(B) \to GF(B)\) and \(G(\epsilon_B) : GFG(B) \to G(B)\) is the identity) and \(\epsilon F \circ F\eta = \text{id}_F\). Then you can check that \(F\) is left adjoint to \(G\). These facts aren’t hard to check, so if you want to use them, you should verify everything for yourself.

**2.5.2. Examples from other fields.** For those familiar with representation theory: Frobenius reciprocity may be understood in terms of adjoints. Suppose \(V\) is a finite-dimensional representation of a finite group \(G\), and \(W\) is a representation of a subgroup \(H < G\). Then induction and restriction are an adjoint pair \((\text{Ind}^G_H, \text{Res}^G_H)\) between the category of \(G\)-modules and the category of \(H\)-modules.

Topologists’ favorite adjoint pair may be the suspension functor and the loop space functor.
2.5.3. **Example: groupification.** Here is another motivating example: getting an abelian group from an abelian semigroup. An abelian semigroup is just like an abelian group, except you don’t require an inverse. One example is the non-negative integers \(0, 1, 2, \ldots\) under addition. Another is the positive integers under multiplication \(1, 2, \ldots\). From an abelian semigroup, you can create an abelian group. Here is a formalization of that notion. If \(S\) is a semigroup, then its **groupification** is a map of semigroups \(\pi : S \rightarrow G\) such that \(G\) is a group, and any other map of semigroups from \(S\) to a group \(G'\) factors uniquely through \(G\).

\[
\begin{array}{ccc}
S & \xrightarrow{\pi} & G \\
\downarrow & & \downarrow \\
& & \exists \uparrow \\
& & G' \\
\end{array}
\]

2.5.E. **Exercise.** Construct groupification \(H\) from the category of abelian semigroups to the category of abelian groups. (One possibility of a construction: given an abelian semigroup \(S\), the elements of its groupification \(H(S)\) are \((a, b)\), which you may think of as \(a - b\), with the equivalence that \((a, b) \sim (c, d)\) if \(a + d + e = b + c + e\) for some \(e \in S\). Describe addition in this group, and show that it satisfies the properties of an abelian group. Describe the semigroup map \(S \rightarrow H(S)\).) Let \(F\) be the forgetful morphism from the category of abelian groups \(Ab\) to the category of abelian semigroups. Show that \(H\) is left-adjoint to \(F\).

(Here is the general idea for experts: We have a full subcategory of a category. We want to “project” from the category to the subcategory. We have

\[
\text{Mor}_{\text{category}}(S, H) = \text{Mor}_{\text{subcategory}}(G, H)
\]

automatically; thus we are describing the left adjoint to the forgetful functor. How the argument worked: we constructed something which was in the smaller category, which automatically satisfies the universal property.)

2.5.F. **Exercise.** Show that if a semigroup is already a group then the identity morphism is the groupification (“the semigroup is groupified by itself”), by the universal property. (Perhaps better: the identity morphism is a groupification — but we don’t want tie ourselves up in knots over categorical semantics.)

2.5.G. **Exercise.** The purpose of this exercise is to give you some practice with “adjoints of forgetful functors”, the means by which we get groups from semigroups, and sheaves from presheaves. Suppose \(A\) is a ring, and \(S\) is a multiplicative subset. Then \(S^{-1}A\)-modules are a fully faithful subcategory of the category of \(A\)-modules (meaning: the objects of the first category are a subset of the objects of the second; and the morphisms between any two objects of the second that are secretly objects of the first are just the morphisms from the first). Then \(\text{Mod}_A \rightarrow \text{Mod}_{S^{-1}A}\) can be interpreted as an adjoint to the forgetful functor \(\text{Mod}_{S^{-1}A} \rightarrow \text{Mod}_A\). Figure out the correct statement, and prove that it holds.

(Here is the larger story. Every \(S^{-1}A\)-module is an \(A\)-module, and this is an injective map, so we have a covariant forgetful functor \(F : \text{Mod}_{S^{-1}A} \rightarrow \text{Mod}_A\). In fact this is a fully faithful functor: it is injective on objects, and the morphisms between any two \(S^{-1}A\)-modules as \(A\)-modules are just the same when they are considered as \(S^{-1}A\)-modules. Then there is a functor \(G : \text{Mod}_A \rightarrow \text{Mod}_{S^{-1}A}\), which
might reasonably be called “localization with respect to S”, which is left-adjoint to the forgetful functor. Translation: If \( M \) is an \( A \)-module, and \( N \) is an \( S^{-1}A \)-module, then \( \text{Mor}(GM, N) \) (morphisms as \( S^{-1}A \)-modules, which are the same as morphisms as \( A \)-modules) are in natural bijection with \( \text{Mor}(M, FN) \) (morphisms as \( A \)-modules).

Here is a table of adjoints that will come up for us.

<table>
<thead>
<tr>
<th>situation</th>
<th>category ( A )</th>
<th>category ( B )</th>
<th>left-adjoint ( F : A \rightarrow B )</th>
<th>right-adjoint ( G : B \rightarrow A )</th>
</tr>
</thead>
<tbody>
<tr>
<td>( A )-modules (Ex. 2.5.D)</td>
<td></td>
<td></td>
<td>( \cdot \otimes_A N )</td>
<td>( \text{Hom}_A[N, \cdot] )</td>
</tr>
<tr>
<td>ring maps ( A \rightarrow B )</td>
<td>( \text{Mod}_A )</td>
<td>( \text{Mod}_B )</td>
<td>( \cdot \otimes_A B ) (extension of scalars)</td>
<td>forgetful (restriction of scalars)</td>
</tr>
<tr>
<td>(pre)sheaves on a topological space ( X ) (Ex. 3.4.K)</td>
<td>presheaves on ( X )</td>
<td>sheaves on ( X )</td>
<td>sheafification</td>
<td>forgetful</td>
</tr>
<tr>
<td>(semi)groups (2.5.3)</td>
<td>semigroups</td>
<td>groups</td>
<td>groupification</td>
<td>forgetful</td>
</tr>
<tr>
<td>sheaves, ( f : X \rightarrow Y ) (Ex. 3.6.B)</td>
<td>sheaves on ( Y )</td>
<td>sheaves on ( X )</td>
<td>( f^{-1} )</td>
<td>( f_* )</td>
</tr>
<tr>
<td>sheaves of abelian groups or ( \mathcal{O} )-modules, open immersions ( f : U \rightarrow Y ) (Ex. 3.6.G)</td>
<td>sheaves on ( U )</td>
<td>sheaves on ( Y )</td>
<td>( f_! )</td>
<td>( f^{-1} )</td>
</tr>
<tr>
<td>quasicoherent sheaves, ( f : X \rightarrow Y ) (Prop. 17.3.5)</td>
<td>quasicoherent sheaves on ( Y )</td>
<td>quasicoherent sheaves on ( X )</td>
<td>( f^* )</td>
<td>( f_* )</td>
</tr>
</tbody>
</table>

Other examples will also come up, such as the adjoint pair \((\sim, \Gamma)\) between graded modules over a graded ring, and quasicoherent sheaves on the corresponding projective scheme (16.4).

2.5.4. Useful comment for experts. One last comment only for people who have seen adjoints before: If \((F, G)\) is an adjoint pair of functors, then \( F \) commutes with colimits, and \( G \) commutes with limits. Also, limits commute with limits and colimits commute with colimits. We will prove these facts (and a little more) in §2.6.10.

2.6 Kernels, cokernels, and exact sequences: A brief introduction to abelian categories

Since learning linear algebra, you have been familiar with the notions and behaviors of kernels, cokernels, etc. Later in your life you saw them in the category of abelian groups, and later still in the category of \( A \)-modules. Each of these notions generalizes the previous one.

We will soon define some new categories (certain sheaves) that will have familiar-looking behavior, reminiscent of that of modules over a ring. The notions of kernels, cokernels, images, and more will make sense, and they will behave “the way we expect” from our experience with modules. This can be made precise through the notion of an abelian category. Abelian categories are the right general setting
in which one can do “homological algebra”, in which notions of kernel, cokernel, and so on are used, and one can work with complexes and exact sequences.

We will see enough to motivate the definitions that we will see in general: monomorphism (and subobject), epimorphism, kernel, cokernel, and image. But in these notes we will avoid having to show that they behave “the way we expect” in a general abelian category because the examples we will see are directly interpretable in terms of modules over rings. In particular, it is not worth memorizing the definition of abelian category.

Two central examples of an abelian category are the category $\mathbf{Ab}$ of abelian groups, and the category $\mathbf{Mod}_A$ of $A$-modules. The first is a special case of the second (just take $A = \mathbb{Z}$). As we give the definitions, you should verify that $\mathbf{Mod}_A$ is an abelian category.

We first define the notion of additive category. We will use it only as a stepping stone to the notion of an abelian category.

2.6.1. Definition. A category $C$ is said to be additive if it satisfies the following properties.

Ad1. For each $A, B \in C$, $\text{Mor}(A, B)$ is an abelian group, such that composition of morphisms distributes over addition. (You should think about what this means — it translates to two distinct statements).

Ad2. $C$ has a zero object, denoted $0$. (This is an object that is simultaneously an initial object and a final object, Definition 2.3.2).

Ad3. It has products of two objects (a product $A \times B$ for any pair of objects), and hence by induction, products of any finite number of objects.

In an additive category, the morphisms are often called homomorphisms, and $\text{Mor}$ is denoted by $\text{Hom}$. In fact, this notation $\text{Hom}$ is a good indication that you’re working in an additive category. A functor between additive categories preserving the additive structure of $\text{Hom}$, is called an additive functor.

2.6.2. Remarks. It is a consequence of the definition of additive category that finite direct products are also finite direct sums (coproducts) — the details don’t matter to us. The symbol $\oplus$ is used for this notion. Also, it is quick to show that additive functors send zero objects to zero objects (show that $\alpha$ is a $0$-object if and only if $\text{id}_\alpha = 0$, additive functors preserve both $\text{id}$ and $0$), and preserves products.

One motivation for the name $0$-object is that the $0$-morphism in the abelian group $\text{Hom}(A, B)$ is the composition $A \to 0 \to B$.

Real (or complex) Banach spaces are an example of an additive category. The category of free $A$-modules is another. The category of $A$-modules $\mathbf{Mod}_A$ is also an example, but it has even more structure, which we now formalize as an example of an abelian category.

2.6.3. Definition. Let $C$ be an additive category. A kernel of a morphism $f : B \to C$ is a map $i : A \to B$ such that $f \circ i = 0$, and that is universal with respect
to this property. Diagramatically:

(Note that the kernel is not just an object; it is a morphism of an object to $B$.) Hence it is unique up to unique isomorphism by universal property nonsense. A **cokernel** is defined dually by reversing the arrows — do this yourself. The kernel of $f : B \to C$ is the limit (2.4) of the diagram

and similarly the cokernel is a colimit.

If $i : A \to B$ is a monomorphism, then we say that $A$ is a **subobject** of $B$, where the map $i$ is implicit. Dually, there is the notion of **quotient object**, defined dually to subobject.

An **abelian category** is an additive category satisfying three additional properties.

1. Every map has a kernel and cokernel.
2. Every monomorphism is the kernel of its cokernel.
3. Every epimorphism is the cokernel of its kernel.

It is a non-obvious (and imprecisely stated) fact that every property you want to be true about kernels, cokernels, etc. follows from these three. (Warning: in part of the literature, additional hypotheses are imposed as part of the definition.)

The **image** of a morphism $f : A \to B$ is defined as $\text{im}(f) = \ker(\text{coker } f)$. It is the unique factorization

$$A \xrightarrow{\text{epi.}} \text{im}(f) \xrightarrow{\text{mono.}} B$$

It is the cokernel of the kernel, and the kernel of the cokernel. The reader may want to verify this as an exercise. It is unique up to unique isomorphism. The cokernel of a monomorphism is called the **quotient**. The quotient of a monomorphism $A \to B$ is often denoted $B/A$ (with the map from $B$ implicit).

We will leave the foundations of abelian categories untouched. The key thing to remember is that if you understand kernels, cokernels, images and so on in the category of modules over a ring $\text{Mod}_A$, you can manipulate objects in any abelian category. This is made precise by Freyd-Mitchell Embedding Theorem. (The Freyd-Mitchell Embedding Theorem: If $\mathcal{A}$ is an abelian category such that $\text{Hom}(a, a')$ is a set for all $a, a' \in \mathcal{A}$, then there is a ring $A$ and an exact, fully faithful functor from $\mathcal{A}$ into $\text{Mod}_A$, which embeds $\mathcal{A}$ as a full subcategory. A proof is sketched in [W] §1.6, and references to a complete proof are given there. The moral is that to prove something about a diagram in some abelian category, we may pretend that it is a diagram of modules over some ring, and we may then “diagram-chase” elements. Moreover, any fact about kernels, cokernels, and so on that holds in $\text{Mod}_A$ holds in any abelian category.) However, the abelian categories
we will come across will obviously be related to modules, and our intuition will clearly carry over, so we needn’t invoke a theorem whose proof we haven’t read. For example, we will show that sheaves of abelian groups on a topological space $X$ form an abelian category (§3.5), and the interpretation in terms of “compatible germs” will connect notions of kernels, cokernels etc. of sheaves of abelian groups to the corresponding notions of abelian groups.

2.6.4. Complexes, exactness, and homology.

We say a sequence

$$(2.6.4.1) \quad A \xrightarrow{f} B \xrightarrow{g} C$$

is a complex if $g \circ f = 0$, and is exact if $\ker g = \text{im} f$. An exact sequence with five terms, the first and last of which are 0, is a short exact sequence. Note that $A \xrightarrow{f} B \xrightarrow{g} C \xrightarrow{} 0$ being exact is equivalent to describing $C$ as a cokernel of $f$ (with a similar statement for $0 \xrightarrow{} A \xrightarrow{g} B \xrightarrow{f} C$).

If you would like practice in playing with these notions before thinking about homology, you can prove the Snake Lemma (stated in Example 2.7.5 with a stronger version in Exercise 2.7.B), or the Five Lemma (stated in Example 2.7.6 with a stronger version in Exercise 2.7.C).

If $(2.6.4.1)$ is a complex, then its homology (often denoted $H$) is $\ker g / \text{im} f$. We say that the $\ker g$ are the cycles, and $\text{im} f$ are the boundaries (so homology is “cycles mod boundaries”). If the complex is indexed in decreasing order, the indices are often written as subscripts, and $H_i$ is the homology at $A_{i+1} \rightarrow A_i \rightarrow A_{i-1}$. If the complex is indexed in increasing order, the indices are often written as superscripts, and the homology $H^i$ at $A^{i-1} \rightarrow A^i \rightarrow A^{i+1}$ is often called cohomology.

An exact sequence

$$(2.6.4.2) \quad A^* : \quad \cdots \rightarrow A^{i-1} \xrightarrow{f^{i-1}} A^i \xrightarrow{f^i} A^{i+1} \xrightarrow{f^{i+1}} \cdots$$

can be “factored” into short exact sequences

$$0 \rightarrow \ker f^i \rightarrow A^i \rightarrow \ker f^{i+1} \rightarrow 0$$

which is helpful in proving facts about long exact sequences by reducing them to facts about short exact sequences.

More generally, if $(2.6.4.2)$ is assumed only to be a complex, then it can be “factored” into short exact sequences.

$$(2.6.4.3) \quad 0 \rightarrow \ker f^i \rightarrow A^i \rightarrow \text{im} f^i \rightarrow 0$$

$$0 \rightarrow \text{im} f^{i-1} \rightarrow \ker f^i \rightarrow H^i(A^*) \rightarrow 0$$
2.6.A. EXERCISE. Describe exact sequences

(2.6.4.4) \[
0 \longrightarrow \text{im } f^i \longrightarrow A^{i+1} \longrightarrow \text{coker } f^i \longrightarrow 0
\]

\[
0 \longrightarrow H^i(A^\bullet) \longrightarrow \text{coker } f^{i-1} \longrightarrow \text{im } f^i \longrightarrow 0
\]

(These are somehow dual to (2.6.4.3). In fact in some mirror universe this might have been given as the standard definition of homology.)

2.6.B. EXERCISE. Suppose

\[
0 \longrightarrow d^0 \longrightarrow A^1 \longrightarrow d^1 \longrightarrow \cdots \longrightarrow d^{n-1} \longrightarrow A^n \longrightarrow d^n \longrightarrow 0
\]

is a complex of finite-dimensional \(k\)-vector spaces (often called \(A^\bullet\) for short). Show that \[\sum (-1)^i \dim A^i = \sum (-1)^i h^i(A^\bullet).\] (Recall that \(h^i(A^\bullet) = \dim \ker(d^i)/\im(d^{i-1}).\)) In particular, if \(A^\bullet\) is exact, then \[\sum (-1)^i \dim A^i = 0.\] (If you haven’t dealt much with cohomology, this will give you some practice.)

2.6.C. IMPORTANT EXERCISE. Suppose \(C\) is an abelian category. Define the category \(\text{Com}_C\) as follows. The objects are infinite complexes

\[
A^\bullet: \quad \cdots \longrightarrow A^{i-1} \overset{f^{i-1}}{\longrightarrow} A^i \overset{f^i}{\longrightarrow} A^{i+1} \overset{f^{i+1}}{\longrightarrow} \cdots
\]

in \(C\), and the morphisms \(A^\bullet \to B^\bullet\) are commuting diagrams

(2.6.4.5) \[
\begin{align*}
A^\bullet: & \quad \cdots \longrightarrow A^{i-1} \overset{f^{i-1}}{\longrightarrow} A^i \overset{f^i}{\longrightarrow} A^{i+1} \overset{f^{i+1}}{\longrightarrow} \cdots \\
B^\bullet: & \quad \cdots \longrightarrow B^{i-1} \overset{g^{i-1}}{\longrightarrow} B^i \overset{g^i}{\longrightarrow} B^{i+1} \overset{g^{i+1}}{\longrightarrow} \cdots
\end{align*}
\]

Show that \(\text{Com}_C\) is an abelian category. (Feel free to deal with the special case \(\text{Mod}_A\).)

2.6.D. IMPORTANT EXERCISE. Show that (2.6.4.5) induces a map of homology \(H(A^i) \to H(B^i).\) (Again, feel free to deal with the special case \(\text{Mod}_A\).)

We will later define when two maps of complexes are homotopic (23.1), and show that homotopic maps induce isomorphisms on cohomology (Exercise 23.1.A), but we won’t need that any time soon.
2.6.5. Theorem (Long exact sequence). — A short exact sequence of complexes

\[
\begin{array}{ccccccccc}
0^* & 
\rightarrow & 0 & 
\rightarrow & 0 & 
\rightarrow & 0 & 
\rightarrow & \cdots \\
A^* & 
\rightarrow & A^{i-1} & 
\rightarrow & f^{i-1} & 
\rightarrow & A^i & 
\rightarrow & A^{i+1} & 
\rightarrow & \cdots \\
B^* & 
\rightarrow & B^{i-1} & 
\rightarrow & g^{i-1} & 
\rightarrow & B^i & 
\rightarrow & B^{i+1} & 
\rightarrow & \cdots \\
C^* & 
\rightarrow & C^{i-1} & 
\rightarrow & h^{i-1} & 
\rightarrow & C^i & 
\rightarrow & C^{i+1} & 
\rightarrow & \cdots \\
0^* & 
\rightarrow & 0 & 
\rightarrow & 0 & 
\rightarrow & 0 & 
\rightarrow & \cdots 
\end{array}
\]

induces a long exact sequence in cohomology

\[
\cdots \rightarrow H^{i-1}(C^*) \rightarrow \\
H^i(A^*) \rightarrow H^i(B^*) \rightarrow H^i(C^*) \rightarrow \\
H^{i+1}(A^*) \rightarrow \cdots
\]

(This requires a definition of the connecting homomorphism \( H^{i-1}(C^*) \rightarrow H^i(A^*) \), which is natural in an appropriate sense.) For a concise proof in the case of complexes of modules, and a discussion of how to show this in general, see [W] §1.3. It will also come out of our discussion of spectral sequences as well (again, in the category of modules over a ring), see Exercise 2.7.E, but this is a somewhat perverse way of proving it.

2.6.6. Exactness of functors. If \( F : \mathcal{A} \rightarrow \mathcal{B} \) is a covariant additive functor from one abelian category to another, we say that \( F \) is right-exact if the exactness of

\[
A' \rightarrow A \rightarrow A'' \rightarrow 0,
\]

in \( \mathcal{A} \) implies that

\[
F(A') \rightarrow F(A) \rightarrow F(A'') \rightarrow 0
\]

is also exact. Dually, we say that \( F \) is left-exact if the exactness of

\[
0 \rightarrow A' \rightarrow A \rightarrow A'' \rightarrow 0
\]

implies

\[
0 \rightarrow F(A') \rightarrow F(A) \rightarrow F(A'')
\]

is exact.
A contravariant functor is **left-exact** if the exactness of

\[
A' \longrightarrow A \longrightarrow A'' \longrightarrow 0
\]

implies

\[
0 \longrightarrow F(A'') \longrightarrow F(A) \longrightarrow F(A')
\]
is exact.

The reader should be able to deduce what it means for a contravariant functor to be **right-exact**.

A covariant or contravariant functor is **exact** if it is both left-exact and right-exact.

**2.6.E. Exercise.** Suppose \( F \) is an exact functor. Show that applying \( F \) to an exact sequence preserves exactness. For example, if \( F \) is covariant, and \( A' \to A \to A'' \) is exact, then \( FA' \to FA \to FA'' \) is exact. (This will be generalized in Exercise 2.6.H(c).)

**2.6.F. Exercise.** Suppose \( A \) is a ring, \( S \subset A \) is a multiplicative subset, and \( M \) is an \( A \)-module.

(a) Show that localization of \( A \)-modules \( \text{Mod}_A \to \text{Mod}_{S^{-1}A} \) is an exact covariant functor.

(b) Show that \( \cdot \otimes M \) is a right-exact covariant functor \( \text{Mod}_A \to \text{Mod}_A \). (This is a repeat of Exercise 2.3.H)

(c) Show that \( \text{Hom}(M,\cdot) \) is a left-exact covariant functor \( \text{Mod}_A \to \text{Mod}_A \).

(d) Show that \( \text{Hom}(\cdot,M) \) is a left-exact contravariant functor \( \text{Mod}_A \to \text{Mod}_A \).

**2.6.G. Exercise.** Suppose \( M \) is a **finitely presented** \( A \)-module: \( M \) has a finite number of generators, and with these generators it has a finite number of relations; or usefully equivalently, fits in an exact sequence

\[
(2.6.6.1)\quad A^{\oplus q} \to A^{\oplus p} \to M \to 0
\]

Use (2.6.6.1) and the left-exactness of \( \text{Hom} \) to describe an isomorphism

\[
S^{-1}\text{Hom}_A(M,N) \cong \text{Hom}_{S^{-1}A}(S^{-1}M, S^{-1}N)
\]

(You might be able to interpret this in light of a variant of Exercise 2.6.H below, for left-exact contravariant functors rather than right-exact covariant functors.)

**2.6.7. ⋆ Two useful facts in homological algebra.**

We now come to two (sets of) facts I wish I had learned as a child, as they would have saved me lots of grief. They encapsulate what is best and worst of abstract nonsense. The statements are so general as to be nonintuitive. The proofs are very short. They generalize some specific behavior it is easy to prove in an ad hoc basis. Once they are second nature to you, many subtle facts will become obvious to you as special cases. And you will see that they will get used (implicitly or explicitly) repeatedly.

**2.6.8. ⋆ Interaction of homology and (right/left-)exact functors.**

You might wait to prove this until you learn about cohomology in Chapter 20 when it will first be used in a serious way.
2.6.H. Important exercise (the FHHF theorem). This result can take you far, and perhaps for that reason it has sometimes been called the fernbahnhof (FernbaHnHoF) theorem. Suppose \( F : A \to B \) is a covariant functor of abelian categories. Suppose \( C^\bullet \) is a complex in \( A \).

(a) (\( F \) right-exact yields \( H^\bullet F \to H^\bullet \)) If \( F \) is right-exact, describe a natural morphism \( F H^\bullet \to H^\bullet F \). (More precisely, for each \( i \), the left side is \( F \) applied to the cohomology at piece \( i \) of \( C^\bullet \), while the right side is the cohomology at piece \( i \) of \( FC^\bullet \).)

(b) (\( F \) left-exact yields \( F \to H^\bullet \)) If \( F \) is left-exact, describe a natural morphism \( H^\bullet F \to F H^\bullet \).

(c) (\( F \) exact yields \( F \to H^\bullet \)) If \( F \) is exact, show that the morphisms of (a) and (b) are inverses and thus isomorphisms.

Hint for (a): use \( C^p \to C^{p+1} \to \text{coker } d^p \to 0 \) to give an isomorphism \( F \text{coker } d^p \cong \text{coker } Fd^p \). Then use the first line of (2.6.4.4) to give a surjection \( \text{im } d^p \to \text{im } Fd^p \). Then use the second line of (2.6.4.4) to give the desired map \( FH^p C^\bullet \to H^p FC^\bullet \). While you are at it, you may as well describe a map for the fourth member of the quartet \( \{ \ker, \text{coker}, \text{im}, H \} : F \ker d^p \to \ker Fd^p \).

2.6.9. If this makes your head spin, you may prefer to think of it in the following specific case, where both \( A \) and \( B \) are the category of \( A \)-modules, and \( F \) is \( \cdot \otimes N \) for some fixed \( N \)-module. Your argument in this case will translate without change to yield a solution to Exercise 2.6.H(a) and (c) in general. If \( \otimes N \) is exact, then \( N \) is called a flat \( A \)-module. (The notion of flatness will turn out to be very important, and is discussed in detail in Chapter [24].)

For example, localization is exact, so \( S^{-1} A \) is a flat \( A \)-algebra for all multiplicative sets \( S \). Thus taking cohomology of a complex of \( A \)-modules commutes with localization — something you could verify directly.

2.6.10. * Interaction of adjoints, (co)limits, and (left- and right-) exactness.

A surprising number of arguments boil down to the statement:

Limits commute with limits and right-adjoints. In particular, because kernels are limits, both right-adjoints and limits are left exact.

as well as its dual:

Colimits commute with colimits and left-adjoints. In particular, because cokernels are colimits, both left-adjoints and colimits are right exact.

These statements were promised in [2.5.4] The latter has a useful extension: In an abelian category, colimits over filtered index categories are exact.

(“Filtered” was defined in [2.4.6]) If you want to use these statements (for example, later in these notes), you will have to prove them. Let’s now make them precise.

2.6.I. Exercise (Kernels commute with limits). Suppose \( C \) is an abelian category, and \( a : I \to C \) and \( b : I \to C \) are two diagrams in \( C \) indexed by \( I \). For convenience, let \( A_i = a(i) \) and \( B_i = b(i) \) be the objects in those two diagrams. Let \( h_i : A_i \to B_i \) be maps commuting with the maps in the diagram. (Translation: \( h_i \) is a natural transformation of functors \( a \to b \), see [2.2.21]) Then the \( \ker h_i \) form
another diagram in \( \mathcal{I} \) indexed by \( \mathcal{I} \). Describe a natural isomorphism \( \varprojlim \ker h_i \cong \ker(\varprojlim A_i \to \varprojlim B_i) \).

2.6.J. **Exercise.** Make sense of the statement that “limits commute with limits” in a general category, and prove it. (Hint: recall that kernels are limits. The previous exercise should be a corollary of this one.)

2.6.11. **Proposition (right-adjoints commute with limits).** — Suppose \( (F : C \to D, G : D \to C) \) is a pair of adjoint functors. If \( A = \varprojlim A_i \) is a limit in \( D \) of a diagram indexed by \( I \), then \( GA = \varprojlim GA_i \) (with the corresponding maps \( GA \to GA_i \)) is a limit in \( C \).

**Proof.** We must show that \( GA \to GA_i \) satisfies the universal property of limits. Suppose we have maps \( W \to GA_i \) commuting with the maps of \( I \). We wish to show that there exists a unique \( W \to GA \) extending the \( W \to GA_i \). By adjointness of \( F \) and \( G \), we can restate this as: Suppose we have maps \( FW \to A_i \) commuting with the maps of \( I \). We wish to show that there exists a unique \( FW \to A \) extending the \( FW \to A_i \). But this is precisely the universal property of the limit. □

Of course, the dual statements to Exercise 2.6.J and Proposition 2.6.11 hold by the dual arguments.

If \( F \) and \( G \) are additive functors between abelian categories, and \( (F, G) \) is an adjoint pair, then (as kernels are limits and cokernels are colimits) \( G \) is left-exact and \( F \) is right-exact.

2.6.K. **Exercise.** Show that in \( \text{Mod}_A \), colimits over filtered index categories are exact. (Your argument will apply without change to any abelian category whose objects can be interpreted as “sets with additional structure”.) Right-exactness follows from the above discussion, so the issue is left-exactness. (Possible hint: After you show that localization is exact, Exercise 2.6.H(a), or sheafification is exact, Exercise 3.5.D in a hands on way, you will be easily able to prove this. Conversely, if you do this exercise, those two will be easy.)

2.6.L. **Exercise.** Show that filtered colimits commute with homology. Hint: use the FHHF Theorem (Exercise 2.6.H), and the previous Exercise.

In light of Exercise 2.6.L you may want to think about how limits (and colimits) commute with homology in general, and which way maps go. The statement of the FHHF Theorem should suggest the answer. (Are limits analogous to left-exact functors, or right-exact functors?) We won’t directly use this insight.

2.6.12. **Dreaming of derived functors.** When you see a left-exact functor, you should always dream that you are seeing the end of a long exact sequence. If 

\[
0 \to M' \to M \to M'' \to 0
\]

is an exact sequence in abelian category \( A \), and \( F : A \to B \) is a left-exact functor, then

\[
0 \to FM' \to FM \to FM''
\]

is exact, and you should always dream that it should continue in some natural way. For example, the next term should depend only on \( M' \), call it \( R^1FM' \), and if it is zero, then \( FM \to FM'' \) is an epimorphism. This remark holds true for left-exact
and contravariant functors too. In good cases, such a continuation exists, and is incredibly useful. We will discuss this in Chapter 23.

2.7 * Spectral sequences

Spectral sequences are a powerful book-keeping tool for proving things involving complicated commutative diagrams. They were introduced by Leray in the 1940’s at the same time as he introduced sheaves. They have a reputation for being abstruse and difficult. It has been suggested that the name ‘spectral’ was given because, like spectres, spectral sequences are terrifying, evil, and dangerous. I have heard no one disagree with this interpretation, which is perhaps not surprising since I just made it up.

Nonetheless, the goal of this section is to tell you enough that you can use spectral sequences without hesitation or fear, and why you shouldn’t be frightened when they come up in a seminar. This will allow you to get some hands-on experience for how to use them. We will also see them only in the special case of double complexes (which is the version by far the most often used in algebraic geometry), and not in the general form usually presented (filtered complexes, exact couples, etc.). See [W, Ch. 5] for more detailed information if you wish.

You should not read this section when you are reading the rest of Chapter 2. Instead, you should read it just before you need it for the first time. When you finally do read this section, you must do the exercises.

For concreteness, we work in the category \( \text{Vec}_k \) of vector spaces over a field \( k \). However, everything we say will apply in any abelian category, such as the category \( \text{Mod}_A \) of \( A \)-modules.

2.7.1. Double complexes.

A **double complex** is a collection of vector spaces \( E^p,q \) (\( p, q \in \mathbb{Z} \)), and “rightward” morphisms \( d^{p,q}_r : E^{p,q} \to E^{p,q+1} \) and “upward” morphisms \( d^{p,q+1} \) : \( E^{p,q} \to E^{p+1,q} \). In the superscript, the first entry denotes the row number, and the second entry denotes the column number, in keeping with the convention for matrices, but opposite to how the \((x, y)\)-plane is labeled. The subscript is meant to suggest the direction of the arrows. We will always write these as \( d_\rightarrow \) and \( d_\uparrow \) and ignore the superscripts. We require that \( d_\rightarrow \) and \( d_\uparrow \) satisfying (a) \( d^2_{\rightarrow} = 0 \), (b) \( d^2_{\uparrow} = 0 \), and one more condition: (c) either \( d_\rightarrow d_\uparrow = d_\uparrow d_\rightarrow \) (all the squares commute) or \( d_\rightarrow d_\uparrow + d_\uparrow d_\rightarrow = 0 \) (they all anticommute). Both come up in nature, and you can switch from one to the other by replacing \( d^{p,q}_\uparrow \) with \( (-1)^q d^{p,q}_\rightarrow \). So I will assume that all the squares anticommute, but that you know how to turn the commuting case into this one. (You will see that there is no difference in the recipe, basically because the image and kernel of a homomorphism \( f \) equal the image and kernel
respectively of $-f$.)

\[
\begin{array}{c}
\mathbb{E}^{p+1,q} \\
\downarrow \quad d_{\rightarrow}^{p,q} \\
\mathbb{E}^{p,q} \\
\downarrow \quad d_{\rightarrow}^{p,q} \\
\mathbb{E}^{p,q+1} \\
\end{array}
\]

\[
\begin{array}{c}
\mathbb{E}^{p+1,q+1} \\
\downarrow \quad d_{\rightarrow}^{p+1,q} \\
\mathbb{E}^{p,q+1} \\
\downarrow \quad d_{\rightarrow}^{p,q+1}
\end{array}
\]

There are variations on this definition, where for example the vertical arrows go downwards, or some different subset of the $\mathbb{E}^{p,q}$ are required to be zero, but I will leave these straightforward variations to you.

From the double complex we construct a corresponding (single) complex $E^\bullet$ with $E^k = \bigoplus_i E^{i,k-i}$, with $d = d_{\rightarrow} + d_{\uparrow}$. In other words, when there is a single superscript $k$, we mean a sum of the $k$th antidiagonal of the double complex. The single complex is sometimes called the total complex. Note that $d^2 = (d_{\rightarrow} + d_{\uparrow})^2 = d_{\rightarrow}^2 + (d_{\rightarrow}d_{\uparrow} + d_{\uparrow}d_{\rightarrow} + d_{\uparrow}^2) = 0$, so $E^\bullet$ is indeed a complex.

The cohomology of the single complex is sometimes called the hypercohomology of the double complex. We will instead use the phrase “cohomology of the double complex”.

Our initial goal will be to find the cohomology of the double complex. You will see later that we secretly also have other goals.

A spectral sequence is a recipe for computing some information about the cohomology of the double complex. I won’t yet give the full recipe. Surprisingly, this fragmentary bit of information is sufficient to prove lots of things.

2.7.2. Approximate Definition. A spectral sequence with rightward orientation is a sequence of tables or pages $\cdots \to E_0^{p,q} \to E_1^{p,q} \to E_2^{p,q} \to \cdots$ ($p, q \in \mathbb{Z}$), where $E_0^{p,q} = E^{p,q}$, along with a differential

\[d_r^{p,q} : E_r^{p,q} \to E_r^{p+r,q-r+1}\]

with $d_r^{p,q} \circ d_r^{p,q} = 0$, and with an isomorphism of the cohomology of $d_r$ at $E_r^{p,q}$ (i.e. $\ker d_r^{p,q} / \text{im } d_r^{p-r,q+r-1}$) with $E_r^{p,q}$.

The orientation indicates that our 0th differential is the rightward one: $d_0 = d_{\rightarrow}$. The left subscript “$\to$” is usually omitted.

The order of the morphisms is best understood visually:

(2.7.2.1)

\[
\begin{array}{c}
\bullet \\
\downarrow d_1 \\
\bullet \\
\downarrow d_2 \\
\bullet \\
\downarrow d_3 \\
\bullet \\
\end{array}
\]

(the morphisms each apply to different pages). Notice that the map always is “degree 1” in the grading of the single complex $E^\bullet$. 

\[\text{anticommutes}\]
The actual definition describes what $E_{r}^{p,q}$ and $d_{r}^{p,q}$ really are, in terms of $E^{p,q}$. We will describe $d_0$, $d_1$, and $d_2$ below, and you should for now take on faith that this sequence continues in some natural way.

Note that $E_{r}^{p,q}$ is always a subquotient of the corresponding term on the 0th page $E_{0}^{p,q} = E^{p,q}$. In particular, if $E_{r}^{p,q} = 0$, then $E_{r+1}^{p,q} = 0$ for all $r$, so $E_{r}^{p,q} = 0$ unless $p, q \in \mathbb{Z}_{\geq 0}$.

Suppose now that $E_{\cdot,\cdot}$ is a **first quadrant double complex**, i.e. $E_{p,q} = 0$ for $p < 0$ or $q < 0$. Then for any fixed $p, q$, once $r$ is sufficiently large, $E_{r+1}^{p,q}$ is computed from $(E_{r}^{\cdot,q}, d_{r})$ using the complex

$$
\begin{array}{c}
0 \\
d_{r}^{p,q} \\
E_{r}^{p,q} \\
d_{r+1}^{p+q-r-1} \\
0
\end{array}
$$

and thus we have canonical isomorphisms

$$
E_{r}^{p,q} \cong E_{r+1}^{p,q} \cong E_{r+2}^{p,q} \cong \ldots
$$

We denote this module $E_{\infty}^{p,q}$. The same idea works in other circumstances, for example if the double complex is only nonzero in a finite number of rows — $E_{p,q} = 0$ unless $p_0 < p < p_q$. This will come up for example in the long exact sequence and mapping cone discussion (Exercises 2.7.E and 2.7.F below).

We now describe the first few pages of the spectral sequence explicitly. As stated above, the differential $d_0$ on $E_{0}^{\cdot,\cdot} = E^{\cdot,\cdot}$ is defined to be $d_{-r}$. The rows are complexes:

$$
\begin{array}{c}
\bullet \\
\bullet \\
\bullet
\end{array}
$$

The 0th page $E_0$:

$$
\begin{array}{c}
\bullet \\
\bullet \\
\bullet
\end{array}
$$

$$
\begin{array}{c}
\bullet \\
\bullet
\end{array}
$$

and so $E_1$ is just the table of cohomologies of the rows. You should check that there are now vertical maps $d_{r}^{p,q} : E_{1}^{p,q} \to E_{1}^{p+1,q}$ of the row cohomology groups, induced by $d_{r}$, and that these make the columns into complexes. (This is essentially the fact that a map of complexes induces a map on homology.) We have
“used up the horizontal morphisms”, but “the vertical differentials live on”.

The 1st page $E_1$:

We take cohomology of $d_1$ on $E_1$, giving us a new table, $E_2^{p,q}$. It turns out that there are natural morphisms from each entry to the entry two above and one to the left, and that the composition of these two is 0. (It is a very worthwhile exercise to work out how this natural morphism $d_2$ should be defined. Your argument may be reminiscent of the connecting homomorphism in the Snake Lemma 2.7.5 or in the long exact sequence in cohomology arising from a short exact sequence of complexes, Exercise 2.6.C. This is no coincidence.)

The 2nd page $E_2$:

This is the beginning of a pattern.

Then it is a theorem that there is a filtration of $H^k(E^*)$ by $E_\infty^{p,q}$ where $p+q = k$. (We can’t yet state it as an official Theorem because we haven’t precisely defined the pages and differentials in the spectral sequence.) More precisely, there is a filtration

\[(2.7.2.2) \quad E_\infty^{0,k} \rightarrow E_\infty^{1,k-1} \rightarrow \cdots \rightarrow E_\infty^{k,0} \rightarrow H^k(E^*)\]

where the quotients are displayed above each inclusion. (I always forget which way the quotients are supposed to go, i.e. whether $E^{k,0}$ or $E^{0,k}$ is the subobject. One way of remembering it is by having some idea of how the result is proved.)

We say that the spectral sequence $\rightarrow E^{*,*}$ converges to $H^*(E^*)$. We often say that $\rightarrow E_2^{*,*}$ (or any other page) abuts to $H^*(E^*)$.

Although the filtration gives only partial information about $H^*(E^*)$, sometimes one can find $H^*(E^*)$ precisely. One example is if all $E_\infty^{i,k-i}$ are zero, or if all but one of them are zero (e.g. if $E_\infty^{i,k-i}$ has precisely one non-zero row or column, in which case one says that the spectral sequence collapses at the $r$th step, although we will not use this term). Another example is in the category of vector spaces over a field, in which case we can find the dimension of $H^k(E^*)$. Also, in lucky circumstances, $E_2$ (or some other small page) already equals $E_\infty$.

2.7.A. Exercise: Information from the second page. Show that $H^0(E^*) = E_\infty^{0,0} = E_2^{0,0}$ and

$0 \rightarrow E_2^{0,1} \rightarrow H^1(E^*) \rightarrow E_2^{1,0} \rightarrow E_2^{1,2} \rightarrow H^2(E^*)$. 
2.7.3. The other orientation.

You may have observed that we could as well have done everything in the opposite direction, i.e., reversing the roles of horizontal and vertical morphisms. Then the sequences of arrows giving the spectral sequence would look like this (compare to (2.7.2.1)).

\[
\begin{array}{cccccc}
\bullet & \rightarrow & \bullet & \rightarrow & \bullet & \rightarrow \\
\uparrow & & \uparrow & & \uparrow & \\
\rightarrow & & \rightarrow & & \rightarrow & \\
\end{array}
\]

This spectral sequence is denoted \( \mathcal{E}^{\bullet, \bullet}_p \) (“with the upwards orientation”). Then we would again get pieces of a filtration of \( H^\bullet(\mathcal{E}^\bullet) \) (where we have to be a bit careful with the order with which \( \mathcal{E}^{p,q}_\infty \) corresponds to the subquotients — it in the opposite order to that of \( \mathcal{E}^{p,q}_\infty \) for \( \mathcal{E}^{p,q}_\infty \)). Warning: in general there is no isomorphism between \( \mathcal{E}^{p,q}_\infty \) and \( \mathcal{E}^{p,q}_\infty \).

In fact, this observation that we can start with either the horizontal or vertical maps was our secret goal all along. Both algorithms compute information about the same thing (\( H^\bullet(\mathcal{E}^\bullet) \)), and usually we don’t care about the final answer — we often care about the answer we get in one way, and we get at it by doing the spectral sequence in the other way.

2.7.4. Examples.

We are now ready to see how this is useful. The moral of these examples is the following. In the past, you may have proved various facts involving various sorts of diagrams, by chasing elements around. Now, you will just plug them into a spectral sequence, and let the spectral sequence machinery do your chasing for you.

2.7.5. Example: Proving the Snake Lemma. Consider the diagram

\[
\begin{array}{cccccc}
0 & \rightarrow & D & \rightarrow & E & \rightarrow & F & \rightarrow & 0 \\
\alpha & & & & \beta & & \gamma & & \\
0 & \rightarrow & A & \rightarrow & B & \rightarrow & C & \rightarrow & 0
\end{array}
\]

where the rows are exact in the middle (at B, C, D, G, H, I) and the squares commute. (Normally the Snake Lemma is described with the vertical arrows pointing downwards, but I want to fit this into my spectral sequence conventions.) We wish to show that there is an exact sequence

\[
0 \rightarrow \ker \alpha \rightarrow \ker \beta \rightarrow \ker \gamma \rightarrow \coker \alpha \rightarrow \coker \beta \rightarrow \coker \gamma \rightarrow 0.
\]

We plug this into our spectral sequence machinery. We first compute the cohomology using the rightwards orientation, i.e., using the order (2.7.2.1). Then because the rows are exact, \( E^{1,p,q}_\infty = 0 \), so the spectral sequence has already converged: \( E^{p,q}_\infty = 0 \).
We next compute this “0” in another way, by computing the spectral sequence using the upwards orientation. Then $\tilde{E}_1^{\cdot \cdot}$ (with its differentials) is:

$$
\begin{array}{cccccc}
0 & \longrightarrow & \text{coker } \alpha & \longrightarrow & \text{coker } \beta & \longrightarrow & \text{coker } \gamma & \longrightarrow & 0 \\
0 & \longrightarrow & \text{ker } \alpha & \longrightarrow & \text{ker } \beta & \longrightarrow & \text{ker } \gamma & \longrightarrow & 0.
\end{array}
$$

Then $\tilde{E}_2^{\cdot \cdot}$ is of the form:

$$
\begin{array}{cccccc}
0 & \longrightarrow & 0 & \longrightarrow & ? & \longrightarrow & ? & \longrightarrow & 0 \\
0 & \longrightarrow & ? & \longrightarrow & ? & \longrightarrow & ? & \longrightarrow & 0 \\
0 & \longrightarrow & ? & \longrightarrow & ? & \longrightarrow & ? & \longrightarrow & 0 \\
0 & \longrightarrow & 0 & \longrightarrow & 0 & \longrightarrow & 0 & \longrightarrow & 0.
\end{array}
$$

We see that after $\tilde{E}_2$, all the terms will stabilize except for the double-question-marks — all maps to and from the single question marks are to and from 0-entries. And after $\tilde{E}_3$, even these two double-question-mark terms will stabilize. But in the end our complex must be the 0 complex. This means that in $\tilde{E}_2$, all the entries must be zero, except for the two double-question-marks, and these two must be isomorphic. This means that $0 \rightarrow \text{ker } \alpha \rightarrow \text{ker } \beta \rightarrow \text{ker } \gamma$ and $\text{coker } \alpha \rightarrow \text{coker } \beta \rightarrow \text{coker } \gamma \rightarrow 0$ are both exact (that comes from the vanishing of the single-question-marks), and

$$
\text{coker}(\text{ker } \beta \rightarrow \text{ker } \gamma) \cong \text{ker}(\text{coker } \alpha \rightarrow \text{coker } \beta)
$$

is an isomorphism (that comes from the equality of the double-question-marks). Taken together, we have proved the exactness of (2.7.5.1), and hence the Snake Lemma! (Notice: in the end we didn’t really care about the double complex. We just used it as a prop to prove the snake lemma.)

Spectral sequences make it easy to see how to generalize results further. For example, if $A \rightarrow B$ is no longer assumed to be injective, how would the conclusion change?

2.7.B. **Unimportant Exercise (Grafting exact sequences, a weaker version of the snake lemma)**. Extend the snake lemma as follows. Suppose we have a commuting diagram

$$
\begin{array}{cccccccc}
0 & \longrightarrow & X' & \longrightarrow & Y' & \longrightarrow & Z' & \longrightarrow & A' & \longrightarrow & \cdots \\
\uparrow \alpha & & \uparrow b & & \uparrow c & & & & \uparrow & \\
0 & \longrightarrow & W & \longrightarrow & X & \longrightarrow & Y & \longrightarrow & Z & \longrightarrow & 0.
\end{array}
$$
where the top and bottom rows are exact. Show that the top and bottom rows can be “grafted together” to an exact sequence

\[ \cdots \rightarrow W \rightarrow \ker a \rightarrow \ker b \rightarrow \ker c \rightarrow \text{coker } a \rightarrow \text{coker } b \rightarrow \text{coker } c \rightarrow A' \rightarrow \cdots \]

2.7.6. Example: the Five Lemma. Suppose

\[ (2.7.6.1) \]

\[ F \rightarrow G \rightarrow H \rightarrow I \rightarrow J \]

\[ \alpha \uparrow \beta \uparrow \gamma \uparrow \delta \uparrow \epsilon \]

\[ A \rightarrow B \rightarrow C \rightarrow D \rightarrow E \]

where the rows are exact and the squares commute.

Suppose \( \alpha, \beta, \delta, \epsilon \) are isomorphisms. We will show that \( \gamma \) is an isomorphism.

We first compute the cohomology of the total complex using the rightwards orientation (2.7.2.1). We choose this because we see that we will get lots of zeros.

Then \( \cdots, E^{1**} \) looks like this:

\[ \begin{array}{ccccc}
? & 0 & 0 & 0 & ? \\
\uparrow & \uparrow & \uparrow & \uparrow & \uparrow \\
? & 0 & 0 & 0 & ?
\end{array} \]

Then \( \cdots, E_2 \) looks similar, and the sequence will converge by \( E_2 \), as we will never get any arrows between two non-zero entries in a table thereafter. We can’t conclude that the cohomology of the total complex vanishes, but we can note that it vanishes in all but four degrees — and most important, it vanishes in the two degrees corresponding to the entries \( C \) and \( H \) (the source and target of \( \gamma \)).

We next compute this using the upwards orientation (2.7.3.1). Then \( \uparrow E_1 \) looks like this:

\[ \begin{array}{ccccc}
0 & \rightarrow & 0 & \rightarrow & ? \\
\rightarrow & \downarrow & \downarrow & \downarrow & \downarrow \\
0 & \rightarrow & 0 & \rightarrow & 0
\end{array} \]

and the spectral sequence converges at this step. We wish to show that those two question marks are zero. But they are precisely the cohomology groups of the total complex that we just showed were zero — so we are done!

The best way to become comfortable with this sort of argument is to try it out yourself several times, and realize that it really is easy. So you should do the following exercises!

2.7.C. Exercise: The Subtle Five Lemma. By looking at the spectral sequence proof of the Five Lemma above, prove a subtler version of the Five Lemma, where one of the isomorphisms can instead just be required to be an injection, and another can instead just be required to be a surjection. (I am deliberately not telling you which ones, so you can see how the spectral sequence is telling you how to improve the result.)
2.7.D. Exercise. If $\beta$ and $\delta$ (in $(2.7.6.1)$) are injective, and $\alpha$ is surjective, show that $\gamma$ is injective. Give the dual statement (whose proof is of course essentially the same).

2.7.E. Exercise. Use spectral sequences to show that a short exact sequence of complexes gives a long exact sequence in cohomology (Exercise 2.6.C).

2.7.F. Exercise (the mapping cone). Suppose $\mu : A^\bullet \to B^\bullet$ is a morphism of complexes. Suppose $C^\bullet$ is the single complex associated to the double complex $A^\bullet \to B^\bullet$. ($C^\bullet$ is called the mapping cone of $\mu$.) Show that there is a long exact sequence of complexes:

$$\cdots \to H^{i-1}(C^\bullet) \to H^i(A^\bullet) \to H^i(B^\bullet) \to H^i(C^\bullet) \to H^{i+1}(A^\bullet) \to \cdots.$$  

(There is a slight notational ambiguity here; depending on how you index your double complex, your long exact sequence might look slightly different.) In particular, we will use the fact that $\mu$ induces an isomorphism on cohomology if and only if the mapping cone is exact. (We won’t use it until the proof of Theorem 20.2.4.)

The Grothendieck (or composition of functor) spectral sequence (Exercise 23.3.D) will be an important example of a spectral sequence that specializes in a number of useful ways.

You are now ready to go out into the world and use spectral sequences to your heart’s content!

2.7.7. ** Complete definition of the spectral sequence, and proof.**

You should most definitely not read this section any time soon after reading the introduction to spectral sequences above. Instead, flip quickly through it to convince yourself that nothing fancy is involved.

We consider the rightwards orientation. The upwards orientation is of course a trivial variation of this.

2.7.8. Goals. We wish to describe the pages and differentials of the spectral sequence explicitly, and prove that they behave the way we said they did. More precisely, we wish to:

(a) describe $E_r^{p,q}$,
(b) verify that $H^k(E^\bullet)$ is filtered by $E_r^{p,k-p}$ as in $(2.7.2.2)$,
(c) describe $d_r$ and verify that $d_r^2 = 0$, and
(d) verify that $E_r^{p,q+r}$ is given by cohomology using $d_r$.

Before tacking these goals, you can impress your friends by giving this short description of the pages and differentials of the spectral sequence. We say that an element of $E^{\bullet,\bullet}$ is a $(p, q)$-strip if it is an element of $\oplus_{i\geq 0} E^{p+i,q-i}$ (see Fig. 2.1). Its non-zero entries lie on a semi-infinite antidiagonal starting with position $(p, q)$. We say that the $(p, q)$-entry (the projection to $E^{p,q}$) is the leading term of the $(p, q)$-strip. Let $S_r^{p,q} \subset E^{\bullet,\bullet}$ be the submodule of all the $(p, q)$-strips. Clearly $S_r^{p,q} \subset E^{p,q}$, and $S_r^{0,k} = E_k$.

Note that the differential $d = d_1 + d_\infty$ sends a $(p, q)$-strip $x$ to a $(p, q + 1)$-strip $dx$. If $dx$ is furthermore a $(p + r, q + r + 1)$-strip ($r \in \mathbb{Z}_{\geq 0}$), we say that $x$ is an $r$-closed $(p, q)$-strip. We denote the set of such $S_r^{p,q}$ (so for example $S_0^{p,q} = S^{p,q}$,
and $S^0_{0,k} = E^k$). An element of $S^0_{p,q}$ may be depicted as:

![Diagram]

**Figure 2.1.** A $(p, q)$-strip (in $S^{p,q} \subset E^{p+q}$). Clearly $S^0_{0,k} = E^k$.

### 2.7.9 Preliminary definition of $E^{p,q}$. We are now ready to give a first definition of $E^{p,q}$, which by construction should be a subquotient of $E^{p,q} = E^0_{p,q}$. We describe it as such by describing two submodules $Y_{r}^{p,q} \subset X^{p,q} \subset E^{p,q}$, and defining $E^{p,q}_r = X^{p,q}_r / Y_{r}^{p,q}$. Let $X^{p,q}_r$ be those elements of $E^{p,q}$ that are the leading terms of $r$-closed $(p, q)$-strips. Note that by definition, $d$ sends $(r-1)$-closed $S^{p-(r-1), q+(r-1)-1}$-strips to $(p, q)$-strips. Let $Y_{r}^{p,q}$ be the leading $(p, q)$-terms of the differential $d$ of $(r-1)$-closed $(p-(r-1), q+(r-1)-1)$-strips (where the differential is considered as a $(p, q)$-strip).

We next give the definition of the differential $d_r$ of such an element $x \in X^{p,q}_r$. We take any $r$-closed $(p, q)$-strip with leading term $x$. Its differential $d$ is a $(p+r, q-r+1)$-strip, and we take its leading term. The choice of the $r$-closed $(p, q)$-strip means that this is not a well-defined element of $E^{p,q}$. But it is well-defined modulo the $(r-1)$-closed $(p+1, r+1)$-strips, and hence gives a map $E^{p,q}_r \to E^{p+r,q-r+1}$. This definition is fairly short, but not much fun to work with, so we will forget it, and instead dive into a snakes’ nest of subscripts and superscripts.
We have completed Goal 2.7.8(a).

2.7.G. Exercise. Verify the following.

(a) $S^{p,q} = S^{p+1,q-1} \oplus E_r^{p,q}$.
(b) Any closed $(p, q)$-strip is $r$-closed for all $r$. Any element $x$ of $S^{p,q} = S_0^{p,q}$ that is a cycle (i.e. $dx = 0$) is automatically in $S_r^{p,q}$ for all $r$. For example, this holds when $x$ is a boundary (i.e. of the form $dy$).
(c) Show that for fixed $p, q$,
\[ S_0^{p,q} \supset S_1^{p,q} \supset \cdots \supset S_r^{p,q} \supset \cdots \]
stabilizes for $r \gg 0$ (i.e. $S_r^{p,q} = S_{r+1}^{p,q} = \cdots$). Denote the stabilized module $S_{\infty}^{p,q}$. Show $S_{\infty}^{p,q}$ is the set of closed $(p, q)$-strips (those $(p, q)$-strips annihilated by $d$, i.e. the cycles). In particular, $S_{\infty}^{p+k}$ is the set of cycles in $E^k$.

2.7.10. Defining $E_r^{p,q}$.

Define $X_r^{p,q} := S_r^{p,q}/S_{r-1}^{p+1,q-1}$ and $Y_r := dS_{r-1}^{p-(r-1),q+(r-1)-1}/S_{r-1}^{p+1,q-1}$. Then $Y_r^{p,q} \subset X_r^{p,q}$ by Exercise 2.7.G(b). We define

\[ E_r^{p,q} = \frac{X_r^{p,q}}{Y_r^{p,q}} = \frac{S_r^{p,q}}{dS_{r-1}^{p-(r-1),q+(r-1)-1} + S_{r-1}^{p+1,q-1}} \]

We have completed Goal 2.7.8(a).

You are welcome to verify that these definitions of $X_r^{p,q}$ and $Y_r^{p,q}$ and hence $E_r^{p,q}$ agree with the earlier ones of 2.7.9 (and in particular $X_r^{p,q}$ and $Y_r^{p,q}$ are both submodules of $E_r^{p,q}$), but we won’t need this fact.

2.7.H. Exercise: $E_{\infty}^{p-k,p}$ gives subquotients of $H^k(E^*)$. By Exercise 2.7.G(c), $E_r^{p,q}$ stabilizes as $r \to \infty$. For $r \gg 0$, interpret $S_r^{p,q}/dS_{r-1}^{p-(r-1),q+(r-1)-1}$ as the cycles in $S_{\infty}^{p,q} \subset E_{\infty}^{p,q}$ modulo those boundary elements of $dE^{p,q}$ contained in $S_{\infty}^{p,q}$. Finally, show that $H^k(E^*)$ is indeed filtered as described in 2.7.2.

We have completed Goal 2.7.8(b).

2.7.11. Definition of $d_r$.

We shall see that the map $d_r : E_r^{p,q} \to E_{r+r}^{p+r,q-r+1}$ is just induced by our differential $d$. Notice that $d$ sends $r$-closed $(p, q)$-strips $S_r^{p,q}$ to $(p+r,q-r+1)$-strips $S_{r+r}^{p+r,q-r+1}$, by the definition “$r$-closed”. By Exercise 2.7.G(b), the image lies in $E_{r+r}^{p+r,q-r+1}$.

2.7.I. Exercise. Verify that $d$ sends
\[ dS_{r-1}^{p-(r-1),q+(r-1)-1} + S_{r-1}^{p+1,q-1} \to dS_{r-1}^{(p+r)-(r-1),(q-r+1)+(r-1)-1} + S_{r-1}^{p+r+1,(q-r+1)-1}. \]
(The first term on the left goes to 0 from $d^2 = 0$, and the second term on the left goes to the first term on the right.)
Thus we may define

\[
d_r : E^r_{p,q} = \frac{S_{r,q}^p}{dS_{r-1}^{p-(r-1),q+(r-1)-1} + S_{r-1}^{p+1,q-1}} \rightarrow \frac{S_{r+1,q-r+1}^p}{dS_{r-1}^{p+1,q-1} + S_{r-1}^{p+1,q-r+1}}
\]

and clearly \(d_r^2 = 0\) (as we may interpret it as taking an element of \(S^p_{r,q}\) and applying \(d\) twice).

We have accomplished Goal 2.7.8(c).

2.7.12. Verifying that the cohomology of \(d_r\) at \(E^r_{p,q}\) is \(E^r_{p,q+1}\). We are left with the unpleasant job of verifying that the cohomology of

\[
(2.7.12.1) \quad \frac{S_{r,q+1}^{p-r} + S_{r-1}^{p+1,q-1}}{dS_{r-1}^{p+1,q-1} + S_{r-1}^{p+1,q-r+1}} \rightarrow \frac{S_{r,q}^p}{dS_{r-1}^{p-(r-1),q+(r-1)-1} + S_{r-1}^{p+1,q-1}}
\]

is naturally identified with

\[
\frac{S_{r+1}^{p-r} + S_{r-1}^{p+1,q-1}}{dS_{r-1}^{p+1,q-1} + S_{r-1}^{p+1,q-r+1}}
\]

and this will conclude our final Goal 2.7.8(d).

We begin by understanding the kernel of the right map of (2.7.12.1). Suppose \(a \in S^p_{r,q}\) is mapped to \(0\). This means that \(da = db + c\), where \(b \in S^{p+1,q-1}_{r-1}\). If \(u = a - b\), then \(u \in S^p_{r-1,q}\), while \(du = c \in S^{p+1,q-1}_{r-1}\). From which \(u\) is \(r\)-closed, i.e. \(u \in S^p_{r,q}\). Hence \(a = b + u + x\), where \(dx = 0\), from which \(a - x = b + c \in S^{p+1,q-1}_{r-1} + S^{p+1,q}_{r+1}\). However, \(x \in S^1_{p,q}\), so \(x \in S^{p+1}_{r-1}\) by Exercise 2.7.G(b). Thus \(a \in S^{p+1,q-1}_{r-1} + S^1_{p,q}\). Conversely, any \(a \in S^{p+1,q-1}_{r-1} + S^1_{p,q}\) satisfies

\[
da \in dS^{p+1,q-1}_{r-1} + dS^1_{r-1} + S^{p+1,q-1}_{r-1} + S^{p+1,q-1}_{r-1}
\]

(using \(dS^1_{r-1} \subset S^1_{p+1,q-1}\) and Exercise 2.7.G(b)) so any such \(a\) is indeed in the kernel of

\[
S^p_{r,q} \rightarrow \frac{S_{r,q}^{p+1,q-r+1}}{dS_{r-1}^{p+1,q-1} + S_{r-1}^{p+1,q-r+1}}
\]

Hence the kernel of the right map of (2.7.12.1) is

\[
\ker = \frac{S_{r-1}^{p+1,q-1} + S^{p,q}_{r+1}}{dS_{r-1}^{p+1,q+r-2} + S_{r-1}^{p+1,q-1}}
\]

Next, the image of the left map of (2.7.12.1) is immediately

\[
\text{im} = \frac{dS_{r-1}^{p-r,q+r+1} + dS_{r-1}^{q+1,q-1} + S_{r-1}^{p+1,q-1}}{dS_{r-1}^{p-r,q+r-1} + dS_{r-1}^{q+1,q-1} + S_{r-1}^{p+1,q-1}} = \frac{dS_{r-1}^{p-r,q+r+1} + S_{r-1}^{p+1,q-1}}{dS_{r-1}^{p-r,q+r-1} + dS_{r-1}^{q+1,q-1} + S_{r-1}^{p+1,q-1}}
\]

(as \(S_{r-1}^{p-r,q+r+1}\) contains \(S_{r-1}^{p-r+1,q+r-1}\)).
Thus the cohomology of (2.7.12.1) is
\[
\frac{\ker}{\text{im}} = \frac{S_{r-1}^{p+1,q-1} + S_{r+1}^{p, q}}{dS^p_{r, q+r-1} + S_{r-1}^{p+1,q-1}} = \frac{S_{r+1}^{p,q}}{S_{r+1}^{p,q} \cap (dS^p_{r,q+r-1} + S_{r-1}^{p+1,q-1})}
\]
where the equality on the right uses the fact that $dS^p_{r,q+r-1} \subset S_{r+1}^{p,q}$ and an isomorphism theorem. We thus must show
\[
S_{r+1}^{p,q} \cap (dS^p_{r,q+r-1} + S_{r-1}^{p+1,q-1}) = dS^p_{r,q+r-1} + S_{r+1}^{p+1,q-1}.
\]
However,
\[
S_{r+1}^{p,q} \cap (dS^p_{r,q+r-1} + S_{r-1}^{p+1,q-1}) = dS^p_{r,q+r-1} + S_{r+1}^{p,q} \cap S_{r-1}^{p+1,q-1}
\]
and $S_{r+1}^{p,q} \cap S_{r-1}^{p+1,q-1}$ consists of $(p, q)$-strips whose differential vanishes up to row $p + r$, from which $S_{r+1}^{p,q} \cap S_{r-1}^{p+1,q-1} = S_r^{p,q}$ as desired.

This completes the explanation of how spectral sequences work for a first-quadrant double complex. The argument applies without significant change to more general situations, including filtered complexes.
CHAPTER 3

Sheaves

It is perhaps surprising that geometric spaces are often best understood in terms of (nice) functions on them. For example, a differentiable manifold that is a subset of $\mathbb{R}^n$ can be studied in terms of its differentiable functions. Because “geometric spaces” can have few (everywhere-defined) functions, a more precise version of this insight is that the structure of the space can be well understood by considering all functions on all open subsets of the space. This information is encoded in something called a sheaf. Sheaves were introduced by Leray in the 1940’s, and Serre introduced them to algebraic geometry. (The reason for the name will be somewhat explained in Remark 3.4.3.) We will define sheaves and describe useful facts about them. We will begin with a motivating example to convince you that the notion is not so foreign.

One reason sheaves are slippery to work with is that they keep track of a huge amount of information, and there are some subtle local-to-global issues. There are also three different ways of getting a hold of them.

- in terms of open sets (the definition §3.2) — intuitive but in some ways the least helpful
- in terms of stalks (see §3.4.1)
- in terms of a base of a topology (§3.7).

Knowing which to use requires experience, so it is essential to do a number of exercises on different aspects of sheaves in order to truly understand the concept.

3.1 Motivating example: The sheaf of differentiable functions.

Consider differentiable functions on the topological space $X = \mathbb{R}^n$ (or more generally on a smooth manifold $X$). The sheaf of differentiable functions on $X$ is the data of all differentiable functions on all open subsets on $X$. We will see how to manage this data, and observe some of its properties. On each open set $U \subset X$, we have a ring of differentiable functions. We denote this ring of functions $\mathcal{O}(U)$.

Given a differentiable function on an open set, you can restrict it to a smaller open set, obtaining a differentiable function there. In other words, if $U \subset V$ is an inclusion of open sets, we have a “restriction map” $\text{res}_{V,U} : \mathcal{O}(V) \to \mathcal{O}(U)$.

Take a differentiable function on a big open set, and restrict it to a medium open set, and then restrict that to a small open set. The result is the same as if you restrict the differentiable function on the big open set directly to the small open set.
In other words, if \( U \rightarrow V \rightarrow W \), then the following diagram commutes:

\[
\begin{array}{ccc}
\mathcal{O}(W) & \xrightarrow{\text{res}_{W,V}} & \mathcal{O}(V) \\
\text{res}_{W,U} & & \text{res}_{V,U} \\
\mathcal{O}(U) & & \\
\end{array}
\]

Next take two differentiable functions \( f_1 \) and \( f_2 \) on a big open set \( U \), and an open cover of \( U \) by some \( \{U_i\} \). Suppose that \( f_1 \) and \( f_2 \) agree on each of these \( U_i \). Then they must have been the same function to begin with. In other words, if \( \{U_i\}_{i \in I} \) is a cover of \( U \), and \( f_1, f_2 \in \mathcal{O}(U) \), and \( \text{res}_{U_i U_j} f_1 = \text{res}_{U_i U_j} f_2 \), then \( f_1 = f_2 \). Thus we can identify functions on an open set by looking at them on a covering by small open sets.

Finally, given the same \( U \) and cover \( \{U_i\} \), take a differentiable function on each of the \( U_i \) — a function \( f_1 \) on \( U_1 \), a function \( f_2 \) on \( U_2 \), and so on — and they agree on the pairwise overlaps. Then they can be “glued together” to make one differentiable function on all of \( U \). In other words, given \( f_i \in \mathcal{O}(U_i) \) for all \( i \), such that \( \text{res}_{U_i U_j} f_i = \text{res}_{U_i U_j} f_j \) for all \( i \) and \( j \), then there is some \( f \in \mathcal{O}(U) \) such that \( \text{res}_{U_i} f = f_i \) for all \( i \).

The entire example above would have worked just as well with continuous function, or smooth functions, or just plain functions. Thus all of these classes of “nice” functions share some common properties. We will soon formalize these properties in the notion of a sheaf.

3.1.1. The germ of a differentiable function. Before we do, we first give another definition, that of the germ of a differentiable function at a point \( p \in X \). Intuitively, it is a “shred” of a differentiable function at \( p \). Germs are objects of the form \( \{ (f, \text{open } U) : p \in U, f \in \mathcal{O}(U) \} \) modulo the relation that \( (f, U) \sim (g, V) \) if there is some open set \( W \subset U, V \) containing \( p \) where \( f|_W = g|_W \) (i.e., \( \text{res}_{U \cap V} f = \text{res}_{U \cap V} g \)). In other words, two functions that are the same in a neighborhood of \( p \) (but may differ elsewhere) have the same germ. We call this set of germs the stalk at \( p \), and denote it \( \mathcal{O}_p \). Notice that the stalk is a ring: you can add two germs, and get another germ: if you have a function \( f \) defined on \( U \), and a function \( g \) defined on \( V \), then \( f + g \) is defined on \( U \cap V \). Moreover, \( f + g \) is well-defined: if \( f' \) has the same germ as \( f \), meaning that there is some open set \( W \) containing \( p \) on which they agree, and \( g' \) has the same germ as \( g \), meaning they agree on some open \( W' \) containing \( p \), then \( f' + g' \) is the same function as \( f + g \) on \( U \cap V \cap W \cap W' \).

Notice also that if \( p \in U \), you get a map \( \mathcal{O}(U) \to \mathcal{O}_p \). Experts may already see that we are talking about germs as colimits.

We can see that \( \mathcal{O}_p \) is a local ring as follows. Consider those germs vanishing at \( p \), which we denote \( m_p \subset \mathcal{O}_p \). They certainly form an ideal: \( m_p \) is closed under addition, and when you multiply something vanishing at \( p \) by any other function, the result also vanishes at \( p \). We check that this ideal is maximal by showing that the quotient map is a field:

\[
\begin{array}{c}
0 \to m_p := \text{ideal of germs vanishing at } p \to \mathcal{O}_p \xrightarrow{f \mapsto f(p)} \mathbb{R} \to 0
\end{array}
\]

3.1.A. Exercise. Show that this is the only maximal ideal of \( \mathcal{O}_p \). (Hint: show that every element of \( \mathcal{O}_p \setminus m \) is invertible.)
Note that we can interpret the value of a function at a point, or the value of a germ at a point, as an element of the local ring modulo the maximal ideal. (We will see that this doesn’t work for more general sheaves, but does work for things behaving like sheaves of functions. This will be formalized in the notion of a locally ringed space, which we will see, briefly, in §7.3.

3.1.2. Aside. Notice that \( \mathfrak{m}/\mathfrak{m}^2 \) is a module over \( \mathcal{O}_p/\mathfrak{m} \cong \mathbb{R} \), i.e. it is a real vector space. It turns out to be naturally (whatever that means) the cotangent space to the manifold at \( p \). This insight will prove handy later, when we define tangent and cotangent spaces of schemes.

3.1.B. Exercise for those with differential geometric background. Prove this.

3.2 Definition of sheaf and presheaf

We now formalize these notions, by defining presheaves and sheaves. Presheaves are simpler to define, and notions such as kernel and cokernel are straightforward. Sheaves are more complicated to define, and some notions such as cokernel require more thought. But sheaves are more useful because they are in some vague sense more geometric; you can get information about a sheaf locally.

3.2.1. Definition of sheaf and presheaf on a topological space \( X \).

To be concrete, we will define sheaves of sets. However, in the definition the category \( \text{Sets} \) can be replaced by any category, and other important examples are abelian groups \( \text{Ab} \), \( k \)-vector spaces \( \text{Vec}_k \), rings \( \text{Rings} \), modules over a ring \( \text{Mod}_A \), and more. (You may have to think more when dealing with a category of objects that aren’t “sets with additional structure”, but there aren’t any new complications. In any case, this won’t be relevant for us.) Sheaves (and presheaves) are often written in calligraphic font. The fact that \( \mathcal{F} \) is a sheaf on a topological space \( X \) is often written as

\[
\mathcal{F}
\]

\[
\downarrow
\]

\[
X
\]

3.2.2. Definition: Presheaf. A presheaf \( \mathcal{F} \) on a topological space \( X \) is the following data.

- To each open set \( U \subset X \), we have a set \( \mathcal{F}(U) \) (e.g. the set of differentiable functions in our motivating example). (Notational warning: Several notations are in use, for various good reasons: \( \mathcal{F}(U) = \Gamma(U, \mathcal{F}) = H^0(U, \mathcal{F}) \). We will use them all.) The elements of \( \mathcal{F}(U) \) are called sections of \( \mathcal{F} \) over \( U \).
- For each inclusion \( U \hookrightarrow V \) of open sets, we have a restriction map \( \text{res}_{V,U} : \mathcal{F}(V) \to \mathcal{F}(U) \) (just as we did for differentiable functions).

The data is required to satisfy the following two conditions.

- The map \( \text{res}_{U,U} \) is the identity: \( \text{res}_{U,U} = \text{id}_{\mathcal{F}(U)} \).
• If $U \hookrightarrow V \hookrightarrow W$ are inclusions of open sets, then the restriction maps commute, i.e.

$$
\begin{array}{ccc}
\mathcal{F}(W) & \xrightarrow{\text{res}_{W,V}} & \mathcal{F}(V) \\
\downarrow & & \downarrow \\
\mathcal{F}(U) & \xrightarrow{\text{res}_{W,U}, \text{res}_{V,U}} &
\end{array}
$$

commutes.

3.2.A. EXERCISE FOR CATEGORY-LOVERS: “A PRESHAEF IS THE SAME AS A CONTRAVARIANT FUNCTOR”. Given any topological space $X$, we have a “category of open sets” (Example 2.2.9), where the objects are the open sets and the morphisms are inclusions. Verify that the data of a presheaf is precisely the data of a contravariant functor from the category of open sets of $X$ to the category of sets. (This interpretation is surprisingly useful.)

3.2.3. Definition: Stalks and germs. We define the stalk of a presheaf at a point in two equivalent ways. One will be hands-on, and the other will be as a colimit.

3.2.4. Define the stalk of a presheaf $\mathcal{F}$ at a point $p$ to be the set of germs of $\mathcal{F}$ at $p$, denoted $\mathcal{F}_p$, as in the example of §3.1.1. Germs correspond to sections over some open set containing $p$, and two of these sections are considered the same if they agree on some smaller open set. More precisely: the stalk is

$$
\{ (f, \text{open } U) : p \in U, f \in \mathcal{F}(U) \}
$$

modulo the relation that $(f, U) \sim (g, V)$ if there is some open set $W \subset U, V$ where $\text{res}_{U,W} f = \text{res}_{V,W} g$.

3.2.5. A useful (and better) equivalent definition of a stalk is as a colimit of all $\mathcal{F}(U)$ over all open sets $U$ containing $p$:

$$
\mathcal{F}_p = \lim_{\rightarrow} \mathcal{F}(U).
$$

The index category is a directed set (given any two such open sets, there is a third such set contained in both), so these two definitions are the same by Exercise 2.4.C. Hence we can define stalks for sheaves of sets, groups, rings, and other things for which colimits exist for directed sets.

If $p \in U$, and $f \in \mathcal{F}(U)$, then the image of $f$ in $\mathcal{F}_p$ is called the germ of $f$ at $p$. (Warning: unlike the example of §3.1.1 in general, the value of a section at a point doesn’t make sense.)

3.2.6. Definition: Sheaf. A presheaf is a sheaf if it satisfies two more axioms. Notice that these axioms use the additional information of when some open sets cover another.

Identity axiom. If $\{U_i\}_{i \in I}$ is an open cover of $U$, and $f_1, f_2 \in \mathcal{F}(U)$, and $\text{res}_{U_i,U} f_1 = \text{res}_{U_i,U} f_2$ for all $i$, then $f_1 = f_2$.

(A presheaf satisfying the identity axiom is called a separated presheaf, but we will not use that notation in any essential way.)

Gluability axiom. If $\{U_i\}_{i \in I}$ is a open cover of $U$, then given $f_i \in \mathcal{F}(U_i)$ for all $i$, such that $\text{res}_{U_i, U \cap U_i} f_i = \text{res}_{U_i, U \cap U_i} f_j$ for all $i, j$, then there is some $f \in \mathcal{F}(U)$ such that $\text{res}_{U_i, U} f = f_i$ for all $i$. 

In mathematics, definitions often come paired: “at most one” and “at least one”. In this case, identity means there is at most one way to glue, and gluability means that there is at least one way to glue.

(For experts and scholars of the empty set only: an additional axiom sometimes included is that $F(\emptyset)$ is a one-element set, and in general, for a sheaf with values in a category, $F(\emptyset)$ is required to be the final object in the category. This actually follows from the above definitions, assuming that the empty product is appropriately defined as the final object.)

Example. If $U$ and $V$ are disjoint, then $\mathcal{F}(U \cup V) = \mathcal{F}(U) \times \mathcal{F}(V)$. Here we use the fact that $F(\emptyset)$ is the final object.

The stalk of a sheaf at a point is just its stalk as a presheaf — the same definition applies — and similarly for the germs of a section of a sheaf.

3.2.B. UNIMPORTANT EXERCISE: PRESHIFVES THAT ARE NOT SHEAVES. Show that the following are presheaves on $\mathbb{C}$ (with the classical topology), but not sheaves:

(a) bounded functions, (b) holomorphic functions admitting a holomorphic square root.

Both of the presheaves in the previous Exercise satisfy the identity axiom. A “natural” example failing even the identity axiom will be given in Remark 3.7.2.

We now make a couple of points intended only for category-lovers.

3.2.7. Interpretation in terms of the equalizer exact sequence. The two axioms for a presheaf to be a sheaf can be interpreted as “exactness” of the “equalizer exact sequence”:

$\mathcal{F}(U) \rightarrow \prod \mathcal{F}(U_i) \rightarrow \prod \mathcal{F}(U_i \cap U_j)$.

Identity is exactness at $\mathcal{F}(U)$, and gluability is exactness at $\prod \mathcal{F}(U_i)$. I won’t make this precise, or even explain what the double right arrow means. (What is an exact sequence of sets?!)

But you may be able to figure it out from the context.

3.2.C. EXERCISE. The gluability axiom may be interpreted as saying that $\mathcal{F}(\bigcup_{i \in I} U_i)$ is a certain limit. What is that limit?

We now give a number of examples of sheaves.

3.2.D. EXERCISE. (a) Verify that the examples of §3.1 are indeed sheaves (of differentiable functions, or continuous functions, or smooth functions, or functions on a manifold or $\mathbb{R}^n$).

(b) Show that real-valued continuous functions on (open sets of) a topological space $X$ form a sheaf.

3.2.8. Important Example: Restriction of a sheaf. Suppose $\mathcal{F}$ is a sheaf on $X$, and $U \subset X$ is an open set. Define the restriction of $\mathcal{F}$ to $U$, denoted $\mathcal{F}|_U$, to be the collection $\mathcal{F}|_U(V) = \mathcal{F}(V)$ for all $V \subset U$. Clearly this is a sheaf on $U$. (Unimportant but fun fact: §3.6 will tell us how to restrict sheaves to arbitrary subsets.)

3.2.9. Important Example: skyscraper sheaf. Suppose $X$ is a topological space, with $p \in X$, and $S$ is a set. Let $i_p : p \to X$ be the inclusion. Then $i_p^*, S$ defined by

$$i_p^*, S(U) = \begin{cases} S & \text{if } p \in U, \\ \{e\} & \text{if } p \not\in U \end{cases}$$
forms a sheaf. Here \((e)\) is any one-element set. (Check this if it isn’t clear to you.) This is called a skewscrapersheaf, because the informal picture of it looks like a skyscraper at \(p\). There is an analogous definition for sheaves of abelian groups, except \(i_p)_*(S)(U) = \{0\} \) if \(p \not\in U\); and for sheaves with values in a category more generically, \(i_p)_*(S)(U)\) should be a final object.

### 3.2.10. Constant presheaves and constant sheaves.

Let \(X\) be a topological space, and \(S\) a set. Define \(\mathcal{S}^\text{pre}(U) = S\) for all open sets \(U\). You will readily verify that \(\mathcal{S}^\text{pre}\) forms a presheaf (with restriction maps the identity). This is called the constant presheaf associated to \(S\). This isn’t (in general) a sheaf. (It may be distracting to say why. Lovers of the empty set will note that the sheaf axioms force the sections over the empty set to be the final object in the category, i.e. a one-element set. But even if we patch the definition by setting \(\mathcal{S}^\text{pre}(\emptyset) = \{e\}\), if \(S\) has more than one element, and \(X\) is the two-point space with the discrete topology, you can check that \(\mathcal{S}^\text{pre}\) fails glubility.)

### 3.2.E. Exercise (Constant sheaves).

Now let \(\mathcal{F}(U)\) be the maps to \(S\) that are locally constant, i.e. for any point \(x\) in \(U\), there is a neighborhood of \(x\) where the function is constant. Show that this is a sheaf. (A better description is this: endow \(S\) with the discrete topology, and let \(\mathcal{F}(U)\) be the continuous maps \(U \to S\).) This is called the constant sheaf (associated to \(S\)); do not confuse it with the constant presheaf. We denote this sheaf \(\mathcal{S}\).

### 3.2.F. Exercise (“Morphisms glue”).

Suppose \(Y\) is a topological space. Show that “continuous maps to \(Y\)” form a sheaf of sets on \(X\). More precisely, to each open set \(U\) of \(X\), we associate the set of continuous maps of \(U\) to \(Y\). Show that this forms a sheaf. (Exercise 3.2.D(b), with \(Y = \mathbb{R}\), and Exercise 3.2.E, with \(Y = S\) with the discrete topology, are both special cases.)

### 3.2.G. Exercise.

This is a fancier version of the previous exercise.

(a) (sheaf of sections of a map) Suppose we are given a continuous map \(f: Y \to X\). Show that “sections of \(f\)” form a sheaf. More precisely, to each open set \(U\) of \(X\), associate the set of continuous maps \(s: U \to Y\) such that \(f \circ s = \text{id}_U\). Show that this forms a sheaf. (For those who have heard of vector bundles, these are a good example.) This is motivation for the phrase “section of a sheaf”.

(b) (This exercise is for those who know what a topological group is. If you don’t know what a topological group is, you might be able to guess.) Suppose that \(Y\) is a topological group. Show that continuous maps to \(Y\) form a sheaf of groups. (Example 3.2.D(b), with \(Y = \mathbb{R}\), is a special case.)

### 3.2.11. * The espace étalé of a (pre)sheaf.

Depending on your background, you may prefer the following perspective on sheaves, which we will not discuss further. Suppose \(\mathcal{F}\) is a presheaf (e.g. a sheaf) on a topological space \(X\). Construct a topological space \(Y\) along with a continuous map \(\pi: Y \to X\) as follows: as a set, \(Y\) is the disjoint union of all the stalks of \(X\). This also describes a natural set map \(\pi: Y \to X\). We topologize \(Y\) as follows. Each section \(s\) of \(\mathcal{F}\) over an open set \(U\) determines a subset \(\{x, s_x\} : x \in U\) of \(Y\). The topology on \(Y\) is the weakest topology such that these subsets are open. (These subsets form a base of the topology. For each \(y \in Y\), there is a neighborhood \(V\) of \(y\) and a neighborhood \(U\) of \(X\) such that \(\pi|_V\) is a homeomorphism from \(V\) to \(U\). Do you see why these facts are true?)
space is called the espace étalé of $F$. The reader may wish to show that (a) if $F$ is a sheaf, then the sheaf of sections of $Y \to X$ (see the previous exercise 3.2.G(a)) can be naturally identified with the sheaf $F$ itself. (b) Moreover, if $F$ is a presheaf, the sheaf of sections of $Y \to X$ is the sheafification of $F$, to be defined in Definition 3.4.5 (see Remark 3.4.7). Example 3.2.E may be interpreted as an example of this construction.

3.2.H. IMPORTANT EXERCISE: THE PUSHFORWARD SHEAF OR DIRECT IMAGE SHEAF. Suppose $f: X \to Y$ is a continuous map, and $F$ is a sheaf on $X$. Then define $f_*F$ by $f_*F(V) = F(f^{-1}(V))$, where $V$ is an open subset of $Y$. Show that $f_*F$ is a sheaf. This is called a direct image sheaf or pushforward sheaf. More precisely, $f_*F$ is called the pushforward of $F$ by $f$.

As the notation suggests, the skyscraper sheaf (Example 3.2.9) can be interpreted as the pushforward of the constant sheaf $\mathbb{S}$ on a one-point space $p$ under the inclusion morphism $i: \{p\} \to X$.

Once we realize that sheaves form a category, we will see that the pushforward is a functor from sheaves on $X$ to sheaves on $Y$ (Exercise 3.3.A).

3.2.I. EXERCISE (PUSHFORWARD INDUCES MAPS OF STALKS). Suppose $f: X \to Y$ is a continuous map, and $F$ is a sheaf of sets (or rings or $A$-modules) on $X$. If $f(x) = y$, describe the natural morphism of stalks $(f_*F)_y \to F_x$. (You can use the explicit definition of stalk using representatives, §3.2.4, or the universal property, §3.2.5. If you prefer one way, you should try the other.) Once we define the category of sheaves of sets on a topological space in §3.3.1 you will see that your construction will make the following diagram commute:

\[
\begin{array}{ccc}
\text{Sets}_X & \xrightarrow{f_*} & \text{Sets}_Y \\
\downarrow & & \downarrow \\
\text{Sets} & \xrightarrow{} & \text{Sets}
\end{array}
\]

3.2.12. Important Example: Ringed spaces, and $\mathcal{O}_X$-modules. Suppose $\mathcal{O}_X$ is a sheaf of rings on a topological space $X$ (i.e. a sheaf on $X$ with values in the category of Rings). Then $(X, \mathcal{O}_X)$ is called a ringed space. The sheaf of rings is often denoted by $\mathcal{O}_X$, pronounced “oh-of-X”. This sheaf is called the structure sheaf of the ringed space. We now define the notion of an $\mathcal{O}_X$-module. The notion is analogous to one we have seen before: just as we have modules over a ring, we have $\mathcal{O}_X$-modules over the structure sheaf (of rings) $\mathcal{O}_X$.

There is only one possible definition that could go with this name. An $\mathcal{O}_X$-module is a sheaf of abelian groups $F$ with the following additional structure. For each $U$, $F(U)$ is an $\mathcal{O}_X(U)$-module. Furthermore, this structure should behave well with respect to restriction maps: if $U \subset V$, then

\[
\begin{array}{ccc}
\mathcal{O}_X(U) \times F(U) & \xrightarrow{\text{action}} & F(V) \\
\downarrow & \downarrow & \downarrow \\
\mathcal{O}_X(U) \times F(U) & \xrightarrow{\text{action}} & F(U)
\end{array}
\]

(3.2.12.1) commutes. (You should convince yourself that I haven’t forgotten anything.)
Recall that the notion of $A$-module generalizes the notion of abelian group, because an abelian group is the same thing as a $\mathbb{Z}$-module. Similarly, the notion of $\mathcal{O}_X$-module generalizes the notion of sheaf of abelian groups, because the latter is the same thing as a $\mathbb{Z}$-module, where $\mathbb{Z}$ is the constant sheaf associated to $\mathbb{Z}$. Hence when we are proving things about $\mathcal{O}_X$-modules, we are also proving things about sheaves of abelian groups.

3.2.13. For those who know about vector bundles. The motivating example of $\mathcal{O}_X$-modules is the sheaf of sections of a vector bundle. If $(X, \mathcal{O}_X)$ is a differentiable manifold (so $\mathcal{O}_X$ is the sheaf of differentiable functions), and $\pi : V \to X$ is a vector bundle over $X$, then the sheaf of differentiable sections $\phi : X \to V$ is an $\mathcal{O}_X$-module. Indeed, given a section $s$ of $\pi$ over an open subset $U \subset X$, and a function $f$ on $U$, we can multiply $s$ by $f$ to get a new section $fs$ of $\pi$ over $U$. Moreover, if $V$ is a smaller subset, then we could multiply $f$ by $s$ and then restrict to $V$, or we could restrict both $f$ and $s$ to $V$ and then multiply, and we would get the same answer. That is precisely the commutativity of (3.2.12.1).

3.3 Morphisms of presheaves and sheaves

3.3.1. Whenever one defines a new mathematical object, category theory teaches to try to understand maps between them. We now define morphisms of presheaves, and similarly for sheaves. In other words, we will describe the category of presheaves (of sets, abelian groups, etc.) and the category of sheaves.

A morphism of presheaves of sets (or indeed of sheaves with values in any category) on $X$, $f : \mathcal{F} \to \mathcal{G}$, is the data of maps $f(U) : \mathcal{F}(U) \to \mathcal{G}(U)$ for all $U$ behaving well with respect to restriction: if $U \rightarrow V$ then

\[
\begin{array}{ccc}
\mathcal{F}(V) & \xrightarrow{f(V)} & \mathcal{G}(V) \\
\downarrow \text{res}_{V,U} & & \downarrow \text{res}_{V,U} \\
\mathcal{F}(U) & \xrightarrow{f(U)} & \mathcal{G}(U)
\end{array}
\]

commutes. (Notice: the underlying space of both $\mathcal{F}$ and $\mathcal{G}$ is $X$.)

Morphisms of sheaves are defined identically: the morphisms from a sheaf $\mathcal{F}$ to a sheaf $\mathcal{G}$ are precisely the morphisms from $\mathcal{F}$ to $\mathcal{G}$ as presheaves. (Translation: The category of sheaves on $X$ is a full subcategory of the category of presheaves on $X$.)

An example of a morphism of sheaves is the map from the sheaf of differentiable functions on $\mathbb{R}$ to the sheaf of continuous functions. This is a “forgetful map”: we are forgetting that these functions are differentiable, and remembering only that they are continuous.

We may as well set some notation: let $\text{Sets}_X$, $\text{Ab}_X$, etc. denote the category of sheaves of sets, abelian groups, etc. on a topological space $X$. Let $\text{Mod}_{\mathcal{O}_X}$ denote the category of $\mathcal{O}_X$-modules on a ringed space $(X, \mathcal{O}_X)$. Let $\text{Sets}_X^{\text{pre}}$, etc. denote the category of presheaves of sets, etc. on $X$. 
3.3.2. Side-remark for category-lovers. If you interpret a presheaf on \(X\) as a contravariant functor (from the category of open sets), a morphism of presheaves on \(X\) is a natural transformation of functors \((\ref{sec:3.2:2})\).

3.3.A. Exercise. Suppose \(f : X \to Y\) is a continuous map of topological spaces (i.e., a morphism in the category of topological spaces). Show that pushforward gives a functor \(\text{Sets}_X \to \text{Sets}_Y\). Here \(\text{Sets}\) can be replaced by many other categories. (Watch out for some possible confusion: a presheaf is a functor, and presheaves form a category. It may be best to forget that presheaves are functors for now.)

3.3.B. Important exercise and definition: “Sheaf \(\text{Hom}\)”. Suppose \(\mathcal{F}\) and \(\mathcal{G}\) are two sheaves of sets on \(X\). (In fact, it will suffice that \(\mathcal{F}\) is a presheaf.) Let \(\mathcal{H}\text{om}(\mathcal{F}, \mathcal{G})\) be the collection of data

\[
\mathcal{H}\text{om}(\mathcal{F}, \mathcal{G})(U) := \text{Mor}(\mathcal{F}|_U, \mathcal{G}|_U).
\]

(Recall the notation \(\mathcal{F}|_U\), the restriction of the sheaf to the open set \(U\), Example \(\ref{ex:2.2:8}\).) Show that this is a sheaf of sets on \(X\). This is called the “sheaf \(\mathcal{H}\text{om}\)”. (Strictly speaking, we should reserve \(\text{Hom}\) for when we are in additive category, so this should possibly be called “sheaf \(\text{Mor}\)”. But the terminology sheaf \(\mathcal{H}\text{om}\) is too established to uproot.) Show that if \(\mathcal{G}\) is a sheaf of abelian groups, then \(\mathcal{H}\text{om}(\mathcal{F}, \mathcal{G})\) is a sheaf of abelian groups. Implicit in this fact is that \(\mathcal{H}\text{om}(\mathcal{F}, \mathcal{G})\) is an abelian group. (This exercise is somewhat tedious, but in the end very rewarding.) The same construction will “obviously” work for sheaves with values in any category.

Warning: \(\mathcal{H}\text{om}\) does not commute with taking stalks. More precisely: it is not true that \(\mathcal{H}\text{om}(\mathcal{F}, \mathcal{G})_p\) is isomorphic to \(\mathcal{H}\text{om}(\mathcal{F}_p, \mathcal{G}_p)\). (Can you think of a counterexample? Does there at least exist a map from one of these to the other?) We will use many variants of the definition of \(\mathcal{H}\text{om}\). For example, if \(\mathcal{F}\) and \(\mathcal{G}\) are sheaves of abelian groups on \(X\), then \(\mathcal{H}\text{om}_{\text{Ab}}(\mathcal{F}, \mathcal{G})\) is defined by taking \(\mathcal{H}\text{om}(\mathcal{F}, \mathcal{G})(U)\) to be the maps as sheaves of abelian groups \(\mathcal{F}|_U \to \mathcal{G}|_U\). Similarly, if \(\mathcal{F}\) and \(\mathcal{G}\) are \(\mathcal{O}_X\)-modules, we define \(\mathcal{H}\text{om}_{\text{Mod}}(\mathcal{F}, \mathcal{G})\) in the analogous way. Obnoxiously, the subscripts \(\text{Ab}\) and \(\text{Mod}\) are essentially always dropped (here and in the literature), so be careful which category you are working in! We call \(\mathcal{H}\text{om}_{\text{Mod}}(\mathcal{F}, \mathcal{O}_X)\) the dual of the \(\mathcal{O}_X\)-module \(\mathcal{F}\), and denoted it \(\mathcal{F}^\vee\).

3.3.C. Unimportant exercise (reality check).

(a) If \(\mathcal{F}\) is a sheaf of sets on \(X\), then show that \(\mathcal{H}\text{om}([p], \mathcal{F}) \cong \mathcal{F}\), where \([p]\) is the constant sheaf associated to the one element set \([p]\).

(b) If \(\mathcal{F}\) is a sheaf of abelian groups on \(X\), then show that \(\mathcal{H}\text{om}_{\text{Ab}}(\mathbb{Z}, \mathcal{F}) \cong \mathcal{F}\).

(c) If \(\mathcal{F}\) is an \(\mathcal{O}_X\)-module, then show that \(\mathcal{H}\text{om}_{\text{Mod}}(\mathcal{O}_X, \mathcal{F}) \cong \mathcal{F}\).

A key idea in (b) and (c) is that 1 “generates” (in some sense) \(\mathbb{Z}\) (in (b)) and \(\mathcal{O}_X\) (in (c)).

3.3.3. Presheaves of abelian groups (and even “presheaf \(\mathcal{O}_X\)-modules”) form an abelian category.

We can make module-like constructions using presheaves of abelian groups on a topological space \(X\). (In this section, all (pre)sheaves are of abelian groups.) For example, we can clearly add maps of presheaves and get another map of presheaves: if \(f, g : \mathcal{F} \to \mathcal{G}\), then we define the map \(f + g\) by \((f + g)(V) = \cdots\).
f(V) + g(V). (There is something small to check here: that the result is indeed a map of presheaves.) In this way, presheaves of abelian groups form an additive category (Definition 2.6.1). For exactly the same reasons, sheaves of abelian groups also form an additive category.

If \( f : \mathcal{F} \to \mathcal{G} \) is a morphism of presheaves, define the presheaf kernel \( \ker_{\text{pre}} f \) by \( (\ker_{\text{pre}} f)(U) = \ker f(U) \).

**3.3.D. Exercise.** Show that \( \ker_{\text{pre}} f \) is a presheaf. (Hint: if \( U \hookrightarrow V \), define the restriction map by chasing the following diagram:

\[
\begin{array}{c}
0 & \longrightarrow & \ker_{\text{pre}} f(V) & \longrightarrow & \mathcal{F}(V) & \longrightarrow & \mathcal{G}(V) \\
& & \downarrow & \downarrow & \downarrow & & \\
0 & \longrightarrow & \ker_{\text{pre}} f(U) & \longrightarrow & \mathcal{F}(U) & \longrightarrow & \mathcal{G}(U)
\end{array}
\]

You should check that the restriction maps compose as desired.)

Define the presheaf cokernel \( \text{coker}_{\text{pre}} f \) similarly. It is a presheaf by essentially the same argument.

**3.3.E. Exercise: The Cokernel Deserves Its Name.** Show that the presheaf cokernel satisfies the universal property of cokernels (Definition 2.6.3) in the category of presheaves.

Similarly, \( \ker_{\text{pre}} f \to \mathcal{F} \) satisfies the universal property for kernels in the category of presheaves.

It is not too tedious to verify that presheaves of abelian groups form an abelian category, and the reader is free to do so. The key idea is that all abelian-categorical notions may be defined and verified “open set by open set”. We needn’t worry about restriction maps — they “come along for the ride”. Hence we can define terms such as subpresheaf, image presheaf, quotient presheaf, cokernel presheaf, and they behave the way one expect. You construct kernels, quotients, cokernels, and images open set by open set. Homological algebra (exact sequences and so forth) works, and also “works open set by open set”. In particular:

**3.3.F. Easy Exercise.** Show (or observe) that for a topological space \( X \) with open set \( U \), \( \mathcal{F} \to \mathcal{F}(U) \) gives a functor from presheaves of abelian groups on \( X \), \( \text{Ab}_{\text{pre}}^X \), to abelian groups, \( \text{Ab} \). Then show that this functor is exact.

**3.3.G. Exercise.** Show that \( 0 \to \mathcal{F}_1 \to \mathcal{F}_2 \to \cdots \to \mathcal{F}_n \to 0 \) is exact if and only if \( 0 \to \mathcal{F}_1(U) \to \mathcal{F}_2(U) \to \cdots \to \mathcal{F}_n(U) \to 0 \) is exact for all \( U \).

The above discussion essentially carries over without change to presheaves with values in any abelian category. (Think this through if you wish.)

However, we are interested in more geometric objects, sheaves, where things can be understood in terms of their local behavior, thanks to the identity and gluing axioms. We will soon see that sheaves of abelian groups also form an abelian category, but a complication will arise that will force the notion of sheafification on us. Sheafification will be the answer to many of our prayers. We just don’t realize it yet.
To begin with, sheaves $\mathbb{A}b_X$ may be easily seen to form an additive category (essentially because presheaves $\mathbb{A}b_{\text{pre}X}$ already do, and sheaves form a full subcategory).

Kernels work just as with presheaves:

**3.3.H. IMPORTANT EXERCISE.** Suppose $f : \mathcal{F} \to \mathcal{G}$ is a morphism of sheaves. Show that the presheaf kernel $\ker_{\text{pre}} f$ is in fact a sheaf. Show that it satisfies the universal property of kernels (Definition 2.6.3). (Hint: the second question follows immediately from the fact that $\ker_{\text{pre}} f$ satisfies the universal property in the category of presheaves.)

Thus if $f$ is a morphism of sheaves, we define

$$\ker f := \ker_{\text{pre}} f.$$  

The problem arises with the cokernel.

**3.3.I. IMPORTANT EXERCISE.** Let $X$ be $\mathbb{C}$ with the classical topology, let $\mathcal{Z}$ be the constant sheaf on $X$ associated to $\mathcal{Z}$, $\mathcal{O}_X$ the sheaf of holomorphic functions, and $\mathcal{F}$ the presheaf of functions admitting a holomorphic logarithm. Describe an exact sequence of presheaves on $X$:

$$0 \to \mathcal{Z} \to \mathcal{O}_X \to \mathcal{F} \to 0$$

where $\mathcal{Z} \to \mathcal{O}_X$ is the natural inclusion and $\mathcal{O}_X \to \mathcal{F}$ is given by $f \mapsto \exp 2\pi i f$. (Be sure to verify exactness.) Show that $\mathcal{F}$ is not a sheaf. (Hint: $\mathcal{F}$ does not satisfy the gluability axiom. The problem is that there are functions that don’t have a logarithm but locally have a logarithm.) This will come up again in Example 3.4.9.

We will have to put our hopes for understanding cokernels of sheaves on hold for a while. We will first learn to understand sheaves using stalks.

### 3.4 Properties determined at the level of stalks, and sheafification

**3.4.1. Properties determined by stalks.** In this section, we will see that lots of facts about sheaves can be checked “at the level of stalks”. This isn’t true for presheaves, and reflects the local nature of sheaves. We will see that sections and morphisms are determined “by their stalks”, and the property of a morphism being an isomorphism may be checked at stalks. (The last one is the trickiest.)

**3.4.A. IMPORTANT EXERCISE (sections are determined by germs).** Prove that a section of a sheaf of sets is determined by its germs, i.e. the natural map

$$\mathcal{F}(U) \to \prod_{p \in U} \mathcal{F}_p$$

(3.4.1.1) is injective. Hint 1: you won’t use the gluability axiom, so this is true for separated presheaves. Hint 2: it is false for presheaves in general, see Exercise 3.4.F, so you will use the identity axiom. (Your proof will also apply to sheaves of groups, rings, etc.)
This exercise suggests an important question: which elements of the right side of (3.4.1) are in the image of the left side?

3.4.2. \textbf{Important definition.} We say that an element $\prod_{p \in U} s_p$ of the right side $\prod_{p \in U} \mathcal{F}_p$ of (3.4.1) consists of \textbf{compatible germs} if for all $p \in U$, there is some representative $(U_p, s'_p \in \mathcal{F}(U_p))$ for $s_p$ (where $p \in U_p \subset U$) such that the germ of $s'_p$ at all $y \in U_p$ is $s_y$. You will have to think about this a little. Clearly any section $s$ of $\mathcal{F}$ over $U$ gives a choice of compatible germs for $U$ — take $(U_p, s'_p) = (U, s)$.

3.4.B. \textbf{IMPORTANT EXERCISE.} Prove that any choice of compatible germs for a sheaf $\mathcal{F}$ over $U$ is the image of a section of $\mathcal{F}$ over $U$. (Hint: you will use glubability.)

We have thus completely described the image of (3.4.1), in a way that we will find useful.

3.4.3. \textbf{Remark.} This perspective is part of the motivation for the agricultural terminology “sheaf”: it is (the data of) a bunch of stalks, bundled together appropriately.

Now we throw morphisms into the mix.

3.4.C. \textbf{EXERCISE.} Show a morphism of (pre)sheaves (of sets, or rings, or abelian groups, or $\mathcal{O}_X$-modules) induces a morphism of stalks. More precisely, if $\phi : \mathcal{F} \rightarrow \mathcal{G}$ is a morphism of (pre)sheaves on $X$, and $p \in X$, describe a natural map $\phi_p : \mathcal{F}_p \rightarrow \mathcal{G}_p$. (You may wish to state this in the language of functors.)

3.4.D. \textbf{EXERCISE (morphisms are determined by stalks).} If $\phi_1$ and $\phi_2$ are morphisms from $\mathcal{F}$ to $\mathcal{G}$ that induce the same maps on each stalk, show that $\phi_1 = \phi_2$. Hint: consider the following diagram.

\[
\begin{array}{ccc}
\mathcal{F}(U) & \rightarrow & \mathcal{G}(U) \\
\downarrow & & \downarrow \\
\prod_{p \in U} \mathcal{F}_p & \rightarrow & \prod_{p \in U} \mathcal{G}_p
\end{array}
\]

3.4.E. \textbf{TRICKY EXERCISE (isomorphisms are determined by stalks).} Show that a morphism of sheaves of sets is an isomorphism if and only if it induces an isomorphism of all stalks. Hint: Use (3.4.3.1). Injectivity of maps of stalks uses the previous exercise 3.4.D. Once you have injectivity, show surjectivity using glubility; this is more subtle.

3.4.F. \textbf{EXERCISE.} (a) Show that Exercise 3.4.A is false for general presheaves.
(b) Show that Exercise 3.4.D is false for general presheaves.
(c) Show that Exercise 3.4.E is false for general presheaves.
(General hint for finding counterexamples of this sort: consider a 2-point space with the discrete topology, i.e. every subset is open.)

3.4.4. \textbf{Sheafification.}

Every sheaf is a presheaf (and indeed by definition sheaves on $X$ form a full subcategory of the category of presheaves on $X$). Just as groupification (§2.5.3...
gives a group that best approximates a semigroup, sheafification gives the sheaf that best approximates a presheaf, with an analogous universal property. (One possible example to keep in mind is the sheafification of the presheaf of holomorphic functions admitting a square root on $\mathbb{C}$ with the classical topology.)

3.4.5. Definition. If $\mathcal{F}$ is a presheaf on $X$, then a morphism of presheaves $\mathcal{F} \to \mathcal{F}^{sh}$ on $X$ is a sheafification of $\mathcal{F}$ if $\mathcal{F}^{sh}$ is a sheaf, and for any other sheaf $\mathcal{G}$, and any presheaf morphism $g : \mathcal{F} \to \mathcal{G}$, there exists a unique morphism of sheaves $f : \mathcal{F}^{sh} \to \mathcal{G}$ making the diagram

\[
\begin{array}{ccc}
\mathcal{F} & \xrightarrow{\text{sh}} & \mathcal{F}^{sh} \\
\downarrow{g} & & \downarrow{f} \\
\mathcal{G} & & \\
\end{array}
\]

commute.

3.4.G. Exercise. Show that sheafification is unique up to unique isomorphism. Show that if $\mathcal{F}$ is a sheaf, then the sheafification is $\mathcal{F} \xrightarrow{id} \mathcal{F}$. (This should be second nature by now.)

3.4.6. Construction. We next show that any presheaf of sets (or groups, rings, etc.) has a sheafification. Suppose $\mathcal{F}$ is a presheaf. Define $\mathcal{F}^{sh}$ by defining $\mathcal{F}^{sh}(U)$ as the set of compatible germs of the presheaf $\mathcal{F}$ over $U$. Explicitly:

\[
\mathcal{F}^{sh}(U) := \{(f_x \in \mathcal{F}_x)_{x \in U} : \text{for all } x \in U, \text{ there exists } x \in V \subset U \text{ and } s \in \mathcal{F}(V) \text{ with } s_y = f_y \text{ for all } y \in V\}.
\]

(Those who want to worry about the empty set are welcome to.)

3.4.H. Easy Exercise. Show that $\mathcal{F}^{sh}$ (using the tautological restriction maps) forms a sheaf.

3.4.I. Easy Exercise. Describe a natural map of presheaves $\text{sh} : \mathcal{F} \to \mathcal{F}^{sh}$.

3.4.J. Exercise. Show that the map $\text{sh}$ satisfies the universal property of sheafification (Definition 3.4.5). (This is easier than you might fear.)

3.4.K. Useful exercise, not just for category-lovers. Show that the sheafification functor is left-adjoint to the forgetful functor from sheaves on $X$ to presheaves on $X$. This is not difficult — it is largely a restatement of the universal property. But it lets you use results from §2.6.10 and can “explain” why you don’t need to sheafify when taking kernel (why the presheaf kernel is already the sheaf kernel), and why you need to sheafify when taking cokernel and (soon, in Exercise 3.5.H) $\otimes$.

3.4.L. Easy Exercise. Use the universal property to show that for any morphism of presheaves $\phi : \mathcal{F} \to \mathcal{G}$, we get a natural induced morphism of sheaves $\phi^{sh} : \mathcal{F}^{sh} \to \mathcal{G}^{sh}$. Show that sheafification is a functor from presheaves on $X$ to sheaves on $X$. 

3.4.M. Exercise. Show $\mathcal{F} \to \mathcal{F}^{\text{sh}}$ induces an isomorphism of stalks. (Possible hint: Use the concrete description of the stalks. Another possibility once you read Remark 3.6.3 judicious use of adjoints.)

3.4.7. Remark. The espace étalé construction (3.2.11) yields a different-sounding description of sheafification which may be preferred by some readers. The fundamental idea is identical. This is essentially the same construction as the one given here. Another construction is described in [EH].

3.4.8. Subsheaves and quotient sheaves.

3.4.N. Exercise. Suppose $\phi : \mathcal{F} \to \mathcal{G}$ is a morphism of sheaves (of sets) on a topological space $X$. Show that the following are equivalent.

(a) $\phi$ is a monomorphism in the category of sheaves.
(b) $\phi$ is injective on the level of stalks: $\phi_x : \mathcal{F}_x \to \mathcal{G}_x$ injective for all $x \in X$.
(c) $\phi$ is injective on the level of open sets: $\phi(U) : \mathcal{F}(U) \to \mathcal{G}(U)$ is injective for all open $U \subset X$.

(Possible hints: for (b) implies (a), recall that morphisms are determined by stalks, Exercise 3.4.D. For (a) implies (c), use the “indicator sheaf” with one section over every open set contained in $U$, and no section over any other open set.)

If these conditions hold, we say that $\mathcal{F}$ is a subsheaf of $\mathcal{G}$ (where the “inclusion” $\phi$ is sometimes left implicit).

3.4.O. Exercise. Continuing the notation of the previous exercise, show that the following are equivalent.

(a) $\phi$ is an epimorphism in the category of sheaves.
(b) $\phi$ is surjective on the level of stalks: $\phi_x : \mathcal{F}_x \to \mathcal{G}_x$ surjective for all $x \in X$.

If these conditions hold, we say that $\mathcal{G}$ is a quotient sheaf of $\mathcal{F}$.

Thus monomorphisms and epimorphisms — subsheafiness and quotient sheafiness — can be checked at the level of stalks.

Both exercises generalize readily to sheaves with values in any reasonable category, where “injective” is replaced by “monomorphism” and “surjective” is replaced by “epimorphism”.

Notice that there was no part (c) to Exercise 3.4.O, and Example 3.4.9 shows why. (But there is a version of (c) that implies (a) and (b): surjectivity on all open sets in the base of a topology implies surjectivity of the map of sheaves, Exercise 3.7.E.)

3.4.9. Example (cf. Exercise 3.3.I). Let $X = \mathbb{C}$ with the classical topology, and define $\mathcal{O}_X$ to be the sheaf of holomorphic functions, and $\mathcal{O}_X^*$ to be the sheaf of invertible (nowhere zero) holomorphic functions. This is a sheaf of abelian groups under multiplication. We have maps of sheaves

$$0 \longrightarrow \mathbb{Z} \xrightarrow{x \cdot 2\pi i} \mathcal{O}_X \xrightarrow{\exp} \mathcal{O}_X^* \longrightarrow 1$$

where $\mathbb{Z}$ is the constant sheaf associated to $\mathbb{Z}$. (You can figure out what the sheaves 0 and 1 mean; they are isomorphic, and are written in this way for reasons that may be clear.) We will soon interpret this as an exact sequence of sheaves of abelian
groups (the exponential exact sequence), although we don’t yet have the language to do so.

3.4.P. ENLIGHTENING EXERCISE. Show that $\mathcal{O}_X \xrightarrow{\exp} \mathcal{O}_X^\times$ describes $\mathcal{O}_X^\times$ as a quotient sheaf of $\mathcal{O}_X$. Show that it is not surjective on all open sets.

This is a great example to get a sense of what “surjectivity” means for sheaves: nowhere vanishing holomorphic functions have logarithms locally, but they need not globally.

3.5 Sheaves of abelian groups, and $\mathcal{O}_X$-modules, form abelian categories

We are now ready to see that sheaves of abelian groups, and their cousins, $\mathcal{O}_X$-modules, form abelian categories. In other words, we may treat them similarly to vector spaces, and modules over a ring. In the process of doing this, we will see that this is much stronger than an analogy; kernels, cokernels, exactness, and so forth can be understood at the level of germs (which are just abelian groups), and the compatibility of the germs will come for free.

The category of sheaves of abelian groups is clearly an additive category (Definition 2.6.1). In order to show that it is an abelian category, we must show that any morphism $\phi : \mathcal{F} \to \mathcal{G}$ has a kernel and a cokernel. We have already seen that $\phi$ has a kernel (Exercise 3.3.H): the presheaf kernel is a sheaf, and is a kernel in the category of sheaves.

3.5.A. EXERCISE. Show that the stalk of the kernel is the kernel of the stalks: there is a natural isomorphism

$$(\ker(\mathcal{F} \to \mathcal{G}))_x \cong \ker(\mathcal{F}_x \to \mathcal{G}_x).$$

We next address the issue of the cokernel. Now $\phi : \mathcal{F} \to \mathcal{G}$ has a cokernel in the category of presheaves; call it $\mathcal{H}^{\text{pre}}$ (where the superscript is meant to remind us that this is a presheaf). Let $\mathcal{H}^{\text{pre}} \xrightarrow{\text{sh}} \mathcal{H}$ be its sheafification. Recall that the cokernel is defined using a universal property: it is the colimit of the diagram

$$\mathcal{F} \xrightarrow{\phi} \mathcal{G} \xrightarrow{\longrightarrow} \mathcal{H}$$

in the category of presheaves. We claim that $\mathcal{H}$ is the cokernel of $\phi$ in the category of sheaves, and show this by proving the universal property. Given any sheaf $\mathcal{E}$ and a commutative diagram

$$\begin{array}{ccc}
\mathcal{F} & \xrightarrow{\phi} & \mathcal{G} \\
\downarrow & & \downarrow \\
\mathcal{0} & & \mathcal{E}
\end{array}$$
We construct

\[
\begin{array}{c}
\mathcal{F} \\
\downarrow \phi \\
0 \\
\mathcal{H}^{\text{pre}} \\
\downarrow \text{sh} \\
\mathcal{H} \\
\downarrow \\
\mathcal{E}
\end{array}
\]

We show that there is a unique morphism \( \mathcal{H} \to \mathcal{E} \) making the diagram commute. As \( \mathcal{H}^{\text{pre}} \) is the cokernel in the category of presheaves, there is a unique morphism of presheaves \( \mathcal{H}^{\text{pre}} \to \mathcal{E} \) making the diagram commute. But then by the universal property of sheafification (Definition 3.4.5), there is a unique morphism of sheaves \( \mathcal{H} \to \mathcal{E} \) making the diagram commute.

3.5.B. **Exercise.** Show that the stalk of the cokernel is naturally isomorphic to the cokernel of the stalk.

We have now defined the notions of kernel and cokernel, and verified that they may be checked at the level of stalks. We have also verified that the properties of a morphism being a monomorphism or epimorphism are also determined at the level of stalks (Exercises 3.4.N and 3.4.O). Hence sheaves of abelian groups on \( X \) form an abelian category.

We see more: all structures coming from the abelian nature of this category may be checked at the level of stalks. For example:

3.5.C. **Exercise.** Suppose \( \phi : \mathcal{F} \to \mathcal{G} \) is a morphism of sheaves of abelian groups. Show that the image sheaf \( \text{im} \phi \) is the sheafification of the image presheaf. (You must use the definition of image in an abelian category. In fact, this gives the accepted definition of image sheaf for a morphism of sheaves of sets.) Show that the stalk of the image is the image of the stalk.

As a consequence, **exactness of a sequence of sheaves may be checked at the level of stalks.** In particular:

3.5.D. **Important Exercise.** Show that taking the stalk of a sheaf of abelian groups is an exact functor. More precisely, if \( X \) is a topological space and \( p \in X \) is a point, show that taking the stalk at \( p \) defines an exact functor \( \text{Ab}_X \to \text{Ab} \).

3.5.E. **Exercise (Left-exactness of the functor of “sections over \( U \)).** Suppose \( U \subset X \) is an open set, and \( 0 \to \mathcal{F} \to \mathcal{G} \to \mathcal{H} \) is an exact sequence of sheaves of abelian groups. Show that

\[
0 \to \mathcal{F}(U) \to \mathcal{G}(U) \to \mathcal{H}(U)
\]

is exact. (You should do this “by hand”, even if you realize there is a very fast proof using the left-exactness of the “forgetful” right-adjoint to the sheafification functor.) Show that the section functor need not be exact: show that if \( 0 \to \mathcal{F} \to \mathcal{G} \to \mathcal{H} \to 0 \) is an exact sequence of sheaves of abelian groups, then

\[
0 \to \mathcal{F}(U) \to \mathcal{G}(U) \to \mathcal{H}(U) \to 0
\]

need not be exact. (Hint: the exponential exact sequence (3.4.9).)
3.5.F. Exercise: Left-exactness of pushforward. Suppose $0 \rightarrow \mathcal{F} \rightarrow \mathcal{G} \rightarrow \mathcal{H}$ is an exact sequence of sheaves of abelian groups on $X$. If $f : X \rightarrow Y$ is a continuous map, show that

$$0 \rightarrow f_* \mathcal{F} \rightarrow f_* \mathcal{G} \rightarrow f_* \mathcal{H}$$

is exact. (The previous exercise, dealing with the left-exactness of the global section functor can be interpreted as a special case of this, in the case where $Y$ is a point.)

3.5.G. Exercise. Show that if $(X, \mathcal{O}_X)$ is a ringed space, then $\mathcal{O}_X$-modules form an abelian category. (There isn’t much more to check!)

We end with a useful construction using some of the ideas in this section.

3.5.H. Important exercise: Tensor products of $\mathcal{O}_X$-modules. (a) Suppose $\mathcal{O}_X$ is a sheaf of rings on $X$. Define (categorically) what we should mean by tensor product of two $\mathcal{O}_X$-modules. Give an explicit construction, and show that it satisfies your categorical definition. Hint: take the “presheaf tensor product” — which needs to be defined — and sheafify. Note: $\otimes_{\mathcal{O}_X}$ is often written $\otimes$ when the subscript is clear from the context. (An example showing sheafification is necessary will arise in Example 15.1.1.)

(b) Show that the tensor product of stalks is the stalk of tensor product. (If you can show this, you may be able to make sense of the phrase “colimits commute with tensor products”.)

3.5.1. Conclusion. Just as presheaves are abelian categories because all abelian-categorical notions make sense open set by open set, sheaves are abelian categories because all abelian-categorical notions make sense stalk by stalk.

3.6 The inverse image sheaf

We next describe a notion that is fundamental, but rather intricate. We will not need it for some time, so this may be best left for a second reading. Suppose we have a continuous map $f : X \rightarrow Y$. If $\mathcal{F}$ is a sheaf on $X$, we have defined the pushforward or direct image sheaf $f_* \mathcal{F}$, which is a sheaf on $Y$. There is also a notion of inverse image sheaf. (We will not call it the pullback sheaf, reserving that name for a later construction for quasicoherent sheaves, §17.3.) This is a covariant functor $f^{-1}$ from sheaves on $Y$ to sheaves on $X$. If the sheaves on $Y$ have some additional structure (e.g. group or ring), then this structure is respected by $f^{-1}$.

3.6.1. Definition by adjoint: elegant but abstract. We define $f^{-1}$ as the left-adjoint to $f_*$. This isn’t really a definition; we need a construction to show that the adjoint exists. Note that we then get canonical maps $f^{-1}f_* \mathcal{F} \rightarrow \mathcal{F}$ (associated to the identity in $\text{Mor}_Y(f_* \mathcal{F}, f_* \mathcal{F})$) and $\mathcal{G} \rightarrow f_*f^{-1} \mathcal{G}$ (associated to the identity in $\text{Mor}_X(f^{-1} \mathcal{G}, f^{-1} \mathcal{G})$).

3.6.2. Construction: concrete but ugly. Define the temporary notation $f^{-1} \mathcal{G}^\text{pre}(U) = \lim_{V \supseteq f(U)} \mathcal{G}(V)$. (Recall the explicit description of colimit: sections are sections on
open sets containing $f(U)$, with an equivalence relation. Note that $f(U)$ won’t be an open set in general.)

3.6.A. Exercise. Show that this defines a presheaf on $X$.

Now define the inverse image of $\mathcal{G}$ by $f^{-1}\mathcal{G} := (f^{-1}\mathcal{G})^{\text{pre}}$. The next exercise shows that this satisfies the universal property. But you may wish to try the later exercises first, and come back to Exercise 3.6.B later. (For the later exercises, try to give two proofs, one using the universal property, and the other using the explicit description.)

3.6.B. Important Tricky Exercise. If $f : X \to Y$ is a continuous map, and $\mathcal{F}$ is a sheaf on $X$ and $\mathcal{G}$ is a sheaf on $Y$, describe a bijection $\text{Mor}_X(f^{-1}\mathcal{G}, \mathcal{F}) \leftrightarrow \text{Mor}_Y(\mathcal{G}, f_*\mathcal{F})$.

Observe that your bijection is “natural” in the sense of the definition of adjoints (i.e. functorial in both $\mathcal{F}$ and $\mathcal{G}$).

3.6.3. Remark. As a special case, if $X$ is a point $p \in Y$, we see that $f^{-1}\mathcal{G}$ is the stalk $\mathcal{G}_p$ of $\mathcal{G}$, and maps from the stalk $\mathcal{G}_p$ to a set $S$ are the same as maps of sheaves on $Y$ from $\mathcal{G}$ to the skyscraper sheaf with set $S$ supported at $p$. You may prefer to prove this special case by hand directly before solving Exercise 3.6.B as it is enlightening. (It can also be useful — can you use it to solve Exercises 3.4.M and 3.4.O?)

3.6.C. Exercise. Show that the stalks of $f^{-1}\mathcal{G}$ are the same as the stalks of $\mathcal{G}$. More precisely, if $f(p) = q$, describe a natural isomorphism $\mathcal{G}_q \cong (f^{-1}\mathcal{G})_p$. (Possible hint: use the concrete description of the stalk, as a colimit. Recall that stalks are preserved by sheafification, Exercise 3.4.M. Alternatively, use adjointness.) This, along with the notion of compatible stalks, may give you a way of thinking about inverse image sheaves.

3.6.D. Exercise (Easy but Useful). If $U$ is an open subset of $Y$, $i : U \to Y$ is the inclusion, and $\mathcal{G}$ is a sheaf on $Y$, show that $i^{-1}\mathcal{G}$ is naturally isomorphic to $\mathcal{G}|_U$.

3.6.E. Exercise. Show that $f^{-1}$ is an exact functor from sheaves of abelian groups on $Y$ to sheaves of abelian groups on $X$ (cf. Exercise 3.5.D). (Hint: exactness can be checked on stalks, and by Exercise 3.6.C the stalks are the same.) The identical argument will show that $f^{-1}$ is an exact functor from $\mathcal{O}_Y$-modules (on $Y$) to $f^{-1}\mathcal{O}_Y$-modules (on $X$), but don’t bother writing that down. (Remark for experts: $f^{-1}$ is a left-adjoint, hence right-exact by abstract nonsense, as discussed in §2.6.10. Left-exactness holds because colimits over directed systems are exact.)

3.6.F. Exercise. (a) Suppose $Z \subset Y$ is a closed subset, and $i : Z \hookrightarrow Y$ is the inclusion. If $\mathcal{F}$ is a sheaf on $Z$, then show that the stalk $(i_*\mathcal{F})_y$ is a one element-set if $y \notin Z$, and $\mathcal{F}_y$ if $y \in Z$.

(b) Definition: Define the support of a sheaf $\mathcal{F}$ of sets, denoted $\text{Supp} \mathcal{F}$, as the locus where the stalks are not the one-element set:

$$\text{Supp} \mathcal{F} := \{x \in X : |\mathcal{F}_x| \neq 1\}.$$ (More generally, if the sheaf has value in some category, the support consists of points where the stalk is not the final object. For sheaves of abelian groups, the
support consists of points with non-zero stalks.) Suppose $\text{Supp} \mathcal{F} \subset Z$ where $Z$ is closed. Show that the natural map $\mathcal{F} \to i_! i^{-1} \mathcal{F}$ is an isomorphism. Thus a sheaf supported on a closed subset can be considered a sheaf on that closed subset. ("Support" is a useful notion, and will arise again in §14.7.C)

3.6.G. **EXERCISE (EXTENSION BY ZERO $i_!$: AN OCCASIONAL LEFT-ADJOINT TO $f^{-1}$).** In addition to always being a left-adjoint, $f^{-1}$ can sometimes be a right-adjoint. Suppose $i : U \hookrightarrow Y$ is an open immersion of ringed spaces. Define extension by zero $i_! : \text{Mod}_{\mathcal{O}_U} \to \text{Mod}_{\mathcal{O}_Y}$ as follows. Suppose $\mathcal{F}$ is an $\mathcal{O}_U$-module. For open $W \subset Y$, $i_! \mathcal{F}(W) = \mathcal{F}(W)$ if $W \subset U$, and $0$ otherwise (with the obvious restriction maps). Note that $i_! \mathcal{F}$ is an $\mathcal{O}_Y$-module, and that this clearly defines a functor. (The symbol "!" is read as “shriek”. I have no idea why. Thus $i_!$ is read as “$i$-lower-shriek”.)

(a) For $y \in Y$, show that $(i_! \mathcal{F})_y = \mathcal{F}_y$ if $y \in U$, and $0$ otherwise.

(b) Show that $i_!$ is an exact functor.

(c) Describe an inclusion $i_! i^{-1} \mathcal{F} \hookrightarrow \mathcal{F}$.

(d) Show that $(i_! i^{-1})$ is an adjoint pair, so there is a natural bijection $\text{Hom}_{\mathcal{O}_Y}(i_! \mathcal{F}, \mathcal{G}) \leftrightarrow \text{Hom}_{\mathcal{O}_U}(\mathcal{F}, \mathcal{G}_|_U)$ for any $\mathcal{O}_Y$-module $\mathcal{G}$. (In particular, the sections of $\mathcal{G}$ over $U$ can be identified with $\text{Hom}_{\mathcal{O}_Y}(i_! \mathcal{O}_U, \mathcal{G})$.)

3.7 Recovering sheaves from a “sheaf on a base”

Sheaves are natural things to want to think about, but hard to get our hands on. We like the identity and gluability axioms, but they make proving things trickier than for presheaves. We have discussed how we can understand sheaves using stalks. We now introduce a second way of getting a hold of sheaves, by introducing the notion of a sheaf on a base. Warning: this way of understanding an entire sheaf from limited information is confusing. It may help to keep sight of the central insight that this limited information lets you understand germs, and the notion of when they are compatible (with nearby germs).

First, we define the notion of a base of a topology. Suppose we have a topological space $X$, i.e. we know which subsets $U_i$ of $X$ are open. Then a base of a topology is a subcollection of the open sets $\{B_i\} \subset \{U_i\}$, such that each $U_i$ is a union of the $B_i$. Here is one example that you have seen early in your mathematical life. Suppose $X = \mathbb{R}^n$. Then the way the usual topology is often first defined is by defining open balls $B_r(x) = \{y \in \mathbb{R}^n : \|y - x\| < r\}$, and declaring that any union of open balls is open. So the balls form a base of the classical topology — we say they generate the classical topology. As an application of how we use them, to check continuity of some map $f : X \to \mathbb{R}^n$, you need only think about the pullback of balls on $\mathbb{R}^n$.

Now suppose we have a sheaf $\mathcal{F}$ on $X$, and a base $\{B_i\}$ on $X$. Then consider the information $(\mathcal{F}(B_i)), (\text{res}_{B_i, B_j} : \mathcal{F}(B_i) \to \mathcal{F}(B_j)))$, which is a subset of the information contained in the sheaf — we are only paying attention to the information involving elements of the base, not all open sets.

We can recover the entire sheaf from this information. This is because we can determine the stalks from this information, and we can determine when germs are compatible.
3.7.A. Exercise. Make this precise.

This suggests a notion, called a sheaf on a base. A sheaf of sets (rings etc.) on a base \( \{B_i\} \) is the following. For each \( B_i \) in the base, we have a set \( F(B_i) \). If \( B_i \subset B_j \), we have maps \( \text{res}_{B_i,B_j} : F(B_j) \to F(B_i) \). (Things called B are always assumed to be in the base.) If \( B_i \subset B_j \subset B_k \), then \( \text{res}_{B_k,B_i} = \text{res}_{B_j,B_i} \circ \text{res}_{B_k,B_j} \). So far we have defined a presheaf on a base \( \{B_i\} \).

We also require the base identity axiom: If \( B = \bigcup B_i \), then if \( f, g \in F(B) \) such that \( \text{res}_{B_i,B_j} : F(B) \to F(B_i) \) for all \( i \), then \( f = g \).

We require the base gluability axiom too: If \( B = \bigcup B_i \), and we have \( f_i \in F(B_i) \) such that \( f_i \) agrees with \( f_j \) on any basic open set contained in \( B_i \cap B_j \) (i.e. \( \text{res}_{B_i,B_k} f_i = \text{res}_{B_j,B_k} f_i \) for all \( B_k \subset B_i \cap B_j \) then there exists \( f \in F(B) \) such that \( \text{res}_{B_i,B_j} f = f_i \) for all \( i \).

3.7.1. Theorem. — Suppose \( \{B_i\} \) is a base on \( X \), and \( F \) is a sheaf of sets on this base. Then there is a sheaf \( \mathcal{F} \) extending \( F \) (with isomorphisms \( \mathcal{F}(B_i) \cong F(B_i) \) agreeing with the restriction maps). This sheaf \( \mathcal{F} \) is unique up to unique isomorphism.

Proof. We will define \( \mathcal{F} \) as the sheaf of compatible germs of \( F \).

Define the stalk of a base presheaf \( F \) at \( p \in X \) by

\[
F_p = \lim_{\longrightarrow} F(B_i)
\]

where the colimit is over all \( B_i \) (in the base) containing \( p \).

We will say a family of germs in an open set \( U \) is compatible near \( p \) if there is a section \( s \) of \( F \) over some \( B_i \) containing \( p \) such that the germs over \( B_i \) are precisely the germs of \( s \). More formally, define

\[
\mathcal{F}(U) := \{ (f_p \in F_p)_{p \in U} : \text{for all } p \in U, \text{ there exists } B \text{ with } p \subset B \subset U, s \in F(B), \text{ with } s_q = f_q \text{ for all } q \in B \}
\]

where each \( B \) is in our base.

This is a sheaf (for the same reasons as the sheaf of compatible germs was earlier, cf. Exercise 3.4.H).

I next claim that if \( B \) is in our base, the natural map \( F(B) \to \mathcal{F}(B) \) is an isomorphism.

3.7.B. Tricky Exercise. Describe the inverse map \( \mathcal{F}(B) \to F(B) \), and verify that it is indeed inverse. Possible hint: elements of \( \mathcal{F}(U) \) are determined by stalks, as are elements of \( F(U) \).

Thus sheaves on \( X \) can be recovered from their “restriction to a base”. This is a statement about objects in a category, so we should hope for a similar statement about morphisms.

3.7.C. Important Exercise: Morphisms of sheaves correspond to morphisms of sheaves on a base. Suppose \( \{B_i\} \) is a base for the topology of \( X \).
(a) Verify that a morphism of sheaves is determined by the induced morphism of sheaves on the base.
(b) Show that a morphism of sheaves on the base (i.e. such that the diagram
\[
\begin{array}{ccc}
F(B_i) & \longrightarrow & G(B_i) \\
\downarrow & & \downarrow \\
F(B_j) & \longrightarrow & G(B_j)
\end{array}
\]
commutes for all \(B_i \hookrightarrow B_i\) gives a morphism of the induced sheaves. (Possible hint: compatible stalks.)

3.7.D. IMPORTANT EXERCISE. Suppose \(X = \bigcup U_i\) is an open cover of \(X\), and we have sheaves \(F_i\) on \(U_i\) along with isomorphisms \(\phi_{ij} : F_i|_{U_i \cap U_j} \rightarrow F_j|_{U_i \cap U_j}\) (with \(\phi_{ii}\) the identity) that agree on triple overlaps (i.e. \(\phi_{jk} \circ \phi_{ij} = \phi_{ik}\) on \(U_i \cap U_j \cap U_k\)). Show that these sheaves can be glued together into a sheaf \(F\) on \(X\) (unique up to
unique isomorphism), such that \(F_i = F|_{U_i}\), and the isomorphisms over \(U_i \cap U_j\) are the obvious ones. (Thus we can “glue sheaves together”, using limited patching information.) Warning: we are not assuming this is a finite cover, so you cannot use induction. For this reason this exercise can be perplexing. (You can use the ideas of this section to solve this problem, but you don’t necessarily need to. Hint: As the base, take those open sets contained in some \(U_i\). Small observation: the hypothesis that \(\phi_{ii}\) is extraneous, as it follows from the cocycle condition.)

3.7.2. Remark for experts. Exercise 3.7.D almost says that the “set” of sheaves forms a sheaf itself, but not quite. Making this precise leads one to the notion of a stack.

3.7.E. UNIMPORTANT EXERCISE. Suppose a morphism of sheaves \(F \rightarrow G\) on a base \(B_i\) is surjective for all \(B_i\) (i.e. \(F(B_i) \rightarrow G(B_i)\) is surjective for all \(i\)). Show that the morphism of sheaves (not on the base) is surjective. The converse is not true, unlike the case for injectivity. This gives a useful criterion for surjectivity (“surjectivity on small enough open sets”).

3.7.3. Observation. In the proof of Theorem 3.7.1, we need even less information than given in the hypotheses. What we are really using is that the opens in the base, and their inclusions, form a filtered set. You will appreciate this observation much later, in the proof of Theorem 14.3.2.
Part II

Schemes
CHAPTER 4

Toward affine schemes: the underlying set, and the underlying topological space

The very idea of scheme is of infantile simplicity — so simple, so humble, that no one before me thought of stooping so low. So childish, in short, that for years, despite all the evidence, for many of my erudite colleagues, it was really “not serious”! — Grothendieck

4.1 Toward schemes

We are now ready to consider the notion of a scheme, which is the type of geometric space central to algebraic geometry. We should first think through what we mean by “geometric space”. You have likely seen the notion of a manifold, and we wish to abstract this notion so that it can be generalized to other settings, notably so that we can deal with non-smooth and arithmetic objects.

The key insight behind this generalization is the following: we can understand a geometric space (such as a manifold) well by understanding the functions on this space. More precisely, we will understand it through the sheaf of functions on the space. If we are interested in differentiable manifolds, we will consider differentiable functions; if we are interested in smooth manifolds, we will consider smooth functions; and so on.

Thus we will define a scheme to be the following data

- The set: the points of the scheme
- The topology: the open sets of the scheme
- The structure sheaf: the sheaf of “algebraic functions” (a sheaf of rings) on the scheme.

Recall that a topological space with a sheaf of rings is called a ringed space (§3.2.12).

We will try to draw pictures throughout. Pictures can help develop geometric intuition, which can guide the algebraic development (and, eventually, vice versa). Some people find pictures very helpful, while others are repulsed or nonplussed or confused.

We will try to make all three notions as intuitive as possible. For the set, in the key example of complex (affine) varieties (roughly, things cut out in $\mathbb{C}^n$ by polynomials), we will see that the points are the “traditional points” (n-tuples of complex numbers), plus some extra points that will be handy to have around. For the topology, we will require that “algebraic functions vanish on closed sets”, and require nothing else. For the sheaf of algebraic functions (the structure sheaf), we will expect that in the complex plane, $(3x^2 + y^2)/(2x + 4xy + 1)$ should be
an algebraic function on the open set consisting of points where the denominator doesn’t vanish, and this will largely motivate our definition.

**4.1.1. Example: Differentiable manifolds.** As motivation, we return to our example of differentiable manifolds, reinterpreting them in this light. We will be quite informal in this discussion. Suppose $X$ is a manifold. It is a topological space, and has a sheaf of differentiable functions $\mathcal{O}_X$ (see §3.1). This gives $X$ the structure of a ringed space. We have observed that evaluation at a point $p \in X$ gives a surjective map from the stalk to $\mathbb{R}$

$$\mathcal{O}_{X,p} \longrightarrow \mathbb{R},$$

so the kernel, the (germs of) functions vanishing at $p$, is a maximal ideal $m_X$ (see §3.1.1).

We could define a differentiable real manifold as a topological space $X$ with a sheaf of rings. We would require that there is a cover of $X$ by open sets such that on each open set the ringed space is isomorphic to a ball around the origin in $\mathbb{R}^n$ (with the sheaf of differentiable functions on that ball). With this definition, the ball is the basic patch, and a general manifold is obtained by gluing these patches together. (Admittedly, a great deal of geometry comes from how one chooses to patch the balls together!) In the algebraic setting, the basic patch is the notion of an affine scheme, which we will discuss soon. (In the definition of manifold, there is an additional requirement that the topological space be Hausdorff, to avoid pathologies. Schemes are often required to be “separated” to avoid essentially the same pathologies. Separatedness will be discussed in Chapter 11.)

**Functions are determined by their values at points.** This is an obvious statement, but won’t be true for schemes in general. We will see an example in Exercise 4.2.A(a), and discuss this behavior further in §4.2.9.

**Morphisms of manifolds.** How can we describe differentiable maps of manifolds $X \rightarrow Y$? They are certainly continuous maps — but which ones? We can pull back functions along continuous maps. Differentiable functions pull back to differentiable functions. More formally, we have a map $f^{-1}\mathcal{O}_Y \rightarrow \mathcal{O}_X$. (The inverse image sheaf $f^{-1}$ was defined in §3.6) Inverse image is left-adjoint to pushforward, so we also get a map $f^* : \mathcal{O}_Y \rightarrow f_*\mathcal{O}_X$.

Certainly given a differentiable map of manifolds, differentiable functions pull back to differentiable functions. It is less obvious that this is a sufficient condition for a continuous function to be differentiable.

**4.1.A. IMPORTANT EXERCISE FOR THOSE WITH A LITTLE EXPERIENCE WITH MANIFOLDS.** Prove that a continuous function of differentiable manifolds $f : X \rightarrow Y$ is differentiable if differentiable functions pull back to differentiable functions, i.e. if pullback by $f$ gives a map $\mathcal{O}_Y \rightarrow f_*\mathcal{O}_X$. (Hint: check this on small patches. Once you figure out what you are trying to show, you’ll realize that the result is immediate.)

**4.1.B. EXERCISE.** Show that a morphism of differentiable manifolds $f : X \rightarrow Y$ with $f(p) = q$ induces a morphism of stalks $f^* : \mathcal{O}_{Y,q} \rightarrow \mathcal{O}_{X,p}$. Show that $f^*(m_{Y,q}) \subset m_{X,p}$. In other words, if you pull back a function that vanishes at $q$, you get a function that vanishes at $p$ — not a huge surprise. (In §7.3 we formalize this by saying that maps of differentiable manifolds are maps of locally ringed spaces.)
4.1.2. Aside. Here is a little more for experts: Notice that this induces a map on tangent spaces (see Aside 3.1.2)

\[(m_{X,p}/m^2_{X,p})^\vee \to (m_{Y,q}/m^2_{Y,q})^\vee.\]

This is the tangent map you would geometrically expect. Again, it is interesting that the cotangent map \(m_{Y,q}/m^2_{Y,q} \to m_{X,p}/m^2_{X,p}\) is algebraically more natural than the tangent map (there are no “duals”).

Experts are now free to try to interpret other differential-geometric information using only the map of topological spaces and map of sheaves. For example: how can one check if \(f\) is a smooth map? How can one check if \(f\) is an immersion? (We will see that the algebro-geometric version of these notions are smooth morphism and locally closed immersion, see Chapter 25 and § 9.1.3 respectively.)

4.1.3. Side Remark. Manifolds are covered by disks that are all isomorphic. This isn’t true for schemes (even for “smooth complex varieties”). There are examples of two “smooth complex curves” (the algebraic version of Riemann surfaces) \(X\) and \(Y\) so that no non-empty open subset of \(X\) is isomorphic to a non-empty open subset of \(Y\). And there is an example of a Riemann surface \(X\) such that no two open subsets of \(X\) are isomorphic. Informally, this is because in the Zariski topology on schemes, all non-empty open sets are “huge” and have more “structure”.

4.1.4. Other examples. If you are interested in differential geometry, you will be interested in differentiable manifolds, on which the functions under consideration are differentiable functions. Similarly, if you are interested in topology, you will be interested in topological spaces, on which you will consider the continuous function. If you are interested in complex geometry, you will be interested in complex manifolds (or possibly “complex analytic varieties”), on which the functions are holomorphic functions. In each of these cases of interesting “geometric spaces”, the topological space and sheaf of functions is clear. The notion of scheme fits naturally into this family.

4.2 The underlying set of affine schemes

For any ring \(A\), we are going to define something called \(\text{Spec } A\), the spectrum of \(A\). In this section, we will define it as a set, but we will soon endow it with a topology, and later we will define a sheaf of rings on it (the structure sheaf). Such an object is called an affine scheme. Later \(\text{Spec } A\) will denote the set along with the topology, and a sheaf of functions. But for now, as there is no possibility of confusion, \(\text{Spec } A\) will just be the set.

4.2.1. The set \(\text{Spec } A\) is the set of prime ideals of \(A\). The point of \(\text{Spec } A\) corresponding to the prime ideal \(p\) will be denoted \([p]\). Elements \(a \in A\) will be called functions on \(\text{Spec } A\), and their value at the point \([p]\) will be \(a \mod p\). This is weird: a function can take values in different rings at different points — the function 5 on \(\text{Spec } \mathbb{Z}\) takes the value 1 (mod 2) at \([2]\) and 2 (mod 3) at \([3]\). “An element \(a\) of the ring lying in a prime ideal \(p\)” translates to “a function \(a\) that is 0 at the point \([p]\)” or “a function \(a\) vanishing at the point \([p]\)”, and we will use these phrases interchangeably. Notice that if you add or multiply two functions, you add or
multiply their values at all points; this is a translation of the fact that \( A \to A/p \) is a ring homomorphism. These translations are important — make sure you are very comfortable with them! They should become second nature.

We now give some examples.

**Example 1 (the complex affine line):** \( \mathbb{A}^1_C := \text{Spec} \mathbb{C}[x] \). Let’s find the prime ideals of \( \mathbb{C}[x] \). As \( \mathbb{C}[x] \) is an integral domain, \( 0 \) is prime. Also, \( (x - a) \) is prime, for any \( a \in \mathbb{C} \): it is even a maximal ideal, as the quotient by this ideal is a field:

\[
0 \longrightarrow (x - a) \longrightarrow \mathbb{C}[x] \xrightarrow{f \mapsto f(a)} \mathbb{C} \longrightarrow 0
\]

(This exact sequence may remind you of (3.1.1.1) in our motivating example of manifolds.)

We now show that there are no other prime ideals. We use the fact that \( \mathbb{C}[x] \) has a division algorithm, and is a unique factorization domain. Suppose \( p \) is a prime ideal. If \( p \neq (0) \), then suppose \( f(x) \in p \) is a non-zero element of smallest degree. It is not constant, as prime ideals can’t contain \( 1 \). If \( f(x) \) is not linear, then factor \( f(x) = g(x)h(x) \), where \( g(x) \) and \( h(x) \) have positive degree. (Here we use that \( \mathbb{C} \) is algebraically closed.) Then \( g(x) \in p \) or \( h(x) \in p \), contradicting the minimality of the degree of \( f \). Hence there is a linear element \( x - a \) of \( p \). Then I claim that \( p = (x - a) \). Suppose \( f(x) \in p \). Then the division algorithm would give \( f(x) = g(x)(x - a) + m \) where \( m \in \mathbb{C} \). Then \( m = f(x) - g(x)(x - a) \in p \). If \( m \neq 0 \), then \( 1 \in p \), giving a contradiction.

Thus we have a picture of \( \mathbb{A}^1_C = \text{Spec} \mathbb{C}[x] \) (see Figure 4.1). There is one point for each complex number, plus one extra point \([0]\). We can mostly picture \( \mathbb{A}^1_C \) as \( \mathbb{C} \): the point \([x - a]\) we will reasonably associate to \( a \in \mathbb{C} \). Where should we picture the point \([0]\)? The best way of thinking about it is somewhat zen. It is somewhere on the complex line, but nowhere in particular. Because \( [0] \) is contained in all of these primes, we will somehow associate it with this line passing through all the other points. \([0]\) is called the “generic point” of the line; it is “generically on the line” but you can’t pin it down any further than that. (We will formally define “generic point” in 4.3.) We will place it far to the right for lack of anywhere better to put it. You will notice that we sketch \( \mathbb{A}^1_C \) as one-(real-)dimensional (even though we picture it as an enhanced version of \( \mathbb{C} \)); this is to later remind ourselves that this will be a one-dimensional space, where dimensions are defined in an algebraic (or complex-geometric) sense. (Dimension will be defined in Chapter 12)

![Figure 4.1. A picture of \( \mathbb{A}^1_C = \text{Spec} \mathbb{C}[x] \)](image)

To give you some feeling for this space, we make some statements that are currently undefined, but suggestive. The functions on \( \mathbb{A}^1_C \) are the polynomials. So \( f(x) = x^2 - 3x + 1 \) is a function. What is its value at \([x - 1]\), which we think of as the point \( 1 \in \mathbb{C} \)? Answer: \( f(1) \)!

Or equivalently, we can evaluate \( f(x) \mod x - 1 \) — this is the same thing by the division algorithm. (What is its value at \([0]\)? It is \( f(x) \mod 0 \), which is just \( f(x) \).)
Here is a more complicated example: \( g(x) = \frac{(x - 3)^3}{(x - 2)} \) is a “rational function”. It is defined everywhere but \( x = 2 \). (When we know what the structure sheaf is, we will be able to say that it is an element of the structure sheaf on the open set \( \mathbb{A}^1_k - \{2\} \).) We want to say that \( g(x) \) has a triple zero at 3, and a single pole at 2, and we will be able to after §13.3.

**Example 2 (the affine line over \( k = \overline{k} \)):** \( \mathbb{A}^1_k := \text{Spec } k[x] \) where \( k \) is an algebraically closed field. This is called the affine line over \( k \). All of our discussion in the previous example carries over without change. We will use the same picture, which is after all intended to just be a metaphor.

**Example 3:** \( \text{Spec } \mathbb{Z} \). An amazing fact is that from our perspective, this will look a lot like the affine line \( \mathbb{A}^1_k \). The integers, like \( k[x] \), form a unique factorization domain, with a division algorithm. The prime ideals are: \((0)\), and \((p)\) where \( p \) is prime. Thus everything from Example 1 carries over without change, even the picture. Our picture of \( \text{Spec } \mathbb{Z} \) is shown in Figure 4.2.

![Figure 4.2](image)

**Figure 4.2.** A “picture” of \( \text{Spec } \mathbb{Z} \), which looks suspiciously like Figure 4.1

Let’s blithely carry over our discussion of functions to this space. 100 is a function on \( \text{Spec } \mathbb{Z} \). Its value at \( 3 \) is “1 \( \pmod{3} \)”. Its value at \( 2 \) is “0 \( \pmod{2} \)”, and in fact it has a double zero. \( 27/4 \) is a rational function on \( \text{Spec } \mathbb{Z} \), defined away from \( 2 \). We want to say that it has a double pole at \( 2 \), and a triple zero at \( 3 \). Its value at \( 5 \) is

\[
27 \times 4^{-1} \equiv 2 \times (-1) \equiv 3 \pmod{5}.
\]

**Example 4:** silly but important examples, and the German word for bacon. The set \( \text{Spec } k \) where \( k \) is any field is boring: one point. \( \text{Spec } 0 \), where \( 0 \) is the zero-ring, is the empty set, as \( 0 \) has no prime ideals.

**4.2.A. A small exercise about small schemes.** (a) Describe the set \( \text{Spec } k[\epsilon]/(\epsilon^2) \). The ring \( k[\epsilon]/(\epsilon^2) \) is called the ring of **dual numbers**, and will turn out to be quite useful. You should think of \( \epsilon \) as a very small number, so small that its square is 0 (although itself it is not 0). It is a non-zero function whose value at all points is zero, thus giving our first example of functions not being determined by their values at points. We will discuss this phenomenon further in §4.2.9.

(b) Describe the set \( \text{Spec } k[x]_{(x)} \) (see §2.3.3 for discussion of localization). We will see this scheme again repeatedly, starting with §4.2.6 and Exercise 4.4.J. You might later think of it as a shred of a particularly nice smooth curve.

In Example 2, we restricted to the case of algebraically closed fields for a reason: things are more subtle if the field is not algebraically closed.

**Example 5 (the affine line over \( \mathbb{R} \)):** \( \mathbb{R}[x] \). Using the fact that \( \mathbb{R}[x] \) is a unique factorization domain, similar arguments to those of Examples 1–3 show that the primes are \((0)\), \((x - a)\) where \( a \in \mathbb{R} \), and \((x^2 + ax + b)\) where \( x^2 + ax + b \) is an
irreducible quadratic. The latter two are maximal ideals, i.e. their quotients are fields. For example: \( \mathbb{R}[x]/(x - 3) \cong \mathbb{R}, \mathbb{R}[x]/(x^2 + 1) \cong \mathbb{C} \).

**4.2.B. Unimportant Exercise.** Show that for the last type of prime, of the form \((x^2 + ax + b),\) the quotient is always isomorphic to \(\mathbb{C}\).

So we have the points that we would normally expect to see on the real line, corresponding to real numbers; the generic point \(0\); and new points which we may interpret as conjugate pairs of complex numbers (the roots of the quadratic). This last type of point should be seen as more akin to the real numbers than to the generic point. You can picture \(\mathbb{A}^1_{\mathbb{R}}\) as the complex plane, folded along the real axis. But the key point is that Galois-conjugate points (such as \(i\) and \(-i\)) are considered glued.

Let’s explore functions on this space. Consider the function \(f(x) = x^3 - 1\). Its value at the point \([x - 2]\) is \(f(x) = 7\), or perhaps better, \(7 \pmod{x - 2}\). How about at \((x^2 + 1)\)? We get

\[
x^3 - 1 \equiv -x - 1 \pmod{x^2 + 1},
\]

which may be profitably interpreted as \(-i - 1\).

One moral of this example is that we can work over a non-algebraically closed field if we wish. It is more complicated, but we can recover much of the information we care about.

**4.2.C. Important Exercise.** Describe the set \(\mathbb{A}^1_{\mathbb{Q}}\). (This is harder to picture in a way analogous to \(\mathbb{A}^1_{\mathbb{R}}\). But the rough cartoon of points on a line, as in Figure 4.1 remains a reasonable sketch.)

**Example 6 (the affine line over \(\mathbb{F}_p\)):** \(\mathbb{A}^1_{\mathbb{F}_p} = \text{Spec} \mathbb{F}_p[x]\). As in the previous examples, \(\mathbb{F}_p[x]\) is a Euclidean domain, so the prime ideals are of the form \((0)\) or \((f(x))\) where \(f(x) \in \mathbb{F}_p[x]\) is an irreducible polynomial, which can be of any degree. Irreducible polynomials correspond to sets of Galois conjugates in \(\mathbb{F}_p\).

Note that \(\text{Spec} \mathbb{F}_p[x]\) has \(p\) points corresponding to the elements of \(\mathbb{F}_p\), but also (infinitely) many more. This makes this space much richer than simply \(p\) points. For example, a polynomial \(f(x)\) is not determined by its values at the \(p\) elements of \(\mathbb{F}_p\), but it is determined by its values at the points of \(\text{Spec} \mathbb{F}_p[x]\). (As we have mentioned before, this is not true for all schemes.)

You should think about this, even if you are a geometric person — this intuition will later turn up in geometric situations. Even if you think you are interested only in working over an algebraically closed field (such as \(\mathbb{C}\)), you will have non-algebraically closed fields (such as \(\mathbb{C}(x)\)) forced upon you.

**Example 7 (the complex affine plane):** \(\mathbb{A}^2_{\mathbb{C}} = \text{Spec} \mathbb{C}[x,y]\). (As with Examples 1 and 2, our discussion will apply with \(\mathbb{C}\) replaced by \(any\) algebraically closed field.) Sadly, \(\mathbb{C}[x,y]\) is not a principal ideal domain: \((x,y)\) is not a principal ideal. We can quickly name some prime ideals. One is \((0)\), which has the same flavor as the \((0)\) ideals in the previous examples. \((x - 2, y - 3)\) is prime, and indeed maximal, because \(\mathbb{C}[x,y]/(x - 2, y - 3) \cong \mathbb{C}\), where this isomorphism is via \(f(x,y) \mapsto f(2,3)\).

More generally, \((x - a, y - b)\) is prime for any \((a,b) \in \mathbb{C}^2\). Also, if \(f(x,y)\) is an irreducible polynomial (e.g. \(y - x^2\) or \(y^2 - x^3\) then \((f(x,y))\) is prime.
4.2.D. Exercise. (We will see a different proof of this in Exercise 12.2.3) Show that we have identified all the prime ideals of \( \mathbb{C}[x, y] \). Hint: Suppose \( p \) is a prime ideal that is not principal. Show you can find \( f(x, y), g(x, y) \in p \) with no common factor. By considering the Euclidean algorithm in the Euclidean domain \( k(x)[y] \), show that you can find a nonzero \( h(x) \in \{ f(x, y), g(x, y) \} \subset p \). Using primality, show that one of the linear factors of \( h(x) \), say \( (x - a) \), is in \( p \). Similarly show there is some \( (y - b) \in p \).

We now attempt to draw a picture of \( \mathbb{A}^2_{\mathbb{C}} \). The maximal primes of \( \mathbb{C}[x, y] \) correspond to the traditional points in \( \mathbb{C}^2 \): \( [(x - a, y - b)] \) corresponds to \( (a, b) \in \mathbb{C}^2 \). We now have to visualize the “bonus points”. \( [(0)] \) somehow lives behind all of the traditional points; it is somewhat on the plane, but nowhere in particular. So for example, it does not lie on the parabola \( y = x^2 \). The point \( [(y - x^2)] \) lies on the parabola \( y = x^2 \), but nowhere in particular on it. You can see from this picture that we already are implicitly thinking about “dimension”. The primes \( (x - a, y - b) \) are somehow of dimension \( 0 \), the primes \( (f(x, y)) \) are of dimension 1, and \( (0) \) is of dimension 2. (All of our dimensions here are complex or algebraic dimensions. The complex plane \( \mathbb{C}^2 \) has real dimension 4, but complex dimension 2. Complex dimensions are in general half of real dimensions.) We won’t define dimension precisely until Chapter 12 but you should feel free to keep it in mind before then.

Note too that maximal ideals correspond to the “smallest” points. Smaller ideals correspond to “bigger” points. “One prime ideal contains another” means that the points “have the opposite containment.” All of this will be made precise once we have a topology. This order-reversal is a little confusing, and will remain so even once we have made the notions precise.

We now come to the obvious generalization of Example 7.

Example 8 (complex affine \( n \)-space): \( \mathbb{A}^n_{\mathbb{C}} := \text{Spec} \mathbb{C}[x_1, \ldots, x_n] \). (More generally, \( \mathbb{A}^n_A := \text{Spec} A[x_1, \ldots, x_n] \), where \( A \) is an arbitrary ring.) For concreteness, let’s consider \( n = 3 \). We now have an interesting question in what at first appears to be pure algebra: What are the prime ideals of \( \mathbb{C}[x, y, z] \)?

Analogously to before, \( (x - a, y - b, z - c) \) is a prime ideal. This is a maximal ideal, with residue field \( \mathbb{C} \); we think of these as “0-dimensional points”. We will often write \( (a, b, c) \) for \( [(x - a, y - b, z - c)] \) because of our geometric interpretation of these ideals. There are no more maximal ideals, by Hilbert’s Weak Nullstellensatz.

4.2.2. Hilbert’s Weak Nullstellensatz. — If \( k \) is an algebraically closed field, then the maximal ideals \( k[x_1, \ldots, x_n] \), are precisely those of the form \( (x_1 - a_1, \ldots, x_n - a_n) \), where \( a_i \in k \).

We may as well state a slightly stronger version now.

4.2.3. Hilbert’s Nullstellensatz. — If \( k \) is any field, the maximal ideals of \( k[x_1, \ldots, x_n] \) are precisely those with residue field a finite extension of \( k \).

The Nullstellensatz 4.2.3 clearly implies the Weak Nullstellensatz 4.2.2. You will prove the Nullstellensatz in Exercise 12.2.3.

There are other prime ideals of \( \mathbb{C}[x, y, z] \) too. We have \( (0) \), which is corresponds to a “3-dimensional point”. We have \( (f(x, y, z)) \), where \( f \) is irreducible. To this we associate the hypersurface \( f = 0 \), so this is “2-dimensional” in nature. But we have not found them all! One clue: we have prime ideals of “dimension” 0,
2, and 3 — we are missing “dimension 1”. Here is one such prime ideal: \((x,y)\).
We picture this as the locus where \(x = y = 0\), which is the \(z\)-axis. This is a
prime ideal, as the corresponding quotient \(\mathbb{C}[x,y,z]/(x,y) \cong \mathbb{C}[z]\) is an integral
domain (and should be interpreted as the functions on the \(z\)-axis). There are lots
of one-dimensional primes, and it is not possible to classify them in a reasonable
way. It will turn out that they correspond to things that we think of as irreducible
curves. Thus remarkably the answer to the purely algebraic question (“what are
the primes of \(\mathbb{C}[x,y,z]\)” is fundamentally geometric!

The fact that the closed points \(\mathbb{A}^1_2\) can be interpreted as points of \(\mathbb{Q}\) where
Galois-conjugates are glued together (Exercise 4.2.C) extends to \(\mathbb{A}^n_2\). For example,
in \(\mathbb{A}^2_2\), \((\sqrt{2}, \sqrt{2})\) is glued to \((-\sqrt{2}, -\sqrt{2})\) but not to \((\sqrt{2}, -\sqrt{2})\). The following exercise
will give you some idea of how this works.

4.2.E. Exercise. Describe the maximal ideal of \(Q[x,y]\) corresponding to \((\sqrt{2}, \sqrt{2})\)
and \((-\sqrt{2}, -\sqrt{2})\). Describe the maximal ideal of \(Q[x,y]\) corresponding to \((\sqrt{2}, -\sqrt{2})\)
and \((-\sqrt{2}, \sqrt{2})\). What are the residue fields in both cases?

The description of closed points of \(\mathbb{A}^2_2\) (and its generalizations) as Galois-orbits
can even be extended to non-closed points, as follows.

4.2.F. Unimportant but fun exercise. Consider the map of sets \(\phi : \mathbb{C}^2 \to \mathbb{A}^2_2\)
defined as follows. \((z_1, z_2)\) is sent to the prime ideal of \(Q[x,y]\) consisting of
polynomials vanishing at \((z_1, z_2)\). (a) What is the image of \((\pi, \pi^2)\)? (b) Show that
\(\phi\) is surjective. (Once we define the Zariski topology on \(\mathbb{A}^2_2\), you will be able to
check that \(\phi\) is continuous, where we give \(\mathbb{C}^2\) the classical topology. This example
generalizes.)

4.2.4. Quotients and localizations. Two natural ways of getting new rings from
old — quotients and localizations — have interpretations in terms of spectra.

4.2.5. Quotients: \(\text{Spec } A/I\) as a subset of \(\text{Spec } A\). It is an important fact that the
primes of \(A/I\) are in bijection with the primes of \(A\) containing \(I\).

4.2.G. Essential algebra exercise (mandatory if you haven’t seen it be-
fore). Suppose \(A\) is a ring, and \(I\) an ideal of \(A\). Let \(\phi : A \to A/I\). Show that
\(\phi^{-1}\) gives an inclusion-preserving bijection between primes of \(A/I\) and primes of
\(A\) containing \(I\). Thus we can picture \(\text{Spec } A/I\) as a subset of \(\text{Spec } A\).

As an important motivational special case, you now have a picture of complex
affine varieties. Suppose \(A\) is a finitely generated \(\mathbb{C}\)-algebra, generated by \(x_1, \ldots, x_n\),
with relations \(f_1(x_1, \ldots, x_n) = \cdots = f_r(x_1, \ldots, x_n) = 0\). Then this description
in terms of generators and relations naturally gives us an interpretation of \(\text{Spec } A\)
as a subset of \(\mathbb{A}^n_\mathbb{C}\), which we think of as “traditional points” (n-tuples of complex
numbers) along with some “bonus” points we haven’t yet fully described. To see
which of the traditional points are in \(\text{Spec } A\), we simply solve the equations \(f_1 =
\cdots = f_r = 0\). For example, \(\text{Spec } \mathbb{C}[x,y,z]/(x^2+y^2-z^2)\) may be pictured as shown in
Figure 4.3 (Admittedly this is just a “sketch of the \(\mathbb{R}\)-points”, but we will still find
it helpful later.) This entire picture carries over (along with the Nullstellensatz)
with \(\mathbb{C}\) replaced by any algebraically closed field. Indeed, the picture of Figure 4.3
can be said to depict \(k[x,y,z]/(x^2+y^2-z^2)\) for most algebraically closed fields \(k\).
(although it is misleading in characteristic 2, because of the coincidence \( x^2 + y^2 - z^2 = (x + y + z)^2 \)).

\[ (x^2 + y^2 - z^2) = (x + y + z)^2. \]

**Figure 4.3.** A “picture” of \( \text{Spec} \mathbb{C}[x, y, z]/(x^2 + y^2 - z^2) \)

### 4.2.6. Localizations: \( \text{Spec} S^{-1}A \) as a subset of \( \text{Spec} A \)

The following exercise shows how prime ideals behave under localization.

**4.2.H. Essential Algebra Exercise (Mandatory if you haven’t seen it before).** Suppose \( S \) is a multiplicative subset of \( A \). The map \( \text{Spec} S^{-1}A \to \text{Spec} A \) gives an order-preserving bijection of the primes of \( S^{-1}A \) with the primes of \( A \) that don’t meet the multiplicative set \( S \).

Recall from [§2.3.3](#) that there are two important flavors of localization. The first is \( A_f = \{1, f, f^2, \ldots\}^{-1}A \) where \( f \in A \). A motivating example is \( A = \mathbb{C}[x, y] \), \( f = y - x^2 \). The second is \( A_p = (A - p)^{-1}A \), where \( p \) is a prime ideal. A motivating example is \( A = \mathbb{C}[x, y], S = A - (x, y) \).

If \( S = \{1, f, f^2, \ldots\} \), the primes of \( S^{-1}A \) are just those primes not containing \( f \) — the points where “\( f \) doesn’t vanish”. (In [§4.5](#), we will call this a distinguished open set, once we know what open sets are.) So to picture \( \text{Spec} \mathbb{C}[x, y]/(y - x^2) \), we picture the affine plane, and throw out those points on the parabola \( y - x^2 \) — the points \((a, a^2)\) for \( a \in \mathbb{C} \) (by which we mean \([x - a, y - a^2])\), as well as the “new kind of point” \(([y - x^2])\).

It can be initially confusing to think about localization in the case where zero-divisors are inverted, because localization \( A \to S^{-1}A \) is not injective (Exercise [2.3.C](#)). Geometric intuition can help. Consider the case \( A = \mathbb{C}[x, y]/(xy) \) and \( f = x \). What is the localization \( A_f \)? The space Spec \( \mathbb{C}[x, y]/(xy) \) “is” the union of the two axes in the plane. Localizing means throwing out the locus where \( x \) vanishes. So we are left with the \( x \)-axis, minus the origin, so we expect \( \text{Spec} \mathbb{C}[x]_x \). So there should be some natural isomorphism \( (\mathbb{C}[x, y]/(xy))_x \cong \mathbb{C}[x]_x \).

**4.2.I. Exercise.** Show that these two rings are isomorphic. (You will see that \( y \) on the left goes to 0 on the right.)

If \( S = A - p \), the primes of \( S^{-1}A \) are just the primes of \( A \) contained in \( p \). In our example \( A = \mathbb{C}[x, y], p = (x, y) \), we keep all those points corresponding to “things through the origin”, i.e. the 0-dimensional point \((x, y)\), the 2-dimensional point \((0)\),
and those 1-dimensional points \((f(x, y))\) where \(f(0, 0) = 0\), i.e. those “irreducible curves through the origin”. You can think of this being a shred of the plane near the origin; anything not actually “visible” at the origin is discarded (see Figure 4.4).

Another example is when \(A = \text{Spec} \, k[x]\), and \(p = (x)\) (or more generally when \(p\) is any maximal ideal). Then \(A_p\) has only two prime ideals (Exercise 4.2.A(b)). You should see this as the germ of a “smooth curve”, where one point is the “classical point”, and the other is the “generic point of the curve”. This is an example of a discrete valuation ring, and indeed all discrete valuation rings should be visualized in such a way. We will discuss discrete valuation rings in §13.3. By then we will have justified the use of the words “smooth” and “curve”. (Reality check: try to picture \(\text{Spec} \, \mathbb{Z}\) localized at \((2)\) and at \((0)\). How do the two pictures differ?)

4.2.7. Important fact: Maps of rings induce maps of spectra (as sets). We now make an observation that will later grow up to be the notion of morphisms of schemes.

4.2.J. IMPORTANT EASY EXERCISE. If \(\phi : B \to A\) is a map of rings, and \(p\) is a prime ideal of \(A\), show that \(\phi^{-1}(p)\) is a prime ideal of \(B\).

Hence a map of rings \(\phi : B \to A\) induces a map of sets \(\text{Spec} \, A \to \text{Spec} \, B\) “in the opposite direction”. This gives a contravariant functor from the category of rings to the category of sets: the composition of two maps of rings induces the composition of the corresponding maps of spectra.

4.2.K. EASY EXERCISE. Let \(B\) be a ring.
(a) Suppose \(I \subset B\) is an ideal. Show that the map \(\text{Spec} \, B/I \to \text{Spec} \, B\) is the inclusion of §4.2.5.
(b) Suppose \(S \subset B\) is a multiplicative set. Show that the map \(\text{Spec} \, S^{-1}B \to \text{Spec} \, B\) is the inclusion of §4.2.6.

4.2.8. An explicit example. In the case of affine complex varieties (or indeed affine varieties over any algebraically closed field), the translation between maps given
by explicit formulas and maps of rings is quite direct. For example, consider a
map from the parabola in $\mathbb{C}^2$ (with coordinates $a$ and $b$) given by $b = a^2$, to the
“curve” in $\mathbb{C}^3$ (with coordinates $x$, $y$, and $z$) cut out by the equations $y = x^2$ and
$z = y^2$. Suppose the map sends the point $(a, b) \in \mathbb{C}^2$ to the point $(a, b, b^2) \in \mathbb{C}^3$. In
our new language, we have map

$$\text{Spec } \mathbb{C}[a, b]/(b - a^2) \longrightarrow \text{Spec } \mathbb{C}[x, y, z]/(y - x^2, z - y^2)$$

given by

$$\mathbb{C}[a, b]/(b - a^2) \hookrightarrow \mathbb{C}[x, y, z]/(y - x^2, z - y^2)$$

$$(a, b, b^2) \hookrightarrow (x, y, z),$$
i.e. $x \mapsto a$, $y \mapsto b$, and $z \mapsto b^2$. If the idea is not yet clear, the following two
exercises may help.

\section*{Figure 4.5. The map $\mathbb{C} \to \mathbb{C}$ given by $y \mapsto y^2$}

\subsection*{4.2.L. Exercise (Special Case)} Consider the map of complex manifolds send-
ing $\mathbb{C} \to \mathbb{C}$ via $y \mapsto y^2$; you can picture it as the projection of the parabola $x = y^2$
in the plane to the $x$-axis (see Figure 4.5). Interpret the corresponding map of rings
as given by $\mathbb{C}[x] \hookrightarrow \mathbb{C}[y]$ by $x \mapsto y^2$. Verify that the preimage (the fiber) above the
point $a \in \mathbb{C}$ is the point(s) $\pm \sqrt{a} \in \mathbb{C}$, using the definition given above. (A more
sophisticated version of this example appears in Example 10.3.3)

\subsection*{4.2.M. Exercise (General Case)} (a) Show that the map

$$\phi : \{y_1, y_2, \ldots, y_n\} \mapsto \{f_1(x_1, \ldots, x_m), f_2(x_1, \ldots, x_m), \ldots, f_n(x_1, \ldots, x_m)\}$$

determines a map

$$\text{Spec } \mathbb{C}[x_1, \ldots, x_m]/I \to \text{Spec } \mathbb{C}[y_1, \ldots, y_n]/J$$

if $\phi(J) \subset I$.

(b) Via the identification of the Nullstellensatz, interpret the map of (a) as a map

$$\mathbb{C}^m \to \mathbb{C}^n$$
given by

$$(x_1, \ldots, x_m) \mapsto \{f_1(x_1, \ldots, x_m), \ldots, f_n(x_1, \ldots, x_m)\}.$$

The converse to (a) isn’t quite true. Once you have more experience and intu-
tion, you can figure out when it is true, and when it can be false. The failure of the
converse to hold has to do with nilpotents, which we come to very shortly (§4.2.9).
4.2.N. Important Exercise. Consider the map of sets $f : \mathbb{A}^n_Z \to \text{Spec } \mathbb{Z}$, given by the ring map $\mathbb{Z} \to \mathbb{Z}[x_1, \ldots, x_n]$. If $p$ is prime, describe a bijection between the fiber $f^{-1}([p])$ and $\mathbb{A}^n_{\mathbb{F}_p}$. (You won’t need to describe either set! Which is good because you can’t.) This exercise may give you a sense of how to picture maps (see Figure 4.6), and in particular why you can think of $\mathbb{A}^n_Z$ as an “$\mathbb{A}^n$-bundle” over $\text{Spec } \mathbb{Z}$. (Can you interpret the fiber over $[0]$ as $\mathbb{A}^n_k$ for some field $k$?)

![Figure 4.6](image)

Figure 4.6. A picture of $\mathbb{A}^n_Z \to \text{Spec } \mathbb{Z}$ as a “family of $\mathbb{A}^n$’s”, or an “$\mathbb{A}^n$-bundle over $\text{Spec } \mathbb{Z}$”. What is $k$?

4.2.9. Functions are not determined by their values at points: the fault of nilpotents. We conclude this section by describing some strange behavior. We are developing machinery that will let us bring our geometric intuition to algebra. There is one serious serious point where your intuition will be false, so you should know now, and adjust your intuition appropriately. As noted by Mumford ([M-CAS p. 12]), “it is this aspect of schemes which was most scandalous when Grothendieck defined them.”

Suppose we have a function (ring element) vanishing at all points. Then it is not necessarily the zero function! The translation of this question is: is the intersection of all prime ideals necessarily just 0? The answer is no, as is shown by the example of the ring of dual numbers $k[\epsilon]/(\epsilon^2)$: $\epsilon \neq 0$, but $\epsilon^2 = 0$. (We saw this scheme in Exercise 4.2.A(a).) Any function whose power is zero certainly lies in the intersection of all prime ideals.

4.2.O. Exercise. Ring elements that have a power that is 0 are called nilpotents. (a) If $I$ is an ideal of nilpotents, show that the inclusion $\text{Spec } B/I \to \text{Spec } B$ of Exercise 4.2.C is a bijection. Thus nilpotents don’t affect the underlying set. (We will soon see in 4.4.5 that they won’t affect the topology either — the difference will be in the structure sheaf.) (b) (easy) Show that the nilpotents of a ring $B$ form an ideal. This ideal is called the nilradical, and is denoted $\mathfrak{n}$.

Thus the nilradical is contained in the intersection of all the prime ideals. The converse is also true:

4.2.10. Theorem. — The nilradical $\mathfrak{n}(A)$ is the intersection of all the primes of $A$. 


4.2.P. EXERCISE. If you don’t know this theorem, then look it up, or even better, prove it yourself. (Hint: Use the fact that any proper ideal of \( A \) is contained in a maximal ideal, which requires the axiom of choice. Possible further hint: Suppose \( x \notin \mathfrak{m}(A) \). We wish to show that there is a prime ideal not containing \( x \). Show that \( A_x \) is not the 0-ring, by showing that \( 1 \neq 0 \).)

4.2.11. In particular, although it is upsetting that functions are not determined by their values at points, we have precisely specified what the failure of this intuition is: two functions have the same values at points if and only if they differ by a nilpotent. You should think of this geometrically: a function vanishes at every point of the spectrum of a ring if and only if it has a power that is zero. And if there are no non-zero nilpotents — if \( \mathfrak{m} = \{0\} \) — then functions are determined by their values at points. If a ring has no non-zero nilpotents, we say that it is reduced.

4.2.Q. FUN UNIMPORTANT EXERCISE: DERIVATIVES WITHOUT DELTAS AND EPSILONS (OR AT LEAST WITHOUT DELTAS). Suppose we have a polynomial \( f(x) \in k[x] \). Instead, we work in \( k[x, \epsilon]/\epsilon^2 \). What then is \( f(x + \epsilon) \)? (Do a couple of examples, then prove the pattern you observe.) This is a hint that nilpotents will be important in defining differential information (Chapter 22).

4.3 Visualizing schemes I: generic points

For years, you have been able to picture \( x^2 + y^2 = 1 \) in the plane, and you now have an idea of how to picture \( \text{Spec } \mathbb{Z} \). If we are claiming to understand rings as geometric objects (through the Spec functor), then we should wish to develop geometric insight into them. To develop geometric intuition about schemes, it is helpful to have pictures in your mind, extending your intuition about geometric spaces you are already familiar with. As we go along, we will empirically develop some idea of what schemes should look like. This section summarizes what we have gleaned so far.

Some mathematicians prefer to think completely algebraically, and never think in terms of pictures. Others will be disturbed by the fact that this is an art, not a science. And finally, this hand-waving will necessarily never be used in the rigorous development of the theory. For these reasons, you may wish to skip these sections. However, having the right picture in your mind can greatly help understanding what facts should be true, and how to prove them.

Our starting point is the example of “affine complex varieties” (things cut out by equations involving a finite number variables over \( \mathbb{C} \)), and more generally similar examples over arbitrary algebraically closed fields. We begin with notions that are intuitive (“traditional” points behaving the way you expect them to), and then add in the two features which are new and disturbing, generic points and non-reduced behavior. You can then extend this notion to seemingly different spaces, such as \( \text{Spec } \mathbb{Z} \).

Hilbert’s Weak Nullstellensatz [4.2.2] shows that the “traditional points” are present as points of the scheme, and this carries over to any algebraically closed
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field. If the field is not algebraically closed, the traditional points are glued together into clumps by Galois conjugation, as in Examples 5 (the real affine line) and 6 (the affine line over \( \mathbb{F}_p \)) in §4.2 above. This is a geometric interpretation of Hilbert’s Nullstellensatz [4.2.3].

But we have some additional points to add to the picture. You should remember that they “correspond” to “irreducible” “closed” (algebraic) subsets. As motivation, consider the case of the complex affine plane (Example 7): we had one for each irreducible polynomial, plus one corresponding to the entire plane. We will make “closed” precise when we define the Zariski topology (in the next section). You may already have an idea of what “irreducible” should mean; we make that precise at the start of §4.6. By “correspond” we mean that each closed irreducible subset has a corresponding point sitting on it, called its generic point (defined in §4.6). It is a new point, distinct from all the other points in the subset. The correspondence is described in Exercise 4.7.E for Spec \( A \), and in Exercise 6.1.B for schemes in general. We don’t know precisely where to draw the generic point, so we may stick it arbitrarily anywhere, but you should think of it as being “almost everywhere”, and in particular, near every other point in the subset.

In §4.2.5, we saw how the points of Spec \( A/I \) should be interpreted as a subset of Spec \( A \). So for example, when you see Spec \( \mathbb{C}[x, y]/(x + y) \), you should picture this not just as a line, but as a line in the \( xy \)-plane; the choice of generators \( x \) and \( y \) of the algebra \( \mathbb{C}[x, y] \) implies an inclusion into affine space.

In §4.2.6, we saw how the points of Spec \( S^{-1}A \) should be interpreted as subsets of Spec \( A \). The two most important cases were discussed. The points of Spec \( A_f \) correspond to the points of Spec \( A \) where \( f \) doesn’t vanish; we will later (§4.5) interpret this as a distinguished open set.

If \( p \) is a prime ideal, then Spec \( A_p \) should be seen as a “shred of the space Spec \( A \) near the subset corresponding to \( p \)”. The simplest nontrivial case of this is \( p = (x) \subset \text{Spec } k[x] = A \) (see Exercise 4.2.A which we discuss again in Exercise 4.4.J).

### 4.4 The Zariski topology: The underlying topological space of an affine scheme

We next introduce the Zariski topology on the spectrum of a ring. For example, consider \( \mathbb{A}^2_k = \text{Spec } \mathbb{C}[x, y] \), the complex plane (with a few extra points). In algebraic geometry, we will only be allowed to consider algebraic functions, i.e. polynomials in \( x \) and \( y \). The locus where a polynomial vanishes should reasonably be a closed set, and the Zariski topology is defined by saying that the only sets we should consider closed should be these sets, and other sets forced to be closed by these. In other words, it is the coarsest topology where these sets are closed.

In particular, although topologies are often described using open subsets, it will more convenient for us to define this topology in terms of closed subsets. If \( S \) is a subset of a ring \( A \), define the Vanishing set of \( S \) by

\[
V(S) := \{ [p] \in \text{Spec } A \mid S \subseteq p \}.
\]
It is the set of points on which all elements of $S$ are zero. (It should now be second nature to equate “vanishing at a point” with “contained in a prime”.) We declare that these — and no other — are the closed subsets.

For example, consider $V(xy, yz) \subset \mathbb{A}^3 = \text{Spec } \mathbb{C}[x, y, z]$. Which points are contained in this locus? We think of this as solving $xy = yz = 0$. Of the “traditional” points (interpreted as ordered triples of complex numbers, thanks to the Hilbert’s Nullstellensatz 4.2.2), we have the points where $y = 0$ or $x = z = 0$: the $xz$-plane and the $y$-axis respectively. Of the “new” points, we have the generic point of the $xz$-plane (also known as the point $[(y)]$), and the generic point of the $y$-axis (also known as the point $[(x, z)]$). You might imagine that we also have a number of “one-dimensional” points contained in the $xz$-plane.

4.4.A. Easy Exercise. Check that the $x$-axis is contained in $V(xy, yz)$.

Let’s return to the general situation. The following exercise lets us restrict attention to vanishing sets of ideals.

4.4.B. Easy Exercise. Show that if $(S)$ is the ideal generated by $S$, then $V(S) = V((S))$.

We define the Zariski topology by declaring that $V(S)$ is closed for all $S$. Let’s check that this is a topology:

4.4.C. Exercise. (a) Show that $\emptyset$ and $\text{Spec } A$ are both open.
(b) If $I_i$ is a collection of ideals (as $i$ runs over some index set), show that $\cap_i V(I_i) = V(\sum_i I_i)$. Hence the union of any collection of open sets is open.
(c) Show that $V(I_1) \cup V(I_2) = V(I_1 I_2)$. Hence the intersection of any finite number of open sets is open.

4.4.1. Properties of the “vanishing set” function $V(\cdot)$. The function $V(\cdot)$ is obviously inclusion-reversing: If $S_1 \subset S_2$, then $V(S_2) \subset V(S_1)$. Warning: We could have equality in the second inclusion without equality in the first, as the next exercise shows.

4.4.D. Exercise/Definition. If $I \subset R$ is an ideal, then define its radical by

$$\sqrt{I} := \{r \in R : r^n \in I \text{ for some } n \in \mathbb{Z}^\geq 0\}.$$ 

For example, the nilradical $\mathfrak{n}$ ([4.2.2]) is $\sqrt{(0)}$. Show that $V(\sqrt{I}) = V(I)$. We say an ideal is radical if it equals its own radical.

Here are two useful consequences. As $(I \cap J)^2 \subset IJ \subset I \cap J$, we have that $V(IJ) = V(I \cap J)$ (= $V(I) \cup V(J)$ by Exercise 4.4.C(b)). Also, combining this with Exercise 4.4.B we see $V(S) = V((S)) = V(\sqrt{S})$.

4.4.E. Exercise (Radicals commute with finite intersection). If $I_1, \ldots, I_n$ are ideals of a ring $A$, show that $\sqrt{\cap_{i=1}^n I_i} = \cap_{i=1}^n \sqrt{I_i}$. We will use this property without referring back to this exercise.

4.4.F. Exercise for later use. Show that $\sqrt{I}$ is the intersection of all the prime ideals containing $I$. (Hint: Use Theorem 4.2.10 on an appropriate ring.)
4.4.2. Examples. Let’s see how this meshes with our examples from the previous section.

Recall that $\mathbb{A}^1_\mathbb{C}$, as a set, was just the “traditional” points (corresponding to maximal ideals, in bijection with $a \in \mathbb{C}$), and one “new” point $\{0\}$. The Zariski topology on $\mathbb{A}^1_\mathbb{C}$ is not that exciting: the open sets are the empty set, and $\mathbb{A}^1_\mathbb{C}$ minus a finite number of maximal ideals. (It “almost” has the cofinite topology. Notice that the open sets are determined by their intersection with the “traditional points”. The “new” point $\{0\}$ comes along for the ride, which is a good sign that it is harmless. Ignoring the “new” point, observe that the topology on $\mathbb{A}^1_\mathbb{C}$ is a coarser topology than the classical topology on $\mathbb{C}$.)

The case Spec $\mathbb{Z}$ is similar. The topology is “almost” the cofinite topology in the same way. The open sets are the empty set, and Spec $\mathbb{Z}$ minus a finite number of “ordinary” $(p)$ where $p$ is prime) primes.

4.4.3. Closed subsets of $\mathbb{A}^2_\mathbb{C}$. The case $\mathbb{A}^2_\mathbb{C}$ is more interesting. You should think through where the “one-dimensional primes” fit into the picture. In Exercise 4.2.D we identified all the primes of $\mathbb{C}[x, y]$ (i.e. the points of $\mathbb{A}^2_\mathbb{C}$) as the maximal ideals $(x-a, y-b)$ (where $a, b \in \mathbb{C}$), the “one-dimensional points” $(f(x, y))$ (where $f(x, y)$ is irreducible), and the “two-dimensional point” $\{0\}$.

Then the closed subsets are of the following form:

(a) the entire space, and

(b) a finite number (possibly zero) of “curves” (each of which is the closure of a “one-dimensional point”) and a finite number (possibly zero) of closed points.

4.4.4. Important fact: Maps of rings induce continuous maps of topological spaces. We saw in §4.2.7 that a map of rings $\phi : B \to A$ induces a map of sets $\pi : \text{Spec } A \to \text{Spec } B$.

4.4.G. Important exercise. By showing that closed sets pull back to closed sets, show that $\pi$ is a continuous map.

Not all continuous maps arise in this way. Consider for example the continuous map on $\mathbb{A}^1_\mathbb{C}$ that is the identity except $0$ and $1$ (i.e. $[(x)]$ and $[(x-1)]$) are swapped; no polynomial can manage this marvellous feat.

In §4.2.7 we saw that Spec $B/I$ and Spec $S^{-1}B$ are naturally subsets of Spec $B$. It is natural to ask if the Zariski topology behaves well with respect to these inclusions, and indeed it does.

4.4.H. Important exercise (cf. Exercise 4.2.K). Suppose that $I, S \subseteq B$ are an ideal and multiplicative subset respectively. Show that Spec $B/I$ is naturally a closed subset of Spec $B$. Show that the Zariski topology on Spec $B/I$ (resp. Spec $S^{-1}B$) is the subspace topology induced by inclusion in Spec $B$. (Hint: compare closed subsets.)

4.4.5. In particular, if $I \subseteq \mathfrak{m}$ is an ideal of nilpotents, the bijection Spec $B/I \to$ Spec $B$ (Exercise 4.2.C) is a homeomorphism. Thus nilpotents don’t affect the topological space. (The difference will be in the structure sheaf.)
4.4.I. Useful exercise for later. Suppose $I \subset B$ is an ideal. Show that $f$ vanishes on $V(I)$ if and only if $f \in \sqrt{I}$ (i.e. $f^n \in I$ for some $n$). (If you are stuck, you will get a hint when you see Exercise 4.5.E.)

4.4.J. Easy Exercise (cf. Exercise 4.2.A). Describe the topological space $\text{Spec } k[x](x)$.

4.5 A base of the Zariski topology on $\text{Spec } A$: Distinguished open sets

If $f \in A$, define the distinguished open set $D(f) = \{|p| \in \text{Spec } A : f \notin p\}$. It is the locus where $f$ doesn't vanish. (I often privately write this as $D(f \neq 0)$ to remind myself of this. I also privately call this a “Doesn't-vanish set” in analogy with $V(f)$ being the Vanishing set.) We have already seen this set when discussing $\text{Spec } A_f$ as a subset of $\text{Spec } A$. For example, we have observed that the Zariski-topology on the distinguished open set $D(f) \subset \text{Spec } A$ coincides with the Zariski topology on $\text{Spec } A_f$ (Exercise 4.4.I).

The reason these sets are important is that they form a particularly nice base for the (Zariski) topology:

4.5.A. Easy Exercise. Show that the distinguished open sets form a base for the (Zariski) topology. (Hint: Given a subset $S \subset A$, show that the complement of $V(S)$ is $\bigcup_{n \in S} D(f)$.)

Here are some important but not difficult exercises to give you a feel for this concept.

4.5.B. Exercise. Suppose $f_i \in A$ as $i$ runs over some index set $I$. Show that $\bigcup_{i \in I} D(f_i) = \text{Spec } A$ if and only if $(f_i) = A$, or equivalently and very usefully, there are $a_i$ ($i \in I$), all but finitely many 0, such that $\sum_{i \in I} a_if_i = 1$. (One of the directions will use the fact that any proper ideal of $A$ is contained in some maximal ideal.)

4.5.C. Exercise. Show that if $\text{Spec } A$ is an infinite union of distinguished open sets $\bigcup_{i \in I} D(f_i)$, then in fact it is a union of a finite number of these, i.e. there is a finite subset $J'$ so that $\text{Spec } A = \bigcup_{j \in J'} D(f_j)$. (Hint: exercise 4.5.B)

4.5.D. Easy Exercise. Show that $D(f) \cap D(g) = D(fg)$.

4.5.E. Important Exercise (cf. Exercise 4.4.I). Show that $D(f) \subset D(g)$ if and only if $f^n \in (g)$ for some $n$, if and only if $g$ is a unit in $A_f$.

We will use Exercise 4.5.E often. You can solve it thinking purely algebraically, but the following geometric interpretation may be helpful. Inside $\text{Spec } A$, we have the closed subset $V(g) = \text{Spec } A/(g)$, where $g$ vanishes, and its complement $D(g)$, where $g$ doesn't vanish. Then $f$ is a function on this closed subset $V(g)$ (or more precisely, on $\text{Spec } A/(g)$), and by assumption it vanishes at all points of the closed subset. Now any function vanishing at every point of the spectrum of a ring must be nilpotent (Theorem 4.2.10). In other words, there is some $n$ such that $f^n = 0$ in $A/(g)$, i.e. $f^n \equiv 0 \pmod{g}$ in $A$, i.e. $f^n \in (g)$.
4.5.F. **Easy Exercise.** Show that $D(f) = \emptyset$ if and only if $f \in \mathfrak{m}$.

**4.6 Topological definitions**

A topological space is said to be **irreducible** if it is nonempty, and it is not the union of two proper closed subsets. In other words, $X$ is irreducible if whenever $X = Y \cup Z$ with $Y$ and $Z$ closed, we have $Y = X$ or $Z = X$.

4.6.A. **Easy Exercise.** Show that in an irreducible topological space, any nonempty open set is dense. (The moral: unlike in the classical topology, in the Zariski topology, non-empty open sets are all “huge”.)

4.6.B. **Easy exercise.** If $A$ is an integral domain, show that $\text{Spec } A$ is irreducible. (Hint: pay attention to the generic point $([0])$.)

A point of a topological space $x \in X$ is said to be **closed** if $\{x\}$ is a closed subset. In the classical topology on $\mathbb{C}^n$, all points are closed.

4.6.C. **Exercise.** Show that the closed points of $\text{Spec } A$ correspond to the maximal ideals.

Thus Hilbert’s Nullstellensatz lets us interpret the closed points of $\mathbb{A}^n_\mathbb{C}$ as the $n$-tuples of complex numbers. Hence from now on we will say “closed point” instead of “traditional point” and “non-closed point” instead of “bonus” or “new-fangled” point when discussing subsets of $\mathbb{A}^n_\mathbb{C}$.

4.6.1. **Quasicompactness.** A topological space $X$ is **quasicompact** if given any cover $X = \bigcup_{i \in I} U_i$ by open sets, there is a finite subset $S$ of the index set $I$ such that $X = \bigcup_{i \in S} U_i$. Informally: every cover has a finite subcover. Depending on your definition of “compactness”, this is the definition of compactness, minus possibly a Hausdorff condition. We will like this condition, because we are afraid of infinity.

4.6.D. **Exercise.** (a) Show that $\text{Spec } A$ is quasicompact. (Hint: Exercise 4.5.C.)

(b) Show that in general $\text{Spec } A$ can have nonquasicompact open sets. (Possible hint: let $A = k[x_1, x_2, x_3, \ldots]$ and $m = (x_1, x_2, \ldots) \subset A$, and consider the complement of $V(m)$. This example will be useful to construct other enlightening examples later, e.g. Exercises 8.1.B and 8.3.E. In Exercise 4.6.M, we see that such weird behavior doesn’t happen for “suitably nice” (Noetherian) rings.)

4.6.E. **Exercise.** (a) If $X$ is a topological space that is a finite union of quasicompact spaces, show that $X$ is quasicompact.

(b) Show that every closed subset of a quasicompact topological space is quasicompact.

4.6.2. **Specialization and generization.** Given two points $x, y$ of a topological space $X$, we say that $x$ is a **specialization** of $y$, and $y$ is a **generization** of $x$, if $x \in \overline{\{y\}}$. This now makes precise our hand-waving about “one point containing another”. It is of course nonsense for a point to contain another. But it is not nonsense to say that the closure of a point contains another. For example, in $\mathbb{A}^2_\mathbb{C} = \text{Spec } \mathbb{C}[x, y]$, ...
\[([y - x^2]) \text{ is a generalization of } (2, 4) = [(x - 2, y - 4)], \text{ and } (2, 4) \text{ is a specialization of } ([y - x^2]).\]

**4.6.F. Exercise.** If \(X = \text{Spec } A\), show that \([p]\) is a specialization of \([q]\) if and only if \(q \subset p\).

We say that a point \(x \in X\) is a **generic point** for a closed subset \(K\) if \(\{x\} = K\). (Recall that if \(S\) is a subset of a topological space, then \(\overline{S}\) denotes its closure.)

**4.6.G. Exercise.** Verify that \([y - x^2]) \in \mathbb{A}^2\) is a generic point for \(V(y - x^2)\).

We will soon see (Exercise [4.7.E]) that there is a natural bijection between points of \(\text{Spec } A\) and irreducible closed subsets of \(\text{Spec } A\). You know enough to prove this now, although we will wait until we have developed some convenient terminology.

**4.6.H. Exercise.** (a) Suppose \(I = (wz - xy, wy - x^2, xz - y^2) \subset k[w, x, y, z]\). Show that \(k[w, x, y, z]/I\) is irreducible, by showing that \(I\) is prime. (Possible hint: Show that the quotient ring is an integral domain, by showing that it is isomorphic to the subring of \(k[a, b]\) generated by monomials of degree divisible by 3. There are other approaches as well, some of which we will see later. This is an example of a hard question: how do you tell if an ideal is prime?) We will later see this as the cone over the **twisted cubic curve** (the twisted cubic curve is defined in Exercise [9.2.A] and is a special case of a Veronese embedding, [§9.2.5]).

(b) Note that the generators of the ideal of part (a) may be rewritten as the equations ensuring that

\[
\text{rank} \begin{pmatrix} w & x & y \\ x & y & z \end{pmatrix} \leq 1,
\]

i.e., as the determinants of the \(2 \times 2\) submatrices. Generalize this to the ideal of rank one \(2 \times n\) matrices. This notion will correspond to the cone (§9.2.10) over the **degree n rational normal curve** (Exercise [9.2.K]).

**4.6.3. Noetherian conditions.**

In the examples we have considered, the spaces have naturally broken up into some obvious pieces. Let’s make that a bit more precise.

A topological space \(X\) is called **Noetherian** if it satisfies the descending chain condition for closed subsets: any sequence \(Z_1 \supseteq Z_2 \supseteq \cdots \supseteq Z_n \supseteq \cdots\) of closed subsets eventually stabilizes: there is an \(r\) such that \(Z_r = Z_{r+1} = \cdots\).

The following exercise may be enlightening.

**4.6.I. Exercise.** Show that any decreasing sequence of closed subsets of \(\mathbb{A}^2_C = \text{Spec } C[x, y]\) must eventually stabilize. Note that it can take arbitrarily long to stabilize. (The closed subsets of \(\mathbb{A}^2_C\) were described in [4.4.3]).

**4.6.4. Noetherian rings.** It turns out that all of the spectra we have considered have this property, but that isn’t true of the spectra of all rings. The key characteristic all of our examples have had in common is that the rings were **Noetherian**. A ring is **Noetherian** if every ascending sequence \(I_1 \subset I_2 \subset \cdots\) of ideals eventually stabilizes: there is an \(r\) such that \(I_r = I_{r+1} = \cdots\). (This is called the ascending chain condition on ideals.)
Here are some quick facts about Noetherian rings. You should be able to prove them all.

- Fields are Noetherian. \( \mathbb{Z} \) is Noetherian.
- If \( A \) is Noetherian, and \( \phi : A \to B \) is any ring homomorphism, then \( \phi(A) \) is Noetherian. Equivalently, quotients of Noetherian rings are Noetherian.
- If \( A \) is Noetherian, and \( S \) is any multiplicative set, then \( S^{-1}A \) is Noetherian.
- Any submodule of a finitely generated module over a Noetherian ring is finitely generated. (Hint: prove it for \( A \oplus n \), and use the next exercise.)

(The notion of a Noetherian module will come up in §14.6.)

4.6.J. IMPORTANT EXERCISE. Show that a ring \( A \) is Noetherian if and only if every ideal of \( A \) is finitely generated.

The next fact is non-trivial.

4.6.5. The Hilbert basis theorem. — If \( A \) is Noetherian, then so is \( A[x] \).

By the results described above, any polynomial ring over any field, or over the integers, is Noetherian — and also any quotient or localization thereof. Hence for example any finitely-generated algebra over \( k \) or \( \mathbb{Z} \), or any localization thereof, is Noetherian: \( k[x_1,x_2,\ldots] \) is not, because \( m = (x_1,x_2,\ldots) \) is not finitely generated (cf. Exercise 4.6.D(b)).

Proof of the Hilbert Basis Theorem 4.6.5. We show that any ideal \( I \subset A[x] \) is finitely-generated. We inductively produce a set of generators \( f_1,\ldots \) as follows. For \( n > 0 \), if \( I \neq (f_1,\ldots,f_{n-1}) \), let \( f_n \) be any non-zero element of \( I - (f_1,\ldots,f_{n-1}) \) of lowest degree. Thus \( f_1 \) is any element of \( I \) of lowest degree, assuming \( I \neq (0) \). If this procedure terminates, we are done. Otherwise, let \( a_n \in A \) be the initial coefficient of \( f_n \) for \( n > 0 \). Then as \( A \) is Noetherian, \( (a_1,a_2,\ldots) = (a_1,\ldots,a_N) \) for some \( N \). Say \( a_{N+1} = \sum_{i=1}^{N} b_i a_i \). Then

\[
f_{N+1} - \sum_{i=1}^{N} b_i f_i x^{\deg f_{N+1} - \deg f_i}
\]

is an element of \( I \) that is nonzero (as \( f_{N+1} \not\in (f_1,\ldots,f_N) \)) of lower degree than \( f_{n+1} \), yielding a contradiction.

4.6.K. UNIMPORTANT EXERCISE. Show that if \( A \) is Noetherian, then so is \( A[[x]] := \lim_{\leftarrow n} A[x]/x^n \), the ring of power series in \( x \). (Possible hint: Suppose \( I \subset A[[x]] \) is an ideal. Let \( I_n \subset A \) be the coefficients of \( t^n \) that appear in the elements of \( I \). Show that \( I_n \) is an ideal. Show that \( I_n \subset I_{n+1} \), and that \( I \) is determined by \( (I_0, I_1, I_2,\ldots) \).)

4.6.L. EXERCISE. If \( A \) is Noetherian, show that \( \text{Spec } A \) is a Noetherian topological space. Describe a ring \( A \) such that \( \text{Spec } A \) is not a Noetherian topological space. (As an aside, we note that if \( \text{Spec } A \) is a Noetherian topological space, \( A \) need not be Noetherian.)
4.6.M. **Exercise** (promised in Exercise 4.6.D). Show that if $A$ is Noetherian, every open subset of $	ext{Spec } A$ is quasicompact.

If $X$ is a topological space, and $Z$ is a maximal irreducible subset (an irreducible closed subset not contained in any larger irreducible closed subset), $Z$ is said to be an **irreducible component** of $X$. We think of these as the "pieces of $X" (see Figure 4.7).

![Figure 4.7. This closed subset of $A^2$ has six irreducible components](image)

4.6.N. **Exercise.** If $A$ is any ring, show that the irreducible components of $	ext{Spec } A$ are in bijection with the minimal primes of $A$. (For example, the only minimal prime of $k[x, y]$ is $(0)$.)

4.6.O. **Exercise.** Show that $	ext{Spec } A$ is irreducible if and only if $A$ has only one **minimal prime** ideal. (Minimality is with respect to inclusion.) In particular, if $A$ is an integral domain, then $	ext{Spec } A$ is irreducible.

4.6.P. **Exercise.** What are the minimal primes of $k[x, y]/(xy)$?

4.6.6. **Proposition.** — Suppose $X$ is a Noetherian topological space. Then every non-empty closed subset $Z$ can be expressed uniquely as a finite union $Z = Z_1 \cup \cdots \cup Z_n$ of irreducible closed subsets, none contained in any other.

Translation: any non-empty closed subset $Z$ has a finite number of pieces. As a corollary, this implies that a Noetherian ring $A$ has only finitely many minimal primes.

**Proof.** The following technique is called **Noetherian induction**, for reasons that will become clear.

Consider the collection of nonempty closed subsets of $X$ that cannot be expressed as a finite union of irreducible closed subsets. We will show that it is empty. Otherwise, let $Y_1$ be one such. If it properly contains another such, then
choose one, and call it \( Y_2 \). If this one contains another such, then choose one, and call it \( Y_3 \), and so on. By the descending chain condition, this must eventually stop, and we must have some \( Y_r \) that cannot be written as a finite union of irreducible closed subsets, but every closed subset properly contained in it can be so written. But then \( Y_r \) is not itself irreducible, so we can write \( Y_r = Y' \cup Y'' \) where \( Y' \) and \( Y'' \) are both proper closed subsets. Both of these by hypothesis can be written as the union of a finite number of irreducible subsets, and hence so can \( Y_r \), yielding a contradiction. Thus each closed subset can be written as a finite union of irreducible closed subsets. We can assume that none of these irreducible closed subsets contain any others, by discarding some of them.

We now show uniqueness. Suppose

\[
Z = Z_1 \cup Z_2 \cup \cdots \cup Z_r = Z'_1 \cup Z'_2 \cup \cdots \cup Z'_s
\]

are two such representations. Then \( Z'_1 \subset Z_1 \cup Z_2 \cup \cdots \cup Z_r \), so \( Z'_1 = (Z_1 \cap Z'_1) \cup \cdots \cup (Z_r \cap Z'_1) \). Now \( Z'_1 \) is irreducible, so one of these is \( Z'_1 \) itself, say (without loss of generality) \( Z_1 \cap Z'_1 \). Thus \( Z'_1 \subset Z_1 \). Similarly, \( Z_1 \subset Z''_a \) for some \( a \); but because \( Z'_1 \subset Z_1 \subset Z''_a \), and \( Z_1 \) is contained in no other \( Z'_j \), we must have \( a = 1 \), and \( Z'_1 = Z_1 \). Thus each element of the list of \( Z \)'s is in the list of \( Z'' \)'s, and vice versa, so they must be the same list. \( \square \)

4.6.7. Definition. A topological space \( X \) is connected if it cannot be written as the disjoint union of two non-empty open sets. A subset \( Y \) of \( X \) is a connected component if it is a maximal connected subset.

4.6.Q. Exercise. Show that an irreducible topological space is connected.

4.6.R. Exercise. Give (with proof!) an example of a scheme that is connected but reducible. (Possible hint: a picture may help. The symbol “×” has two “pieces” yet is connected.)

4.6.S. Exercise. If \( A \) is a Noetherian ring, show that the connected components of \( \text{Spec } A \) are unions of the irreducible components. Show that the connected components of \( \text{Spec } A \) are the subsets that are simultaneously open and closed.

4.6.T. Exercise. If \( A = A_1 \times A_2 \times \cdots \times A_n \), describe a homeomorphism \( \text{Spec } A = \text{Spec } A_1 \coprod \text{Spec } A_2 \coprod \cdots \coprod \text{Spec } A_n \). Show that each \( \text{Spec } A_i \) is a distinguished open subset \( D(f_i) \) of \( \text{Spec } A \). (Hint: let \( f_i = (0, \cdots, 0, 1, 0, \cdots, 0) \) where the 1 is in the \( i \)th component.) In other words, \( \coprod_{i=1}^n \text{Spec } A_i = \text{Spec } \coprod_{i=1}^n A_i \).

An extension of the previous exercise (that you can prove if you wish) is that \( \text{Spec } A \) is not connected if and only if \( A \) is isomorphic to the product of nonzero rings \( A_1 \) and \( A_2 \).

4.6.8. *Fun but irrelevant remark.* The previous exercise shows that \( \coprod_{i=1}^n \text{Spec } A_i \cong \text{Spec } \coprod_{i=1}^n A_i \), but this can’t hold if “\( n \) is infinite” as \( \text{Spec } A \) of any ring is quasicompact (Exercise 1.6.D(a)). This leads to an interesting phenomenon. We show that \( \text{Spec } \coprod_{i=1}^\infty A_i \) is “strictly bigger” than \( \coprod_{i=1}^\infty \text{Spec } A_i \) where each \( A_i \) is isomorphic to the field \( k \). First, we have an inclusion of sets \( \coprod_{i=1}^\infty \text{Spec } A_i \hookrightarrow \text{Spec } \coprod_{i=1}^\infty A_i \), as there is a maximal ideal of \( \coprod A_i \) corresponding to each \( i \) (precisely those elements \( 0 \) in the \( i \)th component.) But there are other maximal ideals of \( \coprod A_i \). Hint:
describe a proper ideal not contained in any of these maximal ideals. (One idea: consider elements \( \prod a_i \) that are “eventually zero”, i.e. \( a_i = 0 \) for \( i \gg 0 \).) This leads to the notion of \textit{ultrafilters}, which are very useful, but irrelevant to our current discussion.

4.6.9. \textit{Remark.} The notion of constructible and locally closed subsets will be discussed later, see Exercise 8.4.A

\section{4.7 The function \( I(\cdot) \), taking subsets of \( \Spec A \) to ideals of \( A \)}

We now introduce a notion that is in some sense “inverse” to the vanishing set function \( V(\cdot) \). Given a subset \( S \subset \Spec A \), \( I(S) \) is the set of functions vanishing on \( S \).

We make three quick observations:

- \( I(S) \) is clearly an ideal.
- \( I(S) = I(\overline{S}) \).
- \( I(\cdot) \) is inclusion-reversing: if \( S_1 \subset S_2 \), then \( I(S_2) \subset I(S_1) \).

\subsection{4.7.A. Exercise} Let \( A = k[x, y] \). If \( S = \{[(x)], [(x - 1, y)]\} \) (see Figure 4.8), then \( I(S) \) consists of those polynomials vanishing on the \( y \) axis, and at the point \((1, 0)\). Give generators for this ideal.

\begin{figure}[h]
\centering
\includegraphics[width=0.3\textwidth]{image}
\caption{The set \( S \) of Exercise/example 4.7.A pictured as a subset of \( A^2 \)}
\end{figure}

\subsection{4.7.B. Tricky exercise} Suppose \( X \subset A^3 \) is the union of the three axes. (The \( x \)-axis is defined by \( y = z = 0 \), and the \( y \)-axis and \( z \)-axis are defined analogously.) Give generators for the ideal \( I(X) \). Be sure to prove it! We will see in Exercise 13.1.F that this ideal is not generated by less than three elements.

\subsection{4.7.C. Exercise} Show that \( \text{V}(I(S)) = \overline{S} \). Hence \( \text{V}(I(S)) = S \) for a closed set \( S \). (Compare this to Exercise 4.7.D.)
Note that $I(S)$ is always a radical ideal — if $f \in \sqrt{I(S)}$, then $f^n$ vanishes on $S$ for some $n > 0$, so then $f$ vanishes on $S$, so $f \in I(S)$.

4.7.D. Exercise. Prove that if $J \subset A$ is an ideal, then $I(V(J)) = \sqrt{J}$.

This exercise and Exercise 4.7.C suggest that $V$ and $I$ are “almost” inverse. More precisely:

4.7.1. Theorem. — $V(\cdot)$ and $I(\cdot)$ give a bijection between closed subsets of $\text{Spec } A$ and radical ideals of $A$ (where a closed subset gives a radical ideal by $I(\cdot)$, and a radical ideal gives a closed subset by $V(\cdot)$).

Theorem 4.7.1 is sometimes called Hilbert’s Nullstellensatz, but we reserve that name for Theorem 4.2.3.

4.7.E. Important exercise. Show that $V(\cdot)$ and $I(\cdot)$ give a bijection between irreducible closed subsets of $\text{Spec } A$ and prime ideals of $A$. From this conclude that in $\text{Spec } A$ there is a bijection between points of $\text{Spec } A$ and irreducible closed subsets of $\text{Spec } A$ (where a point determines an irreducible closed subset by taking the closure). Hence each irreducible closed subset of $\text{Spec } A$ has precisely one generic point — any irreducible closed subset $Z$ can be written uniquely as $[z]$. 
The structure sheaf, and the definition of schemes in general

5.1 The structure sheaf of an affine scheme

The final ingredient in the definition of an affine scheme is the structure sheaf $\mathcal{O}_{\text{Spec } A}$, which we think of as the “sheaf of algebraic functions”. You should keep in your mind the example of “algebraic functions” on $\mathbb{C}^n$, which you understand well. For example, in $\mathbb{A}^2$, we expect that on the open set $D(xy)$ (away from the two axes), $(3x^4 + y + 4)/x^7y^3$ should be an algebraic function.

These functions will have values at points, but won’t be determined by their values at points. But like all sections of sheaves, they will be determined by their germs (see §5.3.3).

It suffices to describe the structure sheaf as a sheaf (of rings) on the base of distinguished open sets (Theorem 3.7.1 and Exercise 4.5.A).

5.1.1. Definition. Define $\mathcal{O}_{\text{Spec } A}(D(f))$ to be the localization of $A$ at the multiplicative set of all functions that do not vanish outside of $V(f)$ (i.e. those $g \in A$ such that $V(g) \subset V(f)$, or equivalently $D(f) \subset D(g)$). This depends only on $D(f)$, and not on $f$ itself.

5.1.A. GREAT EXERCISE. Show that the natural map $A \to \mathcal{O}_{\text{Spec } A}(D(f))$ is an isomorphism. (Possible hint: Exercise 4.5.E)

If $D(f') \subset D(f)$, define the restriction map $\text{res}_{D(f), D(f')} : \mathcal{O}_{\text{Spec } A}(D(f)) \to \mathcal{O}_{\text{Spec } A}(D(f'))$ in the obvious way: the latter ring is a further localization of the former ring. The restriction maps obviously commute: this is a “presheaf on the distinguished base”.

5.1.2. Theorem. — The data just described give a sheaf on the distinguished base, and hence determine a sheaf on the topological space $\text{Spec } A$.

This sheaf is called the structure sheaf, and will be denoted $\mathcal{O}_{\text{Spec } A}$, or sometimes $\mathcal{O}$ if the subscript is clear from the context. Such a topological space, with sheaf, will be called an affine scheme. The notation Spec $A$ will hereafter denote the data of a topological space with a structure sheaf.

Proof. We must show the base identity and base gluability axioms hold (§3.7). We show that they both hold for the open set that is the entire space Spec $A$, and leave
to you the trick which extends them to arbitrary distinguished open sets (Exercises 5.1.B and 5.1.C). Suppose $\text{Spec } A = \bigcup_{i\in I} D(f_i)$, or equivalently (Exercise 4.5.B) the ideal generated by the $f_i$ is the entire ring $A$.

We check identity on the base. Suppose that $\text{Spec } A = \bigcup_{i\in I} D(f_i)$ where $i$ runs over some index set $I$. Then there is some finite subset of $I$, which we name $\{1, \ldots, n\}$, such that $\text{Spec } A = \bigcup_{i=1}^{n} D(f_i)$, i.e. $(f_1, \ldots, f_n) = A$ (quasicompactness of $\text{Spec } A$, Exercise 4.5.C). Suppose we are given $s \in A$ such that $\text{res}_{\text{Spec } A, D(f_i)} s = 0$ in $A_{f_i}$ for all $i$. We wish to show that $s = 0$. The fact that $\text{res}_{\text{Spec } A, D(f_i)} s = 0$ in $A_{f_i}$ implies that there is some $m$ such that for each $i \in \{1, \ldots, n\}$, $f_i^m s = 0$. Now $(f_1^m, \ldots, f_n^m) = A$ (for example, from $\text{Spec } A = \bigcup D(f_i) = \bigcup (f_i^m)$), so there are $r_i \in A$ with $\sum_{i=1}^{n} r_i f_i^m = 1$ in $A$, from which

$$s = \left(\sum_i r_i f_i^m\right) s = \sum_i r_i (f_i^m s) = 0.$$

Thus we have checked the “base identity” axiom for $\text{Spec } A$. (Serre has described this as a “partition of unity” argument, and if you look at it in the right way, his insight is very enlightening.)

**5.1.B. EXERCISE.** Make the tiny changes to the above argument to show base identity for any distinguished open $D(f)$. (Hint: judiciously replace $A$ by $A_f$ in the above argument.)

We next show base gluability. Suppose again $\bigcup_{i\in I} D(f_i) = \text{Spec } A$, where $I$ is a index set (possibly horribly infinite). Suppose we are given elements in each $A_{f_i}$ that agree on the overlaps $A_{f_i, f_j}$. Note that intersections of distinguished open sets are also distinguished open sets.

(Aside: experts might realize that we are trying to show exactness of

$$0 \to A \to \prod_i A_{f_i} \to \prod_{i \neq j} A_{f_i, f_j}.$$  

Do you understand what the right-hand map is? Base identity corresponds to injectivity at $A$. The composition of the right two morphisms is trivially zero, and gluability is exactness at $\prod_i A_{f_i}$.)

Choose a finite subset $\{1, \ldots, n\} \subset I$ with $(f_1, \ldots, f_n) = A$ (or equivalently, use quasicompactness of $\text{Spec } A$ to choose a finite subcover by $D(f_i)$). We have elements $a_i/f_i \in A_{f_i}$ agreeing on overlaps $A_{f_i, f_j}$. Letting $g_i = f_i^m$, using $D(f_i) = D(g_i)$, we can simplify notation by considering our elements as of the form $a_i/g_i \in A_{g_i}$.

The fact that $a_i/g_i$ and $a_j/g_j$ “agree on the overlap” (i.e. in $A_{g_i, g_j}$) means that for some $m_{ij}$,

$$ (g_i g_j)^{m_{ij}} (g_j a_i - g_i a_j) = 0$$

in $A$. By taking $m = \max m_{ij}$ (here we use the finiteness of $I$), we can simplify notation:

$$ (g_i g_j)^m (g_j a_i - g_i a_j) = 0$$

for all $i, j$. Let $b_i = a_i g_i^m$ for all $i$, and $h_i = g_i^{m+1}$ (so $D(h_i) = D(g_i)$). Then we can simplify notation even more: on each $D(h_i)$, we have a function $b_i/h_i$, and the overlap condition is

$$ (5.1.2.2) \quad h_i b_i = h_i b_j.$$
Now $\bigcup_i D(h_i) = \text{Spec } A$, implying that $1 = \sum_{i=1}^n r_i h_i$ for some $r_i \in A$. Define $r = \sum r_i b_i$. This will be the element of $A$ that restricts to each $b_j/h_j$. Indeed, from the overlap condition (5.1.2.2),

$$rh_j = \sum_i r_i b_i h_j = \sum_i r_i h_i b_j = b_j.$$

We are not quite done! We are supposed to have something that $r$ restricts to $a_i/f_i^1$ for all $i \in I$, not just $i = 1, \ldots, n$. But a short trick takes care of this. We now show that for any $\alpha \in I - \{1, \ldots, n\}$, $r$ restricts to the desired element $a_\alpha$ of $A_{f_\alpha}$. Define $r = \sum r_i b_i$. This will be the element of $A$ that restricts to each $b_j/h_j$.

5.1.C. Exercise. Alter this argument appropriately to show base gluability for any distinguished open $D(f)$.

We have now completed the proof of Theorem 5.1.2.

The following generalization of Theorem 5.1.2 will be essential for the definition of a quasicoherent sheaf in Chapter 14.

5.1.D. Important Exercise/Definition. Suppose $M$ is an $A$-module. Show that the following construction describes a sheaf $\tilde{M}$ on the distinguished base. Define $\tilde{M}(D(f))$ to be the localization of $M$ at the multiplicative set of all functions that vanish only in $V(f)$. Define restriction maps $\text{res}_{D(f),D(g)}$ in the analogous way to $O_{\text{Spec } A}$. Show that this defines a sheaf on the distinguished base, and hence a sheaf on $\text{Spec } A$. Then show that this is an $O_{\text{Spec } A}$-module. (This sheaf $\tilde{M}$ will be very important soon; it will be an example of a quasicoherent sheaf.)

5.1.3. Remark (cf. (5.1.2.1)). In the course of answering the previous exercise, you will show that if $(f_1, \ldots, f_r) = A$, $M$ can be identified with a specific submodule of $M_{f_1} \times \cdots \times M_{f_r}$. Even though $M \to M_{f_i}$ may not be an inclusion for any $f_i$, $M \to M_{f_1} \times \cdots \times M_{f_r}$ is an inclusion. This will be useful later: we will want to show that if $M$ has some nice property, then $M_{f_i}$ does too, which will be easy. We will also want to show that if $(f_1, \ldots, f_n) = A$, then if $M_{f_i}$ have this property, then $M$ does too, and we will invoke this.

5.2 Visualizing schemes II: nilpotents

In §4.3 we discussed how to visualize the underlying set of schemes, adding in generic points to our previous intuition of “classical” (or closed) points. Our later discussion of the Zariski topology fit well with that picture. In our definition of the “affine scheme” $(\text{Spec } A, O_{\text{Spec } A})$, we have the additional information of nilpotents, which are invisible on the level of points (§4.2.9), so now we figure out to picture them. We will then readily be able to glue them together to picture
schemes in general, once we have made the appropriate definitions. As we are building intuition, we will not be rigorous or precise.

To begin, we picture $\text{Spec } \mathbb{C}[x]/(x)$ as a closed subset (a point) of $\text{Spec } \mathbb{C}[x]$: to the quotient $\mathbb{C}[x] \rightarrow \mathbb{C}[x]/(x)$, we associate the picture of a closed inclusion. The ring map can be interpreted as restriction of functions: to $\mathbb{C}[x]$, we associate its value at 0 (its residue class modulo $(x)$, by the remainder theorem). The quotient $\mathbb{C}[x]/(x^2)$ should fit in between these rings,

$$
\mathbb{C}[x] \rightarrow \mathbb{C}[x]/(x^2) \rightarrow \mathbb{C}[x]/(x)
$$

$f(x) \rightarrow f(0)$,

and we should picture it in terms of the information the quotient remembers. The image of a polynomial $f(x)$ is the information of its value at 0, and its derivative (cf. Exercise 4.2.Q). We thus picture this as being the point, plus a little bit more — a little bit of “fuzz” on the point (see Figure 5.1). (These will later be examples of closed subschemes, the schematic version of closed subsets, §9.1)

Similarly, $\mathbb{C}[x]/(x^3)$ remembers even more information — the second derivative as well. Thus we picture this as the point 0 plus even more fuzz.

More subtleties arise in two dimensions (see Figure 5.2). Consider $\text{Spec } \mathbb{C}[x, y]/(x, y)^2$, which is sandwiched between two rings we know well:

$$
\mathbb{C}[x, y] \rightarrow \mathbb{C}[x, y]/(x, y)^2 \rightarrow \mathbb{C}[x, y]/(x, y)
$$

$f(x, y) \rightarrow f(0)$.

Again, taking the quotient by $(x, y)^2$ remembers the first derivative, “in both directions”. We picture this as fuzz around the point. Similarly, $(x, y)^3$ remembers the second derivative “in all directions”.

Consider instead the ideal $(x^2, y)$. What it remembers is the derivative only in the x direction — given a polynomial, we remember its value at 0, and the coefficient of x. We remember this by picturing the fuzz only in the x direction.
This gives us some handle on picturing more things of this sort, but now it becomes more an art than a science. For example, Spec $\mathbb{C}[x, y]/(x^2, y^2)$ we might picture as a fuzzy square around the origin. One feature of this example is that given two ideals $I$ and $J$ of a ring $A$ (such as $\mathbb{C}[x, y]$), your fuzzy picture of Spec $A/(I, J)$ should be the “intersection” of your picture of Spec $A/I$ and Spec $A/J$ in Spec $A$. (You will make this precise in Exercise 9.1.G(a).) For example, Spec $\mathbb{C}[x, y]/(x^2, y^2)$ should be the intersection of two thickened lines. (How would you picture Spec $\mathbb{C}[x, y]/(x^5, y^3)$? Spec $\mathbb{C}[x, y, z]/(x^3, y^4, z^5, (x + y + z)^2)$? Spec $\mathbb{C}[x, y]/((x, y)^5, y^3)$?)

This idea captures useful information that you already have some intuition for. For example, consider the intersection of the parabola $y = x^2$ and the x-axis (in the $xy$-plane). See Figure 5.3. You already have a sense that the intersection has multiplicity two. In terms of this visualization, we interpret this as intersecting (in Spec $\mathbb{C}[x, y]$):

$$\text{Spec } \mathbb{C}[x, y]/(y - x^2) \cap \text{Spec } \mathbb{C}[x, y]/(y) = \text{Spec } \mathbb{C}[x, y]/(y - x^2, y) = \text{Spec } \mathbb{C}[x, y]/(y, x^2)$$

which we interpret as the fact that the parabola and line not just meet with multiplicity two, but that the “multiplicity 2” part is in the direction of the x-axis. You will make this example precise in Exercise 9.1.G(b).
We will later make the location of the fuzz somewhat more precise when we discuss associated points (§5.5). We will see that (in reasonable circumstances, when associated points make sense) the fuzz is concentrated on closed subsets.

## 5.3 Definition of schemes

We can now define scheme in general. First, define an isomorphism of ringed spaces \((X, \mathcal{O}_X)\) and \((Y, \mathcal{O}_Y)\) as (i) a homeomorphism \(f : X \to Y\), and (ii) an isomorphism of sheaves \(\mathcal{O}_X\) and \(\mathcal{O}_Y\), considered to be on the same space via \(f\). (Part (ii), more precisely, is an isomorphism \(f^{-1}\mathcal{O}_Y \to \mathcal{O}_X\) of sheaves on \(X\), or equivalently \(\mathcal{O}_Y \to f_*\mathcal{O}_X\) of sheaves on \(Y\).) In other words, we have a “correspondence” of sets, topologies, and structure sheaves. An affine scheme is a ringed space that is isomorphic to \((\text{Spec } A, \mathcal{O}_{\text{Spec } A})\) for some \(A\). A scheme \((X, \mathcal{O}_X)\) is a ringed space such that any point \(x \in X\) has a neighborhood \(U\) such that \((U, \mathcal{O}_{X|U})\) is an affine scheme. The scheme can be denoted \((X, \mathcal{O}_X)\), although it is often denoted \(X\), with the structure sheaf implicit.

An isomorphism of two schemes \((X, \mathcal{O}_X)\) and \((Y, \mathcal{O}_Y)\) is an isomorphism as ringed spaces. If \(U \subseteq X\) is an open subset, then \(\Gamma(\mathcal{O}_X, U)\) are said to be the functions on \(U\); this generalizes in an obvious way the definition of functions on an affine scheme. \([4.2.1]\)

### 5.3.1. Remark

From the definition of the structure sheaf on an affine scheme, several things are clear. First of all, if we are told that \((X, \mathcal{O}_X)\) is an affine scheme, we may recover its ring (i.e. find the ring \(A\) such that \(\text{Spec } A = X\)) by taking the ring of global sections, as \(X = \text{D}(1)\), so:

\[
\Gamma(X, \mathcal{O}_X) = \Gamma(\text{D}(1), \mathcal{O}_{\text{Spec } A}) \quad \text{as } \text{D}(1) = \text{Spec } A
\]

(You can verify that we get more, and can “recognize \(X\) as the scheme \(\text{Spec } A\)”): we get an isomorphism \(f : (\text{Spec } \Gamma(X, \mathcal{O}_X), \mathcal{O}_{\text{Spec } \Gamma(X, \mathcal{O}_X)}) \to (X, \mathcal{O}_X)\). For example, if \(m\) is a maximal ideal of \(\Gamma(X, \mathcal{O}_X)\), \(f([m]) = V(m)\). The following exercise will give you some practice with these notions.

### 5.3.A. Exercise (which can be strangely confusing)

Describe a bijection between the isomorphisms \(\text{Spec } A \to \text{Spec } A'\) and the ring isomorphisms \(A' \to A\).

More generally, given \(f \in A\), \(\Gamma(\text{D}(f), \mathcal{O}_{\text{Spec } A}) \cong A_f\). Thus under the natural inclusion of sets \(\text{Spec } A_f \hookrightarrow \text{Spec } A\), the Zariski topology on \(\text{Spec } A\) restricts to give the Zariski topology on \(\text{Spec } A_f\) (Exercise \([4.4.1]\)), and the structure sheaf of \(\text{Spec } A\) restricts to the structure sheaf of \(\text{Spec } A_f\), as the next exercise shows.

### 5.3.B. Important but easy exercise

Suppose \(f \in A\). Show that under the identification of \(\text{D}(f)\) in \(\text{Spec } A\) with \(\text{Spec } A_f\) (Exercise \([4.4.1]\)), there is a natural isomorphism of sheaves \((\text{D}(f), \mathcal{O}_{\text{Spec } A_{\text{D}(f)}}) \cong (\text{Spec } A_f, \mathcal{O}_{\text{Spec } A_f})\). Hint: notice that distinguished open sets of \(\text{Spec } R\) are already distinguished open sets in \(\text{Spec } R\).

### 5.3.C. Easy Exercise

If \(X\) is a scheme, and \(U\) is any open subset, prove that \((U, \mathcal{O}_X|_U)\) is also a scheme.
5.3.2. **Definitions.** We say \((U, \mathcal{O}_X|_U)\) is an **open subscheme** of \(U\). If \(U\) is also an affine scheme, we often say \(U\) is an **affine open subset**, or an **affine open subscheme**, or sometimes informally just an **affine open**. For example, \(D(f)\) is an affine open subscheme of \(\text{Spec} \ A\).

5.3.D. **EASY EXERCISE.** Show that if \(X\) is a scheme, then the affine open sets form a base for the Zariski topology.

5.3.E. **EASY EXERCISE.** The disjoint union of schemes is defined as you would expect: it is the disjoint union of sets, with the expected topology (thus it is the disjoint union of topological spaces), with the expected sheaf. Once we know what morphisms are, it will be immediate (Exercise 10.1.A) that (just as for sets and topological spaces) disjoint union is the coproduct in the category of schemes.

(a) Show that the disjoint union of a **finite** number of affine schemes is also an affine scheme. (Hint: Exercise 4.6.T.)

(b) (a first example of a non-affine scheme) Show that an infinite disjoint union of (non-empty) affine schemes is not an affine scheme. (Hint: affine schemes are quasicompact, Exercise 4.6.D(a).)

5.3.3. **Stalks of the structure sheaf: germs, values at a point, and the residue field of a point.** Like every sheaf, the structure sheaf has stalks, and we shouldn’t be surprised if they are interesting from an algebraic point of view. In fact, we have seen them before.

5.3.F. **IMPORTANT EXERCISE.** Show that the stalk of \(\mathcal{O}_{\text{Spec} \ A}\) at the point \([p]\) is the local ring \(A_p\).

Essentially the same argument will show that the stalk of the sheaf \(\tilde{\mathcal{M}}\) (defined in Exercise 5.1.D) at \([p]\) is \(M_p\). Here is an interesting consequence, or if you prefer, a geometric interpretation of an algebraic fact. A section is determined by its germs (Exercise 3.4.A), meaning that \(M \rightarrow \prod_p M_p\) is an inclusion. So for example an \(A\)-module is zero if and only if all its localizations at primes are zero.

5.3.4. **Definition.** We say a ringed space is a **locally ringed space** if its stalks are local rings. (The motivation for the terminology comes from thinking of sheaves in terms of stalks. A **ringed space** is a sheaf whose stalks are rings. A **locally ringed space** is a sheaf whose stalks are local rings.) Thus schemes are locally ringed spaces. Manifolds are another example of locally ringed spaces, see §3.1.1. In both cases, taking quotient by the maximal ideal may be interpreted as evaluating at the point. The maximal ideal of the local ring \(\mathcal{O}_{X,p}\) is denoted \(m_{X,p}\) or \(m_p\), and the **residue field** \(\mathcal{O}_{X,p}/m_p\) is denoted \(\kappa(p)\). Functions on an open subset \(U\) of a locally ringed space have **values** at each point of \(U\). The value at \(p\) of such a function lies in \(\kappa(p)\). As usual, we say that a function **vanishes** at a point \(p\) if its value at \(p\) is \(0\).

As an example, consider a point \([p]\) of an affine scheme \(\text{Spec} \ A\). (Of course, this example is “universal”, as all points may be interpreted in this way, by choosing an affine neighborhood.) The residue field at \([p]\) is \(A_p/pA_p\), which is isomorphic to \(K(A/p)\), the fraction field of the quotient. It is useful to note that localization at
p and taking quotient by p “commute”, i.e. the following diagram commutes.

\[
\begin{array}{ccc}
A_p & \rightarrow & A_p/pA_p = K(A/p) \\
\text{localize} & & \text{localize}, \ i.e. K(\cdot)
\end{array}
\]

For example, consider the scheme \(\mathbb{A}^2_k\) = Spec \(k[x, y]\), where \(k\) is a field of characteristic not 2. Then \((x^2 + y^2)/x(y^2 - x^5)\) is a function away from the \(y\)-axis and the curve \(y^2 - x^5\). Its value at \((2, 4)\) (by which we mean \([[x - 2, y - 4]]\)) is 
\[
\left(\frac{2^2 + 4^2}{2(4^2 - 2^5)}\right)
\]
in the residue field — check this if it seems mysterious. And its value at \([[y]]\), the generic point of the \(x\)-axis, is \(\frac{x^2}{x^2} = -1/x^4\), which we see by setting \(y\) to 0. This is indeed an element of the fraction field of \(k[x, y]/(y)\), i.e. \(k(x)\). (If you think you care only about algebraically closed fields, let this example be a first warning: \(A_p/pA_p\) won’t be algebraically closed in general, even if \(A\) is a finitely generated \(\mathbb{C}\)-algebra!)

If anything makes you nervous, you should make up an example to make you feel better. Here is one: 27/4 is a function on Spec \(\mathbb{Z} - \{(2), (7)\}\) or indeed on an even bigger open set. What is its value at \([[5]]\)? Answer: \(2/(-1) \equiv -2 \pmod{5}\). What is its value at the generic point \([[0]]\)? Answer: 27/4. Where does it vanish? At \([[3]]\).

5.3.5. Stray definition: the fiber of an \(\mathcal{O}\)-module at a point. If \(\mathcal{F}\) is an \(\mathcal{O}\)-module on a scheme \(X\) (or more generally, a locally ringed space), define the fiber of \(\mathcal{F}\) at a point \(p \in X\) by

\[
\mathcal{F}_p := \mathcal{F}_p \otimes_{\mathcal{O}_X, p} \kappa(p).
\]

As a reality check, \(\mathcal{O}_p\) is essentially by definition \(\kappa(p)\).

5.4 Three examples

We now give three extended examples. Our short-term goal is to see that we can really work with the structure sheaf, and can compute the ring of sections of interesting open sets that aren’t just distinguished open sets of affine schemes. Our long-term goal is to meet interesting examples that will come up repeatedly in the future.

5.4.1. Example: The plane minus the origin. This example will show you that the distinguished base is something that you can work with. Let \(A = k[x, y]\), so Spec \(A = \mathbb{A}^2_k\). Let’s work out the space of functions on the open set \(U = \mathbb{A}^2 - \{(0, 0)\} = \mathbb{A}^2 - \{[(x, y)]\}\).
You can’t cut out this set with a single equation (can you see why?), so this isn’t a distinguished open set. But in any case, even if we are not sure if this is a distinguished open set, we can describe it as the union of two things which are distinguished open sets: \( U = D(x) \cup D(y) \). We will find the functions on \( U \) by gluing together functions on \( D(x) \) and \( D(y) \).

The functions on \( D(x) \) are, by Definition 5.1.1, \( A_x = k[x, y, 1/x] \). The functions on \( D(y) \) are \( A_y = k[x, y, 1/y] \). Note that \( A \rightarrow A_x, A_y \). This is because \( x \) and \( y \) are not zero-divisors. (The ring \( A \) is an integral domain — it has no zero-divisors, besides 0 — so localization is always an inclusion, Exercise 2.3.C.) So we are looking for functions on \( D(x) \) and \( D(y) \) that agree on \( D(x) \cap D(y) = D(xy) \), i.e. they are just the same Laurent polynomial. Which things of this first form are also of the second form? Just traditional polynomials —

\[
\Gamma(U, \mathcal{O}_{\mathbb{A}^2}) = k[x, y].
\]

In other words, we get no extra functions by throwing out this point. Notice how easy that was to calculate!

5.4.2. Aside. Notice that any function on \( \mathbb{A}^2 - \{0, 0\} \) extends over all of \( \mathbb{A}^2 \). This is an analogue of Hartogs’ Lemma in complex geometry: you can extend a holomorphic function defined on the complement of a set of codimension at least two on a complex manifold over the missing set. This will work more generally in the algebraic setting: you can extend over points in codimension at least 2 not only if they are smooth, but also if they are mildly singular — what we will call normal. We will make this precise in §12.3.10. This fact will be very useful for us.

5.4.3. We now show an interesting fact: \((U, \mathcal{O}_{\mathbb{A}^2|U})\) is a scheme, but it is not an affine scheme. (This is confusing, so you will have to pay attention.) Here’s why: otherwise, if \((U, \mathcal{O}_{\mathbb{A}^2|U}) = (\text{Spec } A, \mathcal{O}_{\text{Spec } A})\), then we can recover \( A \) by taking global sections:

\[ A = \Gamma(U, \mathcal{O}_{\mathbb{A}^2|U}), \]

which we have already identified in (5.4.1.1) as \( k[x, y] \). So if \( U \) is affine, then \( U \cong \mathbb{A}^2_k \). But this bijection between primes in a ring and points of the spectrum is more constructive than that: given the prime ideal \( I \), you can recover the point as the generic point of the closed subset cut out by \( I \), i.e. \( V(I) \), and given the point \( p \), you can recover the ideal as those functions vanishing at \( p \), i.e. \( I(p) \). In particular, the prime ideal \((x, y)\) of \( A \) should cut out a point of \( \text{Spec } A \). But on \( U \), \( V(x) \cap V(y) = \emptyset \). Conclusion: \( U \) is not an affine scheme. (If you are ever looking for a counterexample to something, and you are expecting one involving a non-affine scheme, keep this example in mind!)

5.4.4. Gluing two copies of \( \mathbb{A}^1 \) together in two different ways. We have now seen two examples of non-affine schemes: an infinite disjoint union of non-empty schemes: Exercise 5.3.E and \( \mathbb{A}^2 - \{0, 0\} \). I want to give you two more examples. They are important because they are the first examples of fundamental behavior, the first pathological, and the second central.

First, I need to tell you how to glue two schemes together. Before that, you should review how to glue topological spaces together along isomorphic open sets. Given two topological spaces \( X \) and \( Y \), and open subsets \( U \subset X \) and \( V \subset Y \) along with a homeomorphism \( U \cong V \), we can create a new topological space \( W \),
that we think of as gluing X and Y together along $U \cong V$. It is the quotient of the disjoint union $X \bigsqcup Y$ by the equivalence relation $U \cong V$, where the quotient is given the quotient topology. Then X and Y are naturally (identified with) open subsets of $W$, and indeed cover $W$. Can you restate this cleanly with an arbitrary (not necessarily finite) number of topological spaces?

Now that we have discussed gluing topological spaces, let’s glue schemes together. Suppose you have two schemes $(X, \mathcal{O}_X)$ and $(Y, \mathcal{O}_Y)$, and open sub-
sets $U \subset X$ and $V \subset Y$, along with a homeomorphism $f : U \xrightarrow{\sim} V$, and an
isomorphism of structure sheaves $\mathcal{O}_X \cong f^* \mathcal{O}_Y$ (i.e. an isomorphism of schemes $(U, \mathcal{O}_X|_U) \cong (V, \mathcal{O}_Y|_V)$). Then we can glue these together to get a single scheme.

Reason: let $W$ be $X$ and $Y$ glued together using the isomorphism $U \cong V$. Then Exercise 3.7.D shows that the structure sheaves can be glued to
together to get a sheaf of rings. Note that this is indeed a scheme: any point has a ne-
ghborhood that is

5.4.A. ESSENTIAL EXERCISE (cf. EXERCISE 3.7.D). For later reference, show that you can glue an arbitrary collection of schemes together. Suppose we are given:

- schemes $X_i$ (as i runs over some index set I, not necessarily finite),
- open subschemes $X_{ij} \subset X_i$,
- isomorphisms $f_{ij} : X_{ij} \rightarrow X_{ji}$ with $f_{ii}$ the identity

such that

- (the cocycle condition) the isomorphisms “agree on triple intersections”,
  i.e. $f_{ik}|_{X_{ij} \cap X_{ik}} = f_{jk}|_{X_{ij} \cap X_{ik}} \circ f_{ij}|_{X_{ij} \cap X_{ik}}$.

(The cocycle condition ensures that $f_{ij}$ and $f_{ji}$ are inverses. In fact, the hypothesis that $f_{ii}$ is the identity also follows from the cocycle condition.) Show that there is a unique scheme $X$ (up to unique isomorphism) along with open subset isomorphic to $X_i$ respecting this gluing data in the obvious sense. (Hint: what is X as a set? What is the topology on this set? In terms of your description of the open sets of X, what are the sections of this sheaf over each open set?)

I will now give you two non-affine schemes. In both cases, I will glue together
two copies of the affine line $\mathbb{A}^1_k$. Let $X = \text{Spec } k[t]$, and $Y = \text{Spec } k[u]$. Let
$U = D(t) = \text{Spec } k[t, 1/t] \subset X$ and $V = D(u) = \text{Spec } k[u, 1/u] \subset Y$. We will get both examples by gluing X and Y together along $U$ and $V$. The difference will be in how we glue.

5.4.5. Extended example: the affine line with the doubled origin. Consider the
isomorphism $U \cong V$ via the isomorphism $k[t, 1/t] \cong k[u, 1/u]$ given by $t \leftrightarrow u$ (cf. Exercise 5.3.A). The resulting scheme is called the affine line with doubled origin. Figure 5.4 is a picture of it.
As the picture suggests, intuitively this is an analogue of a failure of Hausdorffness. Now $\mathbb{A}^1$ itself is not Hausdorff, so we can’t say that it is a failure of Hausdorffness. We see this as weird and bad, so we will want to make a definition that will prevent this from happening. This will be the notion of separatedness (to be discussed in Chapter [11]). This will answer other of our prayers as well. For example, on a separated scheme, the “affine base of the Zariski topology” is nice — the intersection of two affine open sets will be affine (Proposition [11.1.8]).

5.4.B. Exercise. Show that the affine line with doubled origin is not affine. Hint: calculate the ring of global sections, and look back at the argument for $\mathbb{A}^2 - \{(0, 0)\}$.

5.4.C. Easy Exercise. Do the same construction with $\mathbb{A}^1$ replaced by $\mathbb{A}^2$. You’ll have defined the affine plane with doubled origin. Describe two affine open subsets of this scheme whose intersection is not an affine open subset.

5.4.6. Example 2: the projective line. Consider the isomorphism $U \cong V$ via the isomorphism $k[t, 1/t] \cong k[u, 1/u]$ given by $t \leftrightarrow 1/u$. Figure 5.5 is a suggestive picture of this gluing. The resulting scheme is called the projective line over the field $k$, and is denoted $\mathbb{P}^1_k$.

Notice how the points glue. Let me assume that $k$ is algebraically closed for convenience. (You can think about how this changes otherwise.) On the first affine line, we have the closed (“traditional”) points $[(t - a)]$, which we think of as “$a$ on the $t$-line”, and we have the generic point $[(0)]$. On the second affine line, we have closed points that are “$b$ on the $u$-line”, and the generic point. Then $a$ on the $t$-line is glued to $1/a$ on the $u$-line (if $a \neq 0$ of course), and the generic point is glued to the generic point (the ideal $(0)$ of $k[t]$ becomes the ideal $(0)$ of $k[t, 1/t]$ upon localization, and the ideal $(0)$ of $k[u]$ becomes the ideal $(0)$ of $k[u, 1/u]$). And $(0)$ in $k[t, 1/t]$ is $(0)$ in $k[u, 1/u]$ under the isomorphism $t \leftrightarrow 1/u$.

5.4.7. If $k$ is algebraically closed, we can interpret the closed points of $\mathbb{P}^1_k$ in the following way, which may make this sound closer to the way you have seen projective space defined earlier. The points are of the form $[a:b]$, where $a$ and $b$ are not both zero, and $[a:b]$ is identified with $[ac:bc]$ where $c \in k^*$. Then if $b \neq 0$, this is identified with $a/b$ on the $t$-line, and if $a \neq 0$, this is identified with $b/a$ on the $u$-line.
5.4.8. Proposition. — $\mathbb{P}^1_k$ is not affine.

Proof. We do this by calculating the ring of global sections. The global sections correspond to sections over $X$ and sections over $Y$ that agree on the overlap. A section on $X$ is a polynomial $f(t)$. A section on $Y$ is a polynomial $g(u)$. If we restrict $f(t)$ to the overlap, we get something we can still call $f(t)$; and similarly for $g(u)$. Now we want them to be equal: $f(t) = g(1/t)$. But the only polynomials in $t$ that are at the same time polynomials in $1/t$ are the constants $k$. Thus $\Gamma(\mathbb{P}^1, O_{\mathbb{P}^1}) = k$. If $\mathbb{P}^1$ were affine, then it would be $\text{Spec} \, \Gamma(\mathbb{P}^1, O_{\mathbb{P}^1}) = \text{Spec} \, k$, i.e. one point. But it isn’t — it has lots of points. □

We have proved an analogue of a theorem: the only holomorphic functions on $\mathbb{C}P^1$ are the constants!

5.4.9. Important example: Projective space. We now make a preliminary definition of projective $n$-space over a field $k$, denoted $\mathbb{P}^n_k$, by gluing together $n + 1$ open sets each isomorphic to $\mathbb{A}^n_k$. Judicious choice of notation for these open sets will make our life easier. Our motivation is as follows. In the construction of $\mathbb{P}^1$ above, we thought of points of projective space as $[x_0 : x_1]$, where $[x_0, x_1]$ are only determined up to scalars, i.e. $(x_0, x_1)$ is considered the same as $(\lambda x_0, \lambda x_1)$. Then the first patch can be interpreted by taking the locus where $x_0 \neq 0$, and then we consider the points $[1; t]$, and we think of $t$ as $x_1/x_0$; even though $x_0$ and $x_1$ are not well-defined, $x_1/x_0$ is. The second corresponds to where $x_1 \neq 0$, and we consider the points $[u; 1]$, and we think of $u$ as $x_0/x_1$. It will be useful to instead use the notation $x_1/0$ for $t$ and $x_0/1$ for $u$.

For $\mathbb{P}^n$, we glue together $n + 1$ open sets, one for each of $i = 0, \ldots, n + 1$. The $i$th open set $U_i$ will have coordinates $x_0/i, \ldots, x_{i-1}/i, x_{i+1}/i, \ldots, x_n/i$. It will be convenient to write this as

$$\text{Spec} \, k[x_0/i, x_1/i, \ldots, x_n/i]/(x_i/i - 1)$$

(so we have introduced a “dummy variable” $x_i/i$ which we set to 1). We glue the distinguished open set $D(x_i/i)$ of $U_i$ to the distinguished open set $D(x_j/i)$ of $U_j$, by identifying these two schemes by describing the identification of rings

$$\text{Spec} \, k[x_0/i, x_1/i, \ldots, x_n/i, 1/x_i/i]/(x_i/i - 1) \cong \text{Spec} \, k[x_0/j, x_1/j, \ldots, x_n/j, 1/x_j/j]/(x_j/j - 1)$$

via $x_k/i = x_k/j/x_j/i$ and $x_k/j = x_k/i/x_i/j$ (which implies $x_i/j x_j/i = 1$). We need to check that this gluing information agrees over triple overlaps.

5.4.D. Exercise. Check this, as painlessly as possible. (Possible hint: the triple intersection is affine; describe the corresponding ring.)

Note that our definition does not use the fact that $k$ is a field. Hence we may as well define $\mathbb{P}^n_A$ for any ring $A$. This will be useful later.

5.4.E. Exercise. Show that the only global sections of the structure sheaf are constants, and hence that $\mathbb{P}^n_k$ is not affine if $n > 0$. (Hint: you might fear that you will need some delicate interplay among all of your affine open sets, but you will only need two of your open sets to see this. There is even some geometric intuition behind this: the complement of the union of two open sets has codimension
2. But “Algebraic Hartogs’ Lemma” (discussed informally in §5.4.2 to be stated rigorously in Theorem 12.3.10) says that any function defined on this union extends to be a function on all of projective space. Because we are expecting to see only constants as functions on all of projective space, we should already see this for this union of our two affine open sets.)

5.4.F. Exercise (Generalizing §5.4.7). Show that if \( k \) is algebraically closed, the closed points of \( \mathbb{P}^n_k \) may be interpreted in the traditional way: the points are of the form \([a_0; \ldots; a_n] \), where the \( a_i \) are not all zero, and \([a_0; \ldots; a_n] \) is identified with \([\lambda a_0; \ldots; \lambda a_n] \) where \( \lambda \in k^* \).

We will later give other definitions of projective space (Definition 5.5.4, §17.4.2). Our first definition here will often be handy for computing things. But there is something unnatural about it — projective space is highly symmetric, and that isn’t clear from our current definition.

5.4.10. Fun aside: The Chinese Remainder Theorem is a geometric fact. The Chinese Remainder theorem is embedded in what we have done, which shouldn’t be obvious. I will show this by example, but you should then figure out the general statement. The Chinese Remainder Theorem says that knowing an integer modulo 60 is the same as knowing an integer modulo 3, 4, and 5. Here’s how to see this in the language of schemes. What is \( \text{Spec } \mathbb{Z}/(60) \)? What are the primes of this ring? Answer: those prime ideals containing \( (60) \), i.e. those primes dividing 60, i.e. \( (2) \), \( (3) \), and \( (5) \). Figure 5.6 is a sketch of \( \text{Spec } \mathbb{Z}/(60) \). They are all closed points, as these are all maximal ideals, so the topology is the discrete topology. What are the stalks? You can check that they are \( \mathbb{Z}/4 \), \( \mathbb{Z}/3 \), and \( \mathbb{Z}/5 \). The nilpotents “at \( (2) \)” are indicated by the “fuzz” on that point. (We discussed visualizing nilpotents with “infinitesimal fuzz” in §5.2.) So what are global sections on this scheme? They are sections on this open set \( (2) \), this other open set \( (3) \), and this third open set \( (5) \). In other words, we have a natural isomorphism of rings

\[
\mathbb{Z}/60 \to \mathbb{Z}/4 \times \mathbb{Z}/3 \times \mathbb{Z}/5.
\]

![Figure 5.6. A picture of the scheme \( \text{Spec } \mathbb{Z}/(60) \)](image)

5.4.11. Example. Here is an example of a function on an open subset of a scheme that is a bit surprising. On \( X = \text{Spec } k[w, x, y, z]/(wx - yz) \), consider the open subset \( D(y) \cup D(w) \). Show that the function \( x/y \) on \( D(y) \) agrees with \( z/w \) on \( D(w) \) on their overlap \( D(y) \cap D(w) \). Hence they glue together to give a section. You may have seen this before when thinking about analytic continuation in complex geometry — we have a “holomorphic” function which has the description \( x/y \) on
an open set, and this description breaks down elsewhere, but you can still “analytically continue” it by giving the function a different definition on different parts of the space.

Follow-up for curious experts: This function has no “single description” as a well-defined expression in terms of $w, x, y, z$. There is lots of interesting geometry here. This example will be a constant source of interesting examples for us. We will later recognize it as the cone over the quadric surface. This means that any expression involving throwing away the (cone over the) line $\ell$ in the same ruling as $C$ also involves throwing away the (cone over the) union of two lines $\ell$ and $m_2$, one in each “ruling” of the surface, and $V(z)$ also involves throwing away the (cone over the) union of two lines $\ell$ and $m_2$. The intersection is the (cone over the) line $\ell$, which is a codimension 1 set. Neat fact: despite being “pure codimension 1”, it is not cut out even set-theoretically by a single equation. (It is hard to get an example of this behavior.
not an open subset of a projective variety. So projective schemes are really ubiquitous. Also a projective k-scheme is a good approximation of the algebro-geometric version of compactness ("properness", see §11.3).

Finally, although projective schemes may be obtained by gluing together affines, and we know that keeping track of gluing can be annoying, there is a simple means of dealing with them without worrying about gluing. Just as there is a rough dictionary between rings and affine schemes, we will have an analogous dictionary between graded rings and projective schemes. Just as one can work with affine schemes by instead working with rings, one can work with projective schemes by instead working with graded rings. To get an initial sense of how this works, consider Example 9.2.1 (which secretly gives the notion of projective A-schemes in full generality). Recall that any collection of homogeneous elements of \( A[x_0, \ldots, x_n] \) describes a closed subscheme of \( \mathbb{P}^n_A \). (The \( x_0, \ldots, x_n \) are called projective coordinates on the scheme. Warning: they are not functions on the scheme. Any closed subscheme of \( \mathbb{P}^n_A \) cut out by a set of homogeneous polynomials will soon be called a projective A-scheme.) Thus if \( I \) is a homogeneous ideal in \( A[x_0, \ldots, x_n] \) (i.e. generated by homogeneous polynomials), we have defined a closed subscheme of \( \mathbb{P}^n_A \) deserving the name \( V(I) \). Conversely, given a closed subset \( S \) of \( \mathbb{P}^n_A \), we can consider those homogeneous polynomials in the projective coordinates, vanishing on \( S \). This homogeneous ideal deserves the name \( I(S) \).

5.5.1. A motivating picture from classical geometry. For geometric intuition, we recall how one thinks of projective space "classically" (in the classical topology, over the real numbers). \( \mathbb{P}^n \) can be interpreted as the lines through the origin in \( \mathbb{R}^{n+1} \). Thus subsets of \( \mathbb{P}^n \) correspond to unions of lines through the origin of \( \mathbb{R}^{n+1} \), and closed subsets correspond to such unions which are closed. (The same is not true with "closed" replaced by "open"!)

One often pictures \( \mathbb{P}^n \) as being the "points at infinite distance" in \( \mathbb{R}^{n+1} \), where the points infinitely far in one direction are associated with the points infinitely far in the opposite direction. We can make this more precise using the decomposition

\[
\mathbb{P}^{n+1} = \mathbb{R}^{n+1} \coprod \mathbb{P}^n
\]

by which we mean that there is an open subset in \( \mathbb{P}^{n+1} \) identified with \( \mathbb{R}^{n+1} \) (the points with last projective coordinate non-zero), and the complementary closed subset identified with \( \mathbb{P}^n \) (the points with last projective coordinate zero).

Then for example any equation cutting out some set \( V \) of points in \( \mathbb{P}^n \) will also cut out some set of points in \( \mathbb{R}^n \) that will be a closed union of lines. We call this the affine cone of \( V \). These equations will cut out some union of \( \mathbb{P}^1 \)'s in \( \mathbb{P}^{n+1} \), and we call this the projective cone of \( V \). The projective cone is the disjoint union of the affine cone and \( V \). For example, the affine cone over \( x^2 + y^2 = z^2 \) in \( \mathbb{R}^2 \) is just the "classical" picture of a cone in \( \mathbb{R}^3 \), see Figure 5.7. We will make this analogy precise in our algebraic setting in §9.2.10. To make a connection with the previous discussion on homogeneous ideals: the homogeneous ideal given by the cone is \( (x^2 + y^2 - z^2) \).

5.5.2. The \textit{Proj} construction.

We will now produce a scheme out of a graded ring. A \textit{graded ring} for us is a ring \( S = \bigoplus_{n \in \mathbb{Z}_{\geq 0}} S_n \) (the subscript is called the \textit{grading}), where multiplication respects the grading, i.e. sends \( S_m \times S_n \) to \( S_{m+n} \). (Our graded rings are indexed...
by $\mathbb{Z}^{\geq 0}$. One can define more general graded rings, but we won’t need them.)

Note that $S_0$ is a subring, and $S_\bullet$ is a $S_0$-algebra. In our examples so far, we have

a graded ring $A[x_0, \ldots, x_n]/I$ where $I$ is a homogeneous ideal. We are taking the usual grading on $A[x_0, \ldots, x_n]$, where each $x_i$ has weight 1. In most of the examples below, $S_0 = A$, and $S_\bullet$ is generated over $S_0$ by $S_1$.

**5.5.3. Graded rings over $A$, and finitely generated graded rings.** Fix a ring $A$ (the base ring). Our motivating example is $S_\bullet = A[x_0, x_1, x_2]$, with the usual grading. If $S_\bullet$ is graded by $\mathbb{Z}^{\geq 0}$, with $S_0 = A$, we say that $S_\bullet$ is a graded ring over $A$. Hence each $S_n$ is an $A$-module. The subset $S_+ := \bigoplus_{i>0} S_i \subset S_\bullet$ is an ideal, called the irrelevant ideal. The reason for the name “irrelevant” will be clearer in a few paragraphs. If the irrelevant ideal $S_+$ is a finitely-generated ideal, we say that $S_\bullet$ is a finitely generated graded ring over $A$. If $S_\bullet$ is generated by $S_1$ as an $A$-algebra, we say that $S_\bullet$ is generated in degree 1.

**5.5.A. Exercise.** Show that $S_\bullet$ is a finitely-generated graded ring if and only if $S_\bullet$ is a finitely-generated graded $A$-algebra, i.e. generated over $A = S_0$ by a finite number of homogeneous elements of positive degree. (Hint for the forward implication: show that the generators of $S_+$ as an ideal are also generators of $S_\bullet$ as an algebra.)

Motivated by our example of $\mathbb{P}^n_A$ and its closed subschemes, we now define a scheme $\text{Proj} S_\bullet$. As we did with $\text{Spec}$ of a ring, we will build it first as a set, then as a topological space, and finally as a ringed space. In our preliminary definition of $\mathbb{P}^n_A$, we glued together $n+1$ well-chosen affine pieces, but we don’t want to make any choices, so we do this by simultaneously consider “all possible” affines. Our affine building blocks will be as follows. For each homogeneous $f \in S_+$, consider

$$(5.5.3.1) \quad \text{Spec}((S_\bullet)_f)_0,$$

where $((S_\bullet)_f)_0$ means the 0-graded piece of the graded ring $(S_\bullet)_f$. The notation $((S_\bullet)_f)_0$ is admittedly horrible — the first and third subscripts refer to the grading, and the second refers to localization.

(Before we begin: another possible way of defining $\text{Proj} S_\bullet$ is by gluing together affines, by jumping straight to Exercises 5.5.G, 5.5.H, and 5.5.I. If you prefer that, by all means do so.)
The points of Proj $S_\bullet$ are set of homogeneous prime ideals of $S_\bullet$ not containing the irrelevant ideal $S_+$ (the “relevant prime ideals”).

**5.5.B. Important and tricky exercise.** Suppose $f \in S_+$ is homogeneous. Give a bijection between the primes of $( (S_\bullet)_f )_0$ and the homogeneous prime ideals of $(S_\bullet)_f$. Describe the latter as a subset of Proj $S_\bullet$. Hint: From the ring map $( (S_\bullet)_f )_0 \rightarrow (S_\bullet)_f$, from each homogeneous prime of $(S_\bullet)_f$ we find a homogeneous prime of $( (S_\bullet)_f )_0$. The reverse direction is the harder one. Given a prime ideal $P_0 \subset ( (S_\bullet)_f )_0$, define $P \subset (S_\bullet)_f$ as generated by the following homogeneous elements: $a \in P$ if and only if $a^{\deg f}/f^{\deg a} \in P_0$. Showing that homogeneous $a$ is in $P$ if and only if $a^2 \in P$; show that if $a_1, a_2 \in P$ then $(a_1 + a_2)^2 \in P$ and hence $a_1 + a_2 \in P$; then show that $P$ is an ideal; then show that $P$ is prime.)

The interpretation of the points of Proj $S_\bullet$ with homogeneous prime ideals helps us picture Proj $S_\bullet$. For example, if $S_\bullet = k[x,y,z]$ with the usual grading, then we picture the homogeneous prime ideal $(z^2 - x^2 - y^2)$ as a subset of Spec $S_\bullet$; it is a cone (see Figure 5.7). As in [5.5.D] we picture $\mathbb{P}_k^2$ as the “plane at infinity”. Thus we picture this equation as cutting out a conic “at infinity”. We will make this intuition somewhat more precise in [9.2.10].

**5.5.C. Exercise (the Zariski topology on Proj $S_\bullet$).** If $I$ is a homogeneous ideal of $S_\bullet$, define the **vanishing set of I, $V(I) \subset$ Proj $S_\bullet$** to be those homogeneous prime ideals containing $I$. As in the affine case, let $V(f)$ be $V((f))$, and let $D(f) = \text{Proj } S_\bullet \setminus V(f)$ (the **projective distinguished open set**) be the complement of $V(f)$ (i.e. the open subscheme corresponding to that open set). Show that $D(f)$ is precisely the subset $( (S_\bullet)_f )_0$ you described in the previous exercise.

As in the affine case, the $V(I)$’s satisfy the axioms of the closed set of a topology, and we call this the **Zariski topology** on Proj $S_\bullet$. Many statements about the Zariski topology on Spec of a ring carry over to this situation with little extra work. Clearly $D(f) \cap D(g) = D(fg)$, by the same immediate argument as in the affine case (Exercise 4.5.D). As in the affine case (Exercise 4.5.B), if $D(f) \subset D(g)$, then $f^n \in (g)$ for some $n$, and vice versa.

**5.5.D. Easy exercise.** Verify that the projective distinguished open sets form a base of the Zariski topology.

**5.5.E. Exercise.** Fix a graded ring $S_\bullet$.

(a) Suppose $I$ is any homogeneous ideal of $S_\bullet$, and $f$ is a homogeneous element. Show that $f$ vanishes on $V(I)$ if and only if $f^n \in I$ for some $n$. (Hint: Mimic the affine case; see Exercise 4.4.1)

(b) If $Z \subset \text{Proj } S_\bullet$, define $I(\cdot)$. Show that it is a homogeneous ideal. For any two subsets, show that $I(Z_1 \cup Z_2) = I(Z_1) \cap I(Z_2)$.

(c) For any subset $Z \subset \text{Proj } S_\bullet$, show that $V(I(Z)) = Z$.

**5.5.F. Exercise (cf. Exercise 4.5.B).** Fix a graded ring $S_\bullet$. Show that the following are equivalent.

(a) $V(I) = \emptyset$.

(b) for any $f_i$ (as $i$ runs through some index set) generating $I$, $\cup D(f_i) = \text{Proj } S_\bullet$. 

We then glue the various Spec\(_{\pi}\) and glue these sheaves together using Exercise 3.7.D on gluing sheaves. This is more motivation for the Spec\(_{\pi}\) being “irrelevant”: any ideal whose radical contains it is “geometrically irrelevant”.

Let’s get back to constructing \(\mathrm{Proj} S_{\bullet}\) as a scheme.

5.5.G. Exercise. Suppose some homogeneous \(f \in S_{\bullet}\) is given. Via the inclusion
\[
D(f) = \mathrm{Spec}((S_{\bullet}_f)_0) \hookrightarrow \mathrm{Proj} S_{\bullet},
\]
show that the Zariski topology on \(\mathrm{Proj} S_{\bullet}\) restricts to the Zariski topology on \(\mathrm{Spec}((S_{\bullet}_f)_0)\).

Now that we have defined \(\mathrm{Proj} S_{\bullet}\) as a topological space, we are ready to define the structure sheaf. On \(D(f)\), we wish it to be the structure sheaf of \(\mathrm{Spec}((S_{\bullet}_f)_0)\). We will glue these sheaves together using Exercise 3.7.D on gluing sheaves.

5.5.H. Exercise. If \(f, g \in S_{\bullet}^+\) are homogeneous, describe an isomorphism between \(\mathrm{Spec}((S_{\bullet}_f)_0)\) and the distinguished open subset \(D(g^{\deg f}/f^{\deg g})\) of \(\mathrm{Spec}((S_{\bullet}_f)_0)\).

Similarly, \(\mathrm{Spec}((S_{\bullet}_f)_0)\) is identified with a distinguished open subset of \(\mathrm{Spec}((S_{\bullet}_g)_0)\). We then glue the various \(\mathrm{Spec}((S_{\bullet}_f)_0)(\text{as } f \text{ varies})\) altogether, using these pairwise gluings.

5.5.I. Exercise. By checking that these gluings behave well on triple overlaps (see Exercise 3.7.D), finish the definition of the scheme \(\mathrm{Proj} S_{\bullet}\).

5.5.J. Exercise (Some will find this essential, others will prefer to ignore it). (Re)interpret the structure sheaf of \(\mathrm{Proj} S_{\bullet}\) in terms of compatible stalks.

5.5.4. Definition. We (re)define projective space (over a ring \(A\)) by \(\mathbb{P}^n_A := \mathrm{Proj} A[x_0, \ldots, x_n]\). This definition involves no messy gluing, or special choice of patches.

5.5.K. Exercise. Check that this agrees with our earlier construction of \(\mathbb{P}^n_A\) (Definition 5.4.9). (How do you know that the \(D(x_i)\) cover \(\mathrm{Proj} A[x_0, \ldots, x_n]?)\)

Notice that with our old definition of projective space, it would have been a nontrivial exercise to show that \(D(x^2 + y^2 - z^2) \subset \mathbb{P}^2_A\) (the complement of a plane conic) is affine; with our new perspective, it is immediate — it is \(\mathrm{Spec}(k[x, y, z]/(x^2 + y^2 - z^2))_0\).

5.5.L. Exercise. Both parts of this problem ask you to figure out the “right definition” of the vanishing scheme, in analogy with \(\mathbb{V}(\cdot)\) defined earlier. In both cases, you will be defining a closed subscheme, a notion we will introduce in 9.4.11 (a (the most important part) If \(S_{\bullet}\) is generated in degree 1, and \(f \in S_{\bullet}^+\) is homogeneous, explain how to define \(\mathbb{V}(f)\) “in” \(\mathrm{Proj} S_{\bullet}\), the vanishing scheme of \(f\). (Warning: \(f\) in general isn’t a function on \(\mathrm{Proj} S_{\bullet}\). We will later interpret it as something close: a section of a line bundle.) Hence define \(\mathbb{V}(f)\) for any homogeneous ideal \(I\) of \(S_{\bullet}^+\).

(b) (harder) If \(S_{\bullet}\) is a graded ring over \(A\), but not necessarily generated in degree 1, explain how to define the vanishing scheme \(\mathbb{V}(f)\) “in” \(\mathrm{Proj} S_{\bullet}\). (Hint: On \(D(g)\), let \(\mathbb{V}(f)\) be cut out by all degree 0 equations of the form \(fh/g^n\), where \(n \in \mathbb{Z}^+\), and
is homogeneous. Show that this gives a well defined closed subscheme. Your calculations will mirror those of Exercise 5.5.H.)

5.5.5. Projective and quasiprojective schemes.
We call a scheme of the form \( \text{Proj} \ S_\bullet \), where \( S_\bullet \) is a finitely generated graded ring over \( A \), a \textit{projective scheme over} \( A \), or a \textit{projective} \( A \)-scheme. A \textit{quasiprojective} \( A \)-scheme is a quasicompact open subscheme of a projective \( A \)-scheme. The "\( A \)" is omitted if it is clear from the context; often \( A \) is a field.

5.5.6. \textit{Unimportant} remarks. (1) Note that \( \text{Proj} \ S_\bullet \) makes sense even when \( S_\bullet \) is not finitely generated. This can — rarely — be useful. But having this more general construction can make things easier. For example, you will later be able to do Exercise 7.4.D without worrying about Exercise 7.4.H.)

(2) The quasicompact requirement in the definition quasiprojectivity is of course redundant in the Noetherian case (cf. Exercise 4.6.M), which is all that matters to most.

5.5.7. \textit{Silly} example. Note that \( \mathbb{P}^0_A = \text{Proj} \ A[T] \cong \text{Spec} A \). Thus "Spec \( A \) is a projective \( A \)-scheme".

5.5.8. \textbf{Example:} \( \mathbb{P}V \). We can make this definition of projective space even more choice-free as follows. Let \( V \) be an \( (n+1) \)-dimensional vector space over \( k \). (Here \( k \) can be replaced by any ring \( A \) as usual.) Define

\[
\text{Sym}^*V^\vee = k \oplus V^\vee \oplus \text{Sym}^2 V^\vee \oplus \cdots
\]

(The reason for the dual is explained by the next exercise.) If for example \( V \) is the dual of the vector space with basis associated to \( x_0, \ldots, x_n \), we would have \( \text{Sym}^* V^\vee = k[x_0, \ldots, x_n] \). Then we can define \( \mathbb{P}V := \text{Proj} \text{Sym}^* V^\vee \). In this language, we have an interpretation for \( x_0, \ldots, x_n \): they are the linear functionals on the underlying vector space \( V \).

5.5.M. \textbf{Unimportant} exercise. Suppose \( k \) is algebraically closed. Describe a natural bijection between one-dimensional subspaces of \( V \) and the points of \( \mathbb{P}V \). Thus this construction canonically (in a basis-free manner) describes the one-dimensional subspaces of the vector space Spec \( V \).

Unimportant remark: you may be surprised at the appearance of the dual in the definition of \( \mathbb{P}V \). This is explained by the previous exercise. Most normal (traditional) people define the projectivization of a vector space \( V \) to be the space of one-dimensional subspaces of \( V \). Grothendieck considered the projectivization to be the space of one-dimensional quotients. One motivation for this is that it gets rid of the annoying dual in the definition above. There are better reasons, that we won’t go into here. In a nutshell, quotients tend to be better-behaved than subobjects for coherent sheaves, which generalize the notion of vector bundle. (We will discuss them in Chapter 14.)

On another note related to Exercise 5.5.M you can also describe a natural bijection between points of \( V \) and the points of Spec \( \text{Sym}^* V^\vee \). This construction respects the affine/projective cone picture of \( 9.2.10 \).

5.5.9. \textit{The Grassmannian}. At this point, we could describe the fundamental geometric object known as the Grassmannian, and give the “wrong” definition of it.
We will instead wait until §7.7 to give the wrong definition, when we will know enough to sense that something is amiss. The right definition will be given in §17.6.
CHAPTER 6

Some properties of schemes

6.1 Topological properties

We will now define some useful properties of schemes. The definitions of irreducible, irreducible component, closed point, specialization, generalization, generic point, connected, connected component, and quasicompact were given in §4.5–4.6. You should have pictures in your mind of each of these notions.

Exercise 4.6.O shows that \( \mathbb{A}^n \) is irreducible (it was easy). This argument “behaves well under gluing”, yielding:

6.1.A. EASY EXERCISE. Show that \( \mathbb{P}^n_k \) is irreducible.

6.1.B. EXERCISE. Exercise 4.7.E showed that there is a bijection between irreducible closed subsets and points. Show that this is true of schemes as well.

6.1.C. EASY EXERCISE. Prove that if \( X \) is a scheme that has a finite cover \( X = \bigcup_{i=1}^n \text{Spec}A_i \) where \( A_i \) is Noetherian, then \( X \) is a Noetherian topological space (§4.6.3). (We will soon call such a scheme a Noetherian scheme, [6.3.4])

Thus \( \mathbb{P}^n_k \) and \( \mathbb{P}^n_Z \) are Noetherian topological spaces: we built them by gluing together a finite number of spectra of Noetherian rings.

6.1.D. EASY EXERCISE. Show that a scheme \( X \) is quasicompact if and only if it can be written as a finite union of affine schemes. (Hence \( \mathbb{P}^n_k \) is quasicompact.)

6.1.E. GOOD EXERCISE: QUASICOMPACT SCHEMES HAVE CLOSED POINTS. Show that if \( X \) is a quasicompact scheme, then every point has a closed point in its closure. In particular, every nonempty quasicompact scheme has a closed point. (Warning: there exist non-empty schemes with no closed points, so your argument had better use the quasicompactness hypothesis! We will see that in good situations, the closed points are dense, Exercise 6.3.E.)

6.1.1. Quasiseparatedness. Quasiseparatedness is a weird notion that comes in handy for certain people. (Warning: we will later realize that this is really a property of morphisms, not of schemes [8.3.1].) Most people, however, can ignore this notion, as the schemes they will encounter in real life will all have this property. A topological space is quasiseparated if the intersection of any two quasicompact open sets is quasicompact. Thus a scheme is quasiseparated if the intersection of any two affine open subsets is a finite union of affine open subsets.
6.1.F. **SHORT EXERCISE.** Prove this equivalence.

We will see later that this will be a useful hypothesis in theorems (in conjunction with quasicompactness), and that various interesting kinds of schemes (affine, locally Noetherian, separated, see Exercises [6.1.G][6.3.B] and [11.1.G] resp.) are quasiseparated, and this will allow us to state theorems more succinctly (e.g. “if X is quasicompact and quasiseparated” rather than “if X is quasicompact, and either this or that or the other thing hold”).

6.1.G. **EXERCISE.** Show that affine schemes are quasiseparated.

“Quasicompact and quasiseparated” means something concrete:

6.1.H. **EXERCISE.** Show that a scheme X is quasicompact and quasiseparated if and only if X can be covered by a finite number of affine open subsets, any two of which have intersection also covered by a finite number of affine open subsets.

So when you see “quasicompact and quasiseparated” as hypotheses in a theorem, you should take this as a clue that you will use this interpretation, and that finiteness will be used in an essential way.

6.1.I. **EASY EXERCISE.** Show that all projective A-schemes are quasicompact and quasiseparated. (Hint: use the fact that the graded ring in the definition is finitely generated — those finite number of generators will lead you to a covering set.)

6.1.J. **Dimension.** One very important topological notion is dimension. (It is amazing that this is a topological idea.) But despite being intuitively fundamental, it is more difficult, so we will put it off until Chapter [12].

### 6.2 Reducedness and integrality

Recall that one of the alarming things about schemes is that functions are not determined by their values at points, and that was because of the presence of nilpotents ([4.2.9]).

6.2.1. **Definition.** A ring is said to be reduced if it has no nonzero nilpotents ([4.2.11]). A scheme X is reduced if \( \mathcal{O}_X(U) \) is reduced for every open set U of X.

An example of a nonreduced affine scheme is \( \text{Spec} \ k[x, y]/(y^2, xy) \). A useful representation of this scheme is given in Figure [6.1], although we will only explain in [6.5] why this is a good picture. The fuzz indicates that there is some nonreducedness going on at the origin. Here are two different functions: x and x + y. Their values agree at all points (all closed points \( [[x - a, y]] = (a, 0) \) and at the generic point \( [[y]] \)). They are actually the same function on the open set D(x), which is not surprising, as D(x) is reduced, as the next exercise shows. (This explains why the fuzz is only at the origin, where \( y = 0 \).)

6.2.A. **EXERCISE.** Show that \( (k[x, y]/(y^2, xy))_x \) has no nilpotents. (Possible hint: show that it is isomorphic to another ring, by considering the geometric picture. Exercise [4.2.1] may give another hint.)
Figure 6.1. A picture of the scheme $\text{Spec } k[x, y]/(y^2, xy)$. The fuzz indicates where “the non-reducedness lives”.

6.2.B. Exercise (reducedness is a stalk-local property, i.e. can be checked at stalks). Show that a scheme is reduced if and only if none of the stalks have nilpotents. Hence show that if $f$ and $g$ are two functions on a reduced scheme that agree at all points, then $f = g$. (Two hints: $\mathcal{O}_X(U) \hookrightarrow \prod_{x \in U} \mathcal{O}_{X, x}$ from Exercise 3.4.A, and the nilradical is intersection of all prime ideals from Theorem 4.2.10.)

We remark that the fuzz in Figure 6.1 indicates the points where there is nonreducedness.

6.2.C. Exercise (cf. Exercise 6.1.E). If $X$ is a quasicompact scheme, show that it suffices to check reducedness at closed points. (Hint: Show that any point of a quasicompact scheme has a closed point in its closure.)

Warning for experts: if a scheme $X$ is reduced, then it is immediate from the definition that its ring of global sections is reduced. However, the converse is not true.

6.2.D. Exercise. Suppose $X$ is quasicompact, and $f$ is a function (a global section of $\mathcal{O}_X$) that vanishes at all points of $x$. Show that there is some $n$ such that $f^n = 0$. Show that this may fail if $X$ is not quasicompact. (This exercise is less important, but shows why we like quasicompactness, and gives a standard pathology when quasicompactness doesn’t hold.) Hint: take an infinite disjoint union of $\text{Spec } A_n$ with $A_n := k[\epsilon]/\epsilon^n$.

Definition. A scheme $X$ is integral if it is nonempty, and $\mathcal{O}_X(U)$ is an integral domain for every nonempty open set $U$ of $X$.

6.2.E. Important Exercise. Show that a scheme $X$ is integral if and only if it is irreducible and reduced.

6.2.F. Exercise. Show that an affine scheme $\text{Spec } A$ is integral if and only if $A$ is an integral domain.

6.2.G. Exercise. Suppose $X$ is an integral scheme. Then $X$ (being irreducible) has a generic point $\eta$. Suppose $\text{Spec } A$ is any non-empty affine open subset of $X$. Show that the stalk at $\eta$, $\mathcal{O}_{X, \eta}$, is naturally $K(A)$, the fraction field of $A$. This is called the function field $K(X)$ of $X$. It can be computed on any non-empty open set of $X$, as any such open set contains the generic point.

6.2.H. Exercise. Suppose $X$ is an integral scheme. Show that the restriction maps $\text{res}_{U, V} : \mathcal{O}_X(U) \to \mathcal{O}_X(V)$ are inclusions so long as $V \neq \emptyset$. Suppose $\text{Spec } A$ is any non-empty affine open subset of $X$ (so $A$ is an integral domain). Show that the natural map $\mathcal{O}_X(U) \to \mathcal{O}_{X, \eta} = K(A)$ (where $U$ is any non-empty open set) is an
inclusion. Thus irreducible varieties (an important example of integral schemes defined later) have the convenient property that sections over different open sets can be considered subsets of the same ring. Thus restriction maps (except to the empty set) are always inclusions, and gluing is easy: functions \( f_i \) on a cover \( U_i \) of \( U \) (as \( i \) runs over an index set) glue if and only if they are the same element of \( K(X) \). This is one reason why (irreducible) varieties are usually introduced before schemes.

Integrality is not stalk-local (the disjoint union of two integral schemes is not integral, as Spec \( A \setminus \bigcup \) Spec \( B = \text{Spec} \times B \), cf. Exercise 4.6.T), but it almost is, as is shown in the following believable exercise.

**6.2.1. Unimportant Exercise.** Show that a locally Noetherian scheme \( X \) is integral if and only if \( X \) is connected and all stalks \( \mathcal{O}_{X,p} \) are integral domains. Thus in “good situations” (when the scheme is Noetherian), integrality is the union of local (stalks are integral domains) and global (connected) conditions.

### 6.3 Properties of schemes that can be checked “affine-locally”

This section is intended to address something tricky in the definition of schemes. We have defined a scheme as a topological space with a sheaf of rings, that can be covered by affine schemes. Hence we have all of the affine open sets in the cover, but we don’t know how to communicate between any two of them. Somewhat more explicitly, if I have an affine cover, and you have an affine cover, and we want to compare them, and I calculate something on my cover, there should be some way of us getting together, and figuring out how to translate my calculation over to your cover. The Affine Communication Lemma 6.3.2 will provide a convenient machine for doing this.

Thanks to this lemma, we can define a host of important properties of schemes. All of these are “affine-local” in that they can be checked on any affine cover, i.e. a covering by open affine sets. We like such properties because we can check them using any affine cover we like. If the scheme in question is quasicompact, then we need only check a finite number of affine open sets.

**6.3.1. Proposition.** Suppose \( \text{Spec} \, A \) and \( \text{Spec} \, B \) are affine open subschemes of a scheme \( X \). Then \( \text{Spec} \, A \cap \text{Spec} \, B \) is the union of open sets that are simultaneously distinguished open subschemes of \( \text{Spec} \, A \) and \( \text{Spec} \, B \).

**Proof.** (See Figure 6.2 for a sketch.) Given any point \( p \in \text{Spec} \, A \cap \text{Spec} \, B \), we produce an open neighborhood of \( p \) in \( \text{Spec} \, A \cap \text{Spec} \, B \) that is simultaneously distinguished in both \( \text{Spec} \, A \) and \( \text{Spec} \, B \). Let \( \text{Spec} \, A_f \) be a distinguished open subset of \( \text{Spec} \, A \) contained in \( \text{Spec} \, A \cap \text{Spec} \, B \) and containing \( p \). Let \( \text{Spec} \, B_g \) be a distinguished open subset of \( \text{Spec} \, B \) contained in \( \text{Spec} \, A_f \) and containing \( p \). Then \( g \in \Gamma(\text{Spec} \, B, \mathcal{O}_X) \) restricts to an element \( g' \in \Gamma(\text{Spec} \, A_f, \mathcal{O}_X) = A_f \). The points of \( \text{Spec} \, A_f \) where \( g \) vanishes are precisely the points of \( \text{Spec} \, A_f \) where \( g' \) vanishes, so

\[
\text{Spec} \, B_g = \text{Spec} \, A_f \setminus \{[g] : g' \in p\} = \text{Spec} \, (A_f)_{g'}.
\]
Figure 6.2. A trick to show that the intersection of two affine open sets may be covered by open sets that are simultaneously distinguished in both affine open sets.

If \( g' = g''/f^n \) (\( g'' \in A \)) then \( \text{Spec}(A_f)g' = \text{Spec}Afg'' \), and we are done. \( \square \)

The following easy result will be crucial for us.

6.3.2. Affine Communication Lemma. — Let \( P \) be some property enjoyed by some affine open sets of a scheme \( X \), such that

(i) if an affine open set \( \text{Spec} A \hookrightarrow X \) has property \( P \) then for any \( f \in A \), \( \text{Spec} A_f \hookrightarrow X \) does too.

(ii) if \( (f_1, \ldots, f_n) = A \), and \( \text{Spec} A_{f_i} \hookrightarrow X \) has \( P \) for all \( i \), then so does \( \text{Spec} A \hookrightarrow X \).

Suppose that \( X = \bigcup_{i \in I} \text{Spec} A_i \) where \( \text{Spec} A_i \) has property \( P \). Then every open affine subset of \( X \) has \( P \) too.

We say such a property is **affine-local**. Note that any property that is stalk-local (a scheme has property \( P \) if and only if all its stalks have property \( Q \)) is necessarily affine-local (a scheme has property \( P \) if and only if all of its affines have property \( R \), where an affine scheme has property \( R \) if and only if and only if all its stalks have property \( Q \)), but it is sometimes not so obvious what the right definition of \( Q \) is; see for example the discussion of normality in the next section.

**Proof.** Let \( \text{Spec} A \) be an affine subscheme of \( X \). Cover \( \text{Spec} A \) with a finite number of distinguished open sets \( \text{Spec} A_{f_i} \), each of which is distinguished in some \( \text{Spec} A_i \). This is possible by Proposition 6.3.1 and the quasicompactness of \( \text{Spec} A \) (Exercise 4.6.D(a)). By (i), each \( \text{Spec} A_{f_i} \) has \( P \). By (ii), \( \text{Spec} A \) has \( P \). \( \square \)

By choosing property \( P \) appropriately, we define some important properties of schemes.

6.3.3. Proposition. — Suppose \( A \) is a ring, and \( (f_1, \ldots, f_n) = A \).

(a) If \( A \) is a Noetherian ring, then so is \( A_{f_i} \). If each \( A_{f_i} \) is Noetherian, then so is \( A \).

(b) If \( A \) is reduced, then \( A_{f_i} \) is also reduced. If each \( A_{f_i} \) is reduced, then so is \( A \).

(c) Suppose \( B \) is a ring, and \( A \) is a \( B \)-algebra. (Hence \( A_g \) is a \( B \)-algebra for all \( g \in A \).) If \( A \) is a finitely generated \( B \)-algebra, then so is \( A_{f_i} \). If each \( A_{f_i} \) is a finitely-generated \( B \)-algebra, then so is \( A \).

We will prove these shortly (6.3.8). But let’s first motivate you to read the proof by giving some interesting definitions assuming Proposition 6.3.3 is true.
6.3.4. **Important Definition.** Suppose $X$ is a scheme. If $X$ can be covered by affine open sets $\text{Spec } A$ where $A$ is Noetherian, we say that $X$ is a **locally Noetherian scheme**. If in addition $X$ can be quasicompact, or equivalently covered by finitely many such affine open sets, we say that $X$ is a **Noetherian scheme**. (We will see a number of definitions of the form “if $X$ has this property, we say that it is locally $Q$; if further $X$ is quasicompact, we say that it is $Q$.”) By Exercise 6.1.C the underlying topological space of a Noetherian scheme is Noetherian.

6.3.A. **Exercise.** Show that all open subsets of a Noetherian topological space (hence a Noetherian scheme) are quasicompact.

6.3.B. **Exercise.** Show that locally Noetherian schemes are quasiseparated.

6.3.C. **Exercise.** Show that a Noetherian scheme has a finite number of irreducible components. Show that a Noetherian scheme has a finite number of connected components, each a finite union of irreducible components.

6.3.D. **Exercise.** Show that $X$ is reduced if and only if $X$ can be covered by affine open sets $\text{Spec } A$ where $A$ is reduced.

Our earlier definition of reducedness required us to check that the ring of functions over any open set is nilpotent-free. Our new definition lets us check a single affine cover. Hence for example $\mathbb{A}^n_k$ and $\mathbb{P}^n_k$ are reduced.

6.3.5. **Schemes over a given field, or more generally over a given ring ($A$-schemes).** You may be particularly interested in working over a particular field, such as $\mathbb{C}$ or $\mathbb{Q}$, or over a ring such as $\mathbb{Z}$. Motivated by this, we define the notion of $A$-scheme, or scheme over $A$, where $A$ is a ring, as a scheme where all the rings of sections of the structure sheaf (over all open sets) are $A$-algebras, and all restriction maps are maps of $A$-algebras. (Like some earlier notions such as quasiseparateness, this will later in Exercise 7.3.G be properly understood as a “relative notion”; it is the data of a morphism $X \to \text{Spec } A$.) Suppose now $X$ is an $A$-scheme. If $X$ can be covered by affine open sets $\text{Spec } B_i$ where each $B_i$ is a finitely generated $A$-algebra, we say that $X$ is locally of finite type over $A$, or that it is a **locally of finite type $A$-scheme**. (This is admittedly cumbersome terminology; it will make more sense later, once we know about morphisms in §8.3.9.) If furthermore $X$ is quasicompact, $X$ is (of) **finite type over $A$**, or a **finite type $A$-scheme**. Note that a scheme locally of finite type over $k$ or $\mathbb{Z}$ (or indeed any Noetherian ring) is locally Noetherian, and similarly a scheme of finite type over any Noetherian ring is Noetherian. As our key “geometric” examples: (i) $\text{Spec } \mathbb{C}[x_1, \ldots, x_n]/I$ is a finite-type $\mathbb{C}$-scheme; and (ii) $\mathbb{P}^n_\mathbb{C}$ is a finite type $\mathbb{C}$-scheme. (The field $\mathbb{C}$ may be replaced by an arbitrary ring $A$.)

6.3.6. **Varieties.** We now make a connection to the classical language of varieties. An affine scheme that is a reduced and of finite type $k$-scheme is said to be an **affine variety (over $k$)**, or an **affine $k$-variety**. A reduced (quasi-)projective $k$-scheme is a **(quasi-)projective variety (over $k$)**, or a **(quasi-)projective $k$-variety**. (Warning: in the literature, it is sometimes also assumed in the definition of variety that the scheme is irreducible, or that $k$ is algebraically closed.) We will not define varieties in general until §11.1.7 we will need the notion of separatedness first, to exclude abominations like the line with the doubled origin (Example 5.4.5). But
many of the statements we will make in this section about affine $k$-varieties will automatically apply more generally to $k$-varieties.

6.3.E. Exercise. Show that a point of a locally finite type $k$-scheme is a closed point if and only if the residue field of the stalk of the structure sheaf at that point is a finite extension of $k$. (Hint: the Nullstellensatz [4.2.3]) Show that the closed points are dense on such a scheme (even though they needn’t be quasicompact, cf. Exercise [6.1.E].) (For another exercise on closed points, see [6.1.E] Warning: closed points need not be dense even on quite reasonable schemes, such as that of Exercise [4.4.G].)

6.3.7. Definition. The degree of a closed point of a locally finite type $k$-scheme is the degree of this field extension. For example, in $A_k^1 = \text{Spec} k[t]$, the point $[k[t]/p(t)]$ (p irreducible) is deg $p$. If $k$ is algebraically closed, the degree of every closed point is 1.

6.3.8. Proof of Proposition 6.3.3 We divide each part into (i) and (ii) following the statement of the Affine Communication Lemma 6.3.2 (a) (i) If $I_1 \subseteq I_2 \subseteq I_3 \subseteq \cdots$ is a strictly increasing chain of ideals of $A_{f_i}$, then we can verify that $J_1 \subseteq J_2 \subseteq J_3 \subseteq \cdots$ is a strictly increasing chain of ideals of $A$, where

$$J_j = \{ r \in A : r \in I_j \}$$

where $r \in I_j$ means “the image in $A_{f_i}$ lies in $I_j$”. (We think of this as $I_j \cap A$, except in general $A$ needn’t inject into $A_{f_i}$.) Clearly $J_j$ is an ideal of $A$. If $x/f^n \in I_{j+1} \setminus I_j$ where $x \in A$, then $x \in J_{j+1}$, and $x \notin J_j$ (or else $x(1/f)^n \in J_j$ as well). (ii) Suppose $I_1 \subseteq I_2 \subseteq I_3 \subseteq \cdots$ is a strictly increasing chain of ideals of $A$. Then for each $1 \leq i \leq n$,

$$I_{i,1} \subset I_{i,2} \subset I_{i,3} \subset \cdots$$

is an increasing chain of ideals in $A_{f_i}$, where $I_{i,j} = I_j \otimes_A A_{f_i}$. It remains to show that for each $j$, $I_{i,j} \subseteq I_{i,j+1}$ for some $i$; the result will then follow.

6.3.F. Exercise. Finish this argument.


(c) (i) is clear: if $A$ is generated over $B$ by $r_1, \ldots, r_n$, then $A_f$ is generated over $B$ by $r_1, \ldots, r_n, 1/f$.

(ii) Here is the idea. As the $f_i$ generate $A$, we can write $1 = \sum c_i f_i$ for $c_i \in A$. We have generators of $A_i$: $r_{ij}/f_i^{l_{ij}}$, where $r_{ij} \in A_i$. I claim that $\{f_i\} \cup \{c_i\} \cup \{r_{ij}\}$ generate $A$ as a $B$-algebra. Here’s why. Suppose you have any $r \in A_i$. Then in $A_{f_i}$, we can write $r$ as some polynomial in the $r_{ij}$’s and $f_i$, divided by some huge power of $f_i$. So “in each $A_{f_i}$, we have described $r$ in the desired way”, except for this annoying denominator. Now use a partition of unity type argument as in the proof of Theorem 5.1.2 to combine all of these into a single expression, killing the denominator. Show that the resulting expression you build still agrees with $r$ in each of the $A_{f_i}$. Thus it is indeed $r$.

6.3.H. Exercise. Make this argument precise.

This concludes the proof of Proposition 6.3.3. □
6.3.1. Easy Exercise. Suppose $S_*$ is a finitely generated graded ring over $A$. Show that $\text{Proj } S_*$ is of finite type over $A = S_0$. If $S_0$ is a Noetherian ring, show that $\text{Proj } S_*$ is a Noetherian scheme, and hence that $\text{Proj } S_*$ has a finite number of irreducible components. Show that any quasiprojective scheme is locally of finite type over $A$. If $A$ is Noetherian, show that $\text{Proj } S_*$ is a Noetherian scheme, and hence that $\text{Proj } S_*$ has a finite number of irreducible components. Show that any quasiprojective scheme is locally of finite type over $A$. If $A$ is Noetherian, show that any quasiprojective $A$-scheme is quasicompact, and hence of finite type over $A$. Show this need not be true if $A$ is not Noetherian. Better: give an example of a quasiprojective $A$-scheme that is not quasicompact, necessarily for some non-Noetherian $A$. (Hint: Silly example 5.5.7)

6.4 Normality and factoriality

6.4.1. Normality.

We can now define a property of schemes that says that they are “not too far from smooth”, called normality, which will come in very handy. We will see later that “locally Noetherian normal schemes satisfy Hartogs’ Lemma” (Algebraic Hartogs’ Lemma 12.3.10 for Noetherian normal schemes): functions defined away from a set of codimension 2 extend over that set. (We saw a first glimpse of this in §5.4.2.) As a consequence, rational functions that have no poles (certain sets of codimension one where the function isn’t defined) are defined everywhere. We need definitions of dimension and poles to make this precise.

A scheme $X$ is normal if all of its stalks $\mathcal{O}_{X,p}$ are normal, i.e. are integral domains, and integrally closed in their fraction fields. (An integral domain $A$ is integrally closed if the only zeros in $K(A)$ to any monic polynomial in $A[x]$ must lie in $A$ itself. The basic example is $\mathbb{Z}$.) As reducedness is a stalk-local property (Exercise 6.2.B), normal schemes are reduced.

6.4.A. Exercise. Show that integrally closed domains behave well under localization: if $A$ is an integrally closed domain, and $S$ is a multiplicative subset, show that $S^{-1}A$ is an integrally closed domain. (Hint: assume that $x^n + a_{n-1}x^{n-1} + \cdots + a_0 = 0$ where $a_i \in S^{-1}A$ has a root in the fraction field. Turn this into another equation in $A[x]$ that also has a root in the fraction field.)

It is no fun checking normality at every single point of a scheme. Thanks to this exercise, we know that if $A$ is an integrally closed domain, then $\text{Spec } A$ is normal. Also, for quasicompact schemes, normality can be checked at closed points, thanks to this exercise, and the fact that for such schemes, any point is a generalization of a closed point (see Exercise 6.1.E).

It is not true that normal schemes are integral. For example, the disjoint union of two normal schemes is normal. Thus $\text{Spec } k \coprod \text{Spec } k \cong \text{Spec } (k \times k) \cong \text{Spec } k[x]/(x(x - 1))$ is normal, but its ring of global sections is not an integral domain.

6.4.B. Unimportant exercise. Show that a Noetherian scheme is normal if and only if it is the finite disjoint union of integral Noetherian normal schemes. (Hint: Exercise 6.2.B)

We are close to proving a useful result in commutative algebra, so we may as well go all the way.
6.4.2. Proposition. — If A is an integral domain, then the following are equivalent.

(1) A is integrally closed.
(2) \( A_p \) is integrally closed for all prime ideals \( p \subset A \).
(3) \( A_m \) is integrally closed for all maximal ideals \( m \subset A \).

Proof. Exercise 6.4.A shows that integral closure is preserved by localization, so (1) implies (2). Clearly (2) implies (3).

It remains to show that (3) implies (1). This argument involves a pretty construction that we will use again. Suppose \( A \) is not integrally closed. We show that there is some \( m \) such that \( A_m \) is also not integrally closed. Suppose

\[
(6.4.2.1) \quad x^n + a_{n-1}x^{n-1} + \cdots + a_0 = 0
\]

(with \( a_i \in A \)) has a solution \( s \) in \( K(A) \setminus A \). Let \( I \) be the ideal of denominators of \( s \):

\[
I := \{ r \in A : rs \in A \}.
\]

(Note that \( I \) is clearly an ideal of \( A \).) Now \( I \neq A \), as \( 1 \notin I \). Thus there is some maximal ideal \( m \) containing \( I \). Then \( s \notin A_m \), so equation (6.4.2.1) in \( A_m[x] \) shows that \( A_m \) is not integrally closed as well, as desired. □

6.4.C. UNIMPORTANT EXERCISE. If \( A \) is an integral domain, show that \( A = \cap A_m \), where the intersection runs over all maximal ideals of \( A \). (We won’t use this exercise, but it gives good practice with the ideal of denominators.)

6.4.D. UNIMPORTANT EXERCISE RELATING TO THE IDEAL OF DENOMINATORS. One might naively hope from experience with unique factorization domains that the ideal of denominators is principal. This is not true. As a counterexample, consider our new friend \( A = k[a, b, c, d]/(ad - bc) \) (which we last saw in Example 5.4.11, and which we will later recognize as the cone over the quadric surface), and \( a/c = b/d \in K(A) \). Show that \( I = \langle c, d \rangle \). We will soon see that it is not principal (Exercise 13.1.C).

6.4.3. Factoriality.

We define a notion which implies normality.

6.4.4. Definition. If all the stalks of a scheme \( X \) are unique factorization domains, we say that \( X \) is factorial.

6.4.E. EXERCISE. Show that any localization of a unique factorization domain is a unique factorization domain.

Thus if \( A \) is a unique factorization domain, then \( \text{Spec} \ A \) is factorial. (The converse need not hold. This property is not affine-local, see Exercise 6.4.1.) In fact, we will see that elliptic curves are factorial, yet no affine open set is the Spec of a unique factorization domain, \( (21.10.1) \). Hence it suffices to check factoriality by finding an appropriate affine cover.

6.4.5. ** How to check if a ring is a unique factorization domain. We note here that there are very few means of checking that a Noetherian integral domain is a unique factorization domain. Some useful ones are: (0) elementary means (rings with a euclidean algorithm such as \( \mathbb{Z} \), \( k[t] \), and \( \mathbb{Z}[i] \); polynomial rings over a unique
factorization domain, by Gauss’ Lemma).  (1) Exercise 6.4.E that the localization of a unique factorization domain is also a unique factorization domain. (2) height 1 primes are principal (Proposition 12.3.5). (3) Nagata’s Lemma (Exercise 15.2.S). (4) normal and Cl = 0 (Exercise 15.2.Q).

One of the reasons we like factoriality is that it implies normality.

6.4.F. IMPORTANT EXERCISE. Show that unique factorization domains are integrally closed. Hence factorial schemes are normal, and if A is a unique factorization domain, then Spec A is normal. (However, rings can be integrally closed without being unique factorization domains, as we will see in Exercise 13.1.D. An example without proof: Exercise 6.4.L.)

6.4.G. EASY EXERCISE. Show that the following schemes are normal: $A^n_k$, $P^n_k$, Spec $\mathbb{Z}$. (As usual, $k$ is a field. Although it is true that if A is integrally closed then $A[x]$ is as well [B Ch. 5, §1, no. 3, Cor. 2], this is not an easy fact, so do not use it here.)

6.4.H. HANDY EXERCISE (YIELDING MANY OF ENLIGHTENING EXAMPLES LATER). Suppose A is a unique factorization domain with 2 invertible, $f \in A$ has no repeated prime factors, and $z^2 - f$ is irreducible in $A[z]$. Show that Spec $A[z]/(z^2 - f)$ is normal. Show that if $f$ is not square-free, then Spec $A[z]/(z^2 - f)$ is not normal. (Hint: B := $A[z]/(z^2 - f)$ is an integral domain, as $(z^2 - f)$ is prime in $A[z]$. Suppose we have monic $F(T) \in B[T]$ so that $F(T) = 0$ has a root $\alpha$ in $K(B)$. Then by replacing $F(T)$ by $\overline{F}(T)F(T)$, we can assume $F(T) \in A[T]$. Also, $\alpha = g + hz$ where $g, h \in K(A)$. Now $\alpha$ is the root of $Q(T) = 0$ for monic $Q(T) = T^2 - 2gT + (g^2 - h^2f) \in K(A)[T]$, so we can factor $F(T) = P(T)Q(T)$ in $K(A)[T]$. By Gauss’ lemma, $2g, g^2 - h^2f \in A$. Say $g = r/2, h = s/t$ (s and t have no common factors, $r, s, t \in A$). Then $g^2 - h^2f = (r^2t^2 - 4s^2f)/4t^2$. Then t is a unit, and r is even.)

6.4.I. EXERCISE. Show that the following schemes are normal:

(a) Spec $\mathbb{Z}[x]/(x^2 - n)$ where n is a square-free integer congruent to 3 (mod 4);
(b) Spec $k[x_1, \ldots, x_n]/(x_1^2 + x_2^2 + \cdots + x_m^2)$ where $\text{char } k \neq 2, m \geq 3$;
(c) Spec $k[w, x, y, z]/(wz - xy)$ where $\text{char } k \neq 2$ and $k$ is algebraically closed.

This is our cone over a quadric surface example from Exercises 5.4.11 and 6.4.D (Hint: the side remark below may help.)

This is a good time to define the rank of a quadratic form.

6.4.J. EXERCISE (DIAGONALIZING QUADRICS). Suppose $k$ is an algebraically closed field of characteristic not 2.

(a) Show that any quadratic form in $n$ variables can be “diagonalized” by changing coordinates to be a sum of at most $n$ squares (e.g. $uw - v^2 = ((u + w)/2)^2 + (w - w)/2^2 + (v^2)$, where the linear forms appearing in the squares are linearly independent. (Hint: use induction on the number of variables, by “completing the square” at each step.)
(b) Show that the number of squares appearing depends only on the quadric. For
example, $x^2 + y^2 + z^2$ cannot be written as a sum of two squares. (Possible approach: given a basis $x_1, \ldots, x_n$ of the linear forms, write the quadratic form as

$$
\begin{pmatrix}
  x_1 \\
  \vdots \\
  x_n
\end{pmatrix}
M
\begin{pmatrix}
  x_1 \\
  \vdots \\
  x_n
\end{pmatrix}
$$

where $M$ is a symmetric matrix. Determine how $M$ transforms under a change of basis, and show that the rank of $M$ is independent of the choice of basis.)

The rank of the quadratic form is the number of (“linearly independent”) squares needed.

6.4.K. Exercise. Suppose $A$ is a $k$-algebra where $\text{char } k = 0$, and $l/k$ is a finite field extension. Show that if $A \otimes_k l$ is normal (and in particular an integral domain) then $A$ is normal. (This is a case of a more general fact, and stated correctly, the converse is true.) Show that $\text{Spec } k[w, x, y, z]/(wz - xy)$ is normal if $k$ has characteristic 0. Possible hint: reduce to the case where $l/k$ is Galois.

6.4.L. Exercise (Factoriality is not affine-local). Let $A = (\mathbb{Q}[x, y]_{x^2 + y^2})_0$ denote the homogeneous degree 0 part of the ring $\mathbb{Q}[x, y]_{x^2 + y^2}$. In other words, it consists of quotients $f(x, y)/(x^2 + y^2)^n$, where $f$ has pure degree $2n$. Show that the distinguished open sets $D(x^2 + y^2)$ and $D(y^2 + x^2)$ cover $\text{Spec } A$. (Hint: the sum of those two fractions is 1.) Show that $A_{x^2 + y^2}$ and $A_{y^2 + x^2}$ are unique factorization domains. (Hint for the first: show that each ring is isomorphic to $\mathbb{Q}[t]_{t^2 + 1}$, where $t = y/x$; this is a localization of the unique factorization domain $\mathbb{Q}[t]$.) Finally, show that $A$ is not a unique factorization domain. Possible hint:

$$
\left( \frac{xy}{x^2 + y^2} \right)^2 = -\left( \frac{y^2}{x^2 + y^2} \right)^2.
$$

(This example didn’t come out of thin air; we will see $\text{Spec } A$ later as an example of a scheme with Picard group — or class group — $\mathbb{Z}/2$.)

6.5 Associated points of (locally Noetherian) schemes, and drawing fuzzy pictures

(This important topic won’t be used in an essential way for some time, certainly until we talk about dimension in Chapter 12, so it may be best skipped on a first reading. Better: read this section considering only the case where $A$ is an integral domain, or possibly a reduced Noetherian ring, thereby bypassing some of the annoyances. Then you will at least be comfortable with the notion of a rational function in these situations.)

Recall from just after Definition 6.2.1 (of reduced) our “fuzzy” pictures of the non-reduced scheme $\text{Spec } k[x, y]/(y^2, xy)$ (see Figure 6.1). When this picture was introduced, we mentioned that the “fuzz” at the origin indicated that the non-reduced behavior was concentrated there. This was verified in Exercise 6.2.A, and indeed the origin is the only point where the stalk of the structure sheaf is non-reduced.
You might imagine that in a bigger scheme, we might have different closed subsets with different amount of “non-reducedness”. This intuition will be made precise in this section. We will define associated points of a scheme, which will be the most important points of a scheme, encapsulating much of the interesting behavior of the structure sheaf. For example, in Figure 6.1, the associated points are the generic point of the x-axis, and the origin (where “the nonreducedness lives”).

The primes corresponding to the associated points of an affine scheme Spec A will be called associated primes of A. In fact this is backwards; we will define associated primes first, and then define associated points.

6.5.1. Properties of associated points. The properties of associated points that it will be most important to remember are as follows. Frankly, it is much more important to remember these facts than it is to remember their proofs. But we will, of course, prove these statements.

(0) They will exist for any locally Noetherian scheme, and for integral schemes. There are a finite number in any affine open set (and hence in any quasicompact open set). This will come for free.

(1) The generic points of the irreducible components of a locally Noetherian. The other associated points are called embedded points. Thus in Figure 6.1 the origin is the only embedded point. (By the way, there are easier analogues of these properties where Noetherian hypotheses are replaced by integral conditions; see Exercise 6.5.C.)

(2) If a locally Noetherian scheme X is reduced, then X has no embedded points. (This jibes with the intuition of the picture of associated points described earlier.) It follows from (1) and (2) that if X is integral (i.e. irreducible and reduced, Exercise 6.2.E), then the generic point is the only associated point.

(3) Recall that one nice property of integral schemes X (such as irreducible affine varieties) not shared by all schemes is that for any non-empty open U ⊂ X, the natural map \( \Gamma(U, \mathcal{O}_X) \to K(X) \) is an inclusion (Exercise 6.2.H). Thus all sections over any non-empty open set, and stalks, can be thought of as lying in a single field K(X), which is the stalk at the generic point.

More generally, if X is a locally Noetherian scheme, then for any U ⊂ X, the natural map

\[
\Gamma(U, \mathcal{O}_X) \to \bigoplus_{\text{associated } p \in U} \mathcal{O}_{X,p}
\]

is an injection.

We define a rational function on a scheme with associated points to be an element of the image of \( \Gamma(U, \mathcal{O}_U) \) in (6.5.1.1) for some U containing all the associated points. Equivalently, the set of rational functions is the colimit of \( \mathcal{O}_X(U) \) over all open sets containing the associated points. Thus if X is integral, the rational functions are the elements of the stalk at the generic point, and even if there is more than one associated point, it is helpful to think of them in this stalk-like manner. For example, in Figure 6.1, we think of \( \frac{x-2}{x(x-1)} \) as a rational function, but not \( \frac{x-2}{x(x-1)} \). The rational functions form a ring, called the total fraction ring of X, denoted \( Q(X) \). If \( X = \text{Spec } A \) is affine, then this ring is called the total fraction ring.
of $A$, and is denoted $Q(A)$. (But we will never use this notation.) If $X$ is integral, this is the function field $K(X)$, so this extends our earlier Definition 6.2.G of $K(\cdot)$. It can be more conveniently interpreted as follows, using the injectivity of (6.5.1.1).

A rational function is a function defined on an open set containing all associated points, i.e. an ordered pair $(U, f)$, where $U$ is an open set containing all associated points, and $f \in \Gamma(U, \mathcal{O}_X)$. Two such data $(U, f)$ and $(U', f')$ define the same open rational function if and only if the restrictions of $f$ and $f'$ to $U \cap U'$ are the same. If $X$ is reduced, this is the same as requiring that they are defined on an open set of each of the irreducible components. A rational function has a maximal domain of definition, because any two actual functions on an open set (i.e. sections of the structure sheaf over that open set) that agree as “rational functions” (i.e. on small enough open sets containing associated points) must be the same function, by the injectivity of (6.5.1.1). We say that a rational function $f$ is regular at a point $p$ if $p$ is contained in this maximal domain of definition (or equivalently, if there is some open set containing $p$ where $f$ is defined). For example, in Figure 6.1 the rational function $\frac{x-1}{(x-1)(x-3)}$ has domain of definition consisting of everything but 1 and 3 (i.e. $[(x-1)]$ and $[(x-3)]$), and is regular away from those two points.

The previous facts are intimately related to the following one.

(4) A function on an affine Noetherian scheme $X$ is a zero-divisor if and only if it vanishes at an associated point of $X$.

Motivated by the above four properties, when sketching (locally Noetherian) schemes, we will draw the irreducible components (the closed subsets corresponding to maximal associated points), and then draw “additional fuzz” precisely at the closed subsets corresponding to embedded points. All of our earlier sketches were of this form. (See Figure 6.3) The fact that these sketches “make sense” implicitly uses the fact that the notion of associated points behaves well with respect to open sets (and localization, cf. Theorem 6.5.3(d)).

Figure 6.3. This scheme has 6 associated points, of which 3 are embedded points. A function is a zero-divisor if it vanishes at one of these six points. It is nilpotent if it vanishes at all six of these points. (In fact, it suffices to vanish at the non-embedded associated points.)
6.5.A. Exercise (first practice with making fuzzy pictures). Assume the properties (1)–(4) of associated points. Suppose \( X \) is a closed subscheme of \( \mathcal{A} = \text{Spec} \mathbb{C}[x, y] \) with associated points at \([y - x^2], [(x - 1, y - 1)], [(x - 2, y - 2)]\).

(a) Sketch \( X \), including fuzz. (b) Do you have enough information to know if \( X \) is reduced? (c) Do you have enough information to know if \( x + y - 3 \)? How about \( y - x^2 \)? (Exercise 6.5.K will verify that such an \( X \) actually exists!)

We now finally define associated points, and show that they have the desired properties (1) through (4).

6.5.2. Definition. We work more generally with modules \( M \) over a ring \( A \). A prime \( p \subset A \) is associated to \( M \) if \( p \) is the annihilator of an element \( m \) of \( M \) (\( p = \{a \in A : am = 0\} \)). The set of primes associated to \( M \) is denoted \( \text{Ass} M \) (or \( \text{Ass}_A M \)).

Awkwardly, if \( I \) is an ideal of \( A \), the associated primes of the module \( A/I \) are said to be the associated primes of \( I \). This is not my fault.

6.5.B. Easy Exercise. Show that \( p \) is associated to \( M \) if and only if \( M \) has a submodule isomorphic to \( A/p \).

6.5.3. Theorem (properties of associated primes). — Suppose \( A \) is a Noetherian ring, and \( M \neq 0 \) is finitely generated.

(a) The set \( \text{Ass} M \) is finite and nonempty.
(b) The natural map \( M \to \prod_{p \in \text{Ass} M} \prod M_p \) is an injection.
(c) The set of zero-divisors of \( M \) is \( \cup_{p \in \text{Ass} M} p \).
(d) (association commutes with localization) If \( S \) is a multiplicative set, then

\[
\text{Ass}_{S^{-1} A} S^{-1} M = \text{Ass}_A M \cap \text{Spec} S^{-1} A
\]

(\( = \{p \in \text{Ass}_A M : p \cap S = \emptyset\} \)).
(e) The set \( \text{Ass} M \) contains the primes minimal among those containing \( \text{ann} M := \{a \in A : aM = 0\} \).

6.5.4. Definition. We define the associated points of a locally Noetherian scheme \( X \) to be those points \( p \in X \) such that, on any affine open set \( \text{Spec} \mathcal{A} \) containing \( p \), \( p \) corresponds to an associated prime of \( \mathcal{A} \). This notion is independent of choice of affine neighborhood \( \mathcal{A} \): if \( p \) has two affine open neighborhoods \( \mathcal{A} \) and \( \mathcal{B} \) (say corresponding to primes \( p \subset \mathcal{A} \) and \( q \subset \mathcal{B} \) respectively), then \( p \) corresponds to an associated prime of \( \mathcal{A} \) if and only if it corresponds to an associated prime of \( \mathcal{B} \).

6.5.C. Straightforward exercise. State and prove the analogues of (1)–(4) for schemes that are integral rather than locally Noetherian. State and prove the analogues of Theorem 6.5.3 where the hypothesis that \( A \) is Noetherian is replaced by the hypothesis that \( A \) is an integral domain.

6.5.D. Important Exercise. Show how Theorem 6.5.3 implies properties (0)–(4). (By (3), I mean the injectivity of (6.5.1.1). The trickiest is probably (2).)

We now prove Theorem 6.5.3.
6.5.E. Exercise. Suppose \( M \neq 0 \) is an \( A \)-module. Show that if \( I \subseteq A \) is maximal among all ideals that are annihilators of elements of \( M \), then \( I \) is prime, and hence \( I \in \text{Ass} \ M \). Thus if \( A \) is Noetherian, then \( \text{Ass} \ M \) is nonempty (part of Theorem 6.5.3(a)).

6.5.F. Exercise. Suppose that \( M \) is a module over a Noetherian ring \( A \). Show that \( m = 0 \) if and only if \( m \) is 0 in \( M_p \) for each of the maximal associated primes of \( M \). (Hint: use the previous exercise.)

This immediately implies Theorem 6.5.3(b). It also implies Theorem 6.5.3(c): Any nonzero element of \( \bigcup_{p \in \text{Ass} \ M} p \) is clearly a zero-divisor. Conversely, if \( a \) annihilates a nonzero element of \( M \), then \( r \) is contained in a maximal annihilator ideal.

6.5.G. Exercise. If \( 0 \rightarrow M' \rightarrow M \rightarrow M'' \rightarrow 0 \) is a short exact sequence of \( A \)-modules, show that

\[ \text{Ass} \ M' \subset \text{Ass} \ M \subset \text{Ass} \ M' \cup \text{Ass} \ M'' . \]

(Possible hint for the second containment: if \( m \in M \) has annihilator \( p \), then \( A_m = A/p \), cf. Exercise 6.5.B.)

6.5.H. Exercise. If \( M \) is a finitely generated module over Noetherian \( A \), show that \( M \) has a filtration

\[ 0 = M_0 \subset M_1 \subset \cdots \subset M_n = M \]

where \( M_{i+1}/M_i \cong R/p_i \) for some prime ideal \( p_i \). Show that the associated primes are among the \( p_i \), and thus prove Theorem 6.5.3(a).

6.5.I. Exercise. Prove Theorem 6.5.3(d) as follows.

(a) Show that

\[ \text{Ass}_A M \cap \text{Spec} S^{-1} A \subset \text{Ass}_{S^{-1} A} S^{-1} M. \]

(Hint: suppose \( p \in \text{Ass}_A M \cap \text{Spec} S^{-1} A \), with \( p = \text{ann} m \) for \( m \in M \).)

(b) Suppose \( q \in \text{Ass}_{S^{-1} A} S^{-1} M \), which corresponds to \( p \in A \) (i.e. \( q = p(S^{-1} A) \)). Then \( q = \text{ann}_{S^{-1} A} m \) (\( m \in S^{-1} M \)), which yields a nonzero element of

\[ \text{Hom}_{S^{-1} A}(S^{-1} A/q, S^{-1} M) . \]

Argue that this group is isomorphic to \( S^{-1} \text{Hom}_A(A/p, M) \) (see Exercise 2.6.G), and hence \( \text{Hom}_A(A/p, M) \neq 0 \).

6.5.J. Exercise. Prove Theorem 6.5.3(e) as follows. If \( p \) is minimal over \( \text{ann} \ M \), localize at \( p \), so that \( p \) is the only prime containing \( \text{ann} \ M \). Use Theorem 6.5.3(d).

6.5.K. Exercise. Let \( I = (y - x^2)^3 \cap (x - 1, y - 1)^{15} \cap (x - 2, y - 2) \). Show that \( X = \text{Spec} \mathbb{C}[x, y]/I \) satisfies the hypotheses of Exercise 6.5.A. (Side question: Is there a “smaller” example? Is there a “smallest”?)

6.5.5. Aside: Primary ideals. The notion of primary ideals is important, although we won’t use it. (An ideal \( I \subseteq A \) is in a ring is primary if \( I \neq A \) and if \( xy \in I \) implies either \( x \in I \) or \( y^n \in I \) for some \( n > 0 \).) The associated primes of an ideal turn out to be precisely those primes appearing in its primary decomposition. See [E] §3.3, for example, for more on this topic.
Part III

Morphisms of schemes
CHAPTER 7

Morphisms of schemes

7.1 Introduction

We now describe the morphisms between schemes. We will define some easy-to-state properties of morphisms, but leave more subtle properties for later.

Recall that a scheme is (i) a set, (ii) with a topology, (iii) and a (structure) sheaf of rings, and that it is sometimes helpful to think of the definition as having three steps. In the same way, the notion of morphism of schemes \( X \rightarrow Y \) may be defined (i) as a map of sets, (ii) that is continuous, and (iii) with some further information involving the sheaves of functions. In the case of affine schemes, we have already seen the map as sets (§4.2.7) and later saw that this map is continuous (Exercise 4.4.G).

Here are two motivations for how morphisms should behave. The first is algebraic, and the second is geometric.

7.1.1. Algebraic motivation. We will want morphisms of affine schemes \( \text{Spec } B \rightarrow \text{Spec } A \) to be precisely the ring maps \( A \rightarrow B \). We have already seen that ring maps \( A \rightarrow B \) induce maps of topological spaces in the opposite direction (Exercise 4.4.G); the main new ingredient will be to see how to add the structure sheaf of functions into the mix. Then a morphism of schemes should be something that “on the level of affines, looks like this”.

7.1.2. Geometric motivation. Motivated by the theory of differentiable manifolds (§4.1.1), which like schemes are ringed spaces, we want morphisms of schemes at the very least to be morphisms of ringed spaces; we now describe what these are. Notice that if \( \pi : X \rightarrow Y \) is a map of differentiable manifolds, then a differentiable function on \( Y \) pulls back to a differentiable function on \( X \). More precisely, given an open subset \( U \subset Y \), there is a natural map \( \Gamma(U, O_Y) \rightarrow \Gamma(\pi^{-1}(U), O_X) \). This behaves well with respect to restriction (restricting a function to a smaller open set and pulling back yields the same result as pulling back and then restricting), so in fact we have a map of sheaves on \( Y \): \( O_Y \rightarrow \pi_* O_X \). Similarly a morphism of schemes \( X \rightarrow Y \) should induce a map \( O_Y \rightarrow \pi_* O_X \). But in fact in the category of differentiable manifolds a continuous map \( X \rightarrow Y \) is a map of differentiable manifolds precisely when differentiable functions on \( Y \) pull back to differentiable functions on \( X \) (i.e. the pullback map from differentiable functions on \( Y \) to functions on \( X \) in fact lies in the subset of differentiable functions, i.e. the continuous map \( X \rightarrow Y \) induces a pullback of differential functions \( O_Y \rightarrow O_X \), so this map of
sheaves characterizes morphisms in the differentiable category. So we could use this as the definition of morphism in the differentiable category.

But how do we apply this to the category of schemes? In the category of differentiable manifolds, a continuous map $X \to Y$ induces a pullback of (the sheaf of) functions, and we can ask when this induces a pullback of differentiable functions. However, functions are odder on schemes, and we can’t recover the pullback map just from the map of topological spaces. A reasonable patch is to hardwire this into the definition of morphism, i.e. to have a continuous map $f : X \to Y$, along with a pullback map $f^* : \mathcal{O}_Y \to f_* \mathcal{O}_X$. This leads to the definition of the category of ringed spaces.

One might hope to define morphisms of schemes as morphisms of ringed spaces. This isn’t quite right, as then Motivation 7.1.1 isn’t satisfied: as desired, to each morphism $A \to B$ there is a morphism $\text{Spec } B \to \text{Spec } A$, but there can be additional morphisms of ringed spaces $\text{Spec } B \to \text{Spec } A$ not arising in this way (see Exercise 7.2.E). A revised definition as morphisms of ringed spaces that locally looks of this form will work, but this is awkward to work with, and we take a different approach. However, we will check that our eventual definition actually is equivalent to this (Exercise 7.3.C).

We begin by formally defining morphisms of ringed spaces.

### 7.2 Morphisms of ringed spaces

#### 7.2.1. Definition. A morphism $\pi : X \to Y$ of ringed spaces is a continuous map of topological spaces (which we unfortunately also call $\pi$) along with a “pullback map” $\pi_* \mathcal{O}_Y \to \pi_* \mathcal{O}_X$. By adjointness (§3.6.1), this is the same as a map $\pi^{-1} \mathcal{O}_Y \to \mathcal{O}_X$.

There is an obvious notion of composition of morphisms, so ringed spaces form a category. Hence we have notion of automorphisms and isomorphisms. You can easily verify that an isomorphism $f : (X, \mathcal{O}_X) \to (Y, \mathcal{O}_Y)$ is a homeomorphism $f : X \to Y$ along with an isomorphism $\mathcal{O}_Y \to f_* \mathcal{O}_X$ (or equivalently $f^{-1} \mathcal{O}_Y \to \mathcal{O}_X$).

If $U \subset Y$ is an open subset, then there is a natural morphism of ringed spaces $(U, \mathcal{O}_Y|_U) \to (Y, \mathcal{O}_Y)$. (Check this! The $f^{-1}$ interpretation is cleaner to use here.) More precisely, if $U \to Y$ is an isomorphism of $U$ with an open subset $V$ of $Y$, and we are given an isomorphism $(U, \mathcal{O}_U) \cong (V, \mathcal{O}_Y|_V)$ (via the isomorphism $U \cong V$), then the resulting map of ringed spaces is called an open immersion of ringed spaces.

#### 7.2.A. Exercise (Morphisms of ringed spaces glue). Suppose $(X, \mathcal{O}_X)$ and $(Y, \mathcal{O}_Y)$ are ringed spaces, $X = \bigcup U_i$ is an open cover of $X$, and we have morphisms of ringed spaces $f_i : U_i \to Y$ that “agree on the overlaps”, i.e. $f_i|_{U_i \cap U_j} = f_j|_{U_i \cap U_j}$. Show that there is a unique morphism of ringed spaces $f : X \to Y$ such that $f_i|_{U_i} = f_i$. (Exercise 3.2.F essentially showed this for topological spaces.)

#### 7.2.B. Easy important exercise: $\mathcal{O}$-modules push forward. Given a morphism of ringed spaces $f : X \to Y$, show that sheaf pushforward induces a functor $\text{Mod}_{\mathcal{O}_X} \to \text{Mod}_{\mathcal{O}_Y}$.
7.2.C. EASY IMPORTANT EXERCISE. Given a morphism of ringed spaces \( f : X \to Y \) with \( f(p) = q \), show that there is a map of stalks \( (\mathcal{O}_Y)_q \to (\mathcal{O}_X)_p \).

7.2.D. KEY EXERCISE. Suppose \( \pi^* : B \to A \) is a morphism of rings. Define a morphism of ringed spaces \( \pi : \text{Spec } A \to \text{Spec } B \) as follows. The map of topological spaces was given in Exercise 4.4.G. To describe a morphism of sheaves \( \mathcal{O}_B \to \pi_* \mathcal{O}_A \) on Spec \( B \), it suffices to describe a morphism of sheaves on the distinguished base of Spec \( B \). On \( D(g) \subset \text{Spec } B \), we define

\[
(\mathcal{O}_B)(D(g)) \to (\mathcal{O}_A)(\pi^{-1} D(g)) = (\mathcal{O}_A)(D(\pi^* g))
\]

by \( B_g \to A_{\pi^* g} \). Verify that this makes sense (e.g. is independent of \( g \)), and that this describes a morphism of sheaves on the distinguished base. (This is the third in a series of exercises. We saw that a morphism of rings induces a map of sets in §4.2.7 a map of topological spaces in Exercise 4.4.G and now a map of ringed spaces here.)

This will soon be an example of morphism of schemes! In fact we could make that definition right now.

7.2.2. Tentative Definition we won’t use (cf. Motivation 7.1.1 in §7.1). A morphism of schemes \( f : (X, \mathcal{O}_X) \to (Y, \mathcal{O}_Y) \) is a morphism of ringed spaces that “locally looks like” the maps of affine schemes described in Key Exercise 7.2.D. Precisely, for each choice of affine open sets \( \text{Spec } A \subset X, \text{Spec } B \subset Y \), such that \( f(\text{Spec } A) \subset \text{Spec } B \), the induced map of ringed spaces should be of the form shown in Key Exercise 7.2.D.

We would like this definition to be checkable on an affine cover, and we might hope to use the Affine Communication Lemma to develop the theory in this way. This works, but it will be more convenient to use a clever trick: in the next section, we will use the notion of locally ringed spaces, and then once we have used it, we will discard it like yesterday’s garbage.

The map of ringed spaces of Key Exercise 7.2.D is really not complicated. Here is an example. Consider the ring map \( \mathbb{C}[x] \to \mathbb{C}[y] \) given by \( x \mapsto y^2 \) (see Figure 4.5). We are mapping the affine line with coordinate \( y \) to the affine line with coordinate \( x \). The map is (on closed points) \( a \mapsto a^2 \). For example, where does \( ([y - 3]) \) go to? Answer: \( ([x - 9]) \), i.e. \( 3 \mapsto 9 \). What is the preimage of \( ([x - 3]) \)? Answer: those prime ideals in \( \mathbb{C}[y] \) containing \( ([y^2 - 4]) \), i.e. \( ([y - 2]) \) and \( ([y + 2]) \), so the preimage of 4 is indeed \( \pm 2 \). This is just about the map of sets, which is old news (§4.2.7), so let’s now think about functions pulling back. What is the pullback of the function \( 3/(x - 4) \) on \( D( ([x - 4]) ) = \mathbb{A}^1 - \{ 4 \} \)? Of course it is \( 3/(y^2 - 4) \) on \( \mathbb{A}^1 - \{ -2, 2 \} \).

We conclude with an example showing that not every morphism of ringed spaces between affine schemes is of the form of Key Exercise 7.2.D. (In the language of the next section, this morphism of ringed spaces is not a morphism of locally ringed spaces.)

7.2.E. UNIMPORTANT EXERCISE. Recall (Exercise 4.4.I) that \( \text{Spec } k[x]/(x) \) has two points, corresponding to \( (0) \) and \( (x) \), where the second point is closed, and the first is not. Consider the map of ringed spaces \( \text{Spec } k(x) \to \text{Spec } k[x]/(x) \) sending the point of \( \text{Spec } k(x) \) to \( (x) \), and the pullback map \( f^* \mathcal{O}_{\text{Spec } k[x]/(x)} \to \mathcal{O}_{\text{Spec } k(x)} \) is induced by \( k \to k(x) \). Show that this map of ringed spaces is not of the form described in Key Exercise 7.2.D.
7.3 From locally ringed spaces to morphisms of schemes

In order to prove that morphisms behave in a way we hope, we will use the notion of a locally ringed space. It will not be used later, although it is useful elsewhere in geometry. The notion of locally ringed spaces (and maps between them) is inspired by what we know about manifolds (see Exercise 4.1.B). If $\pi : X \to Y$ is a morphism of manifolds, with $\pi(p) = q$, and $f$ is a function on $Y$ vanishing at $q$, then the pulled back function $\pi^#f$ on $X$ should vanish on $p$. Put differently: germs of functions (at $q \in Y$) vanishing at $q$ should pull back to germs of functions (at $p \in X$) vanishing at $p$.

7.3.1. Definition. Recall (Definition 5.3.4) that a locally ringed space is a ringed space $(X, O_X)$ such that the stalks $O_{X,x}$ are all local rings. A morphism of locally ringed spaces $f : X \to Y$ is a morphism of ringed spaces such that the induced map of stalks $O_{Y,q} \to O_{X,p}$ (Exercise 7.2.C) sends the maximal ideal of the former into the maximal ideal of the latter (a "morphism of local rings"). This means something rather concrete and intuitive: "if $p \mapsto q$, and $g$ is a function vanishing at $q$, then it will pull back to a function vanishing at $p$." Note that locally ringed spaces form a category. (For completeness, we point out that the notion of a morphism of ringed spaces is the same, without the maximal ideal condition. But this idea won't come up for us.)

To summarize: we use the notion of locally ringed space only to define morphisms of schemes, and to show that morphisms have reasonable properties. The main things you need to remember about locally ringed spaces are (i) that the functions have values at points, and (ii) that given a map of locally ringed spaces, the pullback of where a function vanishes is precisely where the pulled back function vanishes.

7.3.A. EXERCISE. Show that morphisms of locally ringed spaces glue (cf. Exercise 7.2.A). (Hint: your solution to Exercise 7.2.A may work without change.)

7.3.B. EASY IMPORTANT EXERCISE. (a) Show that Spec $A$ is a locally ringed space. (Hint: Exercise 5.3.F) (b) Show that the morphism of ringed spaces $f : \text{Spec } A \to \text{Spec } B$ defined by a ring morphism $f^#: B \to A$ (Exercise 4.4.G) is a morphism of locally ringed spaces.

7.3.2. Key Proposition. — If $f : \text{Spec } A \to \text{Spec } B$ is a morphism of locally ringed spaces then it is the morphism of locally ringed spaces induced by the map $f^#: B = \Gamma(\text{Spec } B, O_{\text{Spec } B}) \to \Gamma(\text{Spec } A, O_{\text{Spec } A}) = A$ as in Exercise 7.3.B.b.

Proof. Suppose $f : \text{Spec } A \to \text{Spec } B$ is a morphism of locally ringed spaces. We wish to show that it is determined by its map on global sections $f^#: B \to A$. We first need to check that the map of points is determined by global sections. Now a point $p$ of Spec $A$ can be identified with the prime ideal of global functions vanishing on it. The image point $f(p)$ in Spec $B$ can be interpreted as the unique point $q$ of Spec $B$, where the functions vanishing at $q$ pull back to precisely those functions vanishing at $p$. (Here we use the fact that $f$ is a map of locally ringed spaces.) This is precisely the way in which the map of sets $\text{Spec } A \to \text{Spec } B$ induced by a ring map $B \to A$ was defined (4.2.7).
Note in particular that if \( b \in B \), \( f^{-1}(D(b)) = D(f^#b) \), again using the hypothesis that \( f \) is a morphism of locally ringed spaces.

It remains to show that \( f^# : O_{\text{Spec } B} \to f_*O_{\text{Spec } A} \) is the morphism of sheaves given by Exercise 7.2.D (cf. Exercise 7.3.B(b)). It suffices to check this on the distinguished base (Exercise 3.7.C(a)). We now want to check that for any map of locally ringed spaces inducing the map of sheaves \( O_{\text{Spec } B} \to f_*O_{\text{Spec } A} \), the map of sections on any distinguished open set \( D(b) \subset \text{Spec } B \) is determined by the map of global sections \( B \to A \).

Consider the commutative diagram

\[
\begin{array}{ccc}
B & \xrightarrow{\Gamma(\text{Spec } B, O_{\text{Spec } B})} & A \\
\downarrow{\text{res}_{\text{Spec } B, D(b)}} & & \downarrow{\text{res}_{\text{Spec } A, D(f^#b)}} \\
B_b & \xrightarrow{\Gamma(D(b), O_{\text{Spec } B})} & A_{f^#b} = A \otimes_B B_b
\end{array}
\]

The vertical arrows (restrictions to distinguished open sets) are localizations by \( b \), so the lower horizontal map \( f^#_{D(b)} \) is determined by the upper map (it is just localization by \( b \)).

We are ready for our definition.

**7.3.3. Definition.** If \( X \) and \( Y \) are schemes, then a morphism \( \pi : X \to Y \) as locally ringed spaces is called a **morphism of schemes**. We have thus defined the category of schemes, which we denote \( \text{Sch} \). (We then have notions of isomorphism — just the same as before, §5.3.4 — and automorphism. We note that the target \( Y \) of \( \pi \) is sometimes called the **base scheme** or the **base**, when we are interpreting \( \pi \) as a family of schemes parametrized by \( Y \) — this may become clearer once we have defined the fibers of morphisms in §10.3.2.)

The definition in terms of locally ringed spaces easily implies Tentative Definition 7.2.2.

**7.3.C. Important Exercise.** Show that a morphism of schemes \( f : X \to Y \) is a morphism of ringed spaces that looks locally like morphisms of affines. Precisely, if \( \text{Spec } A \) is an affine open subset of \( X \) and \( \text{Spec } B \) is an affine open subset of \( Y \), and \( f(\text{Spec } A) \subset \text{Spec } B \), then the induced morphism of ringed spaces is a morphism of affine schemes. (In case it helps, note: if \( W \subset X \) and \( Y \subset Z \) are both open immersions of ringed spaces, then any morphism of ringed spaces \( X \to Y \) induces a morphism of ringed spaces \( W \to Z \), by composition \( W \to X \to Y \to Z \).) Show that it suffices to check on a set \( \{\text{Spec } A_i, \text{Spec } B_i\} \) where the \( \text{Spec } A_i \) form an open cover of \( X \).

In practice, we will use the affine cover interpretation, and forget completely about locally ringed spaces. In particular, put imprecisely, the category of affine schemes is the category of rings with the arrows reversed. More precisely:

**7.3.D. Exercise.** Show that the category of rings and the opposite category of affine schemes are equivalent (see §2.2.21 to read about equivalence of categories).
In particular, here is something surprising: there can be interesting maps from one point to another. For example, there are two different maps from the point \( \text{Spec } \mathbb{C} \) to the point \( \text{Spec } \mathbb{C} \): the identity (corresponding to the identity \( \mathbb{C} \to \mathbb{C} \)), and complex conjugation. (There are even more such maps!)

It is clear (from the corresponding facts about locally ringed spaces) that morphisms glue (Exercise 7.3.A), and the composition of two morphisms is a morphism. Isomorphisms in this category are precisely what we defined them to be earlier (§5.3.4).

7.3.4. The category of complex schemes (or more generally the category of \( k \)-schemes where \( k \) is a field, or more generally the category of \( A \)-schemes where \( A \) is a ring, or more generally the category of \( S \)-schemes where \( S \) is a scheme). The category of \( S \)-schemes (where \( S \) is a scheme) is defined as follows. The objects are morphisms of the form

\[
\begin{array}{ccc}
X & \rightarrow & Y \\
\downarrow \quad \text{structure morphism} & & \downarrow \\
S & \rightarrow & S \\
\end{array}
\]

The morphisms in the category of \( S \)-schemes are commutative diagrams

\[
\begin{array}{ccc}
X & \rightarrow & Y \\
\downarrow & & \downarrow \\
S & \rightarrow & S \\
\end{array}
\]

which is more conveniently written as a commutative diagram

\[
\begin{array}{ccc}
X & \rightarrow & Y \\
\downarrow & & \downarrow \\
\quad & \rightarrow & S \\
\end{array}
\]

When there is no confusion (if the base scheme is clear), simply the top row of the diagram is given. In the case where \( S = \text{Spec } A \), where \( A \) is a ring, we get the notion of an \( A \)-scheme, which is the same as the same definition as in §6.3.5 but in a more satisfactory form. For example, complex geometers may consider the category of \( \mathbb{C} \)-schemes.

The next two examples are important. The first will show you that you can work with these notions in a straightforward, hands-on way. The second will show that you can work with these notions in a formal way.

7.3.E. IMPORTANT EXERCISE. (This exercise will give you some practice with understanding morphisms of schemes by cutting up into affine open sets.) Make sense of the following sentence: “\( \mathbb{A}^{n+1} \setminus \{0\} \to \mathbb{P}^n \) given by 

\[(x_0, x_1, \ldots, x_{n+1}) \mapsto [x_0; x_1; \ldots; x_n]\]

is a morphism of schemes.” Caution: you can’t just say where points go; you have to say where functions go. So you will have to divide these up into affines, and describe the maps, and check that they glue.
7.3.F. **Essential Exercise.** Show that morphisms $X \to \text{Spec } A$ are in natural bijection with ring morphisms $A \to \Gamma(X, \mathcal{O}_X)$. Hint: Show that this is true when $X$ is affine. Use the fact that morphisms glue, Exercise 7.3.E. (This is even true in the category of locally ringed spaces, and you are free to prove it in this generality, although it is notably easier in the category of schemes.)

In particular, there is a canonical morphism from a scheme to Spec of its space of global sections. (Warning: Even if $X$ is a finite-type $k$-scheme, the ring of global sections might be nasty! In particular, it might not be finitely generated, see 21.10.7)

7.3.G. **Easy Exercise.** Show that this definition of $A$-scheme agrees with the earlier definition of 6.3.5.

7.3.H. **Easy Exercise.** Describe a natural “structure morphism” $\text{Proj } S \to \text{Spec } A$.

7.3.I. **Easy Exercise.** Show that $\text{Spec } \mathbb{Z}$ is the final object in the category of schemes. In other words, if $X$ is any scheme, there exists a unique morphism to $\text{Spec } \mathbb{Z}$. (Hence the category of schemes is isomorphic to the category of $\mathbb{Z}$-schemes.) If $k$ is a field, show that $\text{Spec } k$ is the final object in the category of $k$-schemes.

7.3.6. **Definition: The functor of points, and $S$-valued points of a scheme.** If $S$ is a scheme, then $S$-valued points of a scheme $X$ are defined to be maps $S \to X$. If $A$ is a ring, then $A$-valued points of a scheme $X$ are defined to be the $(\text{Spec } A)$-valued points of the scheme. This definition isn’t great, because we earlier defined the notion of points of a scheme, and $S$-valued points are not (necessarily) points! But this definition is well-established in the literature. Here is one reason why it is a reasonable notion: the $A$-valued points of an affine scheme $\text{Spec } \mathbb{Z}[x_1, \ldots, x_n]/(f_1, \ldots, f_r)$ (were $f_i \in \mathbb{Z}[x_1, \ldots, x_n]$ are relations) are precisely the solutions to the equations

$$f_1(x_1, \ldots, x_n) = \cdots = f_r(x_1, \ldots, x_n) = 0$$

in the ring $A$. For example, the rational solutions to $x^2 + y^2 = 16$ are precisely the $\mathbb{Q}$-valued points of $\text{Spec } \mathbb{Z}[x,y]/(x^2 + y^2 - 16)$. The integral solutions are precisely the $\mathbb{Z}$-valued points. So $A$-valued points of an affine scheme (finite type over $\mathbb{Z}$) can be interpreted simply. In the special case where $A$ is local, $A$-valued points of a general scheme have a good interpretation too:

7.3.J. **Exercise (Morphisms from Spec of a Local Ring to $X$).** Suppose $X$ is a scheme, and $(A, m)$ is a local ring. Suppose we have a scheme morphism $\pi: \text{Spec } A \to X$ sending $[m]$ to $x$. Show that any open set containing $x$ contains the image of $\pi$. Show that there is a bijection between $\text{Hom}(\text{Spec } A, X)$ and $\{x \in X, \text{local homomorphisms } \mathcal{O}_{x,X} \to A\}$. 
Another reason this notion is good is that “products of \( S \)-valued points” behave as you might hope, see \([\text{10.1.3}](\text{10.1.3})\).

On the other hand, maps to projective space can be confusing. There are some maps we can write down easily, as shown by applying the next exercise in the case \( X = \text{Spec} \ A \), where \( A \) is a \( B \)-algebra.

**7.3.K. Exercise.** Suppose \( B \) is a ring. If \( X \) is a \( B \)-scheme, and \( f_0, \ldots, f_n \) are \( n \) functions on \( X \) with no common zeros, then show that \([f_0; \ldots; f_n]\) gives a morphism \( X \to \mathbb{P}^n_B \).

You might hope that this gives all morphisms. But this isn’t the case. Indeed, even the identity morphism \( X = \mathbb{P}^1_k \to \mathbb{P}^1_k \) isn’t of this form, as the source \( \mathbb{P}^1 \) has no nonconstant global functions with which to build this map. (And there are similar examples with an affine source.) However, there is a correct generalization (characterizing all maps from schemes to projective schemes) in Theorem \([7.4.1](\text{7.4.1})\). This result roughly states that this works, so longer as the \( f_i \) are not quite functions, but sections of a line bundle. Our desire to understand maps to projective schemes in a clean way will be one important motivation for understanding line bundles.

We will see more ways to describe maps to projective space in the next section. Incidentally, before Grothendieck, it was considered a real problem to figure out the right way to interpret points of projective space with “coordinates” in a ring. These difficulties were due to a lack of functorial reasoning. And the clues to the right answer already existed (the same problems arise for maps from a smooth real manifold to \( \mathbb{R}^n \)) — if you ask such a geometric question (for projective space is geometric), the answer is necessarily geometric, not purely algebraic!

**7.3.7. Visualizing schemes III: picturing maps of schemes when nilpotents are present.** You now know how to visualize the points of schemes (\([4.3](\text{4.3})\)), and nilpotents (\([5.2](\text{5.2})\) and \([6.5](\text{6.5})\)). The following imprecise exercise will give you some sense of how to visualize maps of schemes when nilpotents are involved. Suppose \( a \in \mathbb{C} \). Consider the map of rings \( \mathbb{C}[x] \to \mathbb{C}[\epsilon]/\epsilon^2 \) given by \( x \mapsto a\epsilon \). Recall that \( \text{Spec} \ \mathbb{C}[\epsilon] \) may be pictured as a point with a tangent vector (\([5.2](\text{5.2})\)). How would you picture this map if \( a \neq 0 \)? How does your picture change if \( a = 0 \)? (The tangent vector should be “crushed” in this case.)

Exercise \([13.1.G](\text{13.1.G})\) will extend this considerably; you may enjoy reading its statement now.

### 7.4 Maps of graded rings and maps of projective schemes

As maps of rings correspond to maps of affine schemes in the opposite direction, maps of graded rings (over a base ring \( A \)) sometimes give maps of projective schemes in the opposite direction. This is an imperfect generalization: not every map of graded rings gives a map of projective schemes (\([7.4.1](\text{7.4.1})\)); not every map of projective schemes comes from a map of graded rings (later); and different maps of graded rings can yield the same map of schemes (Exercise \([7.4.C](\text{7.4.C})\)).

**7.4.A. Essential exercise.** Suppose that \( f : S \to R \) is a morphism of finitely-generated graded rings over \( A \). By map of finitely generated graded rings,
we mean a map of rings that preserves the grading as a map of grading semi-groups. In other words, there is a \( d > 0 \) such that \( S_n \) maps to \( R_{dn} \). Show that this induces a morphism of schemes \( \text{Proj} \ R_n \setminus V(f(S_+)) \rightarrow \text{Proj} S_\bullet \). (Hint: Suppose \( x \) is a homogeneous element of \( S_+ \). Define a map \( D(f(x)) \rightarrow D(x) \). Show that they glue together (as \( x \) runs over all homogeneous elements of \( S_+ \)). Show that this defines a map from all of \( \text{Proj} R_\bullet \setminus V(f(S_+)) \).) In particular, if

\[
(7.4.0.1) \quad V(f(S_+)) = \emptyset,
\]
then we have a morphism \( \text{Proj} R_\bullet \rightarrow \text{Proj} S_\bullet \).

**7.4.B. Exercise.** Show that if \( f : S_\bullet \rightarrow R_\bullet \) satisfies \( \sqrt{\{f(S_+)\}} = R_+ \), then hypothesis (7.4.0.1) is satisfied. (Hint: Exercise 5.5.F.) This algebraic formulation of the more geometric hypothesis can sometimes be easier to verify.

Let’s see Exercise 7.4.A in action. We will schematically interpret the map of complex projective manifolds \( \mathbb{P}^1 \) to \( \mathbb{P}^2 \) given by

\[
\begin{array}{ccc}
\mathbb{P}^1 & \longrightarrow & \mathbb{P}^2 \\
[s; t] & \mapsto & [s^{20}; s^9 t^{11}; t^{20}]
\end{array}
\]

Notice first that this is well-defined: \([\lambda s; \lambda t]\) is sent to the same point of \( \mathbb{P}^2 \) as \([s; t]\). The reason for it to be well-defined is that the three polynomials \( s^{20} \), \( s^9 t^{11} \), and \( t^{20} \) are all homogeneous of degree 20.

Algebraically, this corresponds to a map of graded rings in the opposite direction

\[
\mathbb{C}[x, y, z] \rightarrow \mathbb{C}[s, t]
\]
given by \( x \mapsto s^{20}, y \mapsto s^9 t^{11}, z \mapsto t^{20} \). You should interpret this in light of your solution to Exercise 7.4.A and compare this to the affine example of 4.2.8.

**7.4.1.** Notice that there is no map of complex manifolds \( \mathbb{P}^2 \rightarrow \mathbb{P}^1 \) given by \([x; y; z] \rightarrow [x; y]\), because the map is not defined when \( x = y = 0 \). This corresponds to the fact that the map of graded rings \( \mathbb{C}[s, t] \rightarrow \mathbb{C}[x, y, z] \) given by \( s \mapsto x \) and \( t \mapsto y \), doesn’t satisfy hypothesis (7.4.0.1).

**7.4.C. Unimportant exercise.** This exercise shows that different maps of graded rings can give the same map of schemes. Let \( R_\bullet = k[x, y, z]/(xz, yz, z^2) \) and \( S_\bullet = k[a, b, c]/(ac, bc, c^2) \), where every variable has degree 1. Show that \( \text{Proj} R_\bullet \cong \text{Proj} S_\bullet \cong \mathbb{P}^1_k \). Show that the maps \( S_\bullet \rightarrow R_\bullet \) given by \((a, b, c) \mapsto (x, y, z)\) and \((a, b, c) \mapsto (x, y, 0)\) give the same (iso)morphism \( \text{Proj} R_\bullet \rightarrow \text{Proj} S_\bullet \). (The real reason is that all of these constructions are insensitive to what happens in a finite number of degrees. This will be made precise in a number of ways later, most immediately in Exercise 7.4.F)

**7.4.2. Veronese subrings.**

Here is a useful construction. Suppose \( S_\bullet \) is a finitely generated graded ring. Define the *n*th Veronese subring of \( S_\bullet \) by \( S_{n\bullet} = \bigoplus_{j=0}^{\infty} S_{nj} \). (The “old degree” \( n \) is “new degree” \( 1 \))
7.4.D. EXERCISE. Show that the map of graded rings $S_{n\bullet} \hookrightarrow S_{\bullet}$ induces an isomorphism $\text{Proj } S_{\bullet} \rightarrow \text{Proj } S_{n\bullet}$. (Hint: if $f \in S_{\bullet}$ is homogeneous of degree divisible by $n$, identify $D(f)$ on $\text{Proj } S_{\bullet}$ with $D(f)$ on $\text{Proj } S_{n\bullet}$. Why do such distinguished open sets cover $\text{Proj } S_{\bullet}$?)

7.4.E. EXERCISE. If $S_{\bullet}$ is generated in degree 1, show that $S_{n\bullet}$ is also generated in degree 1. (You may want to consider the case of the polynomial ring first.)

7.4.F. EXERCISE. Use the previous exercise to show that if $R_{\bullet}$ and $S_{\bullet}$ are the same finitely generated graded rings except in a finite number of nonzero degrees (make this precise!), then $\text{Proj } R_{\bullet} \cong \text{Proj } S_{\bullet}$.

7.4.G. EXERCISE. Suppose $S_{\bullet}$ is generated over $S_0$ by $f_1, \ldots, f_n$. Find a $d$ such that $S_{d\bullet}$ is generated in “new” degree 1 (= “old” degree $d$). (This is surprisingly tricky, so here is a hint. Suppose there are generators $x_1, \ldots, x_n$ of degrees $d_1, \ldots, d_n$ respectively. Show that any monomial $x_1^{a_1} \cdots x_n^{a_n}$ of degree at least $nd_1 \ldots d_n$ has $a_i \geq \left(\prod_j d_j\right)/d_i$ for some $i$. Show that the $nd_1 \ldots d_n$th Veronese subring is generated by elements in “new” degree 1.) This, in combination with the previous exercise, shows that there is little harm in assuming that finitely generated graded rings are generated in degree 1, as after a regrading, this is indeed the case. This is handy, as it means that, using Exercise 7.4.D, we can assume that any finitely-generated graded ring is generated in degree 1. We will see that as a consequence we can place every Proj in some projective space via the construction of Exercise 9.2.H.

7.4.H. LESS IMPORTANT EXERCISE. Show that $S_{n\bullet}$ is a finitely generated graded ring. (Possible approach: use the previous exercise, or something similar, to show there is some $N$ such that $S_{nN\bullet}$ is generated in degree 1, so the graded ring $S_{nN\bullet}$ is finitely generated. Then show that for each $0 < j < N$, $S_{nN\bullet + nj}$ is a finitely generated module over $S_{nN\bullet}$)

7.5 Rational maps from integral schemes

Informally speaking, a “rational map” is a “a morphism defined almost everywhere”, much as a rational function is a name for a function defined almost everywhere. We will later see that in good situations that where a rational map is defined, it is uniquely defined (the Reduced-to-separated Theorem 11.2.1), and has a largest “domain of definition” (11.2.2). For this section, unless otherwise stated, we assume $X$ and $Y$ to be integral. The reader interested in more general notions should consider first the case where the schemes in question are reduced but not necessarily irreducible. A key example will be irreducible varieties, and the language of rational maps is most often used in this case. Many notions can make sense in more generality (without reducedness hypotheses for example), but I’m not sure if there is a widely accepted definition.

7.5.1. Definition. A rational map from $X$ to $Y$, denoted $X \dasharrow Y$, is a morphism on a dense open set, with the equivalence relation $(f : U \rightarrow Y) \sim (g : V \rightarrow Y)$ if there is a dense open set $Z \subset U \cap V$ such that $f|_Z = g|_Z$. (In 11.2.2 we will improve this to:
if \( f|_{U \cap V} = g|_{U \cap V} \) in good circumstances — when \( Y \) is separated.) People often use the word “map” for “morphism”, which is quite reasonable, except that a rational map need not be a map. So to avoid confusion, when one means “rational map”, one should never just say “map”.

7.5.2. Rational maps more generally. The right generality for the notion of rational map, to a situation where no serious pathologies arise, is where \( X \) has associated points — where it is integral or locally Noetherian (§6.5) — and where \( Y \) is arbitrary. In this case, the dense open set of \( X \) is required to contain the associated points. (We will see in §11.2 that rational maps to separated schemes behave particularly well, and they are usually considered in this situation.)

7.5.3. An obvious example of a rational map is a morphism. Another important example is the projection \( \mathbb{P}^n_A \to \mathbb{P}^{n-1}_A \) given by \([x_0; \cdots; x_n] \to [x_0; \cdots; x_{n-1}]\). (How precisely is this a rational map in the sense of Definition 7.5.1? What is its domain of definition?) A third example is the following.

7.5.A. EASY EXERCISE. Interpret rational functions on an integral scheme (§6.5.1) as rational maps to \( \mathbb{A}^1_Z \). (This is analogous to functions corresponding to morphisms to \( \mathbb{A}^1_Z \), which will be described in §7.6.1.)

7.5.B. EASY EXERCISE. Show that a rational map \( X \to Y \) from an integral scheme \( X \) is the same as a \( \mathbb{K}(X) \)-valued point (§7.3.6) of \( Y \).

A rational map \( f : X \to Y \) is dominant (or in some sources, dominating) if for some (and hence every) representative \( U \to Y \), the image is dense in \( Y \). Equivalently, \( f \) is dominant if it sends the generic point of \( X \) to the generic point of \( Y \). A little thought will convince you that you can compose (in a well-defined way) a dominant map \( f : X \to Y \) with a rational map \( g : Y \to Z \). Integral schemes and dominant rational maps between them form a category which is geometrically interesting.

7.5.C. EASY EXERCISE. Show that dominant rational maps of integral schemes give morphisms of function fields in the opposite direction.

It is not true that morphisms of function fields always give dominant rational maps, or even rational maps. For example, \( \text{Spec} \, k[x] \) and \( \text{Spec} \, k(x) \) have the same function field (\( k(x) \)), but there is no rational map \( \text{Spec} \, k[x] \to \text{Spec} \, k(x) \). Reason: that would correspond to a morphism from an open subset \( U \) of \( \text{Spec} \, k[x] \), say \( \text{Spec} \, k[x, 1/f(x)] \), to \( \text{Spec} \, k(x) \). But there is no map of rings \( k(x) \to k[x, 1/f(x)] \) for any one \( f(x) \). However, maps of function fields indeed give dominant rational maps of integral finite type \( k \)-schemes (and in particular, irreducible varieties, to be defined in §11.17), see Proposition 7.5.5 below.

(If you want more evidence that the topologically-defined notion of dominance is simultaneously algebraic, you can show that if \( \phi : A \to B \) is a ring morphism, then the corresponding morphism \( \text{Spec} \, B \to \text{Spec} \, A \) is dominant if and only if \( \phi \) has nilpotent kernel.)

A rational map \( f : X \to Y \) is said to be birational if it is dominant, and there is another rational map (a “rational inverse”) that is also dominant, such that \( f \circ g \) is (in the same equivalence class as) the identity on \( Y \), and \( g \circ f \) is (in the same
equivalence class as) the identity on $X$. This is the notion of isomorphism in the category of integral schemes and dominant rational maps. We say $X$ and $Y$ are birational (to each other) if there exists a birational map $X \dasharrow Y$. Birational maps induce isomorphisms of function fields. The fact that maps of function fields correspond to rational maps in the opposite direction for integral finite type $k$-schemes, to be proved in Proposition 7.5.5, shows that a map between integral finite type $k$-schemes that induces an isomorphism of function fields is birational. An integral finite type $k$-scheme is said to be rational if it is birational to $A^n_k$ for some $k$. A morphism is birational if it is birational as a rational map. We will later see (Proposition 11.2.3) that two integral affine $k$-varieties $X$ and $Y$ are birational if there are open sets $U \subset X$ and $V \subset Y$ that are isomorphic ($U \cong V$). In particular, an integral affine $k$-variety is rational if “it has a big open subset that is a big open subset of affine space $A^n_k$."

7.5.4. Rational maps of irreducible varieties.

7.5.5. Proposition. — Suppose $X, Y$ are integral finite type $k$-schemes, and we are given $\phi^* : K(Y) \hookrightarrow K(X)$. Then there exists a dominant rational map $\phi : X \dasharrow Y$ inducing $\phi^*$.

Proof. By replacing $Y$ with an affine open set, we may assume $Y$ is affine, say $Y = \text{Spec} \ k[x_1, \ldots, x_n]/(f_1, \ldots, f_r)$. Then we have $\phi^* x_1, \ldots, \phi^* x_n \in K(X)$. Let $U$ be an open subset of the domains of definition of these rational functions. Then we get a morphism $U \to A^n_k$. But this morphism factors through $Y \subset A^n_k$, as $x_1, \ldots, x_n$ satisfy the relations $f_1, \ldots, f_r$.

We see that the morphism is dense as follows. If the set-theoretic image is not dense, it is contained in a proper closed subset. Let $f$ be a function vanishing on the closed subset. Then the pullback of $f$ to $U$ is 0 (as $U$ is reduced), implying that $\phi^*(f) = 0$, and $f$ doesn’t vanish on all of $Y$, so $f$ is not the 0-element of $K(Y)$. But this contradicts the fact that $\phi^*$ is an inclusion. $\Box$

7.5.D. Exercise. Let $K$ be a finitely generated field extension of $k$. (Informal definition: a field extension $K$ over $k$ is finitely generated if there is a finite “generating set” $x_1, \ldots, x_n$ in $K$ such that every element of $K$ can be written as a rational function in $x_1, \ldots, x_n$ with coefficients in $k$.) Show that there exists an irreducible affine $k$-variety with function field $K$. (Hint: Consider the map $k[t_1, \ldots, t_n] \to K$ given by $t_i \mapsto x_i$, and show that the kernel is a prime ideal $p$, and that $k[t_1, \ldots, t_n]/p$ has function field $K$. Interpreted geometrically: consider the map $\text{Spec} K \to \text{Spec} k[t_1, \ldots, t_n]$ given by the ring map $t_i \mapsto x_i$, and take the closure of the one-point image.)

7.5.E. Exercise. Describe an equivalence of categories between (a) finitely generated field extensions of $k$, and inclusions extending the identity on $k$, and (b) integral affine $k$-varieties, and dominant rational maps defined over $k$.

In particular, an integral affine $k$-variety $X$ is rational if its function field $K(X)$ is a purely transcendental extension of $k$, i.e. $K(X) \cong k(x_1, \ldots, x_n)$ for some $n$. (This needs to be said more precisely: the map $k \hookrightarrow K(X)$ induced by $X \to \text{Spec} k$ should agree with the “obvious” map $k \hookrightarrow k(x_1, \ldots, x_n)$ under this isomorphism.)
7.5.6. **Definition: degree of a rational map of varieties.** If \( \pi : X \rightarrow Y \) is a dominant rational map of integral affine \( k \)-varieties of the same dimension, the degree of the field extension \( K(X)/K(Y) \) is called the **degree** of the rational map. We will interpret this degree in terms of counting preimages of points of \( Y \) later.

7.5.7. **More examples of rational maps.**

A recurring theme in these examples is that domains of definition of rational maps to projective schemes extend over nonsingular codimension one points. We will make this precise in the Curve-to-projective Extension Theorem 17.5.1, when we discuss curves.

---

**Figure 7.1.** Finding primitive Pythagorean triples using geometry

The first example is the classical formula for Pythagorean triples. Suppose you are looking for rational points on the circle \( C \) given by \( x^2 + y^2 = 1 \) (Figure 7.1). One rational point is \( p = (1, 0) \). If \( q \) is another rational point, then \( pq \) is a line of rational (non-infinite) slope. This gives a rational map from the conic \( C \) to \( \mathbb{A}^1 \). Conversely, given a line of slope \( m \) through \( p \), where \( m \) is rational, we can recover \( q \) by solving the equations \( y = m(x - 1), \ x^2 + y^2 = 1 \). We substitute the first equation into the second, to get a quadratic equation in \( x \). We know that we will have a solution \( x = 1 \) (because the line meets the circle at \( (x, y) = (1, 0) \)), so we expect to be able to factor this out, and find the other factor. This indeed works:

\[
x^2 + (m(x - 1))^2 = 1 \\
\implies (m^2 + 1)x^2 + (-2m^2)x + (m^2 - 1) = 0 \\
\implies (x - 1)((m^2 + 1)x - (m^2 - 1)) = 0
\]

The other solution is \( x = (m^2 - 1)/(m^2 + 1) \), which gives \( y = -2m/(m^2 + 1) \). Thus we get a birational map between the conic \( C \) and \( \mathbb{A}^1 \) with coordinate \( m \), given by \( f : (x, y) \mapsto y/(x - 1) \) (which is defined for \( x \neq 1 \)), and with inverse rational map given by \( m \mapsto ((m^2 - 1)/(m^2 + 1), -2m/(m^2 + 1)) \) (which is defined away from \( m^2 + 1 = 0 \)).

We can extend this to a rational map \( C \rightarrow \mathbb{P}^1 \) via the inclusion \( \mathbb{A}^1 \rightarrow \mathbb{P}^1 \). Then \( f \) is given by \( (x, y) \mapsto [y; x - 1] \). We then have an interesting question: what
is the domain of definition of \( f\)? It appears to be defined everywhere except for where \( y = x - 1 \neq 0\), i.e. everywhere but \( p\). But in fact it can be extended over \( p!\) Note that \((x, y) \mapsto [x + 1; -y]\) (where \((x, y) \neq (-1, 0)\)) agrees with \( f\) on their common domains of definition, as \([x + 1; -y] = [y; x - 1]\). Hence this rational map can be extended farther than we at first thought. This will be a special case of the Curve-to-projective Extension Theorem \[\text{17.5.1}\] (For the curious: we are working with schemes over \( \mathbb{Q}\). But this works for any scheme over a field of characteristic not 2. What goes wrong in characteristic 2?)

**7.5.F. Exercise.** Use the above to find a “formula” yielding all Pythagorean triples.

**7.5.G. Exercise.** Show that the conic \( x^2 + y^2 = z^2 \) in \( \mathbb{P}^2_k \) is isomorphic to \( \mathbb{P}^1_k \) for any field \( k \) of characteristic not 2. (We did this earlier in the case where \( k \) is algebraically closed, by diagonalizing quadrics, \[\text{9.2.6}\].)

In fact, any conic in \( \mathbb{P}^2_k \) with a \( k \)-valued point (i.e. a point with residue field \( k\) of rank 3 (after base change to \( \overline{k}\), so “rank” makes sense, see Exercise \[\text{6.4.}\]) is isomorphic to \( \mathbb{P}^1_k\). (This hypothesis is certainly necessary, as \( \mathbb{P}^1_k\) certainly has \( k\)-valued points, but \( x^2 + y^2 + z^2 = 0 \) over \( k = \mathbb{R}\) is a conic that is not isomorphic to \( \mathbb{P}^1_k\).)

**7.5.H. Exercise.** Find all rational solutions to \( y^2 = x^3 + x^2 \), by finding a birational map to \( \mathbb{A}^1_k\), mimicking what worked with the conic. (In Exercise \[\text{21.8.}\] we will see that these points form a group, and that this is a degenerate elliptic curve.)

You will obtain a rational map to \( \mathbb{P}^1 \) that is not defined over the node \( x = y = 0\), and cannot be extended over this codimension 1 set. This is an example of the limits of our future result, the Curve-to-projective Extension Theorem \[\text{17.5.1}\] showing how to extend rational maps to projective space over codimension 1 sets: the codimension 1 sets have to be nonsingular.

**7.5.I. Exercise.** Use a similar idea to find a birational map from the quadric \( Q = \{x^2 + y^2 = w^2 + z^2\}\) to \( \mathbb{P}^2\). Use this to find all rational points on \( Q\). (This illustrates a good way of solving Diophantine equations. You will find a dense open subset of \( Q\) that is isomorphic to a dense open subset of \( \mathbb{P}^2\), where you can easily find all the rational points. There will be a closed subset of \( Q\) where the rational map is not defined, or not an isomorphism, but you can deal with this subset in an ad hoc fashion.)

**7.5.J. Exercise (the Cremona transformation, a useful classical construction).** Consider the rational map \( \mathbb{P}^2 \to \mathbb{P}^2\), given by \([x; y; z] \mapsto [1/x; 1/y; 1/z]\).

What is the domain of definition? (It is bigger than the locus where \( xyz \neq 0\)) You will observe that you can extend it over codimension 1 sets. This again foreshadows the Curve-to-projective Extension Theorem \[\text{17.5.1}\]

**7.5.8. Complex curves that are not rational (fun but inessential).**

We now describe two examples of curves \( C\) such that do not admit a nonconstant rational map from \( \mathbb{P}^1_C\). Both proofs are by Fermat’s method of infinite descent. By Exercise \[\text{7.5.8}\] these results can be interpreted as the fact that these curves have no “nontrivial” \( \mathbb{C}(t)\)-valued points, where by “nontrivial” we mean any such point
is secretly a \( \mathbb{C} \)-valued point. You may notice that if you consider the same examples with \( \mathbb{C}[t] \) replaced by \( \mathbb{Q} \) (and where \( \mathbb{C} \) is a curve over \( \mathbb{Q} \) rather than \( \mathbb{C} \)), you get two fundamental questions in number theory and geometry. The analog of Exercise 7.5.L is the question of rational points on elliptic curves, and you may realize that the analog of Exercise 7.5.K is even more famous. Also, the arithmetic analogue of Exercise 7.5.L(a) is the “four squares theorem” (there are not four integer squares in arithmetic progression), first stated by Fermat. These examples will give you a glimpse of how and why facts over number fields are often paralleled by facts over function fields of curves. This parallelism is a recurring deep theme in the subject.

7.5.K. Exercise. If \( n > 2 \), show that \( \mathbb{P}^1_\mathbb{C} \) has no dominant rational maps to the “Fermat curve” \( x^n + y^n = z^n \) in \( \mathbb{P}^2_\mathbb{C} \). Hint: reduce this to showing that there is no “nonconstant” solution \((f(t), g(t), h(t))\) to \( f(t)^n + g(t)^n = h(t)^n \), where \( f(t), g(t), \) and \( h(t) \) are rational functions in \( t \). By clearing denominators, reduce this to showing that there is no nonconstant solution where \( f(t), g(t), \) and \( h(t) \) are relatively prime polynomials. For this, assume there is a solution, and consider one of the lowest positive degree. Then use the fact that \( \mathbb{C}[t] \) is a unique factorization domain, and \( h(t)^n - g(t)^n = \prod_{i=1}^n (h(t) - \zeta^i g(t)) \), where \( \zeta \) is a primitive \( n \)th root of unity. Argue that each \( h(t) - \zeta^i g(t) \) is an \( n \)th power. Then use

\[
(h(t) - g(t)) \alpha (h(t) - \zeta g(t)) = \beta (h(t) - \zeta^2 g(t))
\]

for suitably chosen \( \alpha \) and \( \beta \) to get a solution of smaller degree. (How does this argument fail for \( n = 2 \)?)

7.5.L. Exercise. Suppose \( a, b, \) and \( c \) are distinct complex numbers. By the following steps, show that \( x(t) \) and \( y(t) \) are two rational functions of \( t \) (elements of \( \mathbb{C}(t) \)) such that

\[
y(t)^2 = (x(t) - a)(x(t) - b)(x(t) - c),
\]

then \( x(t) \) and \( y(t) \) are constants \( (x(t), y(t) \in \mathbb{C}) \). (Here \( \mathbb{C} \) may be replaced by any field \( K \); slight extra care is needed if \( K \) is not algebraically closed.)

(a) Suppose \( P, Q \in \mathbb{C}[t] \) are relatively prime polynomials such that four distinct linear combinations of them are perfect squares. Show that \( P \) and \( Q \) are constant (i.e. \( P, Q \in \mathbb{C} \)). Hint: By renaming \( P \) and \( Q \), show that you may assume that the perfect squares are \( P, P - Q, P - \lambda Q \) (for some \( \lambda \in \mathbb{C} \)). Define \( u \) and \( v \) to be square roots of \( P \) and \( Q \) respectively. Show that \( u - v, u + v, u - \sqrt{\lambda} v, u + \sqrt{\lambda} v \) are perfect squares, and that \( u \) and \( v \) are relatively prime. If \( p \) and \( q \) are not both constant, note that \( 0 < \max(\deg u, \deg v) < \max(\deg P, \deg Q) \). Assume from the start that \( P \) and \( Q \) were chosen as a counterexample with minimal \( \max(\deg P, \deg Q) \) to obtain a contradiction. (Aside: It is possible to have three distinct linear combinations that are perfect squares. Such examples essentially correspond to primitive Pythagorean triples in \( \mathbb{C}(t) \) — can you see how?)

(b) Suppose \( (x, y) = (p/q, r/s) \) is a solution to (7.5.8.1), where \( p, q, r, s \in \mathbb{C}[t] \), and \( p/q \) and \( r/s \) are in lowest terms. Clear denominators to show that \( r^2 q^3 = s^2 (p - aq)(p - bq)(p - cq) \). Show that \( s^2 |q^3 \) and \( q^3 |s \), and hence that \( s^2 = \delta q^3 \) for some \( \delta \in \mathbb{C} \). From \( r^2 = \delta(p - aq)(p - bq)(p - cq) \), show...
that \((p - aq), (p - bq), (p - cq)\) are perfect squares. Show that \(q\) is also a perfect square, and then apply part (a).

### 7.6 * Representable functors and group schemes

#### 7.6.1. Maps to \(\mathbb{A}^1\) correspond to functions. If \(X\) is a scheme, there is a bijection between the maps \(X \rightarrow \mathbb{A}^1\) and global sections of the structure sheaf: by Exercise 7.3.F, maps \(f : X \rightarrow \mathbb{A}^1\) correspond to maps to ring maps \(f^\# : \mathbb{Z}[t] \rightarrow \Gamma(X, \mathcal{O}_X)\), and \(f^\#(t)\) is a function on \(X\); this is reversible.

This map is very natural in an informal sense: you can even picture this map to \(\mathbb{A}^1\) as being given by the function. (By analogy, a function on a smooth manifold is a map to \(\mathbb{R}\).) But it is natural in a more precise sense: this bijection is functorial in \(X\). We will ponder this example at length, and see that it leads us to two important advanced notions: representable functors and group schemes.

#### 7.6.2. Representable functors.

We restate the bijection of §7.6.1 as follows. We have two different contravariant functors from \(\text{Sch}\) to \(\text{Sets}\): maps to \(\mathbb{A}^1\) (i.e. \(H : X \mapsto \text{Mor}(X, \mathbb{A}^1)\)), and functions on \(X\) (\(F : X \mapsto \Gamma(X, \mathcal{O}_X)\)). The “naturality” of the bijection — the functoriality in \(X\) — is precisely the statement that the bijection gives a natural isomorphism of functors (§2.2.21): given any \(f : X \rightarrow X'\), the diagram

\[
\begin{array}{ccc}
H(X') & \longrightarrow & H(X) \\
\downarrow & & \downarrow \\
F(X') & \longrightarrow & F(X)
\end{array}
\]

(where the vertical maps are the bijections given in §7.6.1) commutes.

More generally, if \(Y\) is an element of a category \(\mathcal{C}\) (we care about the special case \(\mathcal{C} = \text{Sch}\), recall the contravariant functor \(h_Y : \mathcal{C} \rightarrow \text{Sets}\) defined by \(h_Y(X) = \text{Mor}(X, Y)\) (Example 2.2.20). We say a contravariant functor from \(\mathcal{C}\) to \(\text{Sets}\) is representable by \(Y\) if it is naturally isomorphic to the representable functor \(h_Y\). We say it is representable if it is representable by some \(Y\).

#### 7.6.3. Important Easy Exercise (Representing objects are unique up to unique isomorphism).

Show that if a contravariant functor \(F\) is representable by \(Y\) and by \(Z\), then we have a unique isomorphism \(Y \rightarrow Z\) induced by the natural isomorphism of functors \(h_Y \rightarrow h_Z\). Hint: this is a version of the universal property arguments of §2.3 once again, we are recognizing an object (up to unique isomorphism) by maps to that object. This exercise is essentially Exercise 2.3.Y(b). (This extends readily to Yoneda’s Lemma, Exercise 10.1.C. You are welcome to try that now.)

You have implicitly seen this notion before: you can interpret the existence of products and fibered products in a category as examples of representable functors.
(You may wish to work out how a natural isomorphism $h_{Y \times Z} \cong h_Y \times h_Z$ induces the projection maps $Y \times Z \to Y$ and $Y \times Z \to Z$.)

7.6.C. Exercise. In this exercise, $Z$ may be replaced by any ring.
(a) (affine $n$-space represents the functor of $n$ functions) Show that the functor $X \mapsto \{(f_1, \ldots, f_n) : f_i \in \Gamma(X, O_X)\}$ is represented by $\mathbb{A}^n_Z$. Show that $\mathbb{A}^1_Z \times \mathbb{A}^1_Z \cong \mathbb{A}^2_Z$ (i.e. $\mathbb{A}^2$ satisfies the universal property of $\mathbb{A}^1 \times \mathbb{A}^1$).
(b) (The functor of invertible functions is representable) Show that the functor taking $X$ to invertible functions on $X$ is representable by Spec $\mathbb{Z}[t, t^{-1}]$. Definition: This scheme is called $\mathbb{G}_m$.

7.6.D. Less important exercise. Fix a ring $A$. Consider the functor $H$ from the category of locally ringed spaces to Sets given $H(X) = \{A \to \Gamma(X, O_X)\}$. Show that this functor is representable (by Spec $A$). This gives another (admittedly odd) motivation for the definition of Spec $A$, closely related to that of §7.3.5.

7.6.3. Group schemes (or more generally, group objects in a category).
(The rest of §7.6 is intended to be double-starred, and should be read only for entertainment.) We return again to Example 7.6.1. Function $s$ on $X$ are better than a set: they form a group. (Indeed they even form a ring, but we will worry about this later.) Given a morphism $X \to Y$, pullback of functions $\Gamma(Y, O_Y) \to \Gamma(X, O_X)$ is a group homomorphism. So we should expect $\mathbb{A}^1$ to have some group-like structure. This leads us to the notion of group scheme, or more generally a group object in a category, which we now define.

A group object in $\mathcal{C}$ is an element $X$ along with three morphisms:
- Multiplication: $m : X \times X \to X$
- Inverse: $i : X \to X$
- Identity element: $e : Z \to X$ (not the identity map)

These morphisms are required to satisfy several conditions.
(i) associativity axiom:
$$X \times X \times X \xrightarrow{(m, \text{id})} X \times X \xrightarrow{m} X$$
commutes. (Here id means the equality $X \to X$.)
(ii) identity axiom:
$$X \xrightarrow{e, \text{id}} X \times X \xrightarrow{m} X \quad \text{and} \quad X \xrightarrow{\text{id}, e} X \times X \xrightarrow{m} X$$
are both the identity map $X = X$. (This corresponds to group axiom: multiplication by the identity element is the identity map.)
(iii) inverse axiom:
$$X \xrightarrow{i \times \text{id}} X \times X \xrightarrow{m} X \quad \text{and} \quad X \xrightarrow{\text{id} \times i} X \times X \xrightarrow{m} X$$
are both $e$.

As motivation, you can check that a group object in the category of sets is in fact the same thing as a group. (This is symptomatic of how you take some notion and make it categorical. You write down its axioms in a categorical way, and if
all goes well, if you specialize to the category of sets, you get your original notion. You can apply this to the notion of “rings” in an exercise below.

A **group scheme** is defined to be a group object in the category of schemes. A **group scheme** over a ring \( A \) (or a scheme \( S \)) is defined to be a group object in the category of \( A \)-schemes (or \( S \)-schemes).

**7.6.E. Exercise.** Give \( \mathbb{A}^1_\mathbb{Z} \) the structure of a group scheme, by describing the three structural morphisms, and showing that they satisfy the axioms. (Hint: the morphisms should not be surprising. For example, inverse is given by \( t \mapsto -t \).

Note that we know that the product \( \mathbb{A}^1_\mathbb{Z} \times \mathbb{A}^1_\mathbb{Z} \) exists, by Exercise 7.6.C(a).)

**7.6.F. Exercise.** Show that if \( G \) is a group object in a category \( C \), then for any \( X \in C \), \( \text{Mor}(X, G) \) has the structure of a group, and the group structure is preserved by pullback (i.e. \( \text{Mor}(-, G) \) is a contravariant functor to \( \text{Groups} \)).

**7.6.G. Exercise.** Show that the group structure described by the previous exercise translates the group scheme structure on \( \mathbb{A}^1_\mathbb{Z} \) to the group structure on \( \Gamma(X, \mathcal{O}_X) \), via the bijection of §7.6.1.

**7.6.H. Exercise.** Define the notion of **ring scheme**, and **abelian group scheme**.

**7.6.4. Group schemes, more functorially.** There was something unsatisfactory about our discussion of the group scheme nature of the bijection in §7.6.1: we observed that the right side (functions on \( X \)) formed a group, then we developed the axioms of a group scheme, then we cleverly figured out the maps that made \( \mathbb{A}^1_\mathbb{Z} \) into a group scheme, then we showed that this induced a group structure on the left side of the bijection (\( \text{Mor}(X, \mathbb{A}^1_\mathbb{Z}) \)) that precisely corresponded to the group structure on the right side (functions on \( X \)).

The picture is more cleanly explained as follows.

**7.6.I. Exercise.** Suppose we have a contravariant functor \( F \) from \( \text{Sch} \) (or indeed any category) to \( \text{Groups} \). Suppose further that \( F \) composed with the forgetful functor \( \text{Groups} \to \text{Sets} \) is representable by an object \( Y \). Show that the group operations on \( F(X) \) (as \( X \) varies through \( \text{Sch} \)) uniquely determine \( m : Y \times Y \to Y \), \( i : Y \to Y \), \( e : Z \to Y \) satisfying the axioms defining a group scheme, such that the group operation on \( \text{Mor}(X, Y) \) is the same as that on \( F(X) \).

In particular, the definition of a group object in a category was forced upon us by the definition of group. More generally, you should expect that any category that can be interpreted as sets with additional structure should fit into this picture.

You should apply this exercise to \( \mathbb{A}^1_X \), and see how the explicit formulas you found in Exercise 7.6.E are forced on you.

**7.6.J. Exercise.** Work out the maps \( m \), \( i \), and \( e \) in the group schemes of Exercise 7.6.C.
7.6.K. Exercise. (a) Define morphism of group schemes.
(b) Define the group scheme GL_n, and describe the determinant map \( \det : \text{GL}_n \to \mathbb{G}_m \).
(c) Make sense of the statement: \( \cdot^n : \mathbb{G}_m \to \mathbb{G}_m \) given by \( t \mapsto t^n \) is a morphism of group schemes.

7.6.L. Exercise (Kernels of Maps of Group Schemes). Suppose \( F : G_1 \to G_2 \) is a morphism of group schemes. Consider the contravariant functor \( \text{Sch} \to \text{Groups} \) given by \( X \mapsto \ker(\text{Mor}(X, G_1) \to \text{Mor}(X, G_2)) \). If this is representable, by group scheme \( G_0 \), say, show that \( G_0 \to G_1 \) is the kernel of \( F \) in the category of group schemes.

7.6.M. Exercise. Show that the kernel of \( \cdot^p \) (Exercise 7.6.K) is representable. Show that over a field \( k \) of characteristic \( p \), this group scheme is non-reduced. (Clarification: \( \mathbb{G}_m \) over a field \( k \) means \( \text{Spec} k[t, t^{-1}] \), with the same group operations. Better: it represents the group of invertible functions in the category of \( k \)-schemes. We can similarly define \( \mathbb{G}_m \) over an arbitrary scheme.)

7.6.N. Exercise. Show (as easily as possible) that \( \mathbb{A}_k^1 \) is a ring scheme.

7.6.5. Aside: Hopf algebras. Here is a notion that we won’t use, but it is easy enough to define now. Suppose \( G = \text{Spec} A \) is an affine group scheme, i.e. a group scheme that is an affine scheme. The categorical definition of group scheme can be restated in terms of the ring \( A \). Then these axioms define a Hopf algebra. For example, we have a “comultiplication map” \( A \to A \otimes A \).

7.6.O. Exercise. As \( \mathbb{A}_k^1 \) is a group scheme, \( k[t] \) has a Hopf algebra structure. Describe the comultiplication map \( k[t] \to k[t] \otimes_k k[t] \).

7.7 ** The Grassmannian (initial construction)

The Grassmannian is a useful geometric construction that is “the geometric object underlying linear algebra”. In (classical) geometry over a field \( K = \mathbb{R} \) or \( \mathbb{C} \), just as projective space parametrizes one-dimensional subspaces of a given \( n \)-dimensional vector space, the Grassmannian parametrizes \( k \)-dimensional subspaces of \( n \)-dimensional space. The Grassmannian \( G(k, n) \) is a manifold of dimension \( k(n - k) \) (over the field). The manifold structure is given as follows. Given a basis \( \{v_1, \ldots, v_n\} \) of \( n \)-space, “most” \( k \)-planes can be described as the span of the \( k \) vectors

\[
(v_1 + \sum_{i=k+1}^n a_{1i}v_i, v_2 + \sum_{i=k+1}^n a_{2i}v_i, \ldots, v_k + \sum_{i=k+1}^n a_{ki}v_i).
\]

(Can you describe which \( k \)-planes are not of this form? Hint: row reduced echelon form. Aside: the stratification of \( G(k, n) \) by normal form is the decomposition of the Grassmannian into Schubert cells. You may be able to show using the normal form that each Schubert cell is isomorphic to an affine space.) Any \( k \)-plane of this form can be described in such a way uniquely. We use this to identify those \( k \)-planes of this form with the manifold \( K^{k(n-k)} \) (with coordinates \( a_{ji} \)). This is a large
affine patch on the Grassmannian (called the “open Schubert cell” with respect to this basis). As the $v_i$ vary, these patches cover the Grassmannian (why?), and the manifold structures agree (a harder fact).

We now define the Grassmannian in algebraic geometry, over a ring $A$. Suppose $v = (v_1, \ldots, v_n)$ is a basis for $A^n$. More precisely: $v_i \in A^n$, and the map $A^n \to A^n$ given by $(a_1, \ldots, a_n) \mapsto a_1 v_1 + \cdots + a_n v_n$ is an isomorphism.

**7.7.A. Exercise.** Show that any two bases are related by an invertible $n \times n$ matrix over $A$ — a matrix with entries in $A$ whose determinant is an invertible element of $A$.

For each such $v$, we consider the scheme $U_v \cong \mathbb{A}_A^{k(n-k)}$, with coordinates $a_{ij}$ ($k + 1 \leq i \leq n$, $1 \leq j \leq k$), which we imagine as corresponding to the $k$-plane spanned by the vectors (7.7.0.1).

**7.7.B. Exercise.** Given two bases $v$ and $w$, explain how to glue $U_v$ to $U_w$ along appropriate open sets. You may find it convenient to work with coordinates $a_{ij}$ where $i$ runs from 1 to $n$, not just $k + 1$ to $n$, but imposing $a_{ij} = \delta_{ij}$ (i.e. 1 when $i = j$ and 0 otherwise). This convention is analogous to coordinates $x_{i/j}$ on the patches of projective space (§5.4.9). Hint: the relevant open subset of $U_v$ will be where a certain determinant doesn’t vanish.

**7.7.C. Exercise/Definition.** By checking triple intersections, verify that these patches (over all possible bases) glue together to a single scheme (Exercise 5.4.A). This is the Grassmannian $G(k, n)$ over the ring $A$.

Although this definition is pleasantly explicit (it is immediate that the Grassmannian is covered by $\mathbb{A}_A^{k(n-k)}$'s), and perhaps more “natural” than our original definition of projective space in §5.4.9 (we aren’t making a choice of basis; we use all bases), there are several things unsatisfactory about this definition of the Grassmannian. In fact the Grassmannian is always projective; this isn’t obvious with this definition. Furthermore, the Grassmannian comes with a natural closed immersion into $\mathbb{P}^{\binom{n}{k}-1}$ (the Plücker embedding). We will address these issues in §17.6 by giving a better description, as a moduli space.
CHAPTER 8

Useful classes of morphisms of schemes

We now define an unreasonable number of types of morphisms. Some (often finiteness properties) are useful because every "reasonable morphism" has such properties, and they will be used in proofs in obvious ways. Others correspond to geometric behavior, and you should have a picture of what each means.

One of Grothendieck’s lessons is that things that we often think of as properties of objects are better understood as properties of morphisms. One way of turning properties of objects into properties of morphisms is as follows. If \( P \) is a property of schemes, we say that a morphism \( f : X \to Y \) has \( P \) if for every affine open \( U \subset X \), \( f^{-1}(U) \) has \( P \). We will see this for \( P = \) quasicompact, quasiseparated, affine, and more. (As you might hope, in good circumstances, \( P \) will satisfy the hypotheses of the Affine Communication Lemma 6.3.2.) Informally, you can think of such a morphism as one where all the fibers have \( P \). (You can quickly define the fiber of a morphism as a topological space, but once we define fiber product, we will define the scheme-theoretic fiber, and then this discussion will make sense.) But it means more than that: it means that the “\( P \)-ness” is really not just fiber-by-fiber, but behaves well as the fiber varies. (For comparison, a smooth morphism of manifolds means more than that the fibers are smooth.)

8.1 Open immersions

An open immersion of schemes is defined to be an open immersion as ringed spaces (§7.2.1). In other words, a morphism \( f : (X, \mathcal{O}_X) \to (Y, \mathcal{O}_Y) \) of schemes is an open immersion if \( f \) factors as

\[
\begin{array}{ccc}
(X, \mathcal{O}_X) & \xrightarrow{g} & (U, \mathcal{O}_Y|_U) \xrightarrow{h} & (Y, \mathcal{O}_Y)
\end{array}
\]

where \( g \) is an isomorphism, and \( U \to Y \) is an inclusion of an open set. It is immediate that isomorphisms are open immersions. We say that \( (U, \mathcal{O}_Y|_U) \) is an open subscheme of \( (Y, \mathcal{O}_Y) \), and often sloppily say that \( (X, \mathcal{O}_X) \) is an open subscheme of \( (Y, \mathcal{O}_Y) \).

8.1.A. Important but easy exercise. Suppose \( i : U \to Z \) is an open immersion, and \( f : Y \to Z \) is any morphism. Show that \( U \times_Z Y \) exists. (Hint: I’ll even tell you what it is: \( (f^{-1}(U), \mathcal{O}_Y|_{f^{-1}(U)}) \).) In particular, if \( U \to Z \) and \( V \to Z \) are open immersions, \( U \times_Z V \cong U \cap V \).

8.1.B. Easy exercise. Suppose \( f : X \to Y \) is an open immersion. Show that if \( Y \) is locally Noetherian, then \( X \) is too. Show that if \( Y \) is Noetherian, then \( X \) is too.
However, show that if $Y$ is quasicompact, $X$ need not be. (Hint: let $Y$ be affine but not Noetherian, see Exercise 4.6.10b.)

“Open immersions” are scheme-theoretic analogues of open subsets. “Closed immersions” are scheme-theoretic analogues of closed subsets, but they have a surprisingly different flavor, as we will see in §9.1.

8.2 Algebraic interlude: Integral morphisms, the Lying Over Theorem, and Nakayama’s lemma

To set up our discussion in the next section on integral morphisms, we develop some algebraic preliminaries. A clever trick we use can also be used to show Nakayama’s lemma, so we discuss that as well.

Suppose $\phi : B \to A$ is a ring homomorphism. We say $a \in A$ is integral over $B$ if $a$ satisfies some monic polynomial $a^n + \cdots + a^1 + \cdots = 0$ where the coefficients lie in $\phi(B)$. A ring homomorphism $\phi : B \to A$ is integral if every element of $A$ is integral over $\phi(B)$. An integral ring homomorphism $\phi$ is an integral extension if $\phi$ is an inclusion of rings. You should think of integral homomorphisms and integral extensions as ring-theoretic generalizations of the notion of algebraic extensions of fields.

8.2.A. Exercise. Show that if $\phi : B \to A$ is a ring homomorphism, $(b_1, \ldots, b_n) = 1$ in $B$, and $B_{b_i} \to A_{\phi(b_i)}$ is integral for all $i$, then $\phi$ is integral.

8.2.B. Exercise. (a) Show that the property of a homomorphism $\phi : B \to A$ being integral is well behaved with respect to localization and quotient of $B$, and quotient of $A$, but not localization of $A$. More precisely: suppose $\phi$ is integral. Show that the induced maps $T^{-1}B \to \phi(T)^{-1}A$, $B/J \to A/\phi(J)A$, and $B \to A/I$ are integral (where $T$ is a multiplicative subset of $B$, $I$ is an ideal of $B$, and $I$ is an ideal of $A$), but $B \to S^{-1}A$ need not be integral (where $S$ is a multiplicative subset of $A$). (Hint for the latter: show that $k[t] \to k[t]$ is an integral homomorphism, but $k[t] \to k[t]_I$ is not.)

(b) Show that the property of $f$ being an integral extension is well behaved with respect to localization and quotient of $B$, but not quotient of $A$. (Hint for the latter: $k[t] \to k[t]$ is an integral extension, but $k[t] \to k[t]/(t)$ is not.)

8.2.C. Exercise. Show that if $C \to B$ and $B \to A$ are both integral homomorphisms, then so is their composition.

The following lemma uses a useful but sneaky trick.

8.2.1. Lemma. — Suppose $\phi : B \to A$ is a ring homomorphism. Then $a \in A$ is integral over $B$ if and only if it is contained in a subalgebra of $A$ that is a finitely generated $B$-module.
Proof. If \( a \) satisfies a monic polynomial equation of degree \( n \), then the \( B \)-submodule of \( A \) generated by \( 1, a, \ldots, a^{n-1} \) is closed under multiplication, and hence a subalgebra of \( A \).

Assume conversely that \( a \) is contained in a subalgebra \( A' \) of \( A \) that is a finitely generated \( B \)-module. Choose a finite generating set \( m_1, \ldots, m_n \) of \( A' \) (as a \( B \)-module). Then \( a m_i = \sum b_{ij} m_j \), for some \( b_{ij} \in B \). Thus

\[
\begin{pmatrix}
m_1 \\
\vdots \\
m_n
\end{pmatrix} = (a I_{n \times n} - [b_{ij}])^{-1} \begin{pmatrix}
m_1 \\
\vdots \\
m_n
\end{pmatrix}.
\]

We can’t invert the matrix \( (a I_{n \times n} - [b_{ij}]) \), but we almost can. Recall that an \( n \times n \) matrix \( M \) has an adjugate matrix \( \text{adj}(M) \) such that \( \text{adj}(M)M = \det(M)I_n \).

The coefficients of \( \text{adj}(M) \) are polynomials in the coefficients of \( M \). Multiplying (8.2.1.1) by \( \text{adj}(a I_{n \times n} - [b_{ij}]) \), we get

\[
\det(a I_{n \times n} - [b_{ij}]) \begin{pmatrix}
m_1 \\
\vdots \\
m_n
\end{pmatrix} = \begin{pmatrix}
0 \\
\vdots \\
0
\end{pmatrix}.
\]

So \( \det(a I - [b_{ij}]) \) annihilates every element of \( A' \), i.e. \( \det(a I - [b_{ij}]) = 0 \). But expanding the determinant yields an integral equation for \( a \) with coefficients in \( B \). \( \square \)

8.2.2. Corollary (finite implies integral). — If \( A \) is a finite \( B \)-algebra (a finitely generated \( B \)-module), then \( \phi \) is an integral homomorphism.

The converse is false: integral does not imply finite, as \( \mathbb{Q} \hookrightarrow \overline{\mathbb{Q}} \) is an integral homomorphism, but \( \overline{\mathbb{Q}} \) is not a finite \( \mathbb{Q} \)-module. (A field extension is integral if it is algebraic.)

8.2.D. Exercise. Suppose \( \phi : B \to A \) is a ring homomorphism. Show that the elements of \( A \) integral over \( B \) form a subalgebra of \( A \).

8.2.3. Remark: transcendence theory. These ideas lead to the main facts about transcendence theory we will need for a discussion of dimension of varieties, see Exercise/Definition 12.2.A.

8.2.4. The Lying Over and Going-Up Theorems. The Lying Over Theorem is a useful property of integral extensions.

8.2.5. The Lying Over Theorem (Cohen-Seidenberg). — Suppose \( \phi : B \to A \) is an integral extension. Then for any prime ideal \( q \subset B \), there is a prime ideal \( p \subset A \) such that \( p \cap B = q \).

8.2.6. Geometric translation: \( \text{Spec } A \to \text{Spec } B \) is surjective. (A map of schemes is surjective if the underlying map of sets is surjective.)
Although this is a theorem in algebra, the name can be interpreted geometrically: the theorem asserts that the corresponding morphism of schemes is surjective, and that “above” every prime \( q \) “downstairs”, there is a prime \( q \) “upstairs”, see Figure 8.1. (For this reason, it is often said that \( p \) “lies over” \( q \) if \( p \cap B = q \).) The following exercise sets up the proof.

**Figure 8.1.** A picture of the Lying Over Theorem 8.2.5 if \( \phi : A \to B \) is an integral extension, then \( \text{Spec } A \to \text{Spec } B \) is surjective

### 8.2.E. \* Exercise

Show that the special case where \( A \) is a field translates to: if \( B \subset A \) is a subring with \( A \) integral over \( B \), then \( B \) is a field. Prove this. (Hint: you must show that all nonzero elements in \( B \) have inverses in \( B \). Here is the start: If \( b \in B \), then \( 1/b \in A \), and this satisfies some integral equation over \( B \).)

*Proof of the Lying Over Theorem 8.2.5.* We first make a reduction: by localizing at \( q \) (preserving integrality by Exercise 8.2.B), we can assume that \( (B, q) \) is a local ring. Then let \( p \) be any maximal ideal of \( A \). Consider the following diagram.

\[
\begin{array}{ccc}
A & \to & A/p \\
\downarrow & & \downarrow \\
B & \to & B/(p \cap B)
\end{array}
\]

(Do you see why the right vertical arrow is an integral extension?) By Exercise 8.2.E, \( B/(p \cap B) \) is a field too, so \( p \cap B \) is a maximal ideal, hence it is \( q \).

### 8.2.F. Important Exercise (The Going-Up Theorem)

Suppose \( \phi : B \to A \) is an integral homomorphism (not necessarily an integral extension). Show that if \( q_1 \subset q_2 \subset \cdots \subset q_n \) is a chain of prime ideals of \( B \), and \( p_1 \subset \cdots \subset p_m \) is a chain of prime ideals of \( A \) such that \( p_i \) “lies over” \( q_i \) (and \( m < n \)), then the second chain can be extended to \( p_1 \subset \cdots \subset p_n \) so that this remains true. (Hint: reduce to the case \( m = 1, n = 2 \); reduce to the case where \( q_1 = (0) \) and \( p_1 = (0) \); use the Lying Over Theorem.)
8.2.7. Nakayama’s lemma.

The trick in the proof of Lemma 8.2.1 can be used to quickly prove Nakayama’s lemma. This name is used for several different but related results, which we discuss here. (A geometric interpretation will be given in Exercise 14.7.b.)

8.2.8. Nakayama’s Lemma version 1. — Suppose $A$ is a ring, $I$ is an ideal of $A$, and $M$ is a finitely-generated $A$-module, such that $M = IM$. Then there exists an $a \in A$ with $a \equiv 1 \pmod{I}$ with $aM = 0$.

Proof. Say $M$ is generated by $m_1, \ldots, m_n$. Then as $M = IM$, we have $m_i = \sum_j a_{ij} m_j$ for some $a_{ij} \in I$. Thus

\[(8.2.8.1) \quad (\text{Id}_n - A) \begin{pmatrix} m_1 \\ \vdots \\ m_n \end{pmatrix} = 0\]

where $\text{Id}_n$ is the $n \times n$ identity matrix in $A$, and $A = (a_{ij})$. Multiplying both sides of (8.2.8.1) on the left by $\text{adj}(\text{Id}_n - A)$, we obtain

\[\det(\text{Id}_n - A) \begin{pmatrix} m_1 \\ \vdots \\ m_n \end{pmatrix} = 0.\]

But when you expand out $\det(\text{Id}_n - A)$, you get something that is $1 \pmod{I}$. □

Here is why you care. Suppose $I$ is contained in all maximal ideals of $A$. (The intersection of all the maximal ideals is called the Jacobson radical, but we won’t use this phrase. For comparison, recall that the nilradical was the intersection of the prime ideals of $A$.) Then I claim that any $a \equiv 1 \pmod{I}$ is invertible. For otherwise (a) $\neq A$, so the ideal $(a)$ is contained in some maximal ideal $m$ — but $a \equiv 1 \pmod{m}$, contradiction. Then as $a$ is invertible, we have the following.

8.2.9. Nakayama’s Lemma version 2. — Suppose $A$ is a ring, $I$ is an ideal of $A$ contained in all maximal ideals, and $M$ is a finitely-generated $A$-module. (The most interesting case is when $A$ is a local ring, and $I$ is the maximal ideal.) Suppose $M = IM$. Then $M = 0$.

8.2.G. Exercise (Nakayama’s Lemma version 3). Suppose $A$ is a ring, and $I$ is an ideal of $A$ contained in all maximal ideals. Suppose $M$ is a finitely generated $A$-module, and $N \subseteq M$ is a submodule. If $N/IN \rightarrow M/IM$ an isomorphism, then $M = N$. (This can be useful, although it won’t be relevant for us.)

8.2.H. Important exercise (Nakayama’s Lemma version 4: generators of $M/mM$ lift to generators of $M$). Suppose $(A, m)$ is a local ring. Suppose $M$ is a finitely-generated $A$-module, and $f_1, \ldots, f_n \in M$, with (the images of) $f_1, \ldots, f_n$ generating $M/mM$. Then $f_1, \ldots, f_n$ generate $M$. (In particular, taking $M = m$, if we have generators of $m/m^2$, they also generate $m$.)
8.2.I. **UNIMPORTANT AND EASY EXERCISE (NAKAYAMA’S LEMMA VERSION 5).** Prove Nakayama version 1 (Lemma 8.2.8) without the hypothesis that $M$ is finitely generated, but with the hypothesis that $I^n = 0$ for some $n$. (This argument does not use the trick.) This result is quite useful, although we won’t use it.

8.2.J. **IMPORTANT EXERCISE GENERALIZING LEMMA 8.2.1.** Suppose $S$ is a subring of a ring $A$, and $r \in A$. Suppose there is a faithful $S[r]$-module $M$ that is finitely generated as an $S$-module. Show that $r$ is integral over $S$. (Hint: change a few words in the proof of Nakayama’s Lemma version 1.)

8.2.K. **EXERCISE.** Suppose $A$ is an integral domain, and $\hat{A}$ is the integral closure of $A$ in $K(A)$, i.e. those elements of $K(A)$ integral over $A$, which form a subalgebra by Exercise 8.2.D. Show that $\hat{A}$ is integrally closed in $K(\hat{A}) = K(A)$.

### 8.3 Finiteness conditions on morphisms

#### 8.3.1. Quasicompact and quasiseparated morphisms.

A morphism $f : X \to Y$ of schemes is **quasicompact** if for every open affine subset $U$ of $Y$, $f^{-1}(U)$ is quasicompact. (Equivalently, the preimage of any quasicompact open subset is quasicompact.)

We will like this notion because (i) we know how to take the maximum of a finite set of numbers, and (ii) most reasonable schemes will be quasicompact.

Along with quasicompactness comes the weird notion of quasiseparatedness. A morphism $f : X \to Y$ is **quasiseparated** if for every affine open subset $U$ of $Y$, $f^{-1}(U)$ is a quasiseparated scheme (§6.1.1). This will be a useful hypothesis in theorems (in conjunction with quasicompactness). Various interesting kinds of morphisms (locally Noetherian source, affine, separated, see Exercises 8.3.B, 8.3.D, and 11.1.C resp.) are quasiseparated, and this will allow us to state theorems more succinctly.

8.3.A. **EASY EXERCISE.** Show that the composition of two quasicompact morphisms is quasicompact. (It is also true that the composition of two quasiseparated morphisms is quasiseparated. This is not easy to show directly, but will follow easily once we understand it in a more sophisticated way, see Exercise 11.1.13(b).)

8.3.B. **EASY EXERCISE.** (a) Show that any morphism from a Noetherian scheme is quasicompact.

(b) Show that any morphism from a locally Noetherian scheme is quasiseparated. (Hint: Exercise 6.3.B) Thus those readers working only with locally Noetherian schemes may take quasiseparatedness as a standing hypothesis.

8.3.C. **EXERCISE.** (Obvious hint for both parts: the Affine Communication Lemma 6.3.2)

(a) (quasicompactness is affine-local on the target) Show that a morphism $f : X \to Y$ is quasicompact if there is a cover of $Y$ by open affine sets $U_i$ such that $f^{-1}(U_i)$ is quasicompact.

(b) (quasiseparatedness is affine-local on the target) Show that a morphism $f : X \to Y$
is quasiseparated if there is cover of $Y$ by open affine sets $U_i$ such that $f^{-1}(U_i)$ is quasiseparated.

Following Grothendieck’s philosophy of thinking that the important notions are properties of morphisms, not of objects, we can restate the definition of quasicompact (resp. quasiseparated) scheme as a scheme that is quasicompact (resp. quasiseparated) over the final object Spec $\mathbb{Z}$ in the category of schemes (Exercise 7.3.I).

8.3.2. Affine morphisms.

A morphism $f : X \to Y$ is **affine** if for every affine open set $U$ of $Y$, $f^{-1}(U)$ (interpreted as an open subscheme of $X$) is an affine scheme.

8.3.D. **Fast Exercise.** Show that affine morphisms are quasicompact and quasiseparated. (Hint for the second: Exercise 6.1.G)

8.3.E. **Exercise (A nonquasiseparated scheme).** Let $X = \text{Spec } k[x_1, x_2, \ldots]$, and let $U$ be $X - [m]$ where $m$ is the maximal ideal $(x_1, x_2, \ldots)$. Take two copies of $X$, glued along $U$. Show that the result is not quasiseparated. Hint: This open immersion $U \subset X$ came up earlier in Exercise 4.6.D(b) as an example of a nonquasicompact open subset of an affine scheme.

8.3.3. Proposition (the property of “affineness” is affine-local on the target). — A morphism $f : X \to Y$ is affine if there is a cover of $Y$ by affine open sets $U$ such that $f^{-1}(U)$ is affine.

This proof is the hardest part of this section. For part of the proof (which will start in §8.3.5), it will be handy to have a lemma.

8.3.4. **Lemma.** — If $X$ is a quasicompact quasiseparated scheme and $s \in \Gamma(X, O_X)$, then the natural map $\Gamma(X, O_X)_s \to \Gamma(\text{D}(s), O_X)$ is an isomorphism.

Here $\text{D}(s)$ means the locus on $X$ where $s$ doesn’t vanish. This was earlier defined only in the case where $X$ was affine, and here we don’t yet know that $X$ is affine, so the quotes are intended to warn you about this.

To repeat the brief reassuring comment on the “quasicompact quasiseparated” hypothesis: this just means that $X$ can be covered by a finite number of affine open subsets, any two of which have intersection also covered by a finite number of affine open subsets (Exercise 6.1.H). The hypothesis applies in lots of interesting situations, such as if $X$ is affine (Exercise 6.1.G) or Noetherian (Exercise 6.3.B).

**Proof.** Cover $X$ with finitely many affine open sets $U_i = \text{Spec } A_i$. Let $U_{ij} = U_i \cap U_j$. Then

$$0 \to \Gamma(X, O_X) \to \prod_i A_i \to \prod_{i,j} \Gamma(U_{ij}, O_X)$$

is exact. By the quasiseparated hypotheses, we can cover each $U_{ij}$ with a finite number of affines $U_{ijk} = \text{Spec } A_{ijk}$, so we have that

$$0 \to \Gamma(X, O_X) \to \prod_i A_i \to \prod_{i,j,k} A_{ijk}$$
is exact. Localizing at $s$ (an exact functor, Exercise 2.6.F(a)) gives

$$0 \to \Gamma(X, \mathcal{O}_X)_s \to \left( \prod_i A_i \right)_s \to \left( \prod_{i,j,k} A_{ijk} \right)_s$$

As localization commutes with finite products (Exercise 2.3.L(b)),

(8.3.4.1) $$0 \to \Gamma(X, \mathcal{O}_X)_s \to \prod_i (A_i)_{s_i} \to \prod_{i,j,k} (A_{ijk})_{s_{ijk}}$$

is exact, where the global function $s$ induces functions $s_i \in A_i$ and $s_{ijk} \in A_{ijk}$.

But similarly, the scheme $\text{Spec}(D(s))$ can be covered by affine opens $\text{Spec}(A_i)_{s_i}$, and $\text{Spec}(A_i)_{s_i} \cap \text{Spec}(A_j)_{s_j}$ are covered by a finite number of affine opens $\text{Spec}(A_{ijk})_{s_{ijk}}$, so we have

(8.3.4.2) $$0 \to \Gamma(X, \mathcal{O}_X)_s \to \prod_i (A_i)_{s_i} \to \prod_{i,j,k} (A_{ijk})_{s_{ijk}}.$$

Notice that the maps $\prod_i (A_i)_{s_i} \to \prod_{i,j,k} (A_{ijk})_{s_{ijk}}$ in (8.3.4.1) and (8.3.4.2) are the same, and we have described the kernel of the map in two ways, so $\Gamma(X, \mathcal{O}_X)_s \to \Gamma(\text{Spec}(D(s)), \mathcal{O}_X)$ is indeed an isomorphism. (Notice how the quasicompact and quasiseparated hypotheses were used in an easy way: to obtain finite products, which would commute with localization.)

\[\square\]

8.3.5. Proof of Proposition 8.3.3 As usual, we use the Affine Communication Lemma 6.3.2. We check our two criteria. First, suppose $f : X \to Y$ is affine over $\text{Spec} B$, i.e. $f^{-1}(\text{Spec} B) = \text{Spec} A$. Then $f^{-1}(\text{Spec} B_s) = \text{Spec} A_{rs}$. Second, suppose we are given $f : X \to \text{Spec} B$ and $(s_1, \ldots, s_n) = B$ with $X_{s_i}$ affine (Spec $A_i$, say). We wish to show that $X$ is affine too. Let $A = \Gamma(X, \mathcal{O}_X)$. Then $X \to \text{Spec} B$ factors through the tautological map $g : X \to \text{Spec} A$ (arising from the (iso)morphism $A \to \Gamma(X, \mathcal{O}_X)$, Exercise 7.3.F).

Then $h^{-1}D(s_i) = D(h^s s_i) \cong \text{Spec} A_{rs}$ (the preimage of a distinguished open set is a distinguished open set), and $f^{-1}D(s_i) = \text{Spec} A_i$. Now $X$ is quasicompact and quasiseparated by the affine-locality of these notions (Exercise 8.3.C), so the hypotheses of Lemma 8.3.4 are satisfied. Hence we have an induced isomorphism of $A_{rs} = \Gamma(X, \mathcal{O}_X)_{h^s}$, so $\Gamma(\text{Spec}(s_i), \mathcal{O}_X) = A_i$. Thus $g$ induces an isomorphism $\text{Spec} A_i \to \text{Spec} A_{rs}$, (an isomorphism of rings induces an isomorphism of affine schemes, by strangely confusing exercise 5.3.A). Thus $g$ is an isomorphism over each Spec $A_{rs}$, which cover Spec $A$, and thus $g$ is an isomorphism. Hence $X \cong \text{Spec} A$, so is affine as desired.

The affine-locality of affine morphisms (Proposition 8.3.3) has some non-obvious consequences, as shown in the next exercise.
8.3.F. Exercise. Suppose $Z$ is a closed subset of an affine scheme $X$ locally cut out by one equation. (In other words, $\text{Spec } A$ can be covered by smaller open sets, and on each such set $Z$ is cut out by one equation.) Show that the complement $Y$ of $Z$ is affine. (This is clear if $Y$ is globally cut out by one equation; then if $X = \text{Spec } A$ then $Y = \text{Spec } A/I$. However, $Y$ is not always of this form, see Exercise [6.4.1].)

8.3.6. Finite and integral morphisms.

Before defining finite and integral morphisms, we give an example to keep in mind. If $L/K$ is a field extension, then $\text{Spec } L \to \text{Spec } K$ (i) is always affine; (ii) is integral if $L/K$ is algebraic; and (iii) is finite if $L/K$ is finite.

An affine morphism $f : X \to Y$ is finite if for every affine open set $\text{Spec } B$ of $Y$, $f^{-1}(\text{Spec } B)$ is the spectrum of a $B$-algebra that is a finitely-generated $B$-module. Warning about terminology (finite vs. finitely-generated): Recall that if we have a ring homomorphism $A \to B$ such that $B$ is a finitely-generated $A$-module then we say that $B$ is a finite $A$-algebra. This is stronger than being a finitely-generated $A$-algebra.

By definition, finite morphisms are affine.

8.3.G. Exercise (the property of finiteness is affine-local on the target). Show that a morphism $f : X \to Y$ is finite if there is a cover of $Y$ by affine open sets $\text{Spec } A$ such that $f^{-1}(\text{Spec } A)$ is the spectrum of a finite $A$-algebra.

The following four examples will give you some feeling for finite morphisms. In each example, you will notice two things. In each case, the maps are always finite-to-one (as maps of sets). We will verify this in general in Exercise [8.3.K]. You will also notice that the morphisms are closed as maps of topological spaces, i.e. the images of closed sets are closed. We will show that finite morphisms are always closed in Exercise [8.3.N] and give a second proof in [§ 9.2.4]. Intuitively, you should think of finite as being closed plus finite fibers, although this isn’t quite true. We will make this precise later.

Example 1: Branched covers. Consider the morphism $\text{Spec } k[t] \to \text{Spec } k[u]$ given by $u \mapsto p(t)$, where $p(t) \in k[t]$ is a degree $n$ polynomial (see Figure [8.2]). This is finite: $k[t]$ is generated as a $k[u]$-module by $1, t, t^2, \ldots, t^{n-1}$.

![Figure 8.2](image)

Figure 8.2. The “branched cover” $A_k^1 \to A_k^1$ of the “$u$-line” by the “$t$-line” given by $u \mapsto p(t)$ is finite.

Example 2: Closed immersions (to be defined soon, in § 9.1). If $I$ is an ideal of a ring $A$, consider the morphism $\text{Spec } A/I \to \text{Spec } A$ given by obvious map $A \to A/I$ (see
Figure 8.3. This is a finite morphism \((A/I)\) is generated as a \(A\)-module by the element \(1 \in A/I\).

**Figure 8.3.** The “closed immersion” \(\text{Spec } k \to \text{Spec } k[t]\) is finite.

**Example 3: Normalization (to be defined in \(\S 10.6\).)** Consider the morphism \(\text{Spec } k[t] \to \text{Spec } k[x, y]/(y^2 - x^2 - x^3)\) corresponding to \(k[x, y]/(y^2 - x^2 - x^3) \to k[t]\) given by \((x, y) \mapsto (t^2 - 1, t^3 - t)\) (check that this is a well-defined ring map!), see Figure 8.4. This is a finite morphism, as \(k[t]\) is generated as a \((k[x, y]/(y^2 - x^2 - x^3))\)-module by 1 and \(t\). (The figure suggests that this is an isomorphism away from the “node” of the target. You can verify this, by checking that it induces an isomorphism between \(D(t^2 - 1)\) in the source and \(D(x)\) in the target. We will meet this example again!)

**Figure 8.4.** The “normalization” \(\text{Spec } k[t] \to \text{Spec } k[x, y]/(y^2 - x^2 - x^3)\) given by \((x, y) \mapsto (t^2 - 1, t^3 - t)\) is finite.

**8.3.H. Important Exercise (Example 4, finite morphisms to \(\text{Spec } k\).)** Show that if \(X \to \text{Spec } k\) is a finite morphism, then \(X\) is a discrete finite union of points, each with residue field a finite extension of \(k\), see Figure 8.5. (An example is \(\text{Spec } \mathbb{F}_8 \times \mathbb{F}_4[x, y]/(x^2, y^4) \times \mathbb{F}_4[t]/(t^3) \times \mathbb{F}_2 \to \text{Spec } \mathbb{F}_2\).) Do not just quote some fancy theorem! (Possible approach: Show that any integral domain which is a finite \(k\)-algebra must be a field. Show that every prime \(p\) of \(A\) is maximal. Show that the irreducible components of \(A\) are closed points. Show \(\text{Spec } A\) is discrete and hence finite. Show that the residue fields of \(A/p\) are finite field extensions of \(k\).)
8.3.I. EASY EXERCISE (cf. Exercise 8.2.C). Show that the composition of two finite morphisms is also finite.

8.3.J. EXERCISE: FINITE MORPHISMS TO Spec A ARE PROJECTIVE. If B is a finite A-algebra, define a graded ring $S_\bullet$ by $S_0 = A$, and $S_n = B$ for $n > 0$. (What is the multiplicative structure? Hint: you know how to multiply elements of B together, and how to multiply elements of A with elements of B.) Describe an isomorphism $\text{Proj } S_\bullet \cong \text{Spec } B$.

8.3.K. IMPORTANT EXERCISE. Show that finite morphisms have finite fibers. (This is a useful exercise, because you will have to figure out how to get at points in a fiber of a morphism: given $f : X \rightarrow Y$, and $y \in Y$, what are the points of $f^{-1}(y)$? Hint: if $X = \text{Spec } A$ and $Y = \text{Spec } B$ are both affine, and $y = [p]$, then we can throw out everything in $A$ outside $\mathfrak{g}$ by modding out by $p$; you can show that the preimage is $A/\mathfrak{g}$. Then we have reduced to the case where $Y$ is the Spec of an integral domain, and $[p] = [0]$ is the generic point. We can throw out the rest of the points by localizing at 0. You can show that the preimage is $(A_p)/pA_p$ (cf. (5.3.4.1)).

That finiteness behaves well with respect to the operations you made done, you have reduced the problem to Exercise 8.3.H.)

8.3.7. Example. The open immersion $\mathbb{A}^2 - \{(0, 0)\} \rightarrow \mathbb{A}^2$ has finite fibers, but is not affine (as $\mathbb{A}^2 - \{(0, 0)\}$ isn’t affine, 5.4.1) and hence not finite.

8.3.L. EASY EXERCISE. Show that the open immersion $\mathbb{A}^1 - \{0\} \rightarrow \mathbb{A}^1$ has finite fibers and is affine, but is not finite.

8.3.M. LESS IMPORTANT EXERCISE. (This exercise shows that the seemingly set-theoretic notion of surjectivity is also quite algebraic.) Suppose that $f : \text{Spec } A \rightarrow \text{Spec } B$ is a finite morphism, corresponding to $\phi : B \rightarrow A$. Show that A is surjective if and only if $\phi$ is injective. (One direction is the Lying Over Theorem 8.2.5.)

8.3.8. Definition. A morphism $\pi : X \rightarrow Y$ of schemes is integral if $\pi$ is affine, and for every affine open subset $\text{Spec } B \subset Y$, with $\pi^{-1}(\text{Spec } B) = \text{Spec } A$, the induced
map $A \to B$ is an integral homomorphism of rings. This is an affine-local condition by Exercises 8.2.A and 8.2.B and the Affine Communication Lemma 6.3.2. It is closed under composition by Exercise 8.2.C. Integral morphisms are mostly useful because finite morphisms are integral by Corollary 8.2.2. Note that the converse implication doesn’t hold (witness Spec $\mathbb{Q} \to$ Spec $\mathbb{Q}$, as discussed after the statement of Corollary 8.2.2).

8.3.N. Exercise. Prove that integral morphisms are closed, i.e. that the image of closed subsets are closed. (Hence finite morphisms are closed. A second proof will be given in §9.2.4.) Hint: Reduced to the affine case. If $f^* : B \to A$ is a ring map, inducing finite $f : \text{Spec } A \to \text{Spec } B$, then suppose $I \subset A$ cuts out a closed set of $\text{Spec } A$, and $J = (f^*)^{-1}(1)$, then note that $B/J \subset A/I$, and apply the Lying Over Theorem 8.2.5 here.

8.3.O. Unimportant exercise. Suppose $f : B \to A$ is integral. Show that for any ring homomorphism $B \to C$, $C \to A \otimes B C$ is integral. (Hint: We wish to show that any $\sum_{i=1}^n a_i \otimes c_i \in A \otimes B C$ is integral over $C$. Use the fact that each of the finitely many $a_i$ are integral over $B$, and the Exercise 8.2.D.) Once we know what “base change” is, this will imply that the property of integrality of a morphism is preserved by base change.

8.3.9. Morphisms (locally) of finite type. A morphism $f : X \to Y$ is locally of finite type if for every affine open set $\text{Spec } B$ of $Y$, and every affine open subset $\text{Spec } A$ of $f^{-1}(\text{Spec } B)$, the induced morphism $B \to A$ expresses $A$ as a finitely generated $B$-algebra. By the affine-locality of finite-typeness of $B$-schemes (see Proposition 6.3.3), this is equivalent to: $f^{-1}(\text{Spec } B)$ can be covered by affine open subsets $\text{Spec } A_i$ so that each $A_i$ is a finitely generated $B$-algebra.

A morphism is of finite type if it is locally of finite type and quasicompact. Translation: for every affine open set $\text{Spec } B$ of $Y$, $f^{-1}(\text{Spec } B)$ can be covered with a finite number of open sets $\text{Spec } A_i$ so that the induced morphism $B \to A_i$ expresses $A_i$ as a finitely generated $B$-algebra.

8.3.10. Side remark. It is a common practice to name properties as follows: $P = \text{locally } P \text{ plus } \text{quasicompact}$. Two exceptions are “ringed space” (§7.3) and “finite presentation” (§8.3.13).

8.3.P. Exercise (The notions “locally of finite type” and “finite type” are affine-local on the target). Show that a morphism $f : X \to Y$ is locally of finite type if there is a cover of $Y$ by affine open sets $\text{Spec } B_i$ such that $f^{-1}(\text{Spec } B_i)$ is locally of finite type over $B_i$.

Example: the “structure morphism” $\mathbb{P}^n_A \to \text{Spec } A$ is of finite type, as $\mathbb{P}^n_A$ is covered by $n + 1$ open sets of the form $\text{Spec } A[x_1, \ldots, x_n]$.

Our earlier definition of schemes of “finite type over $k$” (or “finite type $k$-schemes”) from (6.3.2) is now a special case of this more general notion: a scheme $X$ is of finite type over $k$ means that we are given a morphism $X \to \text{Spec } k$ (the “structure morphism”) that is of finite type.

Here are some properties enjoyed by morphisms of finite type.
8.3.Q. Exercise (finite = integral + finite type). (a) (easier) Show that finite morphisms are of finite type.
(b) Show that a morphism is finite if and only if it is integral and of finite type.

8.3.R. Exercises (not hard, but important).

(a) Show that every open immersion is locally of finite type. Show that every open immersion into a locally Noetherian scheme is of finite type. More generally, show that every quasicompact open immersion is of finite type.
(b) Show that the composition of two morphisms locally of finite type is locally of finite type. (Hence as the composition of two quasicompact morphisms is quasicompact, the composition of two morphisms of finite type is of finite type.)
(c) Suppose \( f : X \to Y \) is locally of finite type, and \( Y \) is locally Noetherian. Show that \( X \) is also locally Noetherian. If \( X \to Y \) is a morphism of finite type, and \( Y \) is Noetherian, show that \( X \) is Noetherian.

8.3.11. Definition. A morphism \( f \) is quasifinite if it is of finite type, and for all \( y \in Y \), \( f^{-1}(y) \) is a finite set. The main point of this definition is the “finite fiber” part; the “finite type” hypothesis will ensure that this notion is “preserved by fibered product,” Exercise [0.4.C].

Combining Exercise [8.3.K] with Exercise [8.3.Q]a), we see that finite morphisms are quasifinite. There are quasifinite morphisms which are not finite, such as \( \mathbb{A}^2 - \{(0,0)\} \to \mathbb{A}^2 \) (Example [8.3.7]). A key example of a morphism with finite fibers that is not quasifinite is \( \text{Spec } \mathbb{C}(t) \to \text{Spec } \mathbb{C} \). Another is \( \text{Spec } \overline{\mathbb{Q}} \to \text{Spec } \mathbb{Q} \).

8.3.12. How to picture quasifinite morphisms. If \( X \to Y \) is a finite morphism, then any quasi-compact open subset \( U \subset X \) is quasi-finite over \( Y \). In fact every reasonable quasifinite morphism arises in this way. (This simple-sounding statement is in fact a deep and important result — Zariski’s Main Theorem.) Thus the right way to visualize quasifiniteness is as a finite map with some (closed locus of) points removed.

8.3.13. ** Morphisms (locally) of finite presentation.

There is a variant often useful to non-Noetherian people. A morphism \( f : X \to Y \) is locally of finite presentation (or locally finitely presented) if for each affine open set \( \text{Spec } B \) of \( Y \), \( f^{-1}(\text{Spec } B) = \bigcup_i \text{Spec } A_i \) with \( B \to A_i \) finitely presented (finitely generated with a finite number of relations). A morphism is of finite presentation (or finitely presented) if it is locally of finite presentation and quasiseparated and quasicompact. This is a violation of the general principle that erasing “locally” is the same as adding “quasicompact and” (Remark [8.3.10]). But it is well motivated: finite presentation means “finite in all possible ways” (each affine has a finite number of generators, and a finite number of relations, and a finite number of such affines cover, and their intersections are also covered by a finite number affines) — it is all you would hope for in a scheme without it actually being Noetherian.

If \( X \) is locally Noetherian, then locally of finite presentation is the same as locally of finite type, and finite presentation is the same as finite type. So if you are a Noetherian person, you don’t need to worry about this notion.
8.3.S. Exercise. Show that the notion of “locally of finite presentation” is affine-local on the target.

8.3.T. Exercise. Show that the composition of two finitely presented morphisms is finitely presented.

8.4 Images of morphisms: Chevalley’s theorem and elimination theory

In this section, we will answer a question that you may have wondered about long before hearing the phrase “algebraic geometry”. If you have a number of polynomial equations in a number of variables with indeterminate coefficients, you would reasonably ask what conditions there are on the coefficients for a (common) solution to exist. Given the algebraic nature of the problem, you might hope that the answer should be purely algebraic in nature — it shouldn’t be “random”, or involve bizarre functions like exponentials or cosines. This is indeed the case, and it can be profitably interpreted as a question about images of maps of varieties or schemes, in which guise it is answered by Chevalley’s theorem.

In special cases, the image is nicer still. For example, we have seen that finite morphisms are closed (the image of closed subsets under finite morphisms are closed, Exercise 8.3.N). We will prove a classical result, the Fundamental Theorem of Elimination Theory 8.4.5, which essentially generalizes this (as explained in §9.2.4) to maps from projective space. We will use it repeatedly.

8.4.1. Chevalley’s theorem.

If $f : X \to Y$ is a morphism of schemes, the notion of the image of $f$ as sets is clear: we just take the points in $Y$ that are the image of points in $X$. We know that the image can be open (open immersions), and we will soon see that it can be closed (closed immersions), and hence locally closed (locally closed immersions). But it can be weirder still: consider the morphism $\mathbb{A}^2_k \to \mathbb{A}^2_k$ given by $(x, y) \mapsto (x, xy)$. The image is the plane, with the x-axis removed, but the origin put back in. This isn’t so horrible. We make a definition to capture this phenomenon. A constructible subset of a Noetherian topological space is a subset which belongs to the smallest family of subsets such that (i) every open set is in the family, (ii) a finite intersection of family members is in the family, and (iii) the complement of a family member is also in the family. For example the image of $(x, y) \mapsto (x, xy)$ is constructible. (A generalization of the notion of constructibility to more general topological spaces is mentioned in Exercise 8.4.F)

8.4.A. Exercise: Constructible subsets are finite unions of locally closed subsets. Recall that a subset of a topological space $X$ is locally closed if it is the intersection of an open subset and a closed subset. (Equivalently, it is an open subset of a closed subset, or a closed subset of an open subset. We will later have trouble extending this to open and closed and locally closed subschemes, see Exercise 9.1.K.) Show that a subset of a Noetherian topological space $X$ is constructible
One useful property of constructible subsets of schemes is that there is a short criterion for openness: a constructible subset is open if it is “closed under generalization” (see Exercise 24.2.N).

The image of a morphism of schemes can be stranger than constructible. Indeed if $S$ is any subset of a scheme $Y$, it can be the image of a morphism: let $X$ be the disjoint union of spectra of the residue fields of all the points of $S$, and let $f : X \to Y$ be the natural map. This is quite pathological, but in any reasonable situation, the image is essentially no worse than arose in the previous example of $(x, y) \mapsto (x, xy)$. This is made precise by Chevalley’s theorem.

8.4.2. Chevalley’s theorem. — If $\pi : X \to Y$ is a finite type morphism of Noetherian schemes, the image of any constructible set is constructible. In particular, the image of $\pi$ is constructible.

**Proof.** We begin with a series of reductions.

8.4.B. Exercise.

(a) Reduce to the case where $Y$ is affine, say $Y = \text{Spec } B$.
(b) Reduce further to the case where $X$ is affine.
(c) Reduce further to the case where $X = \mathbb{A}^n_B = \text{Spec } B[t_1, \ldots, t_n]$.
(d) By induction on $n$, reduce further to the case where $X = \mathbb{A}^1_B = \text{Spec } B[t]$.
(e) Reduce to showing that for any Noetherian ring $B$, and any irreducible locally closed subset $Z \subset \mathbb{A}^1_B$, the image of $Z$ under the projection $\pi : \mathbb{A}^1_B \to \text{Spec } B$ is constructible.
(f) Reduce to showing that for any Noetherian integral domain $B$ (with $\pi : \mathbb{A}^1_B \to B$), and any irreducible locally closed subset $Z \subset \mathbb{A}^1_B$, where $\pi|_Z : Z \to \text{Spec } B$ is dominant, $\pi(Z)$ is constructible. (Hint: replace Spec $B$ from (e) by the closure of the image of the generic point of $Z$.)
(g) Use Noetherian induction to show that it suffices to show that for any Noetherian integral domain $B$ (with $\pi : \mathbb{A}^1_B \to B$), and any locally closed subset $Z \subset \mathbb{A}^1_B$, where $\pi(Z)$ contains a non-empty open subset of Spec $B$.

8.4.C. Exercise. Reduce to showing the following statement. Given Noetherian integral domains $B$ and $C$, where $C$ is a $B$-algebra generated by a single element $t$ (possibly with some relations), and the induced map $\pi : \text{Spec } C \to \text{Spec } B$ is dominant (with $\pi$ thus inducing an inclusion $B \hookrightarrow C$), then for any nonzero $g \in C$, $\pi(D(g))$ contains a nonempty open subset of Spec $B$. Hint: choose $g \in C$ so that its set is the closure of $Z$ in $\mathbb{A}^1_B$ in the statement given in Exercise 8.4.B(g), and choose $g \in C$ such that $D(g) \subset \pi(D(g))$. (Optional: draw a picture.)

We now prove this statement. If $C = B[t]/I$, then we deal first with the case $I = 0$, and second with $I \neq 0$.

8.4.D. Exercise. Prove the statement of Exercise 8.4.C in the case $C = B[t]$ as follows. Write $g = \sum_{i=0}^n b_i t^i$, where $b_i \in B$ and $b_n \neq 0$. Show that $D(b_n) \subset \pi(D(g))$. 
where the images of

(a) Suppose that \( \sum_{i=0}^{n} b_i t^i \in I \), where \( b_n \neq 0 \). Show that \( \text{Spec } C \to \text{Spec } B \) is finite over \( D(b_n) \). More precisely, show that \( C_{b_n} \) is generated as a \( B_{b_n} \)-module by (the images of) \( 1, t, \ldots, t^{n-1} \).

Thus by replacing \( B \) by \( B_{b_n} \), we may assume that \( \text{Spec } C \to \text{Spec } B \) is finite. But finite morphisms are closed (Exercise 8.3.N), so the image of \( V(g) \) is closed, and doesn’t contain the generic point of \( \text{Spec } B \) (why?). Thus its complement is dense and open in \( \text{Spec } B \), so in particular \( \pi(D(g)) \) contains a dense open subset of \( \text{Spec } B \).

**8.4.E. Exercise (Chevalley’s theorem for locally finitely presented morphisms).** If you are macho and are embarrassed by Noetherian rings, the following extension of Chevalley’s theorem will give you a sense of one of the standard ways of removing Noetherian hypotheses.

(a) Suppose that \( A \) is a finitely presented \( B \)-algebra (\( B \) not necessarily Noetherian), so \( A = B[x_1, \ldots, x_n]/(f_1, \ldots, f_t) \). Show that the image of \( \text{Spec } A \to \text{Spec } B \) is a finite union of locally closed subsets of \( \text{Spec } B \). Hint: describe \( \text{Spec } A \to \text{Spec } B \) as the base change of

\[
\text{Spec } \mathbb{Z}[x_1, \ldots, x_n, a_1, \ldots, a_n]/(g_1, \ldots, g_n) \to \text{Spec } \mathbb{Z}[a_1, \ldots, a_N],
\]

where the images of \( a_i \) in \( \text{Spec } B \) are the coefficients of the \( f_i \) (there is one \( a_i \) for each coefficient of each \( f_i \)), and \( g_i \mapsto f_i \).

(b) Show that if \( \pi : X \to Y \) is a quasicompact locally finitely presented morphism, and \( Y \) is quasicompact, then \( \pi(X) \) is a finite union of locally closed subsets. (For hardened experts only: [EGA] 0III.9.1 gives a definition of constructibility, and local constructability, in more generality. The general form of Chevalley’s constructibility theorem [EGA] IV,1.8.4 is that the image of a locally constructible set, under a finitely presented map, is also locally constructible.)

**8.4.3. * Elimination of quantifiers.** A basic sort of question that arises in any number of contexts is when a system of equations has a solution. Suppose for example you have some polynomials in variables \( x_1, \ldots, x_n \) over an algebraically closed field \( \mathbb{K} \), some of which you set to be zero, and some of which you set to be nonzero. (This question is of fundamental interest even before you know any scheme theory!) Then there is an algebraic condition on the coefficients which will tell you if this is the case. Define the Zariski topology on \( \mathbb{K}^n \) in the obvious way: closed subsets are cut out by equations.

**8.4.G. Exercise (elimination of quantifiers, over an algebraically closed field).** Fix an algebraically closed field \( \mathbb{K} \). Suppose

\[
f_1, \ldots, f_p, g_1, \ldots, g_q \in \mathbb{K}[A_1, \ldots, A_m, X_1, \ldots X_n]
\]

are given. Show that there is a Zariski-constructible subset \( Y \) of \( \mathbb{K}^m \) such that

\[
\text{(8.4.3.1)} \quad f_1(a_1, \ldots, a_m, X_1, \ldots, X_n) = \cdots = f_p(a_1, \ldots, a_m, X_1, \ldots, X_n) = 0
\]

and

\[
\text{(8.4.3.2)} \quad g_1(a_1, \ldots, a_m, X_1, \ldots, X_n) \neq 0 \quad \cdots \quad g_q(a_1, \ldots, a_m, X_1, \ldots, X_n) \neq 0
\]
has a solution \((X_1, \ldots, X_n) = (x_1, \ldots, x_n) \in \mathbb{K}^n\) if and only if \((a_1, \ldots, a_m) \in Y\).

Hints: if \(Z\) is a finite type scheme over \(k\), and the closed points are denoted \(Z^{\text{cl}}\) ("cl" is for either "closed" or "classical"), then under the inclusion of topological spaces \(Z^{\text{cl}} \hookrightarrow Z\), the Zariski topology on \(Z\) induces the Zariski topology on \(Z^{\text{cl}}\). Note that we can identify \((\mathbb{A}^p_k)^{\text{cl}}\) with \(k^p\) by the Nullstellensatz (Exercise 6.3.E). If \(X\) is the locally closed subset of \(\mathbb{A}^{m+n}\) cut out by the equalities and inequalities (8.4.3.1) and (8.4.3.2), we have the diagram

\[
\begin{array}{ccc}
\mathbb{A}^{m+n} & \xrightarrow{\pi^{\text{cl}}} & \mathbb{A}^m \\
\downarrow & & \downarrow \\
Y & \xrightarrow{\pi} & \mathbb{A}^m
\end{array}
\]

where \(Y = \text{im} \pi^{\text{cl}}\). By Chevalley’s theorem 8.4.2, \(\text{im} \pi\) is constructible, and hence so is \((\text{im} \pi) \cap \mathbb{A}^m\). It remains to show that \((\text{im} \pi) \cap \mathbb{A}^m = Y = (\text{im} \pi^{\text{cl}})\). You might use the Nullstellensatz.

This is called “elimination of quantifiers” because it gets rid of the quantifier “there exists a solution”. The analogous statement for real numbers, where inequalities are also allowed, is a special case of Tarski’s celebrated theorem of elimination of quantifiers for real closed fields.

**8.4.4. The Fundamental Theorem of Elimination Theory.**

**8.4.5. Theorem (Fundamental Theorem of Elimination Theory).** — The morphism \(\pi : \mathbb{P}^n_A \to \text{Spec} A\) is closed (sends closed sets to closed sets).

A great deal of classical algebra and geometry is contained in this theorem as special cases. Here are some examples.

First, let \(A = k[a, b, c, \ldots, i]\), and consider the closed subscheme of \(\mathbb{P}^2_A\) (taken with coordinates \(x, y, z\) corresponding to \(ax + by + cz = 0, dx + ey + fz = 0, gx + hy + iz = 0\). Then we are looking for the locus in \(\text{Spec} A\) where these equations have a non-trivial solution. This indeed corresponds to a Zariski-closed set — where

\[
\det \begin{pmatrix} a & b & c \\ d & e & f \\ g & h & i \end{pmatrix} = 0.
\]

Thus the idea of the determinant is embedded in elimination theory.

As a second example, let \(A = k[a_0, a_1, \ldots, a_m, b_0, b_1, \ldots, b_n]\). Now consider the closed subscheme of \(\mathbb{P}^1_A\) (taken with coordinates \(x, y\) corresponding to \(a_0x^m + a_1x^{m-1}y + \cdots + a_mym = 0\) and \(b_0x^n + b_1x^{n-1}y + \cdots + b_nyn = 0\). Then there is a polynomial in the coefficients \(a_0, \ldots, b_n\) (an element of \(A\)) which vanishes if and only if these two polynomials have a common non-zero root — this polynomial is called the resultant.

More generally, this question boils down to the following question. Given a number of homogeneous equations in \(n + 1\) variables with indeterminate coefficients, Theorem 8.4.5 implies that one can write down equations in the coefficients that will precisely determine when the equations have a nontrivial solution.
Proof of the Fundamental Theorem of Elimination Theory \[8.4.5\]

Suppose \( Z \hookrightarrow \mathbb{P}^n_A \) is a closed subset. We wish to show that \( \pi(Z) \) is closed. (See Figure 8.6.)

\[ \begin{array}{c}
\pi^{-1}y \\
\downarrow \pi \\
D(f) \\
\downarrow \\
y \\
\end{array} \]

\( \mathbb{P}^n_A \)

**Figure 8.6.**

Suppose \( y \notin \pi(Z) \) is a closed point of Spec \( A \). We will check that there is a distinguished open neighborhood \( D(f) \) of \( y \) in Spec \( A \) such that \( D(f) \) doesn’t meet \( \pi(Z) \). (If we could show this for all points of \( \pi(Z) \), we would be done. But I prefer to concentrate on closed points first for simplicity.) Suppose \( y \) corresponds to the maximal ideal \( m \) of \( A \). We seek \( f \in A - m \) such that \( \pi^*f \) vanishes on \( Z \).

Let \( U_0, \ldots, U_n \) be the usual affine open cover of \( \mathbb{P}^n_A \). The closed subsets \( \pi^{-1}y \) and \( Z \) do not intersect. On the affine open set \( U_i \), we have two closed subsets \( Z \cap U_i \) and \( \pi^{-1}y \cap U_i \) that do not intersect, which means that the ideals corresponding to the two closed sets generate the unit ideal, so in the ring of functions \( A[x_0/i, x_1/i, \ldots, x_n/i]/(x_i/i - 1) \) on \( U_i \), we can write

\[ 1 = a_i + \sum m_{ij}g_{ij} \]

where \( m_{ij} \in m \), and \( a_i \) vanishes on \( Z \). Note that \( a_i, g_{ij} \in A[x_0/i, \ldots, x_n/i]/(x_i/i - 1) \), so by multiplying by a sufficiently high power \( x_i^n \) of \( x_i \), we have an equality

\[ x_i^N = a'_i + \sum m_{ij}g'_{ij} \]

in \( S_n = A[x_0, \ldots, x_n] \). We may take \( N \) large enough so that it works for all \( i \). Thus for \( N' \) sufficiently large, we can write any monomial in \( x_1, \ldots, x_n \) of degree \( N' \) as something vanishing on \( Z \) plus a linear combination of elements of \( m \) times other polynomials. Hence

\[ S_{N'} = I(Z)_{N'} + mS_{N'} \]

where \( I(Z) \) is the graded ideal of functions vanishing on \( Z \). By Nakayama’s lemma (version 1, Lemma 8.2.8), taking \( M = S_{N'}/I(Z)_{N'} \), we see that there exists \( f \in A - m \) such that

\[ fS_{N'} \subset I(Z)_{N'} \]

Thus we have found our desired \( f \).
We now tackle Theorem 8.4.5 in general, by simply extending the above argument so that \( y \) need not be a closed point. Suppose \( y = [p] \) not in the image of \( Z \). Applying the above argument in \( \text{Spec} \mathcal{A}_p \), we find \( S_{N'} \otimes_{\mathcal{A}_p} A_p = I(Z)_{N'} \otimes_{\mathcal{A}_p} A_p + mS_{N'} \otimes_{\mathcal{A}_p} A_p \), from which \( g(S_{N'}/I(Z)_{N'}) \otimes_{\mathcal{A}_p} A_p = 0 \) for some \( g \in \mathcal{A}_p - p\mathcal{A}_p \), from which \( (S_{N'}/I(Z)_{N'}) \otimes_{\mathcal{A}_p} A_p = 0 \). As \( S_{N'} \) is a finitely generated \( \mathcal{A} \)-module, there is some \( f \in \mathcal{A} - p \) with \( fS_{N} \subset I(Z) \) (if the module-generators of \( S_{N'} \) are \( h_1, \ldots, h_a \), and \( f_1, \ldots, f_a \) are annihilate the generators \( h_1, \ldots, h_a \), respectively, then take \( f = \prod f_i \)), so once again we have found \( D(f) \) containing \( p \), with (the pullback of) \( f \) vanishing on \( Z \). \( \square \)

Notice that projectivity was crucial to the proof: we used graded rings in an essential way.
CHAPTER 9

Closed immersions and related notions

9.1 Closed immersions and closed subschemes

Just as open immersions (the scheme-theoretic version of open set) are locally modeled on open sets \( U \subset Y \), the analogue of closed subsets also has a local model. This was foreshadowed by our understanding of closed subsets of Spec \( B \) as roughly corresponding to ideals. If \( I \subset B \) is an ideal, then \( \text{Spec} B/I \hookrightarrow \text{Spec} B \) is a morphism of schemes, and we have checked that on the level of topological spaces, this describes Spec \( B/I \) as a closed subset of Spec \( B \), with the subspace topology (Exercise 4.4.H). This morphism is our “local model” of a closed immersion.

9.1.1. Definition. A morphism \( f : X \to Y \) is a closed immersion if it is an affine morphism, and for each open subset \( \text{Spec} B \subset Y \), with \( f^{-1}(\text{Spec} B) \cong \text{Spec} A, B \to A \) is a surjective map (i.e. of the form \( B \to B/I \), our desired local model). If \( X \) is a subset of \( Y \) (and \( f \) on the level of sets is the inclusion), we say that \( X \) is a closed subscheme of \( Y \).

9.1.A. Easy exercise. Show that closed immersions are finite, hence of finite type.

9.1.B. Easy exercise. Show that the composition of two closed immersions is a closed immersion.

9.1.C. Exercise. Show that the property of being a closed immersion is affine-local on the target.

A closed immersion \( f : X \hookrightarrow Y \) determines an ideal sheaf on \( Y \), as the kernel \( \mathcal{I}_{X/Y} \) of the map of \( \mathcal{O}_Y \)-modules

\[
\mathcal{O}_Y \to f_*\mathcal{O}_X
\]

(An ideal sheaf on \( Y \) is what it sounds like: it is a sheaf of ideals. It is a sub-\( \mathcal{O}_Y \)-module \( I \hookrightarrow \mathcal{O}_Y \). On each open subset, it gives an ideal \( I(U) \hookrightarrow \mathcal{O}_Y(U) \).) We thus have an exact sequence (of \( \mathcal{O}_Y \)-modules) \( 0 \to \mathcal{I}_{X/Y} \to \mathcal{O}_Y \to f_*\mathcal{O}_X \to 0 \).

Thus for each affine open subset \( \text{Spec} B \hookrightarrow Y \), we have an ideal \( I_B \subset B \), and we can recover \( X \) from this information: the \( I_B \) (as \( \text{Spec} B \hookrightarrow Y \) varies over the affine opens) defines an \( \mathcal{O} \)-module on the base, hence an \( \mathcal{O}_Y \)-module on \( Y \), and the cokernel of \( I \hookrightarrow \mathcal{O}_Y \) yields \( X \). It will be useful to understand when the information of the \( I_B \) (for all affine opens \( \text{Spec} B \hookrightarrow Y \)) actually determine a closed subscheme.
Our life is complicated by the fact that the answer is “not always”, as shown by the following example.

**9.1.D. UNIMPORTANT EXERCISE.** Let \( X = \text{Spec} \ k[x]/(x^2) \), the germ of the affine line at the origin, which has two points, the closed point and the generic point \( \eta \).

Define \( \mathcal{I}(X) = \{0\} \subseteq \mathcal{O}_X(X) = k[x]/(x^2) \), and \( \mathcal{I}(\eta) = k(\eta) = \mathcal{O}_X(\eta) \). Show that this sheaf of ideals does not correspond to a closed subscheme. (Possible hint: do the next exercise first.)

The next exercise gives a necessary condition.

**9.1.E. EXERCISE.** Suppose \( \mathcal{I}_{X/Y} \) is a sheaf ideals corresponding to a closed immersion \( X \hookrightarrow Y \). Suppose \( \text{Spec} B_f \) is a distinguished open of the affine open \( \text{Spec} B \hookrightarrow Y \). Show that the natural map \( (I_B)_f \rightarrow I_{(B_f)} \) is an isomorphism.

It is an important and useful fact that this is sufficient:

**9.1.F. ESSENTIAL (HARD) EXERCISE: A USEFUL CRITERION FOR WHEN IDEALS IN AFFINE OPEN SETS DEFINE A CLOSED SUBSCHEME.** Suppose \( Y \) is a scheme, and for each affine open subset \( \text{Spec} B \) of \( Y \), \( I_B \subseteq B \) is an ideal. Suppose further that for each affine open subset \( \text{Spec} B \hookrightarrow Y \) and each \( f \in B \), restriction of functions from \( B \rightarrow B_f \) induces an isomorphism \( I_{(B_f)} \rightarrow (I_B)_f \). Show that this data arises from a (unique) closed subscheme \( X \hookrightarrow Y \) by the above construction. In other words, the closed immersions \( \text{Spec} B/1 \hookrightarrow \text{Spec} B \) glue together in a well-defined manner to obtain a closed immersion \( X \hookrightarrow Y \).

This is a hard exercise, so as a hint, here are three different ways of proceeding: some combination of them may work for you. **Approach 1.** For each affine open \( \text{Spec} B \), we have a closed subscheme \( \text{Spec} B/I \hookrightarrow \text{Spec} B \). (i) For any two affine open subschemes \( \text{Spec} A \) and \( \text{Spec} B \), show that the two closed subschemes \( \text{Spec} A/I \hookrightarrow \text{Spec} A \) and \( \text{Spec} B/I \hookrightarrow \text{Spec} B \) restrict to the same closed subscheme of their intersection. (Hint: cover their intersection with open sets simultaneously distinguished in both affine open sets, Proposition 6.3.1.) Thus for example we can glue these two closed subschemes together to get a closed subscheme of \( \text{Spec} A \cup \text{Spec} B \). (ii) Use Exercise 5.3.4 on gluing schemes (or the ideas therein) to glue together the closed immersions in all affine open subschemes simultaneously. You will only need to worry about triple intersections. **Approach 2.** (i) Use the data of the ideals \( I_B \) to define a sheaf of ideals \( \mathcal{I} \hookrightarrow \mathcal{O} \). (ii) For each affine open subscheme \( \text{Spec} B \), show that \( \mathcal{I}(\text{Spec} B) \) is indeed \( I_B \), and \( (\mathcal{O}/\mathcal{I})(\text{Spec} B) \) is indeed \( B/I_B \), so the data of \( \mathcal{I} \) recovers the closed subscheme on each \( \text{Spec} B \) as desired. **Approach 3.** (i) Describe \( X \) first as a subset of \( Y \). (ii) Check that \( X \) is closed. (iii) Define the sheaf of functions \( \mathcal{O}_X \) on this subset, perhaps using compatible stalks. (iv) Check that this resulting ringed space is indeed locally the closed subscheme given by \( \text{Spec} B/I \hookrightarrow \text{Spec} B \).

We will see later (14.5.5) that closed subschemes correspond to quasicoherent sheaves of ideals; the mathematical content of this statement will turn out to be precisely Exercise 9.1.F.

**9.1.G. IMPORTANT EXERCISE.** (a) In analogy with closed subsets, define the notion of a finite union of closed subschemes of \( X \), and an arbitrary (not necessarily finite) intersection of closed subschemes of \( X \).
(b) Describe the scheme-theoretic intersection of \( V(y - x^2) \) and \( V(y) \) in \( \mathbb{A}^2 \). See Figure [S.5] for a picture. (For example, explain informally how this corresponds to two curves meeting at a single point with multiplicity 2 — notice how the 2 is visible in your answer. Alternatively, what is the non-reducedness telling you — both its “size” and its “direction”?) Describe their scheme-theoretic union.

(c) Show that the underlying set of a finite union of closed subschemes is the finite union of the underlying sets, and similarly for arbitrary intersections.

(d) Describe the scheme-theoretic intersection of \((y^2 - x^2)\) and \( y \) in \( \mathbb{A}^2 \). Draw a picture. (Did you expect the intersection to have multiplicity one or multiplicity two?) Hence show that if \( X, Y, \) and \( Z \) are closed subschemes of \( W \), then \((X \cap Z) \cup (Y \cap Z) \neq (X \cup Y) \cap Z \) in general.

9.1.H. **Important Exercise/Definition:** the Vanishing Scheme. (a) Suppose \( Y \) is a scheme, and \( s \in \Gamma(\mathcal{O}_Y, Y) \). Define the closed scheme cut out by \( s \). We call this the **vanishing scheme** \( V(s) \) of \( s \), as it is the scheme theoretical version of our earlier (set-theoretical) version of \( V(s) \). (Hint: on affine open \( \text{Spec } B \), we just take \( \text{Spec } B/(s_B) \), where \( s_B \) is the restriction of \( s \) to \( \text{Spec } B \). Use Exercise 9.1.F to show that this yields a well-defined closed subscheme.) In Exercise 9.1.G(b), you are computing \( V(y - x^2, y) \).

(b) If \( u \) is an invertible function, show that \( V(s) = V(su) \).

(c) If \( S \) is a set of functions, define \( V(S) \).

9.1.2. **Locally principal closed subschemes, and effective Cartier divisors.** (This section is just an excuse to introduce some notation, and is not essential to the current discussion.) A closed subscheme is locally principal if on each open set in a small enough open cover it is cut out by a single equation. Thus each homogeneous polynomial in \( n + 1 \) variables defines a locally principal closed subscheme. (Warning: this is not an affine-local condition, see Exercise 6.4.I.) Also, the example of a projective hypersurface given soon in [9.2.I] shows that a locally principal closed subscheme need not be cut out by a (global) function.) A case that will be important later is when the ideal sheaf is not just locally generated by a function, but is generated by a function that is not a zero-divisor. For reasons that will become clearer later, we call such a closed subscheme an effective Cartier divisor. Warning: We will use this terminology before we explain where it came from!

9.1.I. **Unimportant Exercise.** Suppose \( V(s) = V(s') \subset \text{Spec } A \) is an effective Cartier divisor, with \( s \) and \( s' \) non-zero-divisors in \( A \). Show that \( s \) is a unit times \( s' \).

9.1.J. **Unimportant and Hard Exercise.** In the literature, the usual definition of a closed immersion is a morphism \( f : X \to Y \) such that \( f \) induces a homeomorphism of the underlying topological space of \( X \) onto a closed subset of the topological space of \( Y \), and the induced map \( f^\#: \mathcal{O}_Y \to f_*\mathcal{O}_X \) of sheaves on \( Y \) is surjective. Show that this definition agrees with the one given above. (To show that our definition involving surjectivity on the level of affine open sets implies this definition, you can use the fact that surjectivity of a morphism of sheaves can be checked on a base, Exercise [3.7.E].)

We have now defined the analogue of open subsets and closed subsets in the land of schemes. Their definition is slightly less “symmetric” than in the classical topological setting: the “complement” of a closed subscheme is a unique open
subscheme, but there are many “complementary” closed subschemes to a given open subscheme in general. (We will soon define one that is “best”, that has a reduced structure, \(^{[9.3.8]}\))

### 9.1.3. Locally closed immersions and locally closed subschemes.

Now that we have defined analogues of open and closed subsets, it is natural to define the analogue of locally closed subsets. Recall that locally closed subsets are intersections of open subsets and closed subsets. Hence they are closed subsets of open subsets, or equivalently open subsets of closed subsets. The analog of these equivalences will be a little problematic in the land of schemes.

We say a morphism \( h : X \rightarrow Y \) is a **locally closed immersion** if \( h \) can factored into \( X \xrightarrow{f} Z \xrightarrow{g} Y \) where \( f \) is a closed immersion and \( g \) is an open immersion. If \( X \) is a subset of \( Y \) (and \( h \) on the level of sets is the inclusion), we say \( X \) is a **locally closed subscheme** of \( Y \). (Warning: The term *immersion* is often used instead of *locally closed immersion*, but this is unwise terminology. The differential geometric notion of immersion is closer to the what algebraic geometers call unramified, which we will define in \(^{[22.4.5]}\). The algebro-geometric notion of locally closed immersion is closer to the differential geometric notion of *embedding*.)

For example, \( \text{Spec } k[t, t^{-1}] \rightarrow \text{Spec } k[x, y] \) where \( (x, y) \mapsto (t, 0) \) is a locally closed immersion (see Figure 9.1).

---

**Figure 9.1.** The locally closed immersion \( \text{Spec } k[t, t^{-1}] \rightarrow k[x, y] \)

\( (t \mapsto (t, 0)) = (x, y), \text{ i.e. } (x, y) \mapsto (t, 0) \)

At this point, you could define the intersection of two locally closed immersions in a scheme \( X \) (which is also be a locally closed immersion in \( X \)). But it would be awkward, as you would have to show that your construction is independent of the factorizations of each locally closed immersion into a closed immersion and an open immersion. Instead, we wait until Exercise \(^{[10.2.C]}\) when recognizing the intersection as a fibered product will make this easier.

Clearly an open subscheme \( U \) of a closed subscheme \( V \) of \( X \) can be interpreted as a closed subscheme of an open subscheme: as the topology on \( V \) is induced from the topology on \( X \), the underlying set of \( U \) is the intersection of some open
subset $U'$ on $X$ with $V$. We can take $V' = V \cap U$, and then $V' \to U'$ is a closed immersion, and $U' \to X$ is an open immersion.

It is not clear that a closed subscheme $V'$ of an open subscheme $U'$ can be expressed as an open subscheme $U$ of a closed subscheme $V$. In the category of topological spaces, we would take $V$ as the closure of $V'$, so we are now motivated to define the analogous construction, which will give us an excuse to introduce several related ideas, in the next section. We will then resolve this issue in good cases (e.g. if $X$ is Noetherian) in Exercise 9.3.C.

We formalize our discussion in an exercise.

9.1.K. EXERCISE. Suppose $V \to X$ is a morphism. Consider three conditions:

(i) $V$ is an open subscheme of $X$ intersect a closed subscheme of $X$ (which you will have to define, see Exercise 8.1.A or else see below).

(ii) $V$ is an open subscheme of a closed subscheme of $X$ (i.e. it factors into an open immersion followed by a closed immersion).

(iii) $V$ is a closed subscheme of an open subscheme of $X$, i.e. $V$ is a locally closed immersion.

Show that (i) and (ii) are equivalent, and both imply (iii). (Remark: (iii) does not always imply (i) and (ii), see [Stacks, Tag 01QW].) Hint: It may be helpful to think of the problem as follows. You might hope to think of a locally closed immersion as a fibered diagram

$\array{V \to & Y \\
\downarrow \text{open imm.} & \downarrow \text{closed imm.} \\
K \to & X, \\
\downarrow \text{open imm.} & \downarrow \text{closed imm.}}$

Interpret (i) as the existence of the diagram. Interpret (ii) as this diagram minus the lower left corner. Interpret (iii) as the diagram minus the upper right corner.

9.1.L. EXERCISE. Show that the composition of two locally closed immersions is a locally closed immersion. (Hint: you might use (ii) implies (iii) in the previous exercise.)

9.1.4. Unimportant remark. It may feel odd that in the definition of a locally closed immersions, we had to make a choice (as a composition of a closed followed by an open, rather than vice versa), but this type of issue comes up earlier: a subquotient of a group can be defined as the quotient of a subgroup, or a subgroup of a quotient. Which is the right definition? Or are they the same? (Hint: compositions of two subquotients should certainly be a subquotient, cf. Exercise 9.1.L.)

9.2 Closed immersions of projective schemes, and more projective geometry

9.2.1. Example: Closed immersions of projective space $\mathbb{P}^n_A$. Recall the definition of projective space $\mathbb{P}^n_A$ given in [5.4.D](and the terminology defined there). Any homogeneous polynomial $f$ in $x_0, \ldots, x_n$ defines a closed subscheme. (Thus even if
f doesn’t make sense as a function, its vanishing scheme still makes sense.) On the open set \( U_i \), the closed subscheme is \( V(f(x_0/i, \ldots, x_n/i)) \), which we think of as \( V(f(x_0, \ldots, x_n)/x_i^{\deg f}) \). On the overlap
\[
U_i \cap U_j = \text{Spec } A[x_0/i, \ldots, x_n/i, x_j^{-1}]/(x_i - 1),
\]
these functions on \( U_i \) and \( U_j \) don’t exactly agree, but they agree up to a non-vanishing scalar, and hence cut out the same closed subscheme of \( U_i \cap U_j \) (Exercise 9.1.H(b)):
\[
f(x_0/i, \ldots, f_n/i) = x_j^{\deg f} f(x_0/j, \ldots, x_n/j).
\]
Similarly, a collection of homogeneous polynomials in \( A[x_0, \ldots, x_n] \) cuts out a closed subscheme of \( \mathbb{P}^n_A \).

9.2.2. Definition. A closed subscheme cut out by a single (homogeneous) equation is called a \textit{hypersurface} in \( \mathbb{P}^n_A \). A hypersurface is locally principal. Notice that a hypersurface is not in general cut out by a single global function on \( \mathbb{P}^n_A \). For example, if \( A = k \), there are no nonconstant global functions (Exercise 5.4.E). The degree of a hypersurface is the degree of the polynomial. (Implicit in this is that this notion can be determined from the subscheme itself; we haven’t yet checked this.) A hypersurface of degree 1 (resp. degree 2, 3, \ldots) is called a \textit{hyperplane} (resp. \textit{quadric}, \textit{cubic}, \textit{quartic}, \textit{quintic}, \textit{sextic}, \textit{septic}, \textit{octic}, \ldots \textit{hypersurface}). If \( n = 2 \), a degree 1 hypersurface is called a \textit{line}, and a degree 2 hypersurface is called a \textit{conic curve}, or a \textit{conic} for short. If \( n = 3 \), a hypersurface is called a \textit{surface}. (In Chapter 12, we will justify the terms \textit{curve} and \textit{surface}.)

9.2.A. Exercise. (a) Show that \( wz = xy, x^2 = vy, y^2 = xz \) describes an irreducible subscheme in \( \mathbb{P}^3_k \). In fact it is a curve, a notion we will define once we know what dimension is. This curve is called the \textit{twisted cubic}. (The twisted cubic is a good non-trivial example of many things, so you should make friends with it as soon as possible. It implicitly appeared earlier in Exercise 4.6.H)
(b) Show that the twisted cubic is isomorphic to \( \mathbb{P}^1_k \).

9.2.B. Exercise (a special case of Bézout’s theorem). Suppose \( X \subset \mathbb{P}^n \) is a degree \( d \) hypersurface cut out by \( f = 0 \), and \( L \) is a line not contained in \( H \). A very special case of Bézout’s theorem (Exercise 20.5.L) implies that \( X \) and \( L \) meet with multiplicity \( d \), “counted correctly”. Make sense of this, by restricting the degree \( d \) form \( f \) to the line \( H \), and using the fact that a degree \( d \) polynomial in \( k[x] \) has \( d \) roots, counted properly.

We now extend this discussion to projective schemes in general.

9.2.C. Exercise. Suppose that \( S_\bullet \twoheadrightarrow R_\bullet \) is a surjection of finitely-generated graded rings. Show that the induced morphism \( \text{Proj } R_\bullet \twoheadrightarrow \text{Proj } S_\bullet \) (Exercise 7.4.A) is a closed immersion.

9.2.D. Exercise. Suppose \( X \hookrightarrow \text{Proj } S_\bullet \) is a closed immersion in a projective \( A \)-scheme. Show that \( X \) is projective by describing it as \( \text{Proj } S_\bullet/I \), where \( I \) is a homogeneous prime ideal, of “projective functions” vanishing on \( X \).
9.2.E. Exercise. Show that an injective linear map of \( k \)-vector spaces \( V \hookrightarrow W \) induces a closed immersion \( \mathbb{P}V \hookrightarrow \mathbb{P}W \). (This is another justification for the definition of \( \mathbb{P}V \) in Example 9.5.8 in terms of the dual of \( V \).)

This closed subscheme is called a **linear space**. Once we know about dimension, we will call this a linear space of dimension \( \dim V - 1 = \dim \mathbb{P}V \). A linear space of dimension 1 (resp. 2, \( n \), \( \dim \mathbb{P}W - 1 \)) is called a **line** (resp. **plane**, \( n \)-**plane**), **hyperplane**. (If the linear map in the previous exercise is not injective, then the hypothesis (7.4.0.1) of Exercise 7.4.A fails.)

9.2.F. Exercise. Show that the map of graded rings \( k[w, x, y, z] \rightarrow k[s, t] \) given by \( w \mapsto s^3, x \mapsto s^2t, y \mapsto st^2, z \mapsto t^3 \) induces a closed immersion \( \mathbb{P}^1_k \hookrightarrow \mathbb{P}^3_k \), which yields an isomorphism of \( \mathbb{P}^1_k \) with the twisted cubic (defined in Exercise 9.2.A — in fact, this will solve Exercise 9.2.A).

9.2.G. Exercise. Suppose \( S_* \) is a finitely generated graded ring generated in degree 1. Show that \( S_1 \) is a finitely-generated \( S_* \)-module, and the irrelevant ideal \( S_+ \) is generated in degree 1.

9.2.H. Exercise. Show that if \( S_* \) is generated by \( S_1 \) (as an \( A \)-algebra) by \( n + 1 \) elements \( x_0, \ldots, x_n \), then \( \text{Proj} S_* \) may be described as a closed subscheme of \( \mathbb{P}^n_A \) as follows. Consider \( A^{n+1} \) as a free module with generators \( t_0, \ldots, t_n \) associated to \( x_0, \ldots, x_n \). The surjection of

\[
\text{Sym}^* A^{n+1} = A[t_0, t_1, \ldots, t_n] \twoheadrightarrow S_*
\]

implies \( S_* = A[t_0, t_1, \ldots, t_n]/I \), where \( I \) is a homogeneous ideal. (In particular, by Exercise 7.4.C, \( \text{Proj} S_* \) can always be interpreted as a closed subscheme of some \( \mathbb{P}^n_A \).)

This is analogous to the fact that if \( R \) is a finitely-generated \( A \)-algebra, then choosing \( n \) generators of \( R \) as an algebra is the same as describing \( \text{Spec} R \) as a closed subscheme of \( A^n_A \). In the affine case this is “choosing coordinates”; in the projective case this is “choosing projective coordinates”.

For example, \( \text{Proj} k[x, y, z]/(z^2 - x^2 - y^2) \) is a closed subscheme of \( \mathbb{P}^2_k \). (A picture is shown in Figure 9.3.)

Recall (Exercise 5.4.F) that if \( k \) is algebraically closed, then we can interpret the closed points of \( \mathbb{P}^n \) as the lines through the origin in \( (n + 1) \)-space. The following exercise states this more generally.

9.2.I. Exercise. Suppose \( S_* \) is a finitely-generated graded ring over an algebraically closed field \( k \), generated in degree 1 by \( x_0, \ldots, x_n \), inducing closed immersions \( \text{Proj} S_* \hookrightarrow \mathbb{P}^n \) and \( \text{Spec} S_* \hookrightarrow \mathbb{A}^n \). Give a bijection between the closed points of \( \text{Proj} S_* \) and the “lines through the origin” in \( \text{Spec} S_* \subset \mathbb{A}^n \).

9.2.4. A second proof that finite morphisms are closed. This interpretation of \( \text{Proj} S_* \) as a closed subscheme of projective space (when it is generated in degree 1) yields the
following second proof of the fact (shown in Exercise 8.3.N) that finite morphisms are closed. Suppose \( \phi : X \to Y \) is a finite morphism. The question is local on the target, so it suffices to consider the affine case \( Y = \text{Spec} B \). It suffices to show that \( \phi(X) \) is closed. Then by Exercise 8.3.L, \( X \) is a projective \( B \)-scheme, and hence by the Fundamental Theorem of Elimination Theory 8.4.5, its image is closed.

9.2.5. The Veronese embedding.

Suppose \( S_* = k[x, y] \), so \( \text{Proj} S_* = \mathbb{P}^1_k \). Then \( S_{2*} = k[x^2, xy, y^2] \subset k[x, y] \) (see §7.4.2 on the Veronese subring). We identify this subring as follows.

9.2.J. Exercise. Let \( u = x^2, v = xy, w = y^2 \). Show that \( S_{2*} = k[u, v, w]/(uw - v^2) \).

We have a graded ring generated by three elements in degree 1. Thus we think of it as sitting “in” \( \mathbb{P}^2 \), via the construction of §9.2.H. This can be interpreted as “\( \mathbb{P}^1 \) as a conic in \( \mathbb{P}^2 \”).

9.2.6. Thus if \( k \) is algebraically closed of characteristic not 2, using the fact that we can diagonalize quadrics (Exercise 6.4.J), the conics in \( \mathbb{P}^2 \), up to change of coordinates, come in only a few flavors: sums of 3 squares (e.g. our conic of the previous exercise), sums of 2 squares (e.g. \( y^2 - x^2 = 0 \), the union of 2 lines), a single square (e.g. \( x^2 = 0 \), which looks set-theoretically like a line, and is non-reduced), and 0 (perhaps not a conic at all). Thus we have proved: any plane conic (over an algebraically closed field of characteristic not 2) that can be written as the sum of three squares is isomorphic to \( \mathbb{P}^1 \).

We now soup up this example.

9.2.K. Exercise. Show that \( \text{Proj} S_{d*} \) is given by the equations that

\[
\begin{pmatrix}
y_0 & y_1 & \cdots & y_{d-1} \\
y_1 & y_2 & \cdots & y_d
\end{pmatrix}
\]

is rank 1 (i.e. that all the \( 2 \times 2 \) minors vanish). This is called the degree \( d \) rational normal curve “in” \( \mathbb{P}^d \). You did the twisted cubic case \( d = 3 \) in Exercises 9.2.A and 9.2.I.

9.2.7. Remark. More generally, if \( S_* = k[x_0, \ldots, x_n] \), then \( \text{Proj} S_{d*} \subset \mathbb{P}^{N-1} \) (where \( N \) is the number of degree \( d \) polynomials in \( x_0, \ldots, x_n \)) is called the \( d \)-uple embedding or \( d \)-uple Veronese embedding. The reason for the word “embedding” is historical; we really mean closed immersion. (Combining Exercise 7.4.H with Exercise 9.2.H shows that \( \text{Proj} S_* \to \mathbb{P}^{n-1} \) is a closed immersion.)

9.2.L. Combinatorial exercise. Show that \( N = \binom{n+d}{d} \).

9.2.M. Unimportant exercise. Find five linearly independent quadric equations vanishing on the Veronese surface \( \text{Proj} S_{2*} \) where \( S_* = k[x_0, x_1, x_2] \), which sits naturally in \( \mathbb{P}^5 \). (You needn’t show that these equations generate all the equations cutting out the Veronese surface, although this is in fact true.)

9.2.8. Rulings on the quadric surface. We return to rulings on the quadric surface, which first appeared in the optional section §5.4.11.
9.2.N. USEFUL GEOMETRIC EXERCISE: THE RULINGS ON THE QUADRIC SURFACE

$wz = xy$. This exercise is about the lines on the quadric surface $wz - xy = 0$ in $\mathbb{P}^3_k$.

This construction arises all over the place in nature.

(a) Suppose $a_0$ and $b_0$ are elements of $k$, not both zero. Make sense of the statement: as $\begin{bmatrix} c, d \end{bmatrix}$ varies in $\mathbb{P}^1_k$, $\begin{bmatrix} a_0c; b_0c; a_0d; b_0d \end{bmatrix}$ is a line in the quadric surface. (This describes “a family of lines parametrized by $\mathbb{P}^1_k$”, although we can’t yet make this precise.) Find another family of lines. These are the two rulings of the quadric surface.

(b) Show there are no other lines. (There are many ways of proceeding. At risk of predisposing you to one approach, here is a germ of an idea. Suppose $L$ is a line on the quadric surface, and $\begin{bmatrix} 1; x; y; z \end{bmatrix}$ and $\begin{bmatrix} 1; x'; y'; z' \end{bmatrix}$ are distinct points on it. Because they are both on the quadric, $z = xy$ and $z' = x'y'$. Because all of $L$ is on the quadric, $(1 + t)(z + tz') - (x + tx')(y + ty') = 0$ for all $t$. After some algebraic manipulation, this translates into $(x - x')(y - y') = 0$. How can this be made watertight? Another possible approach uses Bézout’s theorem, in the form of Exercise 9.2.B.)

![Figure 9.2](image)

FIGURE 9.2. The two rulings on the quadric surface $V(wz - xy) \subset \mathbb{P}^3$. One ruling contains the line $V(w, x)$ and the other contains the line $V(w, y)$.

Hence by Exercise 6.4.J if we are working over an algebraically closed field of characteristic not 2, we have shown that all rank 4 quadric surfaces have two rulings of lines.

9.2.9. WEIGHTED PROJECTIVE SPACE. If we put a non-standard weighting on the variables of $k[x_1, \ldots, x_n]$ — say we give $x_i$ degree $d_i$ — then $\text{Proj} k[x_1, \ldots, x_n]$ is called weighted projective space $\mathbb{P}(d_1, d_2, \ldots, d_n)$.

9.2.O. EXERCISE. Show that $\mathbb{P}(m, n)$ is isomorphic to $\mathbb{P}^1$. Show that $\mathbb{P}(1, 1, 2) \cong \text{Proj} k[u, v, w, z]/(uw - v^2)$. Hint: do this by looking at the even-graded parts of $k[x_0, x_1, x_2]$, cf. Exercise 7.4.D (This is a projective cone over a conic curve. Over
an algebraically closed field of characteristic not 2, it is isomorphic to the traditional cone \( x^2 + y^2 = z^2 \) in \( \mathbb{P}^3 \), Figure 9.3).

9.2.10. Affine and projective cones.

If \( S_\bullet \) is a finitely-generated graded ring, then the affine cone of \( \text{Proj} S_\bullet \) is \( \text{Spec} S_\bullet \). Note that this construction depends on \( S_\bullet \), not just of \( \text{Proj} S_\bullet \). As motivation, consider the graded ring \( S_\bullet = \mathbb{C}[x, y, z]/(z^2 - x^2 - y^2) \). Figure 9.3 is a sketch of \( \text{Spec} S_\bullet \). (Here we draw the “real picture” of \( z^2 = x^2 + y^2 \) in \( \mathbb{R}^3 \).) It is a cone in the traditional sense; the origin \((0, 0, 0)\) is the “cone point”.

\[ \text{Figure 9.3. The cone } \text{Spec} k[x, y, z]/(z^2 - x^2 - y^2). \]

This gives a useful way of picturing \( \text{Proj} \) (even over arbitrary rings, not just \( \mathbb{C} \)). Intuitively, you could imagine that if you discarded the origin, you would get something that would project onto \( \text{Proj} S_\bullet \). The following exercise makes that precise.

9.2.P. Exercise (cf. Exercise 7.3.E). If \( \text{Proj} S_\bullet \) is a projective scheme over a field \( k \), describe a natural morphism \( \text{Spec} S_\bullet \setminus \{0\} \to \text{Proj} S_\bullet \).

This has the following generalization to \( A \)-schemes, which you might find geometrically reasonable. This again motivates the terminology “irrelevant”.

9.2.Q. Exercise. If \( S_\bullet \) is a graded ring, describe a natural morphism \( \text{Spec} S_\bullet \setminus V(S_+) \to \text{Proj} S_\bullet \).

In fact, it can be made precise that \( \text{Proj} S_\bullet \) is quotient (by the multiplicative group of scalars) of the affine cone minus the origin.

The projective cone of \( \text{Proj} S_\bullet \) is \( \text{Proj} S_\bullet[T] \), where \( T \) is a new variable of degree 1. For example, the cone corresponding to the conic \( \text{Proj} k[x, y, z]/(z^2 - x^2 - y^2) \) is \( \text{Proj} k[x, y, z, T]/(z^2 - x^2 - y^2) \).

9.2.R. Exercise (cf. 5.5.1). Show that the projective cone of \( \text{Proj} S_\bullet[T] \) has a closed subscheme isomorphic to \( \text{Proj} S_\bullet \) (corresponding to \( T = 0 \)), whose complement (the distinguished open set \( D(T) \)) is isomorphic to the affine cone \( \text{Spec} S_\bullet \).
You can also check that $\text{Proj} S_\bullet$ is a locally principal closed subscheme of the projective cone $\text{Proj} S_\bullet[T]$, and is also locally not a zero-divisor (an **effective Cartier divisor**, §9.1.2).

This construction can be usefully pictured as the affine cone union some points “at infinity”, and the points at infinity form the Proj. The reader may wish to ponder Figure 9.3 and try to visualize the conic curve “at infinity”.

We have thus completely described the algebraic analogue of the classical picture of 5.5.1.

### 9.3 “Smallest closed subschemes such that ...”: scheme-theoretic image, scheme-theoretic closure, induced reduced subscheme, and the reduction of a scheme

We now define a series of notions that are all of the form “the smallest closed subscheme such that something or other is true”. One example will be the notion of scheme-theoretic closure of a locally closed immersion, which will allow us to interpret locally closed immersions in three equivalent ways (open subscheme intersect closed subscheme; open subscheme of closed subscheme; and closed subscheme of open subscheme).

#### 9.3.1. Scheme-theoretic image.

We start with the notion of scheme-theoretic image. Set-theoretic images are badly behaved in general (8.4.1), and even with reasonable hypotheses such as those in Chevalley’s theorem 8.4.2, things can be confusing. For example, there is no reasonable way to impose a scheme structure on the image of $\mathbb{A}^2_k \to \mathbb{A}^2_k$ given by $(x, y) \mapsto (x, xy)$. It will be useful (e.g. Exercise 9.3.C) to define a notion of a closed subscheme of the target that “best approximates” the image. This will incorporate the notion that the image of something with non-reduced structure (“fuzz”) can also have non-reduced structure. As usual, we will need to impose reasonable hypotheses to make this notion behave well (see Theorem 9.3.4 and Corollary 9.3.5).

#### 9.3.2. Definition.

Suppose $i : Z \hookrightarrow Y$ is a closed subscheme, giving an exact sequence $0 \to I_{Z/Y} \to O_Y \to i_* O_Z \to 0$. We say that the image of $f : X \to Y$ lies in $Z$ if the composition $I_{Z/Y} \to O_Y \to f_* O_X$ is zero. Informally, locally functions vanishing on $Z$ pull back to the zero function on $X$. If the image of $f$ lies in some subschemes $Z_i$ (as $i$ runs over some index set), it clearly lies in their intersection (cf. Exercise 9.1.C(a) on intersections of closed subschemes). We then define the **scheme-theoretic image** of $f$, a closed subscheme of $Y$, as the “smallest closed subscheme containing the image”, i.e. the intersection of all closed subschemes containing the image.

**Example 1.** Consider $\text{Spec} k[\epsilon]/\epsilon^2 \to \text{Spec} k[x] = \mathbb{A}^1_k$ given by $x \mapsto \epsilon$. Then the scheme-theoretic image is given by $k[x]/x^2$ (the polynomials pulling back to $0$ are precisely multiples of $x^2$). Thus the image of the fuzzy point still has some fuzz.
Example 2. Consider \( f : \text{Spec} k[e]/e^2 \to \text{Spec} k[x] = \mathbb{A}^1_k \) given by \( x \mapsto 0 \). Then the scheme-theoretic image is given by \( k[x]/x \): the image is reduced. In this picture, the fuzz is “collapsed” by \( f \).

Example 3. Consider \( f : \text{Spec} k[t, t^{-1}] = \mathbb{A}^1 - \{0\} \to \mathbb{A}^1 = \text{Spec} k[u] \) given by \( u \mapsto t \). Any function \( g(u) \) which pulls back to \( 0 \) as a function of \( t \) must be the zero-function. Thus the scheme-theoretic image is everything. The set-theoretic image, on the other hand, is the distinguished open set \( \mathbb{A}^1 - \{0\} \). Thus in not-too-pathological cases, the underlying set of the scheme-theoretic image is not the set-theoretic image. But the situation isn’t terrible: the underlying set of the scheme-theoretic image must be closed, and indeed it is the closure of the set-theoretic image. We might imagine that in reasonable cases this will be true, and in even nicer cases, the underlying set of the scheme-theoretic image will be set-theoretic image. We will later see that this is indeed the case (Exercise 9.3.6).

But sadly pathologies can sometimes happen.

Example 4. Let \( X = \bigsqcup \text{Spec} k[\epsilon_n]/(\epsilon_n^n) \) and \( Y = \text{Spec} k[x] \), and define \( X \to Y \) by \( x \mapsto \epsilon_n \) on the \( n \)th component of \( X \). Then if a function \( g(x) \) on \( Y \) pulls back to \( 0 \) on \( X \), then its Taylor expansion is \( 0 \) to order \( n \) (by examining the pullback to the \( n \)th component of \( X \)), so \( g(x) \) must be \( 0 \). Thus the scheme-theoretic image is \( V(0) \) on \( Y \), i.e. \( Y \) itself, while the set-theoretic image is easily seen to be just the origin.

**9.3.3. Criteria for computing scheme-theoretic images affine-locally.** Example 4 clearly is weird though, and we can show that in “reasonable circumstances” such pathology doesn’t occur. It would be great to compute the scheme-theoretic image affine-locally. On the affine open set \( \text{Spec} B \subset Y \), define the ideal \( I_B \subset B \) of functions which pull back to \( 0 \) on \( X \). Formally, \( I_B := \ker(B \to \Gamma(\text{Spec} B, f_*(\mathcal{O}_X))) \). Then if for each such \( B \), and each \( g \in B \), \( I_B \otimes_B B_g \to I_{B_g} \) is an isomorphism, then we will have defined the scheme-theoretic image as a closed subscheme (see Exercise 9.1.F). Clearly each function on \( \text{Spec} B \) that vanishes when pulled back to \( f^{-1}(\text{Spec} B) \) also vanishes when restricted to \( D(g) \) and then pulled back to \( f^{-1}(D(g)) \). So the question is: given a function \( r/g^n \) on \( D(g) \) that pulls back to \( f^{-1}D(g) \), is it true that for some \( m \), \( rg^m = 0 \) when pulled back to \( f^{-1}(\text{Spec} B) \)? Here are three cases where the answer is “yes”.

(i) (I would like to add a picture here, but I can’t think of one that would enlighten more people than it would confuse. So you should try to draw one that suits you.) In a nutshell, for each affine in the source, there is an \( m \) which works. There is one that works for all affines in a cover if (i) if \( m = 1 \) always works, or (ii) or (iii) if there are only a finite number of affines in the cover.

(i) The answer is yes if \( f^{-1}(\text{Spec} B) \) is reduced: we simply take \( m = 1 \) (as \( r \) vanishes on \( \text{Spec} B_g \) and \( g \) vanishes on \( V(g) \), so \( rg \) vanishes on \( \text{Spec} B = \text{Spec} B_g \cup V(g) \)).

(ii) The answer is also yes if \( f^{-1}(\text{Spec} B) \) is affine, say \( \text{Spec} A: \) if \( r' = f^#r \) and \( g' = f^#g \) in \( A \), then if \( r' = 0 \) on \( D(g') \), then there is an \( m \) such that \( r'(g')^m = 0 \) (as the statement \( r' = 0 \) in \( D(g') \) means precisely this fact — the functions on \( D(g') \) are \( A_{g'} \)).

(iii) More generally, the answer is yes if \( f^{-1}(\text{Spec} B) \) is quasicompact: cover \( f^{-1}(\text{Spec} B) \) with finitely many affine open sets. For each one there will be some \( m_i \) so that \( rg^{m_i} = 0 \) when pulled back to this open set. Then let \( m = \max(m_i) \). (We see again that quasicompactness is our friend!)
In conclusion, we have proved the following (subtle) theorem.

9.3.4. Theorem. — Suppose $f : X \to Y$ is a morphism of schemes. If $X$ is reduced or $f$ is quasicompact, then the scheme-theoretic image of $f$ may be computed affine-locally: on $\text{Spec } A$, it is cut out by the functions that pull back to $0$.

9.3.5. Corollary. — Under the hypotheses of the Theorem 9.3.4, the closure of the set-theoretic image of $f$ is the underlying set of the scheme-theoretic image.

(Example 4 above shows that we cannot excise these hypotheses.)

9.3.6. In particular, if the set-theoretic image is closed (e.g. if $f$ is finite or projective), the set-theoretic image is the underlying set of the scheme-theoretic image, as promised in Example 3 above.

Proof. The set-theoretic image is in the underlying set of the scheme-theoretic image. (Check this!) The underlying set of the scheme-theoretic image is closed, so the closure of the set-theoretic image is contained in underlying set of the scheme-theoretic image. On the other hand, if $U$ is the complement of the closure of the set-theoretic image, $f^{-1}(U) = \emptyset$. As under these hypotheses, the scheme theoretic image can be computed locally, the scheme-theoretic image is the empty set on $U$. □

We conclude with a few stray remarks.

9.3.A. Easy Exercise. If $X$ is reduced, show that the scheme-theoretic image of $f : X \to Y$ is also reduced.

More generally, you might expect there to be no unnecessary non-reduced structure on the image not forced by non-reduced structure on the source. We make this precise in the locally Noetherian case, when we can talk about associated points.

9.3.B. * Unimportant Exercise. If $f : X \to Y$ is a quasicompact morphism of locally Noetherian schemes, show that the associated points of the image subscheme are a subset of the image of the associated points of $X$. (The example of $\prod_{a \in \mathbb{C}} \text{Spec } \mathbb{C}[t]/(t - a) \to \text{Spec } \mathbb{C}[t]$ shows what can go wrong if you give up quasicompactness — note that reducedness of the source doesn’t help.) Hint: reduce to the case where $X$ and $Y$ are affine. (Can you develop your geometric intuition so that this is geometrically plausible?)

9.3.7. Scheme-theoretic closure of a locally closed subscheme.

We define the scheme-theoretic closure of a locally closed immersion $f : X \to Y$ as the scheme-theoretic image of $X$.

9.3.C. Exercise. If $V \to X$ is quasicompact (e.g. if $V$ is Noetherian, Exercise 8.3.B(a)), or if $V$ is reduced, show that (iii) implies (i) and (ii) Exercise 9.1.K. Thus in this fortunate situation, a locally closed immersion can be thought of in three different ways, whichever is convenient.
9.3.D. **UNIMPORTANT EXERCISE, USEFUL FOR INTUITION.** If $f : X \to Y$ is a locally closed immersion into a locally Noetherian scheme (so $X$ is also locally Noetherian), then the associated points of the scheme-theoretic closure are (naturally in bijection with) the associated points of $X$. (Hint: Exercise 9.3.B.) Informally, we get no non-reduced structure on the scheme-theoretic closure not “forced by” that on $X$.

9.3.8. The (reduced) subscheme structure on a closed subset.

Suppose $X^{\text{set}}$ is a closed subset of a scheme $Y$. Then we can define a canonical scheme structure $X$ on $X^{\text{set}}$ that is reduced. We could describe it as being cut out by those functions whose values are zero at all the points of $X^{\text{set}}$. On the affine open set $\text{Spec } B$ of $Y$, if the set $X^{\text{set}}$ corresponds to the radical ideal $I = I(X^{\text{set}})$ (recall the $I(\cdot)$ function from §4.7), the scheme $X$ corresponds to $\text{Spec } B/I$. You can quickly check that this behaves well with respect to any distinguished inclusion $\text{Spec } B_f \hookrightarrow \text{Spec } B$. We could also consider this construction as an example of a scheme-theoretic image in the following crazy way: let $W$ be the scheme that is a disjoint union of all the points of $X^{\text{set}}$, where the point corresponding to $p$ in $X^{\text{set}}$ is $\text{Spec}$ of the residue field of $\mathcal{O}_{Y,p}$. Let $f : W \to Y$ be the “canonical” map sending “$p$ to $p$”, and giving an isomorphism on residue fields. Then the scheme structure on $X$ is the scheme-theoretic image of $f$. A third definition: it is the smallest closed subscheme whose underlying set contains $X^{\text{set}}$.

This construction is called the (induced) **reduced subscheme structure** on the closed subset $X^{\text{set}}$. (Vague exercise: Make a definition of the reduced subscheme structure precise and rigorous to your satisfaction.)

9.3.E. **EXERCISE.** Show that the underlying set of the induced reduced subscheme $X \to Y$ is indeed the closed subset $X^{\text{set}}$. Show that $X$ is reduced.

9.3.9. Reduced version of a scheme.

In the main interesting case where $X^{\text{set}}$ is all of $Y$, we obtain a **reduced closed subscheme** $Y^{\text{red}} \to Y$, called the **reduction** of $Y$. On the affine open subset $\text{Spec } B \hookrightarrow Y$, $Y^{\text{red}} \hookrightarrow Y$ corresponds to the nilradical $\mathfrak{N}(B)$ of $B$. The reduction of a scheme is the “reduced version” of the scheme, and informally corresponds to “shearing off the fuzz”.

An alternative equivalent definition: on the affine open subset $\text{Spec } B \hookrightarrow Y$, the reduction of $Y$ corresponds to the ideal $\mathfrak{N}(B) \subset Y$. As for any $f \in B$, $\mathfrak{N}(B)_f = \mathfrak{N}(B_f)$, by Exercise 9.1.F this defines a closed subscheme.

9.3.F. **EXERCISE (USEFUL FOR VISUALIZATION).** Show that if $Y$ is a locally Noetherian scheme, the “reduced locus” of $Y$ (the points of $Y$ where $Y^{\text{red}} \to Y$ induces an isomorphism of stalks of the structure sheaves) is an open subset of $Y$. (Hint: if $Y$ is affine, show that it is the complement of the closure of the embedded associated points.)
CHAPTER 10

Fibered products of schemes

10.1 They exist

Before we get to products, we note that coproducts exist in the category of schemes: just as with the category of sets (Exercise 2.3.S), coproduct is disjoint union. The next exercise makes this precise (and directly extends to coproducts of an infinite number of schemes).

10.1.A. Easy Exercise. Suppose $X$ and $Y$ are schemes. Let $X \coprod Y$ be the scheme whose underlying topological space is the disjoint union of the topological spaces of $X$ and $Y$, and with structure sheaf on (the part corresponding to) $X$ given by $O_X$, and similarly for $Y$. Show that $X \coprod Y$ is the coproduct of $X$ and $Y$ (justifying the use of the symbol $\coprod$).

We will now construct the fibered product in the category of schemes.

10.1.1. Theorem: Fibered products exist. — Suppose $f : X \to Z$ and $g : Y \to Z$ are morphisms of schemes. Then the fibered product

$$X \times_Z Y \xrightarrow{f'} Y$$

$$\downarrow g' \quad \downarrow g$$

$$X \xrightarrow{f} Z$$

exists in the category of schemes.

Note: if $A$ is a ring, people often write $\times_A$ for $\times_{\text{Spec } A}$.

10.1.2. Warning: products of schemes aren’t products of sets. Before showing existence, here is a warning: the product of schemes isn’t a product of sets (and more generally for fibered products). We have made a big deal about schemes being sets, endowed with a topology, upon which we have a structure sheaf. So you might think that we will construct the product in this order. But we won’t, because products behave oddly on the level of sets. You may have checked (Exercise 7.6.C(a)) that the product of two affine lines over your favorite algebraically closed field $\bar{k}$ is the affine plane: $\mathbb{A}^1_{\bar{k}} \times_{\bar{k}} \mathbb{A}^1_{\bar{k}} \cong \mathbb{A}^2_{\bar{k}}$. But the underlying set of the latter is not the underlying set of the former — we get additional points, corresponding to curves in $\mathbb{A}^2$ that are not lines parallel to the axes!
10.1.3. On the other hand, $S$-valued points (where $S$ is a scheme, Definition 7.3.6) do behave well under (fibered) products. This is just the definition of fibered product: an $S$-valued point of a scheme $X$ is defined as $\text{Hom}(S,X)$, and the fibered product is defined by

\[(10.1.3.1) \quad \text{Hom}(S,X \times Z Y) = \text{Hom}(S,X) \times_{\text{Hom}(S,Z)} \text{Hom}(S,Y).\]

This is one justification for making the definition of $S$-valued point. For this reason, those classical people preferring to think only about varieties over an algebraically closed field $\mathbb{K}$ (or more generally, finite-type schemes over $\mathbb{K}$), and preferring to understand them through their closed points — or equivalently, the $\mathbb{K}$-valued points, by the Nullstellensatz (Exercise 6.3.E) — needn’t worry: the closed points of the product of two finite type $\mathbb{K}$-schemes over $\mathbb{K}$ are (naturally identified with) the product of the closed points of the factors. This will follow from the fact that the product is also finite type over $\mathbb{K}$, which we verify in Exercise 10.2.D. This is one of the reasons that varieties over algebraically closed fields can be easier to work with. But over a nonalgebraically closed field, things become even more interesting; Example 10.2.1 is a first glimpse.

(Fancy remark: You may feel that (i) “products of topological spaces are products on the underlying sets” is natural, while (ii) “products of schemes are not necessarily are products on the underlying sets” is weird. But really (i) is the lucky consequence of the fact that the underlying set of a topological space can be interpreted as set of $p$-valued points, where $p$ is a point, so it is best seen as a consequence of paragraph 10.1.3, which is the “more correct” — i.e. more general — fact.)

10.1.4. Philosophy behind the proof of Theorem 10.1.1. The proof of Theorem 10.1.1 can be confusing. The following comments may help a little.

We already basically know existence of fibered products in two cases: the case where $X$, $Y$, and $Z$ is affine (stated explicitly below), and the case where $Y \to Z$ is an open immersion (Exercise 8.1.A).

10.1.B. EXERCISE. Use Exercise 7.3.F (that $\text{Hom}_{\text{Sch}}(W, \text{Spec } A) = \text{Hom}_{\text{Rings}}(A, \Gamma(W, O_W))$) to show that given ring maps $C \to B$ and $C \to A$, $\text{Spec}(A \otimes_C B) \cong \text{Spec } A \times_{\text{Spec } C} \text{Spec } B$.

(Interpret tensor product as the “cofibered product” in the category of rings.) Hence the fibered product of affine schemes exists (in the category of schemes). (This generalizes the fact that the product of affine lines exist, Exercise 7.6.C(a).)

The main theme of the proof of Theorem 10.1.1 is that because schemes are built by gluing affine schemes along open subsets, these two special cases will be all that we need. The argument will repeatedly use the same ideas — roughly, that schemes glue (Exercise 5.4.A), and that morphisms of schemes glue (Exercise 7.3.A). This is a sign that something more structural is going on; §10.1.5 describes this for experts.

Proof of Theorem 10.1.1. The key idea is this: we cut everything up into affine open sets, do fibered products there, and show that everything glues nicely. The conceptually difficult part of the proof comes from the gluing, and the realization that we
have to check almost nothing. We divide the proof up into a number of bite-sized pieces.

**Step 1: fibered products of affine with almost-affine over affine.** We begin by combining the affine case with the open immersion case as follows. Suppose $X$ and $Z$ are affine, and $Y \to Z$ factors as $Y \xrightarrow{i} Y' \xrightarrow{g} Z$ where $i$ is an open immersion and $Y'$ is affine. Then $X \times_Z Y$ exists. This is because if the two small squares of

\[
\begin{array}{ccc}
W & \xrightarrow{f} & Y \\
\downarrow & & \downarrow \\
W' & \xrightarrow{g} & Y' \\
\downarrow & & \downarrow \\
X & \xrightarrow{h} & Z
\end{array}
\]

are fibered diagrams, then the “outside rectangle” is also a fibered diagram. (This was Exercise 2.3.P, although you should be able to see this on the spot.) It will be important to remember that “open immersions” are “preserved by fibered product”: the fact that $Y \to Y'$ is an open immersion implies that $W \to W'$ is an open immersion.

**Key Step 2: fibered product of affine with arbitrary over affine exists.** We now come to the key part of the argument: if $X$ and $Z$ are affine, and $Y$ is arbitrary. This is confusing when you first see it, so we first deal with a special case, when $Y$ is the union of two affine open sets $Y_1 \cup Y_2$. Let $Y_{12} = Y_1 \cap Y_2$.

Now for $i = 1, 2$, $X \times_Z Y_i$ exists by the affine case, Exercise 10.1.B. Call this $W_i$. Also, $X \times_Z Y_{12}$ exists by Step 1 (call it $W_{12}$), and comes with open immersions into $W_1$ and $W_2$ (by construction of fibered products with open immersion). Thus we can glue $W_1$ to $W_2$ along $W_{12}$; call this resulting scheme $W$.

We check that this is the fibered product by verifying that it satisfies the universal property. Suppose we have maps $f'' : V \to X$, $g'' : V \to Y$ that compose (with $f$ and $g$ respectively) to the same map $V \to Z$. We need to construct a unique map $h : V \to W$, so that $f' \circ h = g''$ and $g' \circ h = f''$.

(10.1.4.1)

\[
\begin{array}{ccc}
V & \xrightarrow{g''} & Y \\
\downarrow & & \downarrow \\
W & \xrightarrow{f} & X \\
\downarrow & & \downarrow \\
Z
\end{array}
\]

For $i = 1, 2$, define $V_i := (g'')^{-1}(Y_i)$. Define $V_{12} := (g'')^{-1}(Y_{12}) = V_1 \cap V_2$. Then there is a unique map $V_i \to W_i$ such that the composed maps $V_i \to X$ and $V_i \to Y_i$ are as desired (by the universal product of the fibered product $X \times_Z Y_i = W_i$), hence a unique map $h_i : V_i \to W$. Similarly, there is a unique map $h_{12} : V_{12} \to W$ such that the composed maps $V_{12} \to X$ and $V_{12} \to Y$ are as desired. But the restriction of $h_i$ to $V_{12}$ is one such map, so it must be $h_{12}$. Thus the maps $h_1$ and $h_2$ agree on $V_{12}$, and glue together to a unique map $h : V \to W$. We have shown existence and uniqueness of the desired $h$. 

We have thus shown that if $Y$ is the union of two affine open sets, and $X$ and $Z$ are affine, then $X \times_Z Y$ exists.

We now tackle the general case. (You may prefer to first think through the case where “two” is replaced by “three”.) We now cover $Y$ with open sets $Y_i$, as $i$ runs over some index set (not necessarily finite!). As before, we define $W_i$ and $W_{ij}$.

We can glue these together to produce a scheme $W$ along with open sets we identify with $W_i$ (Exercise \ref{5.4.A} — you should check the triple intersection “cocycle” condition).

As in the two-affine case, we show that $W$ is the fibered product by showing that it satisfies the universal property. Suppose we have maps $f'' : V \to X$, $g'' : V \to Y$ that compose to the same map $V \to Z$. We construct a unique map $h : V \to W$, so that $f' \circ h = g''$ and $g' \circ h = f''$. Define $V_i = (g'')^{-1}(Y_i)$ and $V_{ij} := (g'')^{-1}(Y_{ij}) = V_i \cap V_j$. Then there is a unique map $V_i \to W_i$ such that the composed maps $V_i \to X$ and $V_i \to Y_i$ are as desired, hence a unique map $h_i : V_i \to W_i$. Similarly, there is a unique map $h_{ij} : V_{ij} \to W_i$ such that the composed maps $V_{ij} \to X$ and $V_{ij} \to Y$ are as desired. But the restriction of $h_i$ to $V_{ij}$ is one such map, so it must be $h_{ij}$. Thus the maps $h_i$ and $h_{ij}$ agree on $V_{ij}$. Thus the $h_i$ glue together to a unique map $h : V \to W$. We have shown existence and uniqueness of the desired $h$, completing this step.

Step 3: $Z$ affine, $X$ and $Y$ arbitrary. We next show that if $Z$ is affine, and $X$ and $Y$ are arbitrary schemes, then $X \times_Z Y$ exists. We just follow Step 2, with the roles of $X$ and $Y$ reversed, using the fact that by the previous step, we can assume that the fibered product with an affine scheme with an arbitrary scheme over an affine scheme exists.

Step 4: $Z$ admits an open immersion into an affine scheme $Z'$, $X$ and $Y$ arbitrary. This is akin to Step 1: $X \times_Z Y$ satisfies the universal property of $X \times_Z Y$.

Step 5: the general case. We again employ the trick from Step 4. Say $f : X \to Z$, $g : Y \to Z$ are two morphisms of schemes. Cover $Z$ with affine open subsets $Z_i$. Let $X_i = f^{-1}Z_i$ and $Y_i = g^{-1}Z_i$. Define $Z_{ij} = Z_i \cap Z_j$, and $X_{ij}$ and $Y_{ij}$ analogously. Then $W_i := X_i \times_{Z_i} Y_i$ exists for all $i$, and has as open sets $W_{ij} := X_{ij} \times_{Z_{ij}} Y_{ij}$ along with gluing information satisfying the cocycle condition (arising from the gluing information for $Z$ from the $Z_i$ and $Z_{ij}$). Once again, we show that this satisfies the universal property. Suppose $V$ is any scheme, along with maps to $X$ and $Y$ that agree when they are composed to $Z$. We need to show that there is a unique morphism $V \to W$ completing the diagram (10.1.41). Now break $V$ up into open sets $V_i = g'' \circ f^{-1}(Z_i)$. Then by the universal property for $W_i$, there is a unique map $V_i \to W_i$ (which we can interpret as $V_i \to W$). Thus we have already shown uniqueness of $V \to W$. These must agree on $V_i \cap V_{ij}$, because there is only one map $V_i \cap V_{ij}$ to $W$ making the diagram commute. Thus all of these morphisms $V_i \to W$ glue together, so we are done.

\section{10.1.5. ** Describing the existence of fibered products using the high-falutin’ language of representable functors.} The proof above can be described more cleanly in the language of representable functors (\ref{10.1.6}). This will be enlightening only after you have absorbed the above argument and meditated on it for a long time. It may be most useful to shed light on representable functors, rather than on the existence of the fibered product.
Until the end of §10.1 only, by functor, we mean contravariant functor from the category Sch of schemes to the category of Sets. For each scheme $X$, we have a functor $h_X$, taking a scheme $Y$ to $\text{Mor}(Y, X)$ (§2.2.20). Recall (§2.3.9, §7.6) that a functor is representable if it is naturally isomorphic to some $h_X$. The existence of the fibered product can be reinterpreted as follows. Consider the functor $h_{X \times_Z Y}$ defined by $h_{X \times_Z Y}(W) = h_X(W) \times_{h_Y(W)} h_Y(W)$. (This isn’t quite enough to define a functor; we have only described where objects go. You should work out where morphisms go too.) Then “$X \times_Z Y$ exists” translates to “$h_{X \times_Z Y}$ is representable”.

If a functor is representable, then the representing scheme is unique up to unique isomorphism (Exercise 7.6.B). This can be usefully extended as follows:

10.1.C. \textsc{Exercise (Yoneda’s lemma).} If $X$ and $Y$ are schemes, describe a bijection between morphisms of schemes $X \to Y$ and natural transformations of functors $h_X \to h_Y$. Hence show that the category of schemes is a fully faithful subcategory of the “functor category” of all functors (contravariant, $\text{Sch} \to \text{Sets}$). Hint: this has nothing to do with schemes; your argument will work in any category. This is the contravariant version of Exercise 2.3.Y(c).

One of Grothendieck’s insights is that we should try to treat such functors as “geometric spaces”, without worrying about representability. Many notions carry over to this more general setting without change, and some notions are easier. For example, fibered products of functors always exist: $h \times h'$ may be defined by

$$(h \times h')(W) = h(W) \times_{h'(W)} h'(W)$$

(where the fibered product on the right is a fibered product of sets, which always exists). We didn’t use anything about schemes; this works with $\text{Sch}$ replaced by any category.

10.1.6. \textit{Representable functors are Zariski sheaves.} Because “morphisms to schemes glue” (Exercise 7.3.A), we have a necessary condition for a functor to be representable. We know that if $\{U_i\}$ is an open cover of $Y$, a morphism $Y \to X$ is determined by its restrictions $U_i \to X$, and given morphisms $U_i \to X$ that agree on the overlap $U_i \cap U_j \to X$, we can glue them together to get a morphism $Y \to X$. In the language of equalizer exact sequences (§3.2.7),

$$
\begin{array}{c}
\text{Hom}(Y, X) \\
\text{Hom(U_i, X)} \\
\text{Hom(U_i \cap U_j, X)}
\end{array}
\xrightarrow{\text{}}
\begin{array}{c}
\prod \text{Hom}(U_i, X) \\
\prod \text{Hom}(U_i \cap U_j, X)
\end{array}
$$

is exact. Thus morphisms to $X$ (i.e. the functor $h_X$) form a sheaf on every scheme $Y$. If this holds, we say that the functor is a Zariski sheaf. (You can impress your friends by telling them that this is a sheaf on the big Zariski site.) We can repeat this discussion with $\text{Sch}$ replaced by the category $\text{Sch}_S$ of schemes over a given base scheme $S$. We have proved (or observed) that in order for a functor to be representable, it is necessary for it to be a Zariski sheaf.

The fiber product passes this test:

10.1.D. \textsc{Exercise.} If $X, Y \to Z$ are schemes, show that $h_{X \times_Z Y}$ is a Zariski sheaf. (Do not use the fact that $X \times_Z Y$ is representable! The point of this section is to recover representability from a more sophisticated perspective.)

We can make some other definitions that extend notions from schemes to functors. We say that a map (i.e. natural transformation) of functors $h \to h'$ expresses $h$
as an open subfunctor of $h'$ if for all representable functors $h_X$ and maps $h_X \to h'$, the fibered product $h_X \times_{h'} h$ is representable, by $U$ say, and $h_U \to h_X$ corresponds to an open immersion of schemes $U \to X$. The following fibered square may help.

\[
\begin{array}{ccc}
h_U & \longrightarrow & h \\
\downarrow & & \downarrow \\
h_X & \longrightarrow & h'
\end{array}
\]

Notice that a map of representable functors $h_W \to h_Z$ is an open subfunctor if and only if $W \to Z$ is an open immersion, so this indeed extends the notion of open immersion to (contravariant) functors ($\text{Sch} \to \text{Sets}$).

10.1.E. Exercise. Suppose $h \to h''$ and $h' \to h''$ are two open subfunctors of $h''$. Define the intersection of these two open subfunctors, which should also be an open subfunctor of $h''$.

10.1.F. Exercise. Suppose $X, Y \to Z$ are schemes, and $U \subset X$, $V \subset Y$, $W \subset Z$ are open subsets, where $U$ and $V$ map to $W$. Interpret $U \times_W V$ as an open subfunctor of $X \times_Z Y$. (Hint: given a map $h_T \to h_{X \times_Z Y}$, what open subset of $T$ should correspond to $U \times W V$?)

A collection $h_i$ of open subfunctors of $h'$ is said to cover $h'$ if for every map $h_X \to h'$ from a representable subfunctor, the corresponding open subsets $U_i \hookrightarrow X$ cover $X$.

Given that functors do not have an obvious underlying set (let alone a topology), it is rather amazing that we are talking about when one is an “open subset” of another, or when some functors “cover” another! (Other notions can be similarly extended. If $P$ is a property of morphisms of schemes that is preserved by base change, then we say that a map of functors $h \to h'$ has $P$ if it is representable, and for each representable $h_X$ mapping to $h'$, the map $h_X \times_{h'} h \to h_X$ — interpreted as a map of schemes via Yoneda’s lemma — has $P$. Note that $h_X \to h_Y$ has $P$ if and only if $X \to Y$ has $P$.)

10.1.G. Exercise. Suppose $\{Z_i\}_i$ is an affine cover of $Z$, $\{X_{ij}\}_i$ is an affine cover of the preimage of $Z_i$ in $X$, and $\{Y_{ik}\}_k$ is an affine cover of the preimage of $Z_i$ in $Y$. Show that $\{h_{X_{ij} \times_{Z_i} Y_{ik}}\}_{ijk}$ is an open cover of the functor $h_{X \times_Z Y}$. (Hint: consider a map $h_T \to h_{X \times_Z Y}$, and extend your solution to the Exercise [10.1.F].)

We now come to a key point: a Zariski sheaf that is “locally representable” must be representable:

10.1.H. Key Exercise. If a functor $h$ is a Zariski sheaf that has an open cover by representable functors (“is covered by schemes”), then $h$ is representable. (Hint: use Exercise [5.4.A] to glue together the schemes representing the open subfunctors.)

This immediately leads to the existence of fibered products as follows. Exercise [10.1.D] shows that $h_{X \times_Z Y}$ is a Zariski sheaf. But $\{h_{X_{ij} \times_{Z_i} Y_{ik}}\}_{ijk}$ is representable (fibered products of affines over an affine exist, Exercise [10.1.B]), and these functors are an open cover of $h_{X \times_Z Y}$ by Exercise [10.1.G], so by Key Exercise [10.1.H] we are done.
10.2 Computing fibered products in practice

Before giving some examples, we first see how to compute fibered products in practice. There are four types of morphisms (1)–(4) that it is particularly easy to take fibered products with, and all morphisms can be built from these four atomic components (see the last paragraph of (1)).

(1) Base change by open immersions.

We have already done this (Exercise 8.1.A), and we used it in the proof that fibered products of schemes exist.

I will describe the remaining three on the level of affine open sets, because we obtain general fibered products by gluing. Theoretically, only (2) and (3) are necessary, as any map of rings \( \phi : B \to A \) can be interpreted by adding variables (perhaps infinitely many) to \( A \), and then imposing relations. But in practice (4) is useful, as we will see in examples.

(2) Adding an extra variable.

10.2.A. EASY BUT SLIGHTLY ANNOYING ALGEBRA EXERCISE. Show that \( B \otimes_A A[t] \cong B[t] \), so the following is a fibered diagram. (Your argument might naturally extend to allow the addition of infinitely many variables, but we won’t need this generality.)

\[
\begin{array}{ccc}
\text{Spec } B[t] & \longrightarrow & \text{Spec } A[t] \\
\downarrow & & \downarrow \\
\text{Spec } B & \longrightarrow & \text{Spec } A
\end{array}
\]

(3) Base change by closed immersions

10.2.B. EXERCISE. Suppose \( \phi : A \to B \) is a ring homomorphism, and \( I \subset A \) is an ideal. Let \( I^e := (\phi(i))_{i \in I} \subset B \) be the extension of \( I \) to \( B \). Describe a natural isomorphism \( B/I^e \cong B \otimes_A (A/I) \). (Hint: consider \( I \to A \to A/I \to 0 \), and use the right-exactness of \( \otimes_A B \), Exercise 2.3.H)

As an immediate consequence: the fibered product with a closed subscheme is a closed subscheme of the fibered product in the obvious way. We say that “closed immersions are preserved by base change”.

10.2.C. EXERCISE. (a) Interpret the intersection of two closed immersions into \( X \) (cf. Exercise 9.1.G) as their fibered product over \( X \).

(b) Show that “locally closed immersions” are preserved by base change.

(c) Define the intersection of a finite number of locally closed immersions in \( X \).

As an application of Exercise 10.2.B, we can compute tensor products of finitely generated \( k \) algebras over \( k \). For example, we have a canonical isomorphism

\[
k[x_1, x_2]/(x_1^2 - x_2) \otimes_k k[y_1, y_2]/(y_1^3 + y_2^3) \cong k[x_1, x_2, y_1, y_2]/(x_1^2 - x_2, y_1^3 + y_2^3).
\]

10.2.D. EXERCISE. Suppose \( X \) and \( Y \) are locally finite type \( k \)-schemes. Show that \( X \times_k Y \) is also locally of finite type over \( k \). Prove the same thing with “locally” removed from both the hypothesis and conclusion.
10.2.1. **Example.** We can use Exercise 10.2.B to compute \(C \otimes_R C\):

\[
C \otimes_R C \cong \frac{\mathbb{R}[x]}{(x^2 + 1)}
\]

by (3)

\[
\cong \frac{\mathbb{C}[x]}{(x^2 + 1)}
\]

by (2)

\[
\cong \frac{\mathbb{C}[x]}{(x - i)(x + i)}
\]

by the Chinese Remainder Theorem

\[
\cong \mathbb{C} \times \mathbb{C}
\]

Thus \(\text{Spec } \mathbb{C} \times_R \text{Spec } \mathbb{C} \cong \text{Spec } \mathbb{C} \sqcup \text{Spec } \mathbb{C}\). This example is the first example of many different behaviors. Notice for example that two points somehow correspond to the Galois group of \(\mathbb{C}\) over \(\mathbb{R}\); for one of them, \(x\) (the “\(i\)” in one of the copies of \(\mathbb{C}\)) equals \(i\) (the “\(i\)” in the other copy of \(\mathbb{C}\)), and in the other, \(x = -i\).

10.2.2. **Remark.** Here is a clue that there is more going on. If \(L/K\) is a Galois extension with Galois group \(G\), then \(L \otimes_K L\) is isomorphic to \(L^{|G|}\) (the product of \(|G|\) copies of \(L\)). This turns out to be a restatement of the classical form of linear independence of characters! In the language of schemes, \(\text{Spec } L \times_K \text{Spec } L\) is a union of a number of copies of \(L\) that naturally form a torsor over the Galois group \(G\).

10.2.E. **Hard but fascinating exercise for those familiar with \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})\).** Show that the points of \(\text{Spec } \overline{\mathbb{Q}} \otimes_{\mathbb{Q}} \overline{\mathbb{Q}}\) are in natural bijection with \(\text{Gal}(\overline{\mathbb{Q}}/\mathbb{Q})\), and the Zariski topology on the former agrees with the profinite topology on the latter. (Some hints: first do the case of finite Galois extensions. Relate the topology on \(\text{Spec}\) of a direct limit of rings to the inverse limit of Specs. Can you see which point corresponds to the identity of the Galois group?)

10.2.G. **Exercise:** The three important types of monomorphisms of schemes. Show that the following are monomorphisms (Definition 2.3.8): open immersions, closed immersions, and localization of affine schemes. As monomorphisms are closed under composition, Exercise 2.3.U, compositions of the above are also monomorphisms (e.g. locally closed immersions, or maps from “Spec of stalks at points of \(X\)” to \(X\)).

10.2.H. **Exercise.** If \(X, Y \hookrightarrow Z\) are two locally closed immersions, show that \(X \times_Z Y\) is canonically isomorphic to \(X \cap Y\).
10.2.1. **Exercise.** Prove that $\mathbb{A}^n_A \cong \mathbb{A}^n_P \times_{\text{Spec } Z} \text{Spec } A$. Prove that $\mathbb{P}^n_A \cong \mathbb{P}^n_P \times_{\text{Spec } Z} \text{Spec } A$. Thus affine space and projective space are pulled back from their universal manifestation over the final object $\text{Spec } Z$.

10.2.3. **Extending the base field.** One special case of base change is called extending the base field: if $X$ is a $k$-scheme, and $k'$ is a field extension (often $k'$ is the algebraic closure of $k$), then $X \times_{\text{Spec } k} \text{Spec } k'$ (sometimes informally written $X \times_k k'$ or $X_k$) is a $k'$-scheme. Often properties of $X$ can be checked by verifying them instead on $X_{k'}$. This is the subject of descent — certain properties “descend” from $X_{k'}$ to $X$. We have already seen that the property of being normal descends in this way (in characteristic 0, Exercise 6.4.K).

10.2.3. **Extending the base field.** One special case of base change is called extending the base field: if $X$ is a $k$-scheme, and $k'$ is a field extension (often $k'$ is the algebraic closure of $k$), then $X \times_{\text{Spec } k} \text{Spec } k'$ (sometimes informally written $X \times_k k'$ or $X_k$) is a $k'$-scheme. Often properties of $X$ can be checked by verifying them instead on $X_{k'}$. This is the subject of descent — certain properties “descend” from $X_{k'}$ to $X$. We have already seen that the property of being normal descends in this way (in characteristic 0, Exercise 6.4.K).

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10.2.4. **Unimportant but fun exercise.** Show that $\text{Spec } \mathbb{Q}(t) \otimes_{\mathbb{Q}} \mathbb{C}$ has closed points in natural correspondence with the transcendental complex numbers. (If the description $\text{Spec } \mathbb{Q}(t) \otimes_{\mathbb{Q}[t]} \mathbb{C}[t]$ is more striking, you can use that instead.)

This scheme doesn’t come up in nature, but it is certainly neat!

10.2.5. **Important concrete exercise.** (A first view of a blow-up, see Figure 10.1.) Consider the rational map $\mathbb{A}^2_k \dashrightarrow \mathbb{P}^1_k$ given by $(x, y) \mapsto [x; y]$. Show that this rational map cannot be extended over the origin. (A similar argument arises in Exercise 7.5.J on the Cremona transformation.) Consider the graph of the birational map, which we denote $\text{Bl}_{(0,0)} \mathbb{A}^2_k$. It is a subscheme of $\mathbb{A}^2_k \times \mathbb{P}^1_k$. Show that if the coordinates on $\mathbb{A}^2_k$ are $x, y$, and the projective coordinates on $\mathbb{P}^1_k$ are $u, v$, this subscheme is cut out in $\mathbb{A}^2_k \times \mathbb{P}^1_k$ by the single equation $xv = yu$. Describe the fiber of the morphism $\text{Bl}_{(0,0)} \mathbb{A}^2_k \to \mathbb{P}^1_k$ over each closed point of $\mathbb{P}^1_k$. Show that the morphism $\text{Bl}_{(0,0)} \mathbb{A}^2_k \to \mathbb{A}^2_k$ is an isomorphism away from $\{(0, 0)\} \in \mathbb{A}^2_k$. Show that the fiber over $(0, 0)$ is a closed subscheme that is locally principal and not locally a zero-divisor (what we will call an effective Cartier divisor, §9.1.2). It is called the exceptional divisor. We will discuss blow-ups in Chapter 19. This particular example will come up in the motivating example of §19.1, and in Exercise 20.6.1.

![Figure 10.1. A first example of a blow-up](image-url)
the smooth surface \( \mathbb{A}^2_k \), cuts out a point, and glues back in a \( \mathbb{P}^1 \), in such a way that the outcome is another smooth surface.”

10.2.4. The graph of a rational map.
Define the graph \( \Gamma_f \) of a rational map \( f : X \to Y \) as follows. Let \( (U, f') \) be any representative of this rational map (so \( f' : U \to Y \) is a morphism). Let \( \Gamma_f \) be the scheme-theoretic closure of \( \Gamma_{f'} \hookrightarrow U \times Y \to X \times Y \), where the first map is a closed immersion, and the second is an open immersion. Equivalently, it is the scheme-theoretic image of the morphism \( U \xrightarrow{(I, f')} X \times Y \). The product here should be taken in the category you are working in. For example, if you are working with \( k \)-schemes, the fibered product should be taken over \( k \).

10.2.L. Exercise. Show that these definitions are indeed equivalent. Show that the graph of a rational map is independent of the choice of representative of the rational map.

In analogy with graphs of morphisms (e.g. Figure 11.3), the following diagram of a graph of a rational map can be handy:

\[ \xymatrix{ \Gamma \ar[r] \ar[d] & X \times Y \ar[r] & Y \ar[d] \\ X & & \ar[lu]_f \ar[lu]_Z } \]

10.3 Pulling back families and fibers of morphisms

10.3.1. Pulling back families.
We can informally interpret fibered product in the following geometric way. Suppose \( Y \to Z \) is a morphism. We interpret this as a “family of schemes parametrized by a base scheme” (or just plain base) \( Z \).” Then if we have another morphism \( f : X \to Z \), we interpret the induced map \( X \times_Z Y \to X \) as the “pulled back family” (see Figure 10.2).

\[ \xymatrix{ X \times_Z Y \ar[r] \ar[d] & Y \ar[d] \\ X \ar[r]_f & Z } \]

We sometimes say that \( X \times_Z Y \) is the scheme-theoretic pullback of \( Y \), scheme-theoretic inverse image, or inverse image scheme of \( Y \). (Our forthcoming discussion of fibers may give some motivation for this.) For this reason, fibered product is often called base change or change of base or pullback. In addition to the various names for a Cartesian diagram given in §2.3.5 in algebraic geometry it is often called a base change diagram or a pullback diagram, and \( X \times_Z Y \to X \) is called the pullback of \( Y \to Z \) by \( f \), and \( X \times_Z Y \) is called the pullback of \( Y \) by \( f \).
Before making any definitions, we give a motivating informal example. Consider the “family of curves” \( y^2 = x^3 + tx \) in the \( xy \)-plane parametrized by \( t \). Translation: consider Spec \( k[x, y, t]/(y^2 - x^3 - tx) \rightarrow \text{Spec } k[t] \). If we pull back to a family parametrized by the \( uv \)-plane via \( uv = t \) (i.e. Spec \( k[u, v] \rightarrow \text{Spec } k[t] \) given by \( t \mapsto uv \)), we get \( \text{Spec } k[x, y, u, v]/(y^2 - x^3 - uvx) \rightarrow \text{Spec } k[u, v] \).

If instead we set \( t \) to 3 (i.e. pull back by Spec \( k[t]/(t - 3) \rightarrow \text{Spec } k[t] \)), we get the curve \( y^2 = x^3 + 3x \) (i.e. Spec \( k[x, y]/(y^2 - x^3 - 3x) \rightarrow \text{Spec } k \)), which we interpret as the fiber of the original family above \( t = 3 \). We will soon be able to interpret these constructions in terms of fiber products.

**10.3.2. Fibers of morphisms.**

A special case of pullback is the notion of a fiber of a morphism. We motivate this with the notion of fiber in the category of topological spaces.

**10.3.A. Exercise.** Show that if \( Y \rightarrow Z \) is a continuous map of topological spaces, and \( X \) is a point \( p \) of \( Z \), then the fiber of \( Y \) over \( p \) (the set-theoretic fiber, with the induced topology) is naturally identified with \( X \times_Z Y \).

More generally, for general \( X \rightarrow Z \), the fiber of \( X \times_Z Y \rightarrow X \) over a point \( p \) of \( X \) is naturally identified with the fiber of \( Y \rightarrow Z \) over \( f(p) \).

Motivated by topology, we return to the category of schemes. Suppose \( p \rightarrow Z \) is the inclusion of a point (not necessarily closed). More precisely, if \( p \) is a point, with residue field \( K \), consider the map Spec \( K \rightarrow Z \) sending Spec \( K \) to \( p \), with the natural isomorphism of residue fields. Then if \( g : Y \rightarrow Z \) is any morphism, the base change with \( p \rightarrow Z \) is called the (scheme-theoretic) fiber of \( g \) above \( p \) or the (scheme-theoretic) preimage of \( p \), and is denoted \( g^{-1}(p) \). If \( Z \) is irreducible, the fiber above the generic point is called the generic fiber. In an affine open subscheme Spec \( A \) containing \( p \), \( p \) corresponds to some prime ideal \( p \), and the
morphism corresponds to the ring map \( A \to A_p/pA_p \). This is the composition of localization and closed immersion, and thus can be computed by the tricks above. (Note that \( p \to \mathbb{Z} \) is a monomorphism, by Exercise 10.2.C.)

10.3.B. Exercise. Show that the underlying topological space of the (scheme-theoretic) fiber \( X \to Y \) above a point \( p \) is naturally identified with the topological fiber of \( X \to Y \) above \( p \).

10.3.C. Exercise (analog of Exercise 10.3.A). Suppose that \( \pi : Y \to Z \) and \( f : X \to Z \) are morphisms, and \( x \in X \) is a point. Show that the fiber of \( X \times_Z Y \to X \) over \( x \) is (isomorphic to) the base change to \( x \) of the fiber of \( \pi : Y \to Z \) over \( f(x) \).

10.3.3. Example (enlightening in several ways). Consider the projection of the parabola \( y^2 = x \) to the \( x \) axis over \( \mathbb{Q} \), corresponding to the map of rings \( \mathbb{Q}[x] \to \mathbb{Q}[y] \), with \( x \mapsto y^2 \). If \( \mathbb{Q} \) alarms you, replace it with your favorite field and see what happens. (You should look at Figure 10.5 and figure out how to edit it to reflect what we glean here.) Writing \( \mathbb{Q}[y] \) as \( \mathbb{Q}[x,y]/(y^2-x) \) helps us interpret the morphism conveniently.

(i) Then the preimage of 1 is two points:
\[
\text{Spec } \mathbb{Q}[x,y]/(y^2-x) \otimes_{\mathbb{Q}[x]} \mathbb{Q}[x]/(x-1) \cong \Spec \mathbb{Q}[x,y]/(y^2-x, x-1) \\
\cong \Spec \mathbb{Q}[y]/(y^2-1) \\
\cong \Spec \mathbb{Q}[y]/(y-1) \coprod \Spec \mathbb{Q}[y]/(y+1).
\]

(ii) The preimage of 0 is one nonreduced point:
\[
\Spec \mathbb{Q}[x,y]/(y^2-x, x) \cong \Spec \mathbb{Q}[y]/(y^2).
\]

(iii) The preimage of \(-1\) is one reduced point, but of “size 2 over the base field”:
\[
\Spec \mathbb{Q}[x,y]/(y^2-x, x+1) \cong \Spec \mathbb{Q}[y]/(y^2+1) \cong \Spec \mathbb{Q}[i].
\]

(iv) The preimage of the generic point is again one reduced point, but of “size 2 over the residue field”, as we verify now.
\[
\Spec \mathbb{Q}[x,y]/(y^2-x) \otimes_{\mathbb{Q}[x]} \mathbb{Q}(x) \cong \Spec \mathbb{Q}[y] \otimes \mathbb{Q}(y^2)
\]
i.e. (informally) the Spec of the ring of polynomials in \( y \) divided by polynomials in \( y^2 \). A little thought shows you that in this ring you may invert any polynomial in \( y \), as if \( f(y) \) is any polynomial in \( y \), then
\[
\frac{1}{f(y)} = \frac{f(-y)}{f(y)f(-y)},
\]
and the latter denominator is a polynomial in \( y^2 \). Thus
\[
\Spec \mathbb{Q}[x,y]/(y^2-x) \otimes \mathbb{Q}(x) \cong \mathbb{Q}(y)
\]
which is a degree 2 field extension of \( \mathbb{Q}(x) \).

Notice the following interesting fact: in each of the four cases, the number of preimages can be interpreted as 2, where you count to two in several ways: you can count points (as in the case of the preimage of 1); you can get non-reduced behavior (as in the case of the preimage of 0); or you can have a field extension of degree 2 (as in the case of the preimage of \(-1\) or the generic point). In each case, the fiber is an affine scheme whose dimension as a vector space over the residue field
of the point is 2. Number theoretic readers may have seen this behavior before. We will discuss this example again in §18.4.8. This is going to be symptomatic of a very special and important kind of morphism (a finite flat morphism).

Try to draw a picture of this morphism if you can, so you can develop a pictorial shorthand for what is going on. A good first approximation is the parabola of Figure 4.5, but you will want to somehow depict the peculiarities of (iii) and (iv).

10.3.D. Exercise (important for those with more arithmetic background). What is the scheme-theoretic fiber of $\text{Spec} \mathbb{Z}[i] \to \text{Spec} \mathbb{Z}$ over the prime $(p)$? Your answer will depend on $p$, and there are four cases, corresponding to the four cases of Example 10.3.3. (Can you draw a picture?)

10.3.E. Exercise. Consider the morphism of schemes $X = \text{Spec} \ k[t] \to Y = \text{Spec} \ k[u]$ corresponding to $k[u] \to k[t], u \mapsto t^2$, where $\text{char} k \neq 2$. Show that $X \times_\gamma X$ has 2 irreducible components. (This exercise will give you practice in computing a fibered product over something that is not a field.)

(What happens if $\text{char} k = 2$? See Exercise 10.4.G for a clue.)

10.4 Properties preserved by base change

All reasonable properties of morphisms are preserved under base change. (In fact, one might say that a property of morphisms cannot be reasonable if it is not preserved by base change!) We discuss this, and explain how to fix those that don’t fit this pattern.

We have already shown that the notion of “open immersion” is preserved by base change (Exercise 8.1.A). We did this by explicitly describing what the fibered product of an open immersion is: if $Y \to Z$ is an open immersion, and $f : X \to Z$ is any morphism, then we checked that the open subscheme $f^{-1}(Y)$ of $X$ satisfies the universal property of fibered products.

We have also shown that the notion of “closed immersion” is preserved by base change (§10.2(3)). In other words, given a fiber diagram

\[
\begin{array}{ccc}
W & \rightarrow & X \\
\downarrow & & \downarrow \\
Y \leftarrow \text{cl. imm.} & \rightarrow & Z
\end{array}
\]

where $Y \to Z$ is a closed immersion, $W \to X$ is as well.

10.4.A. Easy Exercise. Show that locally principal closed subschemes pull back to locally principal closed subschemes.

Similarly, other important properties are preserved by base change.

10.4.B. Exercise. Show that the following properties of morphisms are preserved by base change.

(a) quasicompact
(b) quasiseparated
(c) affine morphism
(d) finite
(e) locally of finite type
(f) finite type
(g) locally of finite presentation
(h) finite presentation

10.4.C. ⋆ HARD EXERCISE. Show that the notion of “quasifinite morphism” (finite type + finite fibers, Definition 8.3.11) is preserved by base change. (Warning: the notion of “finite fibers” is not preserved by base change. Spec $\mathbb{Q} \to$ Spec $\mathbb{Q}$ has finite fibers, but Spec $\mathbb{Q} \otimes_{\mathbb{Q}} \mathbb{Q} \to$ Spec $\mathbb{Q}$ has one point for each element of Gal($\overline{\mathbb{Q}}/\mathbb{Q}$), see Exercise [10.2.E]) Hint: reduce to the case Spec $A \to$ Spec $B$. Reduce to the case $\phi :$ Spec $A \to$ Spec $k$. Show that if $\phi$ is quasifinite then $\phi$ is finite.

10.4.D. EXERCISE. Show that surjectivity is preserved by base change. (Surjectivity has its usual meaning: surjective as a map of sets.) You may end up showing that for any fields $k_1$ and $k_2$ containing $k_3$, $k_1 \otimes_{k_3} k_2$ is non-zero, and using the axiom of choice to find a maximal ideal in $k_1 \otimes_{k_3} k_2$.

10.4.1. On the other hand, injectivity is not preserved by base change — witness the bijection Spec $\mathbb{C} \to$ Spec $\mathbb{R}$, which loses injectivity upon base change by Spec $\mathbb{C} \to$ Spec $\mathbb{R}$ (see Example 10.2.1). This can be rectified (§10.4.5).

10.4.E. EXERCISE. Suppose $X$ and $Y$ are integral finite type $k$-schemes. Show that $X \times_S Y$ is an integrable finite type $k$-scheme. (Once we define “variety”, this will become the important fact that the product of irreducible varieties over an algebraically closed field is an irreducible variety, Exercise 11.1.E) The hypothesis that $k$ is algebraically closed is essential, see {10.4.2}. Hint: reduce to the case where $X$ and $Y$ are both affine, say $X = \text{Spec} A$ and $Y = \text{Spec} B$ with $A$ and $B$ integral domains. Suppose $\left( \sum a_i \otimes b_i \right) \left( a'_i \otimes b'_i \right) = 0$ in $A \otimes_S B$ with $a_i, a'_i \in A$, $b_i, b'_i \in B$, where both $\{b_i\}$ and $\{b'_i\}$ are linearly independent over $\overline{k}$ and $a_1$ and $a'_1$ are nonzero. Show that $D(a_1 a'_1) \subset \text{Spec} A$ is nonempty. By the Weak Nullstellensatz $[4.2.2]$ there is a maximal $m \subset A$ in $D(a_1 a'_1)$ with $A/m = \mathbb{F}$. By reducing modulo $m$, deduce $\left( \sum \overline{a}_i \otimes \overline{b}_i \right) \left( \overline{a'}_i \otimes \overline{b'}_i \right) = 0$ in $B$, where the overline indicates residue modulo $m$. Show that this contradicts the fact that $B$ is a domain.

10.4.F. EXERCISE. If $P$ is a property of morphisms preserved by base change and composition, and $X \to Y$ and $X' \to Y'$ are two morphisms of $S$-schemes with property $P$, show that $X \times_S X' \to Y \times_S Y'$ has property $P$ as well.

10.4.2. ⋆ Properties not preserved by base change, and how to fix (some of) them.

There are some notions that you should reasonably expect to be preserved by pullback based on your geometric intuition. Given a family in the topological category, fibers pull back in reasonable ways. So for example, any pullback of a family in which all the fibers are irreducible will also have this property; ditto for connected. Unfortunately, both of these fail in algebraic geometry, as Example [10.2.1]
The family on the right (the vertical map) has irreducible and connected fibers, and the one on the left doesn’t. The same example shows that the notion of “integral fibers” also doesn’t behave well under pullback.

10.4.G. Exercise. Suppose $k$ is a field of characteristic $p$, so $k(u^p)/k(u)$ is an inseparable extension. By considering $k(u^p) \otimes_{k(u)} k(u^p)$, show that the notion of “reduced fibers” does not necessarily behave well under pullback. (The fact that I am giving you this example should show that this happens only in characteristic $p$, in the presence of something as strange as inseparability.)

We rectify this problem as follows.

10.4.3. A geometric point of a scheme $X$ is defined to be a morphism $\text{Spec } k \to X$ where $k$ is an algebraically closed field. Awkwardly, this is now the third kind of “point” of a scheme! There are just plain points, which are elements of the underlying set; there are $S$-valued points, which are maps $S \to X$, §7.3.6 and there are geometric points. Geometric points are clearly a flavor of an $S$-valued point, but they are also an enriched version of a (plain) point: they are the data of a point with an inclusion of the residue field of the point in an algebraically closed field.

A geometric fiber of a morphism $X \to Y$ is defined to be the fiber over a geometric point of $Y$. A morphism has connected (resp. irreducible, integral, reduced) geometric fibers if all its geometric fibers are connected (resp. irreducible, integral, reduced). One usually says that the morphism has geometrically connected (resp. irreducible, integral, reduced) fibers. A $k$-scheme $X$ is geometrically connected (resp. irreducible, integral, reduced) if the structure morphism $X \to \text{Spec } k$ has geometrically connected (resp. irreducible, integral, reduced) fibers.

10.4.H. Exercise. Show that the notion of “connected (resp. irreducible, integral, reduced)” geometric fibers behaves well under base change.

10.4.I. Exercise for the Arithmetically-Minded. Show that for the morphism $\text{Spec } \mathbb{C} \to \text{Spec } \mathbb{R}$, all geometric fibers consist of two reduced points. (Cf. Example 10.2.1) Thus $\text{Spec } \mathbb{C}$ is a geometrically reduced but not geometrically irreducible $\mathbb{R}$-scheme.

10.4.J. Exercise. Recall Example 10.3.3 the projection of the parabola $y^2 = x$ to the x-axis, corresponding to the map of rings $\mathbb{Q}[x] \to \mathbb{Q}[y]$, with $x \mapsto y^2$. Show that the geometric fibers of this map are always two points, except for those geometric fibers “over $0 = [(x)]$”. (Note that $\text{Spec } \mathbb{C} \to \mathbb{Q}[x]$ and $\text{Spec } \mathbb{C}[x] \to \mathbb{Q}[x]$, both with $x \mapsto 0$, are both geometric points “above 0”.)

Checking whether a $k$-scheme is geometrically connected etc. seems annoying: you need to check every single algebraically closed field containing $k$. However, in each of these four cases, the failure of nice behavior of geometric fibers can already be detected after a finite field extension. For example, $\text{Spec } \mathbb{Q}(i) \to \text{Spec } \mathbb{Q}$ is not...
geometrically connected, and in fact you only need to base change by Spec \( \mathbb{Q}(i) \) to see this. We make this precise as follows.

Suppose \( X \) is a \( k \)-scheme. If \( K/k \) is a field extension, define \( X_K = X \times_k \text{Spec} K \). Consider the following twelve statements.

- \( X_K \) is reduced:
  - (Ra) for all fields \( K \),
  - (Rb) for all algebraically closed fields \( K \) (\( X \) is geometrically reduced),
  - (Rc) for \( K = \overline{k} \),
  - (Rd) for \( K = k^p \) (\( k^p \) is the perfect closure of \( k \)).

- \( X_K \) is irreducible:
  - (Ia) for all fields \( K \),
  - (Ib) for all algebraically closed fields \( K \) (\( X \) is geometrically irreducible),
  - (Ic) for \( K = \overline{k} \),
  - (Id) for \( K = k^s \) (\( k^s \) is the separable closure of \( k \)).

- \( X_K \) is connected:
  - (Ca) for all fields \( K \),
  - (Cb) for all algebraically closed fields \( K \) (\( X \) is geometrically connected),
  - (Cc) for \( K = \overline{k} \),
  - (Cd) for \( K = k^s \).

Trivially (Ra) implies (Rb) implies (Rc), and (Ra) implies (Rd), and similarly with “reduced” replaced by “irreducible” and “connected”.

**10.4.K. Exercise.**

(a) Suppose that \( E/F \) is a field extension, and \( A \) is an \( F \)-algebra. Show that \( A \) is a subalgebra of \( A \otimes_F E \). (Hint: think of these as vector spaces over \( F \).)

(b) Show that: (Ra) implies (Rb) and (Rc) implies (Rd).

(c) Show that: (Ia) implies (Ib) and (Ic) implies (Id).

(d) Show that: (Ca) implies (Cb) and (Cc) implies (Cd).

Notice: you may use the fact that if \( Y \) is a nonempty \( F \)-scheme, then \( Y \times_F \text{Spec} E \) is nonempty, cf. Exercise 10.4.D.

Thus for example a \( k \)-scheme is geometrically integral if and only if it remains integral under any field extension.

**10.4.4. Hard fact.**

In fact, (Rd) implies (Ra), and (Ra) through (Rd) are all equivalent, and similarly for the other two rows. You may try to find this fact in some commutative algebra text.

**10.4.5. Universally injective (radicial) morphisms.**

As remarked in [10.4.1] injectivity is not preserved by base change. A better notion is that of universally injective morphisms: morphisms that are injections of sets after any base change. In keeping with the traditional agricultural terminology (sheaves, germs, ..., cf. Remark 3.4.3), these morphisms were named radicial after one of the lesser vegetables. This notion is more useful in positive characteristic, as the following exercise makes clear.

**10.4.L. Exercise.**

(a) Show that locally closed immersions (and in particular open and closed immersions) are universally injective. (a) Show that \( f : X \to Y \) is universally injective only if \( f \) is injective, and for each \( x \in X \), the field extension \( \kappa(x)/\kappa(f(x)) \) is purely inseparable.
(b) Show that the class of universally injective morphisms are stable under composition, products, and base change.
(c) If \( g : Y \rightarrow Z \) is another morphism, show that if \( g \circ f \) is radicial, then \( f \) is radicial.

### 10.5 Products of projective schemes: The Segre embedding

We next describe products of projective \( A \)-schemes over \( A \). (The case of greatest initial interest is if \( A = k \).) To do this, we need only describe \( \mathbb{P}_A^m \times_A \mathbb{P}_A^n \), because any projective \( A \)-scheme has a closed immersion in some \( \mathbb{P}_A^m \), and closed immersions behave well under base change, so if \( X \hookrightarrow \mathbb{P}_A^m \) and \( Y \hookrightarrow \mathbb{P}_A^n \) are closed immersions, then \( X \times_A Y \hookrightarrow \mathbb{P}_A^m \times_A \mathbb{P}_A^n \) is also a closed immersion, cut out by the equations of \( X \) and \( Y \) (Exercise 10.2.3). We will describe \( \mathbb{P}_A^m \times_A \mathbb{P}_A^n \), and see that it too is a projective \( A \)-scheme. (Hence if \( X \) and \( Y \) are projective \( A \)-schemes, then their product \( X \times_A Y \) over \( A \) is also a projective \( A \)-scheme.)

Before we do this, we will get some motivation from classical projective spaces (non-zero vectors modulo non-zero scalars, Exercise 5.4.1) in a special case. Our map will send \([x_0; x_1; x_2] \times [y_0; y_1]\) to a point in \( \mathbb{P}^5 \), whose coordinates we think of as being entries in the "multiplication table"

\[
\begin{bmatrix}
  x_0y_0 & x_1y_0 & x_2y_0 \\
  x_0y_1 & x_1y_1 & x_2y_1
\end{bmatrix}.
\]

This is indeed a well-defined map of sets. Notice that the resulting matrix is rank one, and from the matrix, we can read off \([x_0; x_1; x_2]\) and \([y_0; y_1]\) up to scalars. For example, to read off the point \([x_0; x_1; x_2] \in \mathbb{P}^2\), we take the first row, unless it is all zero, in which case we take the second row. (They can’t both be all zero.)

In conclusion: in classical projective geometry, given a point in \( \mathbb{P}^m \) and \( \mathbb{P}^n \), we have produced a point in \( \mathbb{P}^{m+n} \), and from this point in \( \mathbb{P}^{m+n} \), we can recover the points of \( \mathbb{P}^m \) and \( \mathbb{P}^n \).

Suitably motivated, we return to algebraic geometry. We define a map

\[ \mathbb{P}_A^m \times_A \mathbb{P}_A^n \rightarrow \mathbb{P}_A^{m+n} \]

by

\[
([x_0; \ldots; x_m], [y_0; \ldots; y_n]) \mapsto [z_{00}; z_{01}; \ldots; z_i; \ldots; z_{mn}]
\]

\[ = [x_0y_0; x_0y_1; \ldots; x_iy_j; \ldots; x_my_n].\]

More explicitly, we consider the map from the affine open set \( U_i \times V_j \) (where \( U_i = D(x_i) \) and \( V_j = D(y_j) \)) to the affine open set \( W_{ij} = D(z_{ij}) \) by

\[
[x_0/i, \ldots, x_m/i, y_0/j, \ldots, y_n/j] \mapsto (x_0/i, x_1/i, \ldots, x_i/i, y_j/j, \ldots, x_m/i, y_n/j)
\]

or, in terms of algebras, \( z_{ab}/ij \mapsto x_a/iy_b/j \).

#### 10.5.A. Exercise

Check that these maps glue to give a well-defined morphism \( \mathbb{P}_A^m \times_A \mathbb{P}_A^n \rightarrow \mathbb{P}_A^{m+n} \).

#### 10.5.1. Exercise

We next show that this morphism is a closed immersion. We can check this on an open cover of the target (the notion of being a closed immersion is affine-local, Exercise 9.1.1). Let’s check this on the open set where \( z_{ij} \neq 0 \). The preimage of this open set in \( \mathbb{P}_A^m \times \mathbb{P}_A^n \) is the locus where \( x_i \neq 0 \) and \( y_j \neq 0 \), i.e. \( U_i \times V_j \). As
described above, the map of rings is given by \(z_{ab/ij} \mapsto x_{a/i}y_{b/j}\); this is clearly a surjection, as \(z_{aj/ij} \mapsto x_{a/i}\) and \(z_{ib/ij} \mapsto y_{b/j}\). (A generalization of this ad hoc description will be given in Exercise 17.4.D.)

This map is called the **Segre morphism** or **Segre embedding**. If \(A\) is a field, the image is called the **Segre variety**.

10.5.B. **Exercise**. Show that the Segre scheme (the image of the Segre morphism) is cut out (scheme-theoretically) by the equations corresponding to

\[
\text{rank } \begin{pmatrix}
a_{00} & \cdots & a_{0n} \\
\vdots & \ddots & \vdots \\
a_{m0} & \cdots & a_{mn}
\end{pmatrix} = 1,
\]

i.e. that all \(2 \times 2\) minors vanish. Hint: suppose you have a polynomial in the \(a_{ij}\) that becomes zero upon the substitution \(a_{ij} = x_iy_j\). Give a recipe for subtracting polynomials of the form “monomial times \(2 \times 2\) minor” so that the end result is 0. (The analogous question for the Veronese embedding in special cases is the content of Exercises 9.2.K and 9.2.M.)

10.5.2. **Important Example**. Let’s consider the first non-trivial example, when \(m = n = 1\). We get \(\mathbb{P}^1 \times \mathbb{P}^1 \hookrightarrow \mathbb{P}^3\). We get a single equation

\[
\text{rank } \begin{pmatrix}
a_{00} & a_{01} \\
a_{10} & a_{11}
\end{pmatrix} = 1,
\]

i.e. \(a_{00}a_{11} - a_{01}a_{10} = 0\). We again meet our old friend, the quadric surface (§9.2.8)! Hence: the nonsingular quadric surface \(wz - xy = 0\) is isomorphic to \(\mathbb{P}^1 \times \mathbb{P}^1\) (Figure 9.2). One family of lines corresponds to the image of \(\{x\} \times \mathbb{P}^1\) as \(x\) varies, and the other corresponds to the image \(\mathbb{P}^1 \times \{y\}\) as \(y\) varies.

If we are working over an algebraically closed field of characteristic not 2, then by diagonalizability of quadratics (Exercise 6.4.J), all rank 4 (“full rank”) quadratics are isomorphic, so all rank 4 quadric surfaces over an algebraically closed field of characteristic not 2 are isomorphic to \(\mathbb{P}^1 \times \mathbb{P}^1\).

Note that this is not true over a field that is not algebraically closed. For example, over \(\mathbb{R}\), \(w^2 + x^2 + y^2 + z^2 = 0\) is not isomorphic to \(\mathbb{P}^1_{\mathbb{R}} \times \mathbb{P}^1_{\mathbb{R}}\). Reason: the former has no real points, while the latter has lots of real points.

You may wish to do the next two exercises in either order.

10.5.C. **Exercise**: A **coordinate-free description of the Segre embedding**. Show that the Segre embedding can be interpreted as \(\mathbb{P}V \times \mathbb{P}W \to \mathbb{P}(V \otimes W)\) via the surjective map of graded rings

\[
\text{Sym}^* (V^\vee \otimes W^\vee) \longrightarrow \sum_{i=0}^{\infty} \left(\text{Sym}^i V^\vee\right) \otimes \left(\text{Sym}^i W^\vee\right)
\]

“in the opposite direction”.

10.5.D. **Exercise**: A **coordinate-free description of products of projective \(A\)-schemes in general**. Suppose that \(S_*\) and \(T_*\) are finitely-generated graded rings over \(A\). Describe an isomorphism

\[
(\text{Proj} S_*) \times_A (\text{Proj} T_*) \cong \text{Proj} \bigoplus_{n=0}^{\infty} (S_n \otimes_A T_n)
\]
10.6 Normalization

Normalization is a means of turning a reduced scheme into a normal scheme. A normalization of a scheme X is a morphism \( \nu : \tilde{X} \to X \) from a normal scheme, where \( \nu \) induces a bijection of irreducible components of \( \tilde{X} \) and \( X \), and \( \nu \) gives a birational morphism on each of the irreducible components. (We need the scheme to have irreducible components for this to make sense, so we will often impose hypotheses such as Noetherianness to keep our scheme from being pathological.) It will satisfy a universal property, and hence it is unique up to unique isomorphism.

Figure 8.4 is an example of a normalization. We discuss normalization now because the argument for its existence follows that for the existence of the fibered product.

We begin with the case where \( X \) is irreducible, and hence integral. (We will then deal with a more general case, and also discuss normalization in a function field extension.) In this case of irreducible \( X \), the normalization \( \nu : \tilde{X} \to X \) is a dominant morphism from an irreducible normal scheme to \( X \), such that any other such morphism factors through \( \nu \):

\[
\begin{array}{ccc}
\text{normal} & \xrightarrow{3!} & \tilde{X} \\
\downarrow & \text{f dominant} & \downarrow \text{\nu dominant} \\
X & & \rightarrow \\
\end{array}
\]

Thus if the normalization exists, then it is unique up to unique isomorphism. We now have to show that it exists, and we do this in a way that will look familiar. We deal first with the case where \( X \) is affine, say \( X = \text{Spec} A \), where \( A \) is an integral domain. Then let \( \tilde{A} \) be the integral closure of \( A \) in its fraction field \( K(A) \). (Recall that the integral closure of \( A \) in its fraction field consists of those elements of \( K(A) \) that are solutions to monic polynomials in \( A[x] \). It is a ring extension by Exercise 8.2.D and integrally closed by Exercise 8.2.K)

**10.6.A. Exercise.** Show that \( \nu : \text{Spec} \tilde{A} \to \text{Spec} A \) satisfies the universal property. (En route, you might show that the global sections of a normal scheme are also normal.)

**10.6.B. Important (but surprisingly easy) Exercise.** Show that normalizations of integral schemes exist in general. (Hint: Ideas from the existence of fiber products, §10.1 may help.)

**10.6.C. Easy Exercise.** Show that normalizations are integral and surjective. (Hint for surjectivity: the Lying Over Theorem, see §8.2.6)

**10.6.D. Exercise.** Explain how to extend the notion of normalization to the case where \( X \) is a reduced Noetherian scheme, with possibly more than one component. (We add the Noetherian hypotheses to ensure that we have irreducible components, Proposition 4.6.6) This basically requires defining a universal property. I’m
not sure what the “perfect” definition is, but all reasonable universal properties should be equivalent.

Here are some examples.

10.6.E. Exercise. Show that \( \text{Spec} \, k[t] \rightarrow \text{Spec} \, k[x, y]/(y^2 - x^2(x + 1)) \) given by \((x, y) \mapsto (t^2 - 1, t(t^2 - 1))\) (see Figure 5.4) is a normalization. (Hint: show that \( k[t] \) and \( k[x, y]/(y^2 - x^2(x + 1)) \) have the same fraction field. Show that \( k[t] \) is integrally closed. Show that \( k[t] \) is contained in the integral closure of \( k[x, y]/(y^2 - x^2(x + 1)) \)).

You will see from the previous exercise that once we guess what the normalization is, it isn’t hard to verify that it is indeed the normalization. Perhaps a few words are in order as to where the polynomials \( t^2 - 1 \) and \( t(t^2 - 1) \) arose in the previous exercise. The key idea is to guess \( t = y/x \). (Then \( t^2 = x + 1 \) and \( y = xt \) quickly.) This idea comes from three possible places. We begin by sketching the curve, and noticing the node at the origin. (a) The function \( y/x \) is well-defined away from the node, and at the node, the two branches have “values” \( y/x = 1 \) and \( y/x = -1 \). (b) We can also note that if \( t = y/x \), then \( t^2 \) is a polynomial, so we will need to adjoin \( t \) in order to obtain the normalization. (c) The curve is cubic, so we expect a general line to meet the cubic in three points, counted with multiplicity. (We will make this precise when we discuss Bézout’s Theorem, Exercise 20.5.L, but in this case we have already gotten a hint of this in Exercise 7.5.H) There is a \( P_1 \) parametrizing lines through the origin (with coordinate equal to the slope of the line, \( y/x \)), and most such lines meet the curve with multiplicity two at the origin, and hence meet the curve at precisely one other point of the curve. So this “co-ordinatizes” most of the curve, and we try adding in this coordinate.

10.6.F. Exercise. Find the normalization of the cusp \( y^2 = x^3 \) (see Figure 10.3).

10.6.G. Exercise. Find the normalization of the tacnode \( y^2 = x^4 \), and draw a picture analogous to Figure 10.3.

(Although we haven’t defined “singularity”, “smooth”, “curve”, or “dimension”, you should still read this.) Notice that in the previous examples, normalization “resolves” the singularities (non-smooth points) of the curve. In general, it
will do so in dimension one (in reasonable Noetherian circumstances, as normal Noetherian integral domains of dimension one are all Discrete Valuation Rings, \[ \text{§13.3} \]), but won’t do so in higher dimension (the cone \( z^2 = x^2 + y^2 \) over a field \( k \) of characteristic not 2 is normal, Exercise \[6.4.1b\]).

**10.6.H. Exercise.** Suppose \( X = \text{Spec } \mathbb{Z}[15i] \). Describe the normalization \( \tilde{X} \to X \). (Hint: \( \mathbb{Z}[i] \) is a unique factorization domain, \[6.4.5\]), and hence is integrally closed by Exercise \[6.4.H\].) Over what points of \( X \) is the normalization not an isomorphism?

Another exercise in a similar vein is the normalization of the “knotted plane”, Exercise \[13.3.I\].

**10.6.I. Exercise (Normalization in a function field extension, an important generalization).** Suppose \( X \) is an integral scheme. The **normalization of \( X \)**, \( \nu : \tilde{X} \to X \), **in a given finite field extension \( L \) of the function field \( K(\mathcal{X}) \) of \( X \)** is a dominant morphism from a normal scheme \( \tilde{X} \) with function field \( L \), such that \( \nu \) induces the inclusion \( K(\mathcal{X}) \to L \), and that is universal with respect to this property.

\[
\begin{array}{ccc}
\text{Spec } L = K(Y) & \longrightarrow & Y \\
\downarrow & & \downarrow \exists ! \\
\text{Spec } L = K(\tilde{X}) & \longrightarrow & \tilde{X} \\
\downarrow & & \downarrow \\
K(\mathcal{X}) & \longrightarrow & X
\end{array}
\]

Show that the normalization in a finite field extension exists.

The following two examples, one arithmetic and one geometric, show that this is an interesting construction.

**10.6.J. Exercise.** Suppose \( X = \text{Spec } \mathbb{Z} \) (with function field \( \mathbb{Q} \)). Find its integral closure in the field extension \( \mathbb{Q}(i) \). (There is no “geometric” way to do this; it is purely an algebraic problem, although the answer should be understood geometrically.)

**10.6.1. Remark: rings of integers in number fields.** A finite extension \( K \) of \( \mathbb{Q} \) is called a **number field**, and the integral closure of \( \mathbb{Z} \) in \( K \) the **ring of integers in \( K \)**, denoted \( \mathcal{O}_K \). (This notation is awkward given our other use of the symbol \( \mathcal{O} \).)

\[
\begin{array}{ccc}
\text{Spec } K & \longrightarrow & \text{Spec } \mathcal{O}_K \\
\downarrow & & \downarrow \\
\text{Spec } \mathbb{Q} & \longrightarrow & \text{Spec } \mathbb{Z}
\end{array}
\]

By the previous exercises, \( \text{Spec } \mathcal{O}_K \) is a Noetherian normal integral domain of dimension 1. This is an example of a **Dedekind domain**, see \[13.3.14\]. We will think of it as a smooth curve as soon as we know what “smooth” and “curve” mean.
10.6.K. Exercise. (a) Suppose \( X = \text{Spec} \, k[x] \) (with function field \( k(x) \)). Find its integral closure in the field extension \( k(y) \), where \( y^2 = x^2 + x \). (Again we get a Dedekind domain.) Hint: this can be done without too much pain. Show that \( \text{Spec} \, k[x, y]/(x^2 + x - y^2) \) is normal, possibly by identifying it as an open subset of \( \mathbb{P}^1_k \), or possibly using Exercise 6.4.I.

(b) Suppose \( X = \mathbb{P}^1 \), with distinguished open \( \text{Spec} \, k[x] \). Find its integral closure in the field extension \( k(y) \), where \( y^2 = x^2 + x \). (Part (a) involves computing the normalization over one affine open set; now figure out what happens over the “other” affine open set.)

10.6.2. Fancy fact: finiteness of integral closure.

The following fact is useful.

10.6.3. Theorem (finiteness of integral closure). — Suppose \( A \) is a Noetherian integral domain, \( K = k(A) \), \( L/K \) is a finite separable field extension, and \( B \) is the integral closure of \( A \) in \( L \) (“the integral closure of \( A \) in the field extension \( L/K \), i.e. those elements of \( L \) integral over \( A \)).

(a) If \( A \) is integrally closed, then \( B \) is a finitely generated \( A \)-module.

(b) If \( A \) is a finitely generated \( k \)-algebra and \( L = K \), then \( B \) is a finitely generated \( A \)-module.

Eisenbud gives a proof in a page and a half: (a) is [E] Prop. 13.14] and (b) is [E] Cor. 13.13. A sketch is given in §10.6.4.

Warning: (b) does not hold for Noetherian \( A \) in general. In fact, the integral closure of a Noetherian ring need not be Noetherian (see [E] p. 299] for some discussion). This is alarming. The existence of such an example is a sign that Theorem 10.6.3 is not easy.

10.6.L. Exercise. (a) Show that if \( X \) is an integral finite-type \( k \)-scheme, then its normalization \( \nu : \tilde{X} \to X \) is a finite morphism.

(b) Suppose \( X \) is an integral scheme. Show that if either \( X \) is normal, or \( X \) is a finite type \( k \)-scheme, then the normalization in a finite field extension is a finite morphism. In particular, the normalization of a variety (including in a finite separable field extension) is a variety.

10.6.M. Exercise. Show that if \( X \) is an integral finite type \( k \)-scheme. Show that the normalization map is an isomorphism on an open dense subset of \( X \). Hint: reduce to the case \( X = \text{Spec} \, A \). By Theorem 10.6.3, \( \tilde{A} \) is generated over \( A \) by a finite number of elements of \( K(A) \). Let \( I \) be the ideal generated by their denominators. Show that \( \text{Spec} \, \tilde{A} \to \text{Spec} \, A \) is an isomorphism away from \( V(I) \). (Alternatively, the ideas of Proposition 11.2.3 can also be applied.)

10.6.4. ** Sketch of proof of finiteness of integral closure, Theorem 10.6.3. ** Here is a sketch to show the structure of the argument. It uses commutative algebra ideas from Chapter 12 so you should only glance at this to see that nothing fancy is going on. Part (a): reduce to the case where \( L/K \) is Galois, with group \( \{\sigma_1, \ldots, \sigma_n\} \). Choose \( b_1, \ldots, b_n \in B \) forming a \( K \)-vector space basis of \( L \). Let \( M \) be the matrix (familiar from Galois theory) with \( ij \)th entry \( \sigma_i b_j \), and let \( d = \det M \). Show that the entries of \( M \) lie in \( B \), and that \( d^2 \in K \) (as \( d^2 \) is Galois-fixed). Show that \( d \neq 0 \) using linear independence of characters. Then complete the proof by showing...
that $B \subset d^{-2}(\text{Ab}_1 + \cdots + \text{Ab}_n)$ (submodules of finitely generated modules over Noetherian rings are also Noetherian, §4.6.4) as follows. Suppose $b \in B$, and write $b = \sum c_i b_i$ ($c_i \in K$). If $c$ is the column vector with entries $c_i$, show that the $i$th entry of the column vector $Mc$ is $\sigma_i b \in B$. Multiplying $Mc$ on the left by $\text{adj } M$ (see the trick of the proof of Lemma 8.2.1), show that $d c_i \in B$. Thus $d^2 c_i \in B \cap K = A$ (as $A$ is integrally closed), as desired.

For (b), use the Noether Normalization Lemma [12.2.7] to reduce to the case $A = k[x_1, \ldots, x_n]$. Reduce to the case where $L$ is normally closed over $K$. Let $L'$ be the subextension of $L/K$ so that $L/L'$ is Galois and $L'/K$ is purely inseparable. Use part (a) to reduce to the case $L = L'$. If $L' \neq K$, then for some $q$, $L'$ is generated over $K$ by the $q$th root of a finite set of rational functions. Reduce to the case $L' = k'(x_1^{1/q}, \ldots, x_n^{1/q})$ where $k'/k$ is a finite purely inseparable extension. In this case, show that $B = k'[x_1^{1/q}, \ldots, x_n^{1/q}]$, which is indeed finite over $k[x_1, \ldots, x_n]$. 

CHAPTER 11

Separated and proper morphisms, and (finally!) varieties

11.1 Separated morphisms (and quasiseparatedness done properly)

Separatedness is a fundamental notion. It is the analogue of the Hausdorff condition for manifolds (see Exercise [11.1.A]), and as with Hausdorffness, this geometrically intuitive notion ends up being just the right hypothesis to make theorems work. Although the definition initially looks odd, in retrospect it is just perfect.

11.1.1. Motivation. Let’s review why we like Hausdorffness. Recall that a topological space is Hausdorff if for every two points \(x\) and \(y\), there are disjoint open neighborhoods of \(x\) and \(y\). The real line is Hausdorff, but the “real line with doubled origin” is not (of which Figure 5.4 may be taken as a sketch). Many proofs and results about manifolds use Hausdorffness in an essential way. For example, the classification of compact one-dimensional smooth manifolds is very simple, but if the Hausdorff condition were removed, we would have a very wild set.

So once armed with this definition, we can cheerfully exclude the line with doubled origin from civilized discussion, and we can (finally) define the notion of a variety, in a way that corresponds to the classical definition.

With our motivation from manifolds, we shouldn’t be surprised that all of our affine and projective schemes are separated: certainly, in the land of smooth manifolds, the Hausdorff condition comes for free for “subsets” of manifolds. (More precisely, if \(Y\) is a manifold, and \(X\) is a subset that satisfies all the hypotheses of a manifold except possibly Hausdorffness, then Hausdorffness comes for free. Similarly, locally closed immersions in something separated are also separated: combine Exercise [11.1.B] and Proposition [11.1.13(a)].)

As an unexpected added bonus, a separated morphism to an affine scheme has the property that the intersection of two affine open sets in the source is affine (Proposition [11.1.8]). This will make Čech cohomology work very easily on (quasicompact) schemes (Chapter 20). You might consider this an analogue of the fact that in \(\mathbb{R}^n\), the intersection of two convex sets is also convex. As affine schemes are trivial from the point of view of quasicoherent cohomology, just as convex sets in \(\mathbb{R}^n\) have no cohomology, this metaphor is apt.

A lesson arising from the construction is the importance of the diagonal morphism. More precisely given a morphism \(X \to Y\), good consequences can be leveraged from good behavior of the diagonal morphism \(\delta : X \to X \times_Y X\), usually
through fun diagram chases. This lesson applies across many fields of geometry.
(Another nice gift of the diagonal morphism: it will give us a good algebraic definition of differentials, in Chapter 22.)

Grothendieck taught us that one should try to define properties of morphisms, not of objects; then we can say that an object has that property if its morphism to the final object has that property. We discussed this briefly at the start of Chapter 8. In this spirit, separatedness will be a property of morphisms, not schemes.

11.1.2. Defining separatedness.

Before we define separatedness, we make an observation about all diagonal morphisms.

11.1.3. Proposition. — Let \( \pi : X \to Y \) be a morphism of schemes. Then the diagonal morphism \( \delta : X \to X \times_Y X \) is a locally closed immersion.

We will often use \( \delta \) to denote a diagonal morphism. This locally closed subscheme of \( X \times_Y X \) (which we also call the diagonal) will be denoted \( \Delta \).

Proof. We will describe a union of open subsets of \( X \times_Y X \) covering the image of \( X \), such that the image of \( X \) is a closed immersion in this union.

Say \( Y \) is covered with affine open sets \( V_i \) and \( X \) is covered with affine open sets \( U_{ij} \), with \( \pi : U_{ij} \to V_i \). Note that \( U_{ij} \times_Y U_{ij} \) is an affine open subscheme of the product \( X \times_Y X \) (basically this is how we constructed the product, by gluing together affine building blocks). Then the diagonal is covered by these affine open subsets \( U_{ij} \times_Y U_{ij} \). (Any point \( p \in X \) lies in some \( U_{ij} \); then \( \delta(p) \in U_{ij} \times_Y U_{ij} \). Figure 11.1 may be helpful.) Note that \( \delta^{-1}(U_{ij} \times_Y U_{ij}) = U_{ij} \); clearly \( U_{ij} \subset \delta^{-1}(U_{ij} \times_Y U_{ij}) \), and because \( \text{pr}_1 \circ \delta = \text{id}_X \) (where \( \text{pr}_1 \) is the first projection), \( \delta^{-1}(U_{ij} \times_Y U_{ij}) \subset U_{ij} \). Finally, we check that \( U_{ij} \to U_{ij} \times_Y U_{ij} \) is a closed immersion. Say \( V_i = \text{Spec} B \) and \( U_{ij} = \text{Spec} A \). Then this corresponds to the natural ring map \( A \otimes_B A \to A \) \((a_1 \otimes a_2 \mapsto a_1 a_2)\), which is obviously surjective. \( \square \)

![Figure 11.1](image-url)

**Figure 11.1.** A neighborhood of the diagonal is covered by \( U_{ij} \times_Y U_{ij} \).
The open subsets we described may not cover $X \times Y$, so we have not shown that $\delta$ is a closed immersion.

11.1.4. Definition. A morphism $X \to Y$ is separated if the diagonal morphism $\delta : X \to X \times Y$ is a closed immersion. An $A$-scheme $X$ is said to be separated over $A$ if the structure morphism $X \to \text{Spec} A$ is separated. When people say that a scheme (rather than a morphism) $X$ is separated, they mean implicitly that some “structure morphism” is separated. For example, if they are talking about $A$-schemes, they mean that $X$ is separated over $A$.

Thanks to Proposition 11.1.3, a morphism is separated if and only if the diagonal $\Delta$ is a closed subset — a purely topological condition on the diagonal. This is reminiscent of a definition of Hausdorff, as the next exercise shows.

11.1.A. Unimportant Exercise (for those seeking topological motivation). Show that a topological space $X$ is Hausdorff if and only if the diagonal is a closed subset of $X \times X$. (The reason separatedness of schemes doesn’t give Hausdorffness — i.e. that for any two open points $x$ and $y$ there aren’t necessarily disjoint open neighborhoods — is that in the category of schemes, the topological space $X \times X$ is not in general the product of the topological space $X$ with itself, see §10.1.2.)

11.1.B. Important Easy Exercise. Show that open immersions, closed immersions, and hence locally closed immersions are separated. (Hint: Do this by hand. Alternatively, show that monomorphisms are separated. Open and closed immersions are monomorphisms, by Exercise 10.2.G.)

11.1.C. Important Easy Exercise. Show that every morphism of affine schemes is separated. (Hint: this was essentially done in the proof of Proposition 11.1.3.)

11.1.D. Exercise. Show that the line with doubled origin $X$ (Example 5.4.5) is not separated, by verifying that the image of the diagonal morphism is not closed. (Another argument is given below, in Exercise 11.1.E. A fancy argument is given in Exercise 13.4.C.)

We next come to our first example of something separated but not affine. The following single calculation will imply that all quasiprojective $A$-schemes are separated (once we know that the composition of separated morphisms are separated, Proposition 11.1.13).

11.1.5. Proposition. — $\mathbb{P}^n_A \to \text{Spec } A$ is separated.

We give two proofs. The first is by direct calculation. The second requires no calculation, and just requires that you remember some classical constructions described earlier.

Proof 1: direct calculation. We cover $\mathbb{P}^n_A \times_A \mathbb{P}^n_A$ with open sets of the form $U_i \times_A U_j$, where $U_0, \ldots, U_n$ form the “usual” affine open cover. The case $i = j$ was taken care of before, in the proof of Proposition 11.1.3. If $i \neq j$ then

$$U_i \times_A U_j \cong \text{Spec } A[x_0/i, \ldots, x_n/i, y_0/j, \ldots, y_n/j]/(x_i/i - 1, y_j/j - 1).$$
Now the restriction of the diagonal $\Delta$ is contained in $U_i$ (as the diagonal morphism composed with projection to the first factor is the identity), and similarly is contained in $U_j$. Thus the diagonal morphism over $U_i \times_A U_j$ is $U_i \cap U_j \to U_i \times_A U_j$. This is a closed immersion, as the corresponding map of rings

$$A[x_{0/i}, \ldots, x_{n/i}, y_{0/j}, \ldots, y_{n/j}] \to A[x_{0/i}, \ldots, x_{n/i}, x_{j/i}^{-1}]/(x_{t/i} - 1)$$

(given by $x_{k/i} \mapsto x_{k/i}$, $y_{k/j} \mapsto x_{k/i}/x_{j/i}$) is clearly a surjection (as each generator of the ring on the right is clearly in the image — note that $x_{j/i}^{-1}$ is the image of $y_{1/i}$).

**Proof 2: classical geometry.** Note that the diagonal morphism $\delta : \mathbb{P}^n_A \to \mathbb{P}^n_A \times_A \mathbb{P}^n_A$ followed by the Segre embedding $S : \mathbb{P}^n_A \times_A \mathbb{P}^n_A \to \mathbb{P}^{n^2+n}$ ([10.5] a closed immersion) can also be factored as the second Veronese embedding $\nu_2 : \mathbb{P}^n_A \to \mathbb{P}^{(n+2)^2-1}$ ([9.25] followed by a linear map $L : \mathbb{P}^{(n+2)^2-1} \to \mathbb{P}^{n^2+n}$ (another closed immersion, Exercise [9.2.1]), both of which are closed immersions.

$$
\begin{array}{cc}
\mathbb{P}^n_A & \mathbb{P}^{n+2+n} \\
\mathbb{P}^n_A \times_A \mathbb{P}^n_A & \\
\delta & S \\
\mathbb{P}^{(n+2)^2-1} & L \\
\end{array}
$$

Informally, in coordinates:

$$([x_0; x_1; \cdots; x_n], [x_0; x_1; \cdots; x_n])$$

$$\delta$$

$$S$$

$$\begin{bmatrix}
x_0^2; & x_0x_1; & \cdots; & x_0x_n; \\
x_1x_0; & x_1^2; & \cdots; & x_1x_n; \\
\vdots; & \vdots; & \ddots; & \vdots; \\
x_nx_0; & x_nx_1; & \cdots; & x_n^2
\end{bmatrix}$$

$$\nu_2$$

$$\begin{bmatrix}
x_0^2; & x_0x_1; & \cdots; & x_0x_n; \\
x_1x_0; & x_1^2; & \cdots; & x_1x_n; \\
\vdots; & \vdots; & \ddots; & \vdots; \\
x_nx_0; & x_nx_1; & \cdots; & x_n^2
\end{bmatrix}$$

The composed map $\mathbb{P}^n_A$ may be written as $[x_0; \cdots; x_n] \mapsto [x_0x_0; x_0x_1; \cdots; x_nx_n]$, where the subscripts on the right run over all ordered pairs $(i,j)$ where $0 \leq i, j \leq n$. This forces $\delta$ to send closed sets to closed sets (or else $S \circ \delta$ won’t, but $L \circ \nu_2$ does).

We note for future reference a minor result proved in the course of Proof 1.

**11.1.6. Small Proposition.** — If $U$ and $V$ are open subsets of an $A$-scheme $X$, then $\Delta \cap (U \times_A V) \cong U \cap V$. 

Figure 11.2 may help show why this is natural. You could also interpret this statement as

\[ X \times (X \times X) \{ U \times_A V \} \cong U \times_X V \]

which follows from the magic diagram, Exercise 2.3.R.

**Figure 11.2. Small Proposition 11.1.6**

We finally define variety!

11.1.7. **Definition.** A variety over a field \( k \), or \( k \)-variety, is a reduced, separated scheme of finite type over \( k \). For example, a reduced finite-type affine \( k \)-scheme is a variety. We will soon know that the composition of separate morphisms is separated (Exercise 11.1.13(a)), and then to check if \( \text{Spec} \ k[x_1, \ldots, x_n]/(f_1, \ldots, f_r) \) is a variety, you need only check reducedness. This generalizes our earlier notion of affine variety (§6.3.6) and projective variety (§6.3.6, see Proposition 11.1.14). (Notational caution: In some sources, the additional condition of irreducibility is imposed. Also, it is often assumed that \( k \) is algebraically closed.)

11.1.8. **Exercise (Products of Irreducible Varieties over \( k \) are Irreducible Varieties).** Use Exercise 10.4.E and properties of separatedness to show that the product of two irreducible \( k \)-varieties is a an irreducible \( k \)-variety.

Here is a very handy consequence of separatedness.

11.1.9. **Proposition.** — Suppose \( X \to \text{Spec} \ A \) is a separated morphism to an affine scheme, and \( U \) and \( V \) are affine open subsets of \( X \). Then \( U \cap V \) is an affine open subset of \( X \).

Before proving this, we state a consequence that is otherwise nonobvious. If \( X = \text{Spec} \ A \), then the intersection of any two affine open subsets is an affine open subset (just take \( A = \mathbb{Z} \) in the above proposition). This is certainly not an obvious fact! We know the intersection of two distinguished affine open sets is affine (from \( D(f) \cap D(g) = D(fg) \)), but we have little handle on affine open sets in general.
Warning: this property does not characterize separatedness. For example, if $A = \text{Spec } k$ and $X$ is the line with doubled origin over $k$, then $X$ also has this property.

**Proof.** By Proposition 11.1.6, $(U \times_A V) \cap \Delta \cong U \cap V$, where $\Delta$ is the diagonal. But $U \times_A V$ is affine (the fibered product of two affine schemes over an affine scheme is affine, Step 1 of our construction of fibered products, Theorem 10.1.1), and $\Delta$ is a closed subscheme of an affine scheme, and hence $U \cap V$ is affine. \hfill \square

### 11.1.9. Redefinition: Quasiseparated morphisms.

We say a morphism $f : X \to Y$ is **quasiseparated** if the diagonal morphism $\delta : X \to X \times Y X$ is quasicompact.

### 11.1.F. Exercise.

Show that this agrees with our earlier definition of quasiseparated (§8.3.1), show that $f : X \to Y$ is quasiseparated if and only if for any affine open $\text{Spec } A$ of $Y$, and two affine open subsets $U$ and $V$ of $X$ mapping to $\text{Spec } A$, $U \cap V$ is a finite union of affine open sets. (Possible hint: compare this to Proposition 11.1.8. Another possible hint: the magic diagram, Exercise 2.3.R.)

Here are two large classes of morphisms that are quasiseparated.

### 11.1.G. Easy exercise.

Show that separated morphisms are quasiseparated. (Hint: closed immersions are affine, hence quasicompact.)

Second, if $X$ is a Noetherian scheme, then any morphism to another scheme is quasicompact (easy, see Exercise 8.3.B(a)), so any $X \to Y$ is quasiseparated. Hence those working in the category of Noetherian schemes need never worry about this issue.

We now give four quick propositions showing that separatedness and quasiseparatedness behave well, just as many other classes of morphisms did.

### 11.1.10. Proposition.

*Both separatedness and quasiseparatedness are preserved by base change.*

**Proof.** Suppose

$$
\begin{array}{ccc}
W & \longrightarrow & X \\
\downarrow & & \downarrow \\
Y & \longrightarrow & Z
\end{array}
$$

is a fiber diagram. We will show that if $Y \to Z$ is separated or quasiseparated, then so is $W \to X$. Then you can quickly verify that

$$
\begin{array}{ccc}
W & \overset{\delta_W}{\longrightarrow} & W \times_X W \\
\downarrow & & \downarrow \\
Y & \overset{\delta_Y}{\longrightarrow} & Y \times_Z Y
\end{array}
$$

is a fiber diagram. (This is true in any category with fibered products.) As the property of being a closed immersion is preserved by base change ($\text{[10.2](3)}$), if $\delta_Y$ is a closed immersion, so is $\delta_X$. 

Quasiseparatedness follows in the identical manner, as quasicompactness is also preserved by base change (Exercise 10.4.B).

11.1.11. Proposition. — The condition of being separated is local on the target. Precisely, a morphism \( f : X \to Y \) is separated if and only if for any cover of \( Y \) by open subsets \( U_i \), \( f^{-1}(U_i) \to U_i \) is separated for each \( i \).

11.1.12. Hence affine morphisms are separated, as every morphism of affine schemes is separated (Exercise 11.1.C). In particular, finite morphisms are separated.

Proof. If \( X \to Y \) is separated, then for any \( U_i \to Y \), \( f^{-1}(U_i) \to U_i \) is separated, as separatedness is preserved by base change (Theorem 11.1.10). Conversely, to check if \( \Delta \to X \times_Y Y \) is a closed subset, it suffices to check this on an open cover of \( X \times_Y Y \). Let \( g : X \times_Y X \to Y \) be the natural map. We will use the open cover \( g^{-1}(U_i) \), which by construction of the fiber product is \( f^{-1}(U_i) \times_{U_i} f^{-1}(U_i) \). As \( f^{-1}(U_i) \to U_i \) is separated, \( f^{-1}(U_i) \) is a closed immersion by definition of separatedness.

11.1.13. Proposition. — (a) The condition of being separated is closed under composition. In other words, if \( f : X \to Y \) is separated and \( g : Y \to Z \) is separated, then \( g \circ f : X \to Z \) is separated.

(b) The condition of being quasiseparated is closed under composition.

Proof. (a) We are given that \( \delta_f : X \hookrightarrow X \times_Y X \) and \( \delta_g : Y \hookrightarrow Y \times_Z Y \) are closed immersions, and we wish to show that \( \delta_h : X \hookrightarrow X \times_Z X \) is a closed immersion. Consider the diagram

\[
\begin{array}{ccc}
X & \xrightarrow{\delta_f} & X \times_Y X \\
& \downarrow & \downarrow \\
Y & \xrightarrow{\delta_g} & Y \times_Z Y.
\end{array}
\]

The square is the magic diagram (Exercise 2.3.R). As \( \delta_g \) is a closed immersion, \( c \) is too (closed immersions are preserved by base change, [10.2](3)). Thus \( c \circ \delta_f \) is a closed immersion (the composition of two closed immersions is also a closed immersion, Exercise 9.1.B).

(b) The identical argument (with “closed immersion” replaced by “quasicompact”) shows that the condition of being quasiseparated is closed under composition.

11.1.14. Corollary. — Any quasiprojective \( A \)-scheme is separated over \( A \). In particular, any reduced quasiprojective \( k \)-scheme is a \( k \)-variety.
Proof. Suppose $X \to \text{Spec} A$ is a quasiprojective $A$-scheme. The structure morphism can be factored into an open immersion composed with a closed immersion followed by $\mathbb{P}^n_A \to A$. Open immersions and closed immersions are separated (Exercise 11.1.8), and $\mathbb{P}^n_A \to A$ is separated (Proposition 11.1.5). Compositions of separated morphisms are separated (Proposition 11.1.13), so we are done. □

11.1.15. **Proposition.** — Suppose $f : X \to Y$ and $f' : X' \to Y'$ are separated (resp. quasiseparated) morphisms of $S$-schemes (where $S$ is a scheme). Then the product morphism $f \times f' : X \times_S X' \to Y \times_S Y'$ is separated (resp. quasiseparated).

Proof. Apply Exercise 10.4.1. □

11.1.16. **Applications.**

As a first application, we define the graph morphism.

11.1.17. **Definition.** Suppose $f : X \to Y$ is a morphism of $Z$-schemes. The morphism $\Gamma_f : X \to X \times_Z Y$ given by $\Gamma_f = (\text{id}_X, f)$ is called the **graph morphism**. Then $f$ factors as $\text{pr}_2 \circ \Gamma_f$, where $\text{pr}_2$ is the second projection (see Figure 11.3).

![Figure 11.3. The graph morphism](image)

11.1.18. **Proposition.** — The graph morphism $\Gamma$ is always a locally closed immersion. If $Y$ is a separated $Z$-scheme (i.e. the structure morphism $Y \to Z$ is separated), then $\Gamma$ is a closed immersion. If $Y$ is a quasiseparated $Z$-scheme, then $\Gamma$ is quasicompact.

This will be generalized in Exercise 11.1.1.
Proof by Cartesian diagram. A special case of the magic diagram (Exercise 2.3.R) is:

\[(11.1.18.1)\]

The notions of locally closed immersion and closed immersion are preserved by base change, so if the bottom arrow $\delta$ has one of these properties, so does the top. The same argument establishes the last sentence. \(\square\)

We now come to a very useful, but bizarre-looking, result. Like the magic diagram, I find this result unexpected useful and ubiquitous.

**11.1.19. Cancellation Theorem for a Property $P$ of Morphisms.** — Let $P$ be a class of morphisms that is preserved by base change and composition. Suppose

is a commuting diagram of schemes. Suppose that the diagonal morphism $\delta_g : Y \to Y \times Z$ $Y$ is in $P$ and $h : X \to Z$ is in $P$. Then $f : X \to Y$ is in $P$. In particular:

(i) Suppose that locally closed immersions are in $P$. If $h$ is in $P$, then $f$ is in $P$.
(ii) Suppose that closed immersions are in $P$ (e.g. $P$ could be finite morphisms, morphisms of finite type, closed immersions, affine morphisms). If $h$ is in $P$ and $g$ is separated, then $f$ is in $P$.
(iii) Suppose that quasicompact morphisms are in $P$. If $h$ is in $P$ and $g$ is quasiseparated, then $f$ is in $P$.

The following diagram summarizes this important theorem:

When you plug in different $P$, you get very different-looking (and non-obvious) consequences. For example, if you factor a locally closed immersion $X \to Z$ into $X \to Y \to Z$, then $X \to Y$ must be a locally closed immersion.

**Proof.** By the graph Cartesian diagram \[(11.1.18.1)\]
we see that the graph morphism $\Gamma_f : X \to X \times Z Y$ is in $P$ (Definition 11.1.17), as $P$ is closed under base change. By the fibered square

\[
\begin{array}{ccc}
X \times Z Y & \xrightarrow{h'} & Y \\
\downarrow & & \downarrow \\
X & \xrightarrow{h} & Z \\
\end{array}
\]

the projection $h' : X \times Z Y \to Y$ is in $P$ as well. Thus $f = h' \circ \Gamma_f$ is in $P$. \hfill \Box

Here now are some fun and useful exercises.

11.1.I. Exercise. Suppose $\pi : Y \to X$ is a morphism, and $s : X \to Y$ is a section of a morphism, i.e. $\pi \circ s$ is the identity on $X$. Show that $s$ is a locally closed immersion.

Show that if $\pi$ is separated, then $s$ is a closed immersion. (This generalizes Proposition 11.1.18) Give an example to show that $s$ needn’t be a closed immersion if $\pi$ isn’t separated.

11.1.J. Less important Exercise. Show that an $A$-scheme is separated (over $A$) if and only if it is separated over $Z$. In particular, a complex scheme is separated over $C$ if and only if it is separated over $Z$, so complex geometers and arithmetic geometers can communicate about separated schemes without confusion.

11.1.K. Useful exercise: The locus where two morphisms agree. Suppose $f : X \to Y$ and $g : X \to Y$ are two morphisms over some scheme $Z$. We can now give meaning to the phrase ‘the locus where $f$ and $g$ agree’, and that in particular there is a largest locally closed subscheme where they agree — and even a closed immersion if $Y$ is separated over $Z$. Suppose $h : W \to X$ is some morphism (perhaps a locally closed immersion). We say that $f$ and $g$ agree on $h$ if $f \circ h = g \circ h$. Show that there is a locally closed subscheme $i : V \to X$ such that any morphism $h : W \to X$ on which $f$ and $g$ agree factors uniquely through $i$, i.e. there is a unique $j : W \to V$ such that $h = i \circ j$. Show further that if $Y \to Z$ is separated, then $i : V \to X$ is a closed immersion. Hint: define $V$ to be the following fibered product:

\[
\begin{array}{ccc}
V & \xrightarrow{(f,g)} & Y \\
\downarrow & & \downarrow \\
X & \xrightarrow{\delta} & Y \times Z Y. \\
\end{array}
\]

As $\delta$ is a locally closed immersion, $V \to X$ is too. Then if $h : W \to X$ is any scheme such that $g \circ h = f \circ h$, then $h$ factors through $V$.

Minor Remarks. 1) In the previous exercise, we are describing $V \to X$ by way of a universal property. Taking this as the definition, it is not a priori clear that $V$ is a locally closed subscheme of $X$, or even that it exists.

2) Warning: consider two maps from $\text{Spec} C$ to itself $\text{Spec} C$ over $\text{Spec} \mathbb{R}$, the identity and complex conjugation. These are both maps from a point to a point, yet they do not agree despite agreeing as maps of sets. (If you do not find this reasonable, this might help: after base change $\text{Spec} C \to \text{Spec} \mathbb{R}$, they do not agree as maps of sets.)
3) More generally, in the case of reduced finite type $k$-schemes, the locus where $f$ and $g$ agree can be interpreted as follows: $f$ and $g$ agree at $x$ if $f(x) = g(x)$ and the two maps of residue fields are the same.

11.1.L. LESS IMPORTANT EXERCISE. Show that the line with doubled origin $X$ (Example 5.4.5) is not separated, by finding two morphisms $f_1: W \to X$, $f_2: W \to X$ whose domain of agreement is not a closed subscheme (cf. Proposition 11.1.3). (Another argument was given above, in Exercise 11.1.D. A fancy argument will be given in Exercise 13.4.C.)

11.1.M. LESS IMPORTANT EXERCISE. Suppose $P$ is a class of morphisms such that closed immersions are in $P$, and $P$ is closed under fibered product and composition. Show that if $f: X \to Y$ is in $P$ then $f_{\text{red}}: X_{\text{red}} \to Y_{\text{red}}$ is in $P$. (Two examples are the classes of separated morphisms and quasiseparated morphisms.) Hint:

\[ X_{\text{red}} \to X \times_Y Y_{\text{red}} \to Y_{\text{red}} \]

11.2 Rational maps to separated schemes

When we introduced rational maps in §7.5 we promised that in good circumstances, a rational map has a “largest domain of definition”. We are now ready to make precise what “good circumstances” means.

11.2.1. Reduced-to-separated Theorem (important!). — Two $S$-morphisms $f_1: U \to Z$, $f_2: U \to Z$ from a reduced scheme to a separated $S$-scheme agreeing on a dense open subset of $U$ are the same.

Proof. Let $V$ be the locus where $f_1$ and $f_2$ agree. It is a closed subscheme of $U$ by Exercise 11.1.K, which contains a dense open set. But the only closed subscheme of a reduced scheme $U$ whose underlying set is dense is all of $U$. \[\square\]

11.2.2. Consequence 1. Hence (as $X$ is reduced and $Y$ is separated) if we have two morphisms from open subsets of $X$ to $Y$, say $f: U \to Y$ and $g: V \to Y$, and they agree on a dense open subset $Z \subset U \cap V$, then they necessarily agree on $U \cap V$.

Consequence 2. A rational map has a largest domain of definition on which $f: U \dashrightarrow Y$ is a morphism, which is the union of all the domains of definition. In particular, a rational function on a reduced scheme has a largest domain of definition. For example, the domain of definition of $\mathbb{A}^2_k \dashrightarrow \mathbb{P}^1_k$ given by $(x, y) \mapsto [x; y]$ has domain of definition $\mathbb{A}^2_k \setminus \{(0, 0)\}$ (cf. §7.5.3).

11.2.A. Exercise. Show that the Reduced-to-separated Theorem 11.2.1 is false if we give up reducedness of the source or separatedness of the target. Here are some possibilities. For the first, consider the two maps from $\text{Spec} \, k[x, y]/(y^2, xy)$ to $\text{Spec} \, k[t]$, where we take $f_1$ given by $t \mapsto x$ and $f_2$ given by $t \mapsto x + y$; $f_1$
and \( f_2 \) agree on the distinguished open set \( D(x) \), see Figure 11.4. For the second, consider the two maps from \( \text{Spec } k[t] \) to the line with the doubled origin, one of which maps to the “upper half”, and one of which maps to the “lower half”. These two morphisms agree on the dense open set \( D(f) \), see Figure 11.5.

\[ \begin{array}{ccc}
\ast & \rightarrow & \ast \\
n_1 & \downarrow & \downarrow & n_2 \\
& \ast & \ast 
\end{array} \]

**Figure 11.4.** Two different maps from a nonreduced scheme agreeing on a dense open set

\[ \begin{array}{ccc}
\ast & \rightarrow & \ast \\
n_1 & \downarrow & \downarrow & n_2 \\
& \ast & \ast 
\end{array} \]

**Figure 11.5.** Two different maps to a nonseparated scheme agreeing on a dense open set

11.2.3. **Proposition.** — Suppose \( Y \) and \( Z \) are integral separated schemes. Then \( Y \) and \( Z \) are birational if and only if there is a dense (=non-empty) open subscheme \( U \) of \( Y \) and a dense open subscheme \( V \) of \( Z \) such that \( U \cong V \).

This gives you a good idea of how to think of birational maps. For example, a variety is rational if it has a dense open subset isomorphic to a subset \( \mathbb{A}^n \).

**Proof.** I find this proof surprising and unexpected. Is there a better way to explain it?

Clearly if \( Y \) and \( Z \) have isomorphic open sets \( U \) and \( V \) respectively, then they are birational (with birational maps given by the isomorphisms \( U \rightarrow V \) and \( V \rightarrow U \) respectively).

For the other direction, assume that \( f : Y \rightarrow Z \) is a birational map, with inverse birational map \( g : Z \rightarrow Y \). Choose representatives for these rational maps \( F : W \rightarrow Z \) (where \( W \) is an open subscheme of \( Y \)) and \( G : X \rightarrow Y \) (where \( X \) is an
open subscheme of \( Z \). We will see that \( F^{-1}(G^{-1}(W)) \subset Y \) and \( G^{-1}(F^{-1}(X)) \subset Z \) are isomorphic open subschemes.

The key observation is that the two morphisms

\[
\begin{array}{ccc}
F^{-1}(G^{-1}(W)) & \to & G^{-1}(F^{-1}(X)) \\
\downarrow & & \downarrow \\
F^{-1}(X) & \to & G^{-1}(W) \\
\downarrow & & \downarrow \\
W & \to & X \\
\downarrow & & \downarrow \\
Y & \to & Z
\end{array}
\]

are two maps of varieties over \( k \), such that \( f \) represents the same rational map, so by the Reduced-to-separated Theorem\( \text{[11.2.1]} \) they are the same morphism on \( F^{-1}(G^{-1}(W)) \). Thus \( G \circ F \) gives the identity map from \( F^{-1}(G^{-1}(W)) \) to itself. Similarly \( F \circ G \) gives the identity map on \( G^{-1}(F^{-1}(X)) \).

All that remains is to show that \( F \) maps \( F^{-1}(G^{-1}(W)) \) into \( G^{-1}(F^{-1}(X)) \), and that \( G \) maps \( G^{-1}(F^{-1}(X)) \) into \( F^{-1}(G^{-1}(W)) \), and by symmetry it suffices to show the former. Suppose \( q \in F^{-1}(G^{-1}(W)) \). Then \( F(G(F(q))) = F(q) \in X \), from which \( F(q) \in G^{-1}(F^{-1}(X)) \). (Another approach is to note that each “parallelogram” in the diagram above is a fibered diagram, and to use the key observation of the previous paragraph to construct a morphism \( G^{-1}(F^{-1}(X)) \to F^{-1}(G^{-1}(X)) \) and vice versa, and showing that they are inverses.)

\[ \square \]

11.2.B. Exercise: Maps of varieties are determined by the maps on closed points. Suppose \( f_1 : X \to Y \) and \( f_2 : X \to Y \) are two maps of varieties over \( k \), such that \( f_1(p) = f_2(p) \) for all closed points. Show that \( f_1 = f_2 \). (This implies that the functor from the category of “classical varieties over \( k \)”, which we won’t define here, to the category of \( k \)-schemes, is fully faithful.)

As noted in \( \text{[7.5.2]} \) rational maps can be defined from any \( X \) that has associated points to any \( Y \). The Reduced-to-separated Theorem\( \text{[11.2.1]} \) can be extended to this setting, as follows.

11.2.4. Associated-to-separated Theorem. — Two \( S \)-morphisms \( f_1 : U \to Z \) and \( f_2 : U \to Z \) from a locally Noetherian scheme \( X \) to a separated \( S \)-scheme, agreeing on a dense open subset of \( U \) containing the associated points of \( X \), are the same.

11.2.C. Exercise. Adjust the proof of the Reduced-to-separated Theorem\( \text{[11.2.1]} \) to prove the Associated-to-separated Theorem\( \text{[11.2.4]} \)

11.3 Proper morphisms

Recall that a map of topological spaces (also known as a continuous map!) is said to be proper if the preimage of any compact set is compact. Properness of
morphisms is an analogous property. For example, a variety over $\mathbb{C}$ will be proper if it is compact in the classical topology. Alternatively, we will see that projective $A$-schemes are proper over $A$ — this is the hardest thing we will prove — so you can see this as a nice property satisfied by projective schemes, and quite convenient to work with.

Recall (§8.3.6) that a (continuous) map of topological spaces $f : X \to Y$ is closed if for each closed subset $S \subset X$, $f(S)$ is also closed. A morphism of schemes is closed if the underlying continuous map is closed. We say that a morphism of schemes $f : X \to Y$ is universally closed if for every morphism $g : Z \to Y$, the induced morphism $Z \times_Y X \to Z$ is closed. In other words, a morphism is universally closed if it remains closed under any base change. (More generally, if $P$ is some property of schemes, then a morphism of schemes is said to be universally $P$ if it remains $P$ under any base change.)

To motivate the definition of properness, we remark that a map $f : X \to Y$ of locally compact Hausdorff spaces which have countable bases for their topologies is universally closed if and only if it is proper in the usual topology. (You are welcome to prove this as an exercise.)

11.3.1. Definition. A morphism $f : X \to Y$ is proper if it is separated, finite type, and universally closed. A scheme $X$ is often said to be proper if some implicit structure morphism is proper. For example, a $k$-scheme $X$ is often described as proper if $X \to \text{Spec} \ k$ is proper. (A $k$-scheme is often said to be complete if it is proper. We will not use this terminology.)

Let’s try this idea out in practice. We expect that $\mathbb{A}^1_{\mathbb{C}} \to \text{Spec} \ \mathbb{C}$ is not proper, because the complex manifold corresponding to $\mathbb{A}^1_{\mathbb{C}}$ is not compact. However, note that this map is separated (it is a map of affine schemes), finite type, and (trivially) closed. So the “universally” is what matters here.

11.3.A. Exercise. Show that $\mathbb{A}^1_{\mathbb{C}} \to \text{Spec} \ \mathbb{C}$ is not proper, by finding a base change that turns this into a non-closed map. (Hint: Consider a well-chosen map $\mathbb{A}^1_{\mathbb{C}} \times \mathbb{A}^1_{\mathbb{C}} \to \mathbb{A}^1_{\mathbb{C}}$ or $\mathbb{A}^1_{\mathbb{C}} \times \mathbb{P}^1_{\mathbb{C}} \to \mathbb{P}^1_{\mathbb{C}}$.)

11.3.2. As a first example: closed immersions are proper. They are clearly separated, as affine morphisms are separated, §11.1.12. They are finite type. After base change, they remain closed immersions, and closed immersions are always closed. This easily extends further as follows.

11.3.3. Proposition. — Finite morphisms are proper.

Proof. Finite morphisms are separated (as they are affine by definition, and affine morphisms are separated, §11.1.12), and finite type (basically because finite modules over a ring are automatically finitely generated). To show that finite morphism are closed after any base change, we note that they remain finite after any base change (finiteness is preserved by base change, Exercise 10.4.B(d)), and finite morphisms are closed (Exercise 8.3.N).

11.3.4. Proposition. —

(a) The notion of “proper morphism” is stable under base change.
(b) The notion of “proper morphism” is local on the target (i.e. \( f : X \to Y \) is proper if and only if for any affine open cover \( U_i \to Y \), \( f^{-1}(U_i) \to U_i \) is proper). Note that the “only if” direction follows from (a) — consider base change by \( U_i \).

(c) The notion of “proper morphism” is closed under composition.

(d) The product of two proper morphisms is proper: if \( f : X \to Y \) and \( g : X' \to Y' \) are proper, where all morphisms are morphisms of \( Z \)-schemes, then \( f \times g : X \times_Z X' \to Y \times_Z Y' \) is proper.

(e) Suppose

\[
\begin{array}{ccc}
X & \xrightarrow{f} & Y \\
\downarrow{g} & & \downarrow{h} \\
Z & \xrightarrow{\ } & \\
\end{array}
\]

is a commutative diagram, and \( g \) is proper, and \( h \) is separated. Then \( f \) is proper.

A sample application of (e): a morphism (over Spec \( k \)) from a proper \( k \)-scheme to a separated \( k \)-scheme is always proper.

**Proof.** (a) The notions of separatedness, finite type, and universal closedness are all preserved by fibered product. (Notice that this is why universal closedness is better than closedness — it is automatically preserved by base change!)

(b) We have already shown that the notions of separatedness and finite type are local on the target. The notion of closedness is local on the target, and hence so is the notion of universal closedness.

(c) The notions of separatedness, finite type, and universal closedness are all preserved by composition.

(d) By (a) and (c), this follows from Exercise 10.4.1

(e) Closed immersions are proper, so we invoke the Cancellation Theorem 11.1.19 for proper morphisms. \( \square \)

We now come to the most important example of proper morphisms.

**11.3.5. Theorem.** — **Projective \( A \)-schemes are proper over \( A \).**

(As finite morphisms to Spec \( A \) are projective \( A \)-schemes, Exercise 8.3.1, Theorem 11.3.5 can be used to give a second proof that finite morphisms are proper, Proposition 11.3.3.)

It is not easy to come up with an example of an \( A \)-scheme that is proper but not projective! We will see a simple example of a proper but not projective surface, later. Once we discuss blow-ups, we will see Hironaka’s example of a proper but not projective nonsingular (“smooth”) threefold over \( C \).

**Proof.** The structure morphism of a projective \( A \)-scheme \( X \to \text{Spec} \ A \) factors as a closed immersion followed by \( \mathbb{P}^n_A \to \text{Spec} \ A \). Closed immersions are proper, and compositions of proper morphisms are proper, so it suffices to show that \( \mathbb{P}^n_A \to \text{Spec} \ A \) is proper. We have already seen that this morphism is finite type (Easy Exercise 6.5.1) and separated (Prop. 11.1.5), so it suffices to show that \( \mathbb{P}^n_A \to \text{Spec} \ A \) is universally closed. As \( \mathbb{P}^n_A = \mathbb{P}^n_Z \times_Z \text{Spec} \ A \), it suffices to show that \( \mathbb{P}^n_X := \mathbb{P}^n_{\mathcal{X}} \times_{\mathcal{X}} X \to X \) is closed for any scheme \( X \). But the property of being closed is local on the target on \( X \), so by covering \( X \) with affine open subsets, it suffices to show that
$\mathbb{P}^n_A \to \text{Spec } A$ is closed. This is the Fundamental Theorem of Elimination Theory (Theorem 8.4.5).

11.3.6. Unproved facts that may help you correctly think about finiteness.

We conclude with some interesting facts that we will prove later. They may shed some light on the notion of finiteness.

A morphism is finite if and only if it is proper and affine, if and only if it is proper and quasifinite. We have verified the “only if” parts of this statement; the “if” parts are harder (and involve Zariski’s Main Theorem, cf. §8.3.12).

As an application: quasifinite morphisms from proper schemes to separated schemes are finite. Here is why: suppose $f : X \to Y$ is a quasifinite morphism over $Z$, where $X$ is proper over $Z$. Then by the Cancellation Theorem 11.1.19 for proper morphisms, $X \to Y$ is proper. Hence as $f$ is quasifinite and proper, $f$ is finite.

As an explicit example, consider the map $\pi : \mathbb{P}^1_k \to \mathbb{P}^1_k$ given by $[x:y] \mapsto [f(x,y); g(x,y)]$, where $f$ and $g$ are homogeneous polynomials of the same degree with no common roots in $\mathbb{P}^1$. The fibers are finite, and $\pi$ is proper (from the Cancellation Theorem 11.1.19 for properties of morphisms, as discussed after the statement of Theorem 11.3.4), so $\pi$ is finite. This could be checked directly as well, but now we can save ourselves the annoyance.
Part IV

Harder properties of schemes
CHAPTER 12

Dimension

12.1 Dimension and codimension

At this point, you know a fair bit about schemes, but there are some fundamental notions you cannot yet define. In particular, you cannot use the phrase “smooth surface”, as it involves the notion of dimension and of smoothness. You may be surprised that we have gotten so far without using these ideas. You may also be disturbed to find that these notions can be subtle, but you should keep in mind that they are subtle in all parts of mathematics.

In this chapter, we will address the first notion, that of dimension of schemes. This should agree with, and generalize, our geometric intuition. Although we think of dimension as a basic notion in geometry, it is a slippery concept, as it is throughout mathematics. Even in linear algebra, the definition of dimension of a vector space is surprising the first time you see it, even though it quickly becomes second nature. The definition of dimension for manifolds is equally nontrivial. For example, how do we know that there isn’t an isomorphism between some 2-dimensional manifold and some 3-dimensional manifold? Your answer will likely use topology, and hence you should not be surprised that the notion of dimension is often quite topological in nature.

A caution for those thinking over the complex numbers: our dimensions will be algebraic, and hence half that of the “real” picture. For example, we will see very shortly that $\mathbb{A}^1_C$, which you may picture as the complex numbers (plus one generic point), has dimension 1.

12.1.1. Definition(s): dimension. Surprisingly, the right definition is purely topological — it just depends on the topological space, and not on the structure sheaf. We define the dimension of a topological space $X$ (denoted $\dim X$) as the supremum of lengths of chains of closed irreducible sets, starting the indexing with 0. (The dimension may be infinite.) Scholars of the empty set can take the dimension of the empty set to be $-\infty$. Define the dimension of a ring as the Krull dimension of its spectrum — the supremum of the lengths of the chains of nested prime ideals (where indexing starts at zero). These two definitions of dimension are sometimes called Krull dimension. (You might think a Noetherian ring has finite dimension because all chains of prime ideals are finite, but this isn’t necessarily true — see Exercise 12.1.F)

As we have a natural homeomorphism between Spec $A$ and Spec $A/\mathfrak{m}(A)$ ($\mathfrak{m}$ the Zariski topology disregards nilpotents), we have $\dim A = \dim A/\mathfrak{m}(A)$. 239
Examples. We have identified all the prime ideals of $k[t]$ (they are $0$, and $(f(t))$ for irreducible polynomials $f(t)$, $\mathbb{Z}$ ($(0)$ and $(p)$), $k$ (only $0$), and $k[x]/(x^2)$ (only $0$), so we can quickly check that $\dim \mathbb{A}^1_k = \dim \text{Spec } \mathbb{Z} = 1$, $\dim \text{Spec } k = 0$, $\dim \text{Spec } k[x]/(x^2) = 0$.

We must be careful with the notion of dimension for reducible spaces. If $Z$ is the union of two closed subsets $X$ and $Y$, then $\dim Z = \max(\dim X, \dim Y)$. Thus dimension is not a "local" characteristic of a space. This sometimes bothers us, so we try to only talk about dimensions of irreducible topological spaces. If a topological space can be expressed as a finite union of irreducible subsets, then we say that it is equidimensional or pure dimensional (resp. equidimensional of dimension $n$ or pure dimension $n$) if each of its components has the same dimension (resp. they are all of dimension $n$).

An equidimensional dimension 1 (resp. 2, $n$) topological space is said to be a curve (resp. surface, $n$-fold).

12.1.A. IMPORTANT EXERCISE. Show that if $f : \text{Spec } A \to \text{Spec } B$ corresponds to an integral extension of rings, then $\dim \text{Spec } A = \dim \text{Spec } B$. Hint: show that a chain of prime ideals downstairs gives a chain upstairs of the same length, by the Going-up Theorem (Exercise 8.2.F). Conversely, a chain upstairs gives a chain downstairs. We need to check that no two elements of the chain upstairs goes to the same element $[q] \in \text{Spec } B$ of the chain downstairs. As integral extensions are well-behaved by localization and quotients of Spec $B$ (Exercise 8.2.B), we can replace $B$ by $B_q/qB_q$ (and $A$ by $A \otimes_B (B_q/qB_q)$). Thus we can assume $B$ is a field. Hence we must show that if $\phi : k \to A$ is an integral extension, then $\dim A = 0$. Outline of proof: Suppose $p \subset m$ are two prime ideals of $A$. Mod out by $p$, so we can assume that $A$ is a domain. I claim that any non-zero element is invertible: Say $x \in A$, and $x \neq 0$. Then the minimal monic polynomial for $x$ has non-zero constant term. But then $x$ is invertible — recall the coefficients are in a field.

12.1.B. EXERCISE. Show that if $\tilde{X} \to X$ is the normalization of a scheme (possibly in a finite field extension), then $\dim \tilde{X} = \dim X$.

12.1.C. EXERCISE. Show that $\dim \mathbb{Z}[x] = 2$. (Hint: The primes of $\mathbb{Z}[x]$ were implicitly determined in Exercise 1.2.1N)

12.1.2. Codimension. Because dimension behaves oddly for disjoint unions, we need some care when defining codimension, and in using the phrase. For example, if $Y$ is a closed subset of $X$, we might define the codimension to be $\dim X - \dim Y$, but this behaves badly. For example, if $X$ is the disjoint union of a point $Y$ and a curve $Z$, then $\dim X - \dim Y = 1$, but this has nothing to do with the local behavior of $X$ near $Y$.

A better definition is as follows. In order to avoid excessive pathology, we define the codimension of $Y$ in $X$ only when $Y$ is irreducible. (Use extreme caution when using this word in any other setting.) Define the codimension of an irreducible closed subset $Y \subset X$ of a topological space as the supremum of lengths of increasing chains of irreducible closed subsets starting with $Y$ (where indexing starts at 0). So the codimension of a point is the codimension of its closure.

We say that a prime ideal $p$ in a ring has codimension (denoted $\text{codim}$) equal to the supremum of lengths of the chains of decreasing prime ideals starting at $p$,
with indexing starting at 0. Thus in an integral domain, the ideal \((0)\) has codimension 0; and in \(\mathbb{Z}\), the ideal \((23)\) has codimension 1. Note that the codimension of the prime ideal \(p\) in \(A\) is \(\dim A_p\) (see \[12.6\]). (This notion is often called \textbf{height}.) Thus the codimension of \(p\) in \(A\) is the codimension of \([p]\) in \(\text{Spec} A\).

**12.1.D. Exercise.** Show that if \(Y\) is an irreducible closed subset of a scheme \(X\) with generic point \(y\), then the codimension of \(Y\) is the dimension of the local ring \(O_{X,y}\) (cf. \[12.6\]).

Notice that \(Y\) is codimension 0 in \(X\) if it is an irreducible component of \(X\). Similarly, \(Y\) is codimension 1 if it is strictly contained in an irreducible component \(Y'\), and there is no irreducible subset strictly between \(Y\) and \(Y'\). (See Figure 12.1 for examples.) An closed subset all of whose irreducible components are codimension 1 in some ambient space \(X\) is said to be a \textbf{hypersurface} in \(X\).

**Figure 12.1.** Behavior of codimension

**12.1.E. Easy Exercise.** Show that

\[(12.1.2.1) \quad \text{codim}_X Y + \dim Y \leq \dim X.\]

We will soon see that equality always holds if \(X\) and \(Y\) are varieties (Exercise \[12.2.D\]), but equality doesn’t hold in general (\[12.3.8\]).

\textbf{Warning.} The notion of codimension still can behave slightly oddly. For example, consider Figure 12.1 (You should think of this as an intuitive sketch.) Here the total space \(X\) has dimension 2, but point \(p\) is dimension 0, and codimension 1. We also have an example of a codimension 2 subset \(q\) contained in a codimension 0 subset \(C\) with no codimension 1 subset “in between”.

Worse things can happen; we will soon see an example of a closed point in an \textit{irreducible} surface that is nonetheless codimension 1, not 2, in \[12.3.8\]. However, for
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irreducible varieties this can’t happen, and inequality [12.1.2.1] must be an equality (Proposition [12.2.1]).

12.1.3. A fun but unimportant counterexample. We end this introductory section with a fun pathology. As a Noetherian ring has no infinite chain of prime ideals, you may think that Noetherian rings must have finite dimension. Nagata, the master of counterexamples, shows you otherwise with the following example.

12.1.F. ★ Exercise: An infinite-dimensional Noetherian ring. Let \( A = k[x_1, x_2, \ldots] \). Choose an increasing sequence of positive integers \( m_1, m_2, \ldots \) whose differences are also increasing (\( m_{i+1} - m_i > m_i - m_{i-1} \)). Let \( p_i = (x_{m_i+1}, \ldots, x_{m_{i+1}}) \) and \( S = A - \cup_i p_i \). Show that \( S \) is a multiplicative set. Show that \( S^{-1}A \) is Noetherian. Show that each \( S^{-1}p \) is the smallest prime ideal in a chain of prime ideals of length \( m_{i+1} - m_i \). Hence conclude that \( \dim S^{-1}A = \infty \).

12.2 Dimension, transcendence degree, and Noether normalization

We now prove a powerful alternative interpretation for dimension for irreducible varieties, in terms of transcendence degree. In case you haven’t seen transcendence theory, here is a lightning introduction.

12.2.A. Exercise/Definition. An element of a field extension \( E/F \) is algebraic over \( F \) if it is integral over \( F \). A field extension is algebraic if it is integral. The composition of two algebraic extensions is algebraic, by Exercise [8.2.C]. If \( E/F \) is a field extension, and \( F' \) and \( F'' \) are two intermediate field extensions, then we write \( F' \sim F'' \) if \( F'/F'' \) is algebraic over both \( F' \) and \( F'' \). Here \( F'/F'' \) is the compositum of \( F' \) and \( F'' \), the smallest field extension in \( E \) containing \( F' \) and \( F'' \). (a) Show that \( \sim \) is an equivalence relation on subextensions of \( E/F \). A transcendence basis of \( E/F \) is a set of elements \( \{x_i\} \) that are algebraically independent over \( F \) (there is no nontrivial polynomial relation among the \( x_i \) with coefficients in \( F \)) such that \( F(\{x_i\}) \sim E \). (b) Show that if \( E/F \) has two transcendence bases, and one has cardinality \( n \), then both have cardinality \( n \). (Hint: show that you can substitute elements from the one basis into the other one at a time.) The size of any transcendence basis is called the transcendence degree (which may be \( \infty \)), and is denoted \( \text{tr.deg} \). Any finitely generated field extension necessarily has finite transcendence degree.

12.2.1. Theorem (dimension = transcendence degree). — Suppose \( A \) is a finitely-generated integral domain over a field \( k \). Then \( \dim \text{Spec} A = \text{tr.deg} K(A)/k \).

By “finitely generated domain over \( k \)”, we mean “a finitely generated \( k \)-algebra that is an integral domain”.

We will prove Theorem [12.2.1] shortly (12.2.10). But we first show that it is useful by giving some immediate consequences. We seem to have immediately \( \dim \mathbb{A}^n_k = n \). However, our proof of Theorem [12.2.1] will go through this fact, so it isn’t really a Corollary. Instead, we begin with a proof of the Nullstellensatz, promised earlier.
12.2.B. Exercise: Nullstellensatz from dimension theory. Prove Hilbert’s Nullstellensatz 4.2.3: Suppose \( A = k[x_1, \ldots, x_n]/I \). Show that the residue field of any maximal ideal of \( A \) is a finite extension of \( k \). (Hint: the maximal ideals correspond to dimension 0 points, which correspond to transcendence degree 0 extensions of \( k \), i.e. finite extensions of \( k \).)

For further applications, we make a short observation.

12.2.2. Lemma. — In a unique factorization domain \( A \), all codimension 1 prime ideals are principal.

We will see that the converse (when \( A \) is a Noetherian integral domain) holds as well (Proposition 12.3.5).

Proof. Suppose \( p \) is a codimension 1 prime. Choose any \( f \neq 0 \) in \( p \), and let \( g \) be any irreducible/prime factor of \( f \) that is in \( p \) (there is at least one). Then \( (g) \) is a prime ideal contained in \( p \), so \( (0) \subset (g) \subset p \). As \( p \) is codimension 1, we must have \( p = (g) \), and thus \( p \) is principal. \( \square \)

12.2.3. Points of \( A_k^2 \). We can find a second proof that we have named all the primes of \( k[x, y] \) where \( k \) is algebraically closed (promised in Exercise 4.2.D when \( k = \mathbb{C} \)). Recall that we have discovered the primes \( (0) \), \( f(x, y) \) where \( f \) is irreducible, and \( (x - a, y - b) \) where \( a, b \in k \). As \( A_k^2 \) is irreducible, there is only one irreducible subset of codimension 0. By Lemma 12.2.2, all codimension 1 primes are principal. By inequality (12.1.2.1), there are no primes of codimension greater than 2, and any prime of codimension 2 must be maximal. We have identified all the maximal ideals of \( k[x, y] \) by the Nullstellensatz.

12.2.C. Exercise. Suppose \( X \) is an irreducible variety. Show that \( \dim X \) is the transcendence degree of the function field (the stalk at the generic point) \( O_{X, \eta} \) over \( k \). Thus (as the generic point lies in all non-empty open sets) the dimension can be computed in any open set of \( X \). (This is not true in general, see 12.3.8)

12.2.D. Exercise. Suppose \( Y \subset X \) is an inclusion of irreducible \( k \)-varieties, and \( \eta \) is the generic point of \( Y \). Show that \( \dim Y + \dim O_{X, \eta} = \dim X \). Hence by Exercise 12.1.D \( \dim Y + \text{codim}_X Y = \dim X \). Thus for varieties, the inequality (12.1.2.1) is always an equality.

12.2.E. Exercise. Show that the equations \( wz - xy = 0 \), \( wy - x^2 = 0 \), \( xz - y^2 = 0 \) cut out an integral surface \( S \) in \( A_k^4 \). You may recognize these equations from Exercises 1.6.F and 9.2.A. You might expect \( S \) to be a curve, because it is cut out by three equations in 4-space. One of many ways to proceed: cut \( S \) into pieces. For example, show that \( D(w) \cong \text{Spec} k[x, w]_{w^0} \). (You may recognize \( S \) as the affine cone over the twisted cubic. The twisted cubic was defined in Exercise 9.2.A.) It turns out that you need three equations to cut out this surface. The first equation cuts out a threefold in \( A_k^4 \) (by Krull’s Principal Ideal Theorem 12.3.3 which we will meet soon). The second equation cuts out a surface: our surface, along with another surface. The third equation cuts out our surface, and removes the “extraneous component”. One last aside: notice once again that the cone over the quadric surface \( k[w, x, y, z]/(wz - xy) \) makes an appearance.)
12.2.4. **A first example of the utility of dimension theory.** Although dimension theory is not central to the following statement, it is essential to the proof.

12.2.5. **Enlightening strenuous exercise.** For any \(d > 3\), show that most degree \(d\) surfaces in \(\mathbb{P}^3_k\) contain no lines. Here, “most” means “all closed points of a Zariski-open subset of the parameter space for degree \(d\) homogeneous polynomials in 4 variables, up to scalars. As there are \(\binom{d+3}{3}\) such monomials, the degree \(d\) hypersurfaces are parametrized by \(\mathbb{P}^{\binom{d+3}{3}-1}_k\). Hint: Construct an incidence correspondence

\[
X = \{([\ell], H) : [\ell] \in G(1, 3), [H] \in \mathbb{P}^{\binom{d+3}{3}-1}_k, \ell \subset H\},
\]

parametrizing lines in \(\mathbb{P}^3\) contained in a hypersurface: define a closed subscheme \(X\) of \(G(1, 3) \times \mathbb{P}^{\binom{d+3}{3}-1}_k\) that makes this notion precise. Show that \(X\) is a \(\mathbb{P}^{\binom{d+3}{3}-1-(d+1)}_k\) bundle over \(G(1, 3)\). (Possible hint for this: how many degree \(d\) hypersurfaces contain the line \(x = y = 0\)?) Show that \(\dim G(1, 3) = 4\) (see \(\S 7.7\) : \(G(1, 3)\) is covered by \(A^4\)’s). Show that \(\dim X = \binom{d+3}{3} - 1 - (d + 1) + 4\). Show that the image of the projection \(X \to \mathbb{P}^{d+33} - 1\) must lie in a proper closed subset. The following diagram may help.

\[
\text{dim} \binom{d+3}{3} - 1 - (d + 1) + 4 \\
\begin{array}{ccc}
\mathbb{P}^{\binom{d+3}{3}-1} & \overset{p^{\binom{d+3}{3}-1-(d+1)}}{\longrightarrow} & G(1, 3) \\
X \quad & \quad & \quad \text{dim} 4
\end{array}
\]

12.2.6. **Side Remark.** If you do the previous Exercise, your dimension count will suggest the true facts that degree 1 hypersurfaces — i.e. hyperplanes — have 2-dimensional families of lines, and that most degree 2 hypersurfaces have 1-dimensional families of lines, as shown in Exercise 9.2.N. They will also suggest that most degree 3 hypersurfaces contain a finite number of lines, which reflects the celebrated fact that nonsingular cubic surfaces over an algebraically closed field always contain 27 lines.) The statement about quartics generalizes to the Noether-Lefschetz theorem implying that a very general surface of degree \(d\) at least 4 contains no curves that are not the intersection of the surface with a hypersurface. “**Very general**” means that in the parameter space (in this case, the projective space parametrizing surfaces of degree \(d\)), the statement is true away from a countable union of proper Zariski-closed subsets. It is a weaker version of the phrase “almost every” than “general”.

12.2.7. **Noether Normalization.**

To set up the proof of Theorem 12.2.1 on dimension and transcendence degree, we introduce another important classical notion, Noether Normalization.

12.2.7. **Noether Normalization Lemma.** — Suppose \(A\) is an integral domain, finitely generated over a field \(k\). If \(\text{tr. deg}_k K(A) = n\), then there are elements \(x_1, \ldots, x_n \in A\),
algebraically independent over $k$, such that $A$ is a finite (hence integral by Corollary 8.2.2) extension of $k[x_1, \ldots, x_n]$.

The geometric content behind this result is that given any integral affine $k$-scheme $X$, we can find a surjective finite morphism $X \to \mathbb{A}^n_k$, where $n$ is the transcendence degree of the function field of $X$ (over $k$). Surjectivity follows from the Lying Over Theorem 8.2.5 in particular Exercise [12.1.A].

* Nagata's proof of Noether normalization. Suppose we can write $A = k[y_1, \ldots, y_m]/p$, i.e. that $A$ can be chosen to have $m$ generators. Note that $m \geq n$. We show the result by induction on $m$. The base case $m = n$ is immediate.

Assume now that $m > n$, and that we have proved the result for smaller $m$. We will find $m - 1$ elements $z_1, \ldots, z_{m-1}$ of $A$ such that $A$ is finite over $A' := k[z_1, \ldots, z_{m-1}]$ (i.e. the subring of $A$ generated by $z_1, \ldots, z_{m-1}$). Then by the inductive hypothesis, $A'$ is finite over some $k[x_1, \ldots, x_n]$, and $A$ is finite over $A'$, so by Exercise [8.3.1] $A$ is finite over $k[x_1, \ldots, x_n]$.

\[
\begin{array}{c}
A \\
\text{finite} \\
A' = k[z_1, \ldots, z_{m-1}]/p \\
\text{finite} \\
k[x_1, \ldots, x_n]
\end{array}
\]

As $y_1, \ldots, y_m$ are algebraically dependent, there is some non-zero algebraic relation $f(y_1, \ldots, y_m) = 0$ among them (where $f$ is a polynomial in $m$ variables).

Let $z_1 = y_1 - y_m^{r_1}$, $z_2 = y_2 - y_m^{r_2}$, \ldots, $z_{m-1} = y_{m-1} - y_m^{r_{m-1}}$, where $r_1, \ldots, r_{m-1}$ are positive integers to be chosen shortly. Then

\[
f(z_1 + y_m^{r_1}, z_2 + y_m^{r_2}, \ldots, z_{m-1} + y_m^{r_{m-1}}, y_m) = 0.
\]

Then upon expanding this out, each monomial in $f$ (as a polynomial in $m$ variables) will yield a single term in that is a constant times a power of $y_m$ (with no $z_i$ factors). By choosing the $r_i$ so that $0 \ll r_1 \ll r_2 \ll \cdots \ll r_{m-1}$, we can ensure that the powers of $y_m$ appearing are all distinct, and so that in particular there is a leading term $y_m^N$, and all other terms (including those with $z_i$-factors) are of smaller degree in $y_m$. Thus we have described an integral dependence of $y_m$ on $z_1, \ldots, z_{m-1}$ as desired. 

\[ \square \]

12.2.8. Geometric interpretations and consequences.

12.2.9. Aside: the geometry behind Nagata’s proof. Here is the geometric intuition behind Nagata’s argument. Suppose we have an $m$-dimensional variety in $\mathbb{A}^n_k$ with $m < n$, for example $xy = 1$ in $\mathbb{A}^2$. One approach is to hope the projection to a hyperplane is a finite morphism. In the case of $xy = 1$, if we projected to the $x$-axis, it wouldn’t be finite, roughly speaking because the asymptote $x = 0$ prevents the map from being closed (cf. Exercise [8.3.1]). If we instead projected to a random line, we might hope that we would get rid of this problem, and indeed we usually can: this problem arises for only a finite number of directions. But we might have a problem if the field were finite: perhaps the finite number of
directions in which to project each have a problem. (You can show that if \( k \) is an
infinite field, then the substitution in the above proof \( z_i = y_i - y_{ir} \) can be replaced
by the linear substitution \( z_i = y_i - a_i y_m \) where \( a_i \in k \), and that for a non-empty
Zariski-open choice of \( a_i \), we indeed obtain a finite morphism.) Nagata’s trick in
general is to “jiggle” the variables in a non-linear way, and this is jiggling kills the
non-finiteness of the map.

12.2.G. Exercise. Show that every dimension \( n \) irreducible variety over \( k \) is
birational to a hyperplane in \( \mathbb{A}^n_k \).

12.2.H. Exercise (Geometric Noether Normalization). If \( X \) is an affine
irreducible variety of dimension \( n \) over \( k \), show that there is a dominant finite
morphism \( X \to \mathbb{A}^n_k \) (over \( k \)).

12.2.I. Exercise (Dimension is additive for fibered products of finite type
\( k \)-schemes). Suppose \( X \) and \( Y \) are finite type \( k \)-schemes. Show that \( \dim X \times_k Y = \dim X + \dim Y \). (Hint: Use Noether normalization to find dominant finite
morphisms \( X \to \mathbb{A}^{\dim X}_k \) and \( Y \to \mathbb{A}^{\dim Y}_k \), and use this to construct a dominant
finite morphism \( X \times_k Y \to \mathbb{A}^{\dim X + \dim Y}_k \).

12.2.10. Proof of Theorem 12.2.1 on dimension and transcendence degree. Suppose \( X \) is
an integral affine \( k \)-scheme. We show that \( \dim X \) equals the transcendence degree
\( n \) of its function field, by induction on \( n \). (The idea is that we reduce from \( X \) to \( \mathbb{A}^n_k \)
to a hypersurface in \( \mathbb{A}^n_k \) to \( \mathbb{A}^{n-1}_k \).) Assume the result is known for all transcendence
degrees less than \( n \).

By Noether normalization, there exists a surjective finite morphism \( X \to \mathbb{A}^n_k \).
By Exercise 12.1.A, \( \dim X = \dim \mathbb{A}^n_k \). If \( n = 0 \), we are done, as \( \dim \mathbb{A}^0_k = 0 \).
We now show that \( \dim \mathbb{A}^n_k = n \) for \( n > 0 \), by induction. Clearly \( \dim \mathbb{A}^n_k \geq n \),
as we can describe a chain of irreducible subsets of length \( n + 1 \): if \( x_1, \ldots, x_n \) are
coordinates on \( \mathbb{A}^n_k \), consider the chain of ideals
\[
\{(0) \subset (x_1) \subset \cdots \subset (x_1, \ldots, x_n)\}
\]
in \( k[x_1, \ldots, x_n] \). Suppose we have a chain of prime ideals of length at least \( n 
\)
\[
\{0\} = p_0 \subset \cdots \subset p_m.
\]
where \( p_1 \) is a codimension 1 prime ideal. Then \( p_1 \) is principal (as \( k[x_1, \ldots, x_n] \) is a
unique factorization domain, Lemma 12.2.2) say \( p_1 = (f(x_1, \ldots, x_n)) \), where \( f \) is an
irreducible polynomial. Then \( K[k[x_1, \ldots, x_n]/(f(x_1, \ldots, x_n))] \) has transcendence
degree \( n - 1 \), so by induction,
\[
\dim k[x_1, \ldots, x_n]/(f) = n - 1.
\]

\( \Box \)

12.3 Codimension one miracles: Krull and Hartogs

In this section, we will explore a number of results related to codimension
one. We introduce two results that apply in more general situations, and link
functions and the codimension one points where they vanish, Krull’s Principal
Ideal Theorem 12.3.3 and Algebraic Hartogs' Lemma 12.3.10. We will find these two theorems very useful. For example, Krull’s Principal Ideal Theorem will help us compute codimensions, and will show us that codimension can behave oddly, and Algebraic Hartogs’ Lemma will give us a useful characterization of unique factorization domains (Proposition 12.3.5). The results in this section will require (locally) Noetherian hypotheses.

12.3.1. Krull’s Principal Ideal Theorem. The Principal Ideal Theorem generalizes the linear algebra fact that in a vector space, a single linear equation cuts out a subspace of codimension 0 or 1 (and codimension 0 occurs only when the equation is 0).

12.3.2. Krull’s Principal Ideal Theorem (geometric version). — Suppose $X$ is a locally Noetherian scheme, and $f$ is a function. The irreducible components of $V(f)$ are codimension 0 or 1. This is clearly a consequence of the following algebraic statement. You know enough to prove it for varieties (see Exercise 12.3.G), which is where we will use it most often. The full proof is technical, and included in §12.4 (see §12.4.2) only to show you that it isn’t long.

12.3.3. Krull’s Principal Ideal Theorem (algebraic version). — Suppose $A$ is a Noetherian ring, and $f \in A$. Then every prime $p$ minimal among those containing $f$ has codimension at most 1. If furthermore $f$ is not a zero-divisor, then every minimal prime $p$ containing $f$ has codimension precisely 1.

For example, the scheme $\text{Spec} \ k[w, x, y, z]/(wz - xy)$ (the cone over the quadric surface) is cut out by one non-zero equation $wz - xy$ in $\mathbb{A}^4$, so it is a threefold.

12.3.A. Exercise. What is the dimension of $\text{Spec} \ k[w, x, y, z]/(wz - xy, y^{17} + z^{17})$? (Check the hypotheses before invoking Krull!)

12.3.B. Exercise. Show that an irreducible homogeneous polynomial in $n + 1$ variables over a field $k$ describes an integral scheme of dimension $n - 1$ in $\mathbb{P}^n_k$.

12.3.C. Exercise (very important for later). This is a pretty cool argument. (a) (Hypersurfaces meet everything of dimension at least 1 in projective space, unlike in affine space.) Suppose $X$ is a closed subset of $\mathbb{P}^n_k$ of dimension at least 1, and $H$ is a nonempty hypersurface in $\mathbb{P}^n_k$. Show that $H$ meets $X$. (Hint: note that the affine cone over $H$ contains the origin in $\mathbb{A}^{n+1}_k$. Apply Krull’s Principal Ideal Theorem 12.3.3 to the cone over $X$.)

(b) Suppose $X \hookrightarrow \mathbb{P}^n_k$ is a closed subset of dimension $r$. Show that any codimension $r$ linear space meets $X$. Hint: Refine your argument in (a). (In fact any two things in projective space that you might expect to meet for dimensional reasons do in fact meet. We won’t prove that here.)

(c) Show further that there is an intersection of $r + 1$ nonempty hypersurfaces missing $X$. (The key step: show that there is a hypersurface of sufficiently high degree that doesn’t contain every generic point of $X$. Show this by induction on the number of generic points. To get from $n$ to $n + 1$: take a hypersurface not vanishing on $p_1, \ldots, p_n$. If it doesn’t vanish on $p_{n+1}$, we are done. Otherwise, call this hypersurface $f_{n+1}$. Do something similar with $n + 1$ replaced by $i$ ($1 \leq i \leq n$). Then
consider $\sum_i f_1 \cdots f_i \cdots f_{n+1}$.

(d) If $k$ is an infinite field, show that there is an intersection of $r$ hyperplanes meeting $X$ in a finite number of points. (We will see in Exercise [22.6.C] that if $k = \mathbb{K}$, the number of points for “most” choices of these $r$ hyperplanes, the number of points is the degree of $X$. But first of course we must define “degree”.)

12.3.D. EXERCISE (PRIME AVOIDANCE). As an aside, here is an exercise of a similar flavor to the previous one. Suppose $I \subseteq \bigcup_{i=1}^n p_i$. (The right side is not an ideal!) Show that $I \subset p_i$ for some $i$. (Can you give a geometric interpretation of this result?) Hint: by induction on $n$. Don’t look in the literature — you might find a much longer argument! (See Exercise [12.3.C] for a related problem.)

12.3.E. USEFUL EXERCISE. Suppose $f$ is an element of a Noetherian ring $A$, contained in no codimension 1 primes. Show that $f$ is a unit. (Hint: show that if a function vanishes nowhere, it is a unit.)

12.3.4. A useful characterization of unique factorization domains.

We can use Krull’s Principal Ideal Theorem to prove one of the four useful criteria for unique factorization domains, promised in [6.4.5].

12.3.5. Proposition. — Suppose that $A$ is a Noetherian integral domain. Then $A$ is a unique factorization domain if and only if all codimension 1 primes are principal.

This contains Lemma [12.2.2] and (in some sense) its converse.

Proof. We have already shown in Lemma [12.2.2] that if $A$ is a unique factorization domain, then all codimension 1 primes are principal. Assume conversely that all codimension 1 primes of $A$ are principal. I claim that the generators of these ideals are irreducible, and that we can uniquely factor any element of $A$ into these irreducibles, and a unit. First, suppose $(f)$ is a codimension 1 prime ideal $p$. Then if $f = gh$, then either $g \in p$ or $h \in p$. As codim $p > 0$, $p \neq (0)$, so by Nakayama’s Lemma [5.2.11] (as $p$ is finitely generated), $p \neq p^2$. Thus $g$ and $h$ cannot both be in $p$. Say $g \notin p$. Then $g$ is contained in no codimension 1 primes (as $f$ was contained in only one, namely $p$), and hence is a unit by Exercise [12.3.E].

We next show that any non-zero element $f$ of $A$ can be factored into irreducibles. Now $V(f)$ is contained in a finite number of codimension 1 primes, as $(f)$ has a finite number of associated primes ([6.5]), and hence a finite number of minimal primes. We show that any nonzero $f$ can be factored into irreducibles by induction on the number of codimension 1 primes containing $f$. In the base case where there are none, then $f$ is a unit by Exercise [12.3.E]. For the general case where there is at least one, say $f \in p = (g)$. Then $f = g^n h$ for some $h \notin (g)$. (Reason: otherwise, we have an ascending chain of ideals $(f) \subset (f/g) \subset (f/g^2) \subset \cdots$, contradicting Noetherianess.) Thus $f/g^n \in A$, and is contained in one fewer codimension 1 primes.

12.3.F. EXERCISE. Conclude the proof by showing that this factorization is unique. (Possible hint: the irreducible components of $V(f)$ give you the prime factors, but not the multiplicities.)
12.3.6. Generalizing Krull to more equations. The following generalization of Krull’s Principal Ideal Theorem looks like it might follow by induction from Krull, but it is more subtle.

12.3.7. Theorem. — Suppose $X = \text{Spec } A$ where $A$ is Noetherian, and $Z$ is an irreducible component of $V(r_1, \ldots, r_n)$, where $r_1, \ldots, r_n \in A$. Then the codimension of $Z$ is at most $n$.

A proof is given in §12.4.3. But you already know enough to prove it for varieties:

12.3.G. Exercise. — Prove Theorem 12.3.7 in the special case where $X$ is an affine variety, i.e. if $A$ is finitely generated over some field $k$. Show that $\dim Z \geq \dim X - n$. Hint: Exercise 12.2.D.

12.3.8. * Pathologies of the notion of “codimension”. We can use Krull’s Principal Ideal Theorem to produce the example of pathology promised earlier this chapter. Let $A = k[x]_{(x)}[t]$. In other words, elements of $A$ are polynomials in $t$, whose coefficients are quotients of polynomials in $x$, where no factors of $x$ appear in the denominator. (Warning: $A$ is not $k[x, t]_{(x)}$.) Clearly, $A$ is an integral domain, and $(xt - 1)$ is not a zero divisor. You can verify that $A/(xt - 1) \cong k[x]_{(x)}[1/x] \cong k(x)$ — “in $k[x]_{(x)}$, we may divide by everything but $x$, and now we are allowed to divide by $x$ as well” — so $A/(xt - 1)$ is a field. Thus $(xt - 1)$ is not just prime but also maximal. By Krull’s theorem, $(xt - 1)$ is codimension 1. Thus $(0) \subset (xt - 1)$ is a maximal chain. However, $A$ has dimension at least 2: $(0) \subset (t) \subset (x, t)$ is a chain of primes of length 2. (In fact, $A$ has dimension precisely 2, although we don’t need this fact in order to observe the pathology.) Thus we have a codimension 1 prime in a dimension 2 ring that is dimension 0.

Here is a picture of this poset of ideals.

```
(x, t)
   |
  (t)     (xt - 1)
   |
(0)
```

This example comes from geometry, and it is enlightening to draw a picture, see Figure 12.2. Spec $k[x]_{(x)}$ corresponds to a “germ” of $\mathbb{A}^1_k$ near the origin, and Spec $k[x]_{(x)}[t]$ corresponds to “this $\times$ the affine line”. You may be able to see from the picture some motivation for this pathology — $V(xt - 1)$ doesn’t meet $V(x)$, so it can’t have any specialization on $V(x)$, and there is nowhere else for $V(xt - 1)$ to specialize. It is disturbing that this misbehavior turns up even in a relatively benign-looking ring.

12.3.H. Unimportant exercise. — Show that it is false that if $X$ is an integral scheme, and $U$ is a non-empty open set, then $\dim U = \dim X$.

12.3.9. Algebraic Hartogs’ Lemma for Noetherian normal schemes.
Hartogs’ Lemma in several complex variables states (informally) that a holomorphic function defined away from a codimension two set can be extended over that. We now describe an algebraic analog, for Noetherian normal schemes. We will use this repeatedly and relentlessly when connecting line bundles and divisors.

**12.3.10. Algebraic Hartogs’ Lemma.** — Suppose $A$ is a Noetherian normal integral domain. Then

$$A = \cap \text{codimension 1 } \mathcal{A}_p.$$

The equality takes place inside $K(A)$; recall that any localization of an integral domain $A$ is naturally a subset of $K(A)$ (Exercise 2.3.C). Warning: few people call this Algebraic Hartogs’ Lemma. I call it this because it parallels the statement in complex geometry. The proof is technical and the details are less enlightening, so we postpone it to §12.3.11.

One might say that if $f \in K(A)$ does not lie in $A_p$ where $p$ has codimension 1, then $f$ has a pole at $[p]$, and if $f \in \text{FF}(A)$ lies in $pA_p$ where $p$ has codimension 1, then $f$ has a zero at $[p]$. It is worth interpreting Algebraic Hartogs’ Lemma as saying that a rational function on a normal scheme with no poles is in fact regular (an element of $A$). Informally: “Noetherian normal schemes have the Hartogs property.” (We will properly define zeros and poles in §13.3.7; see also Exercise 13.3.H.)

One can state Algebraic Hartogs’ Lemma more generally in the case that Spec $A$ is a Noetherian normal scheme, meaning that $A$ is a product of Noetherian normal integral domains; the reader may wish to do so.
12.3.11. **Proof of Algebraic Hartogs’ Lemma [12.3.10]** This proof sheds little light on the rest of this section, and thus should not be read. However, you should sleep soundly at night knowing that the proof is this short. The left side is obviously contained in the right. So assume we have some \( x \) in all \( A_p \) but not in \( A \). Let \( I \) be the “ideal of denominators” of \( x \) (cf. the proof of Proposition [6.4.2]):

\[ I := \{ r \in A : rx \in A \} \]

As \( 1/\in I \), we have \( I \neq A \), so choose a minimal prime \( q \) containing \( I \).

This construction behaves well with respect to localization — if \( p \) is any prime, then the ideal of denominators \( x \) in \( A_p \) is \( I_p \), and it again measures "the failure of Algebraic Hartogs’ Lemma for \( x \)," this time in \( A_p \). But Algebraic Hartogs’ Lemma is vacuously true for dimension 1 rings, so no codimension 1 prime contains \( I \).

Thus \( q \) has codimension at least 2. By localizing at \( q \), we can assume that \( A \) is a local ring with maximal ideal \( q \), and that \( q \) is the only prime containing \( I \). Thus \( \sqrt{I} = q \) (Exercise [4.4.F]), so as \( q \) is finitely generated, there is some \( n \) with \( I \supset q^n \) (do you see why?). Take the minimal such \( n \), so \( I \nsubseteq q^{n-1} \), and choose any \( y \in q^{n-1} - I \). Let \( z = yx \). Now \( qy \subseteq q^n \subseteq I \), so \( qz \subseteq Ix \subseteq A \), so \( qz \) is an ideal of \( A \).

I claim \( qz \) is not contained in \( q \). Otherwise, we would have a finitely-generated \( A \)-module (namely \( q \)) with a faithful \( A[z] \)-action, forcing \( z \) to be integral over \( A \) (and hence in \( A \), as \( A \) is integrally closed) by Exercise [8.2.J].

Thus \( qz \) is an ideal of \( A \) not contained in the unique maximal ideal \( q \), so it must be \( A \)! Thus \( qz = A \) from which \( q = A[1/z] \), from which \( q \) is principal. But then \( \text{codim } q = \dim A \leq \dim_{A/q} q/q^2 \leq 1 \) by Nakayama’s lemma [8.2.H] contradicting the fact that \( q \) has codimension at least 2.

□

12.4 **Proof of Krull’s Principal Ideal Theorem [12.3.3]**

The details of this proof won’t matter to us, so you should probably not read it. It is included so you can glance at it and believe that the proof is fairly short, and you could read it if you really wanted to.

If \( A \) is a ring, an **Artinian** \( A \)-module is an \( A \)-module satisfying the descending chain condition for submodules (any infinite descending sequence of submodules must stabilize, [4.6.3]). A **ring** is Artinian ring if it is Artinian over itself as a module. The notion of Artinian rings is very important, but we will get away without discussing it much.

If \( m \) is a maximal ideal of \( A \), then any finite-dimensional \( (A/m) \)-vector space (interpreted as an \( A \)-module) is clearly Artinian, as any descending chain

\[ M_1 \supset M_2 \supset \cdots \]

must eventually stabilize (as \( \dim_{A/m} M_i \) is a non-increasing sequence of non-negative integers).

12.4.A. **Exercise.** Suppose \( m \) is finitely generated. Show that for any \( n \), \( m^n / m^{n+1} \) is a finite-dimensional \( (A/m) \)-vector space. (Hint: show it for \( n = 0 \) and \( n = 1 \). Show surjectivity of \( \text{Sym}^n m/m^2 \to m^n / m^{n+1} \) to bound the dimension for general \( n \).) Hence \( m^n / m^{n+1} \) is an Artinian \( A \)-module.
12.4.B. Exercise. Suppose $A$ is a Noetherian ring with one prime ideal $m$. Suppose $m$ is finitely generated. Prove that $m^n = (0)$ for some $n$. (Hint: As $\sqrt{0}$ is prime, it must be $m$. Suppose $m$ can be generated by $r$ elements, each of which has $k$th power $0$, and show that $m^{r(k-1)+1} = 0$.)

12.4.C. Exercise. Show that if $0 \to M' \to M \to M'' \to 0$ is an exact sequence of modules, then $M$ is Artinian if and only if $M'$ and $M''$ are Artinian. (Hint: given a descending chain in $M$, produce descending chains in $M'$ and $M''$.)

12.4.1. Lemma. — If $A$ is a Noetherian ring with one prime ideal $m$, then $A$ is Artinian, i.e., it satisfies the descending chain condition for ideals.

Proof. As we have a finite filtration

$$A \supset m \supset \cdots \supset m^n = (0)$$

all of whose quotients are Artinian, $A$ is Artinian as well. \qed

12.4.2. Proof of Krull’s Principal Ideal Theorem [12.3.3] Suppose we are given $x \in A$, with $p$ a minimal prime containing $x$. By localizing at $p$, we may assume that $A$ is a local ring, with maximal ideal $p$. Suppose $q$ is another prime strictly contained in $p$.

For the first part of the theorem, we must show that $A_q$ has dimension 0. The second part follows from our earlier work: if any minimal primes are height 0, $f$ is a zero-divisor, by Theorem 6.5.3(c) and (e).

Now $p$ is the only prime ideal containing $(x)$, so $A/(x)$ has one prime ideal. By Lemma 12.4.1, $A/(x)$ is Artinian.

We invoke a useful construction, the $n$th symbolic power of a prime ideal: if $A$ is a ring, and $q$ is a prime ideal, then define

$$q^{(n)} := \{ r \in A : rs \in q^n \text{ for some } s \in A - q \}.$$

We have a descending chain of ideals in $A$

$$q^{(1)} \supset q^{(2)} \supset \cdots,$$

so we have a descending chain of ideals in $A/(x)$

$$q^{(1)} + (x) \supset q^{(2)} + (x) \supset \cdots$$

which stabilizes, as $A/(x)$ is Artinian. Say $q^{(n)} + (x) = q^{(n+1)} + (x)$, so

$$q^{(n)} \subset q^{(n+1)} + (x).$$
Hence for any \( f \in q^{(n)} \), we can write \( f = ax + g \) with \( g \in q^{(n+1)} \). Hence \( ax \in q^{(n)} \). As \( p \) is minimal over \( x, x \notin q \), so \( a \in q^{(n)} \). Thus
\[
q^{(n)} = (x)q^{(n)} + q^{(n+1)}.
\]
As \( x \) is in the maximal ideal \( p \), the second version of Nakayama’s lemma gives \( q^{(n)} = q^{(n+1)} \).

We now shift attention to the local ring \( A_q \), which we are hoping is dimension 0. We have \( q^{(n)}A_q = q^{(n+1)}A_q \) (the symbolic power construction clearly construction commutes with localization). For any \( r \in q^nA_q \subset q^{(n)}A_q \), there is some \( s \in A_q - qA_q \) such that \( rs \in q^{n+1}A_q \). As \( s \) is invertible, \( r \in q^{n+1}A_q \) as well. Thus \( q^nA_q \subset q^{n+1}A_q \), but as \( q^{n+1}A_q \subset q^nA_q \), we have \( q^nA_q = q^{n+1}A_q \). By Nakayama’s Lemma version 4 (Exercise 8.2.H),
\[
q^nA_q = 0.
\]
Finally, any local ring \((R, m)\) such that \( m^n = 0 \) has dimension 0, as \( \text{Spec } R\) consists of only one point: \([m] = V(m) = V(m^n) = V(0) = \text{Spec } R\).

12.4.3. Proof of Theorem 12.3.7, following [E] Thm. 10.2. We argue by induction on \( n \). The case \( n = 1 \) is Krull’s Principal Ideal Theorem [12.3.3]. Assume \( n > 1 \). Suppose \( p \) is a minimal prime containing \( r_1, \ldots, r_n \in A \). We wish to show that \( \text{codim } p \leq n \).

By localizing at \( p \), we may assume that \( p \) is the unique maximal ideal of \( A \). Let \( q \neq p \) be a prime ideal of \( A \) with no prime between \( p \) and \( q \). We shall show that \( q \) is minimal over an ideal generated by \( c - 1 \) elements. Then \( \text{codim } q \leq c - 1 \) by the inductive hypothesis, so we will be done.

Now \( q \) cannot contain every \( r_i \) (as \( V(r_1, \ldots, r_n) = \{[p]\}) \), so say \( r_1 \notin q \). Then \( V(q, r_1) = \{[p]\} \). As each \( r_i \in p \), there is some \( N \) such that \( r_i^N \in (q, r_1) \) (Exercise 4.4.1), so write \( r_i^N = q_i + a_i r_1 \) where \( q_i \in q \) (\( 2 \leq i \leq n \)) and \( a_i \in A \). Note that
\[
V(r_1, q_2, \ldots, q_n) = V(r_1, r_2^N, \ldots, r_n^N) = V(r_1, r_2, \ldots, r_n) = \{[p]\}.
\]

We shall show that \( q \) is minimal among primes containing \( q_2, \ldots, q_n \), completing the proof. In the ring \( A/\langle q_2, \ldots, q_n \rangle \), \( V(r_1) = \{[p]\} \) by (12.4.3.1). By Krull’s principal ideal theorem 12.3.3 \([p]\) is codimension at most 1, so \( [q] \) must be codimension 0 in \( \text{Spec } A/\langle q_2, \ldots, q_n \rangle \), as desired.
Nonsingularity ("smoothness") of Noetherian schemes

One natural notion we expect to see for geometric spaces is the notion of when an object is "smooth". In algebraic geometry, this notion, called nonsingularity (or regularity, although we won’t use this term) is easy to define but a bit subtle in practice. We will soon define what it means for a scheme to be nonsingular (or regular) at a point. The Jacobian criterion will show that this corresponds to smoothness as you may have seen it before. A point that is not nonsingular is (not surprisingly) called singular ("not smooth"). A scheme is said to be nonsingular if all its points are nonsingular, and singular if one of its points is singular.

The notion of nonsingularity is less useful than you might think. Grothendieck taught us that the more important notions are properties of morphisms, not of objects, and there is indeed a "relative notion" that applies to a morphism of schemes $f : X \to Y$ that is much better-behaved (corresponding to the notion of smooth map or submersion in differential geometry). For this reason, the word "smooth" is reserved for these morphisms. We will discuss smooth morphisms in Chapter 25. However, nonsingularity is still useful, especially in (co)dimension 1, and we shall discuss this case (of discrete valuation rings) in §13.3.

13.1 The Zariski tangent space

We first define the tangent space of a scheme at a point. It behaves like the tangent space you know and love at "smooth" points, but also makes sense at other points. In other words, geometric intuition at the "smooth" points guides the definition, and then the definition guides the algebra at all points, which in turn lets us refine our geometric intuition.

This definition is short but surprising. The main difficulty is convincing yourself that it deserves to be called the tangent space. This is tricky to explain, because we want to show that it agrees with our intuition, but our intuition is worse than we realize. So I will just define it for you, and later try to convince you that it is reasonable.

13.1.1. Definition. Suppose $p$ is a prime ideal of a ring $A$, so $[p]$ is a point of $\text{Spec} A$. Then $[pA_p]$ is a point of the scheme $\text{Spec} A_p$. For convenience, we let $m := pA_p \subset A_p =: B$. Let $\kappa = B/m$ be the residue field. Then $m/m^2$ is a vector space over the residue field $\kappa$; it is a $B$-module, and elements of $m$ acts like 0. This is defined to be the Zariski cotangent space. The dual vector space is the Zariski tangent space. Elements of the Zariski cotangent space are called cotangent vectors or differentials; elements of the tangent space are called tangent vectors.
Note that this definition is intrinsic. It does not depend on any specific description of the ring itself (such as the choice of generators over a field \( k \), which is equivalent to the choice of embedding in affine space). Notice that the cotangent space is more algebraically natural than the tangent space (the definition is shorter). There is a moral reason for this: the cotangent space is more naturally determined in terms of functions on a space, and we are very much thinking about schemes in terms of “functions on them”. This will come up later.

Here are two plausibility arguments that this is a reasonable definition. Hopefully one will catch your fancy.

In differential geometry, the tangent space at a point is sometimes defined as the vector space of derivations at that point. A derivation is a function that takes in functions near the point that vanish at the point, and gives elements of the field \( k \), and satisfies the Leibniz rule

\[(fg)' = f'g + g'f.\]

(We will later define derivations in more general settings, \( \S 22.2.14 \) Translation: a derivation is a map \( m \to k \). But \( m^2 \) maps to \( 0 \), as if \( f(p) = g(p) = 0 \), then

\[(fg)'(p) = f'(p)g(p) + g'(p)f(p) = 0.\]

Thus we have a map \( m/m^2 \to k \), i.e. an element of \( (m/m^2)^\vee \).

13.1.A. Exercise. Check that this is reversible, i.e. that any map \( m/m^2 \to k \) gives a derivation. In other words, verify that the Leibniz rule holds.

Here is a second vaguer motivation that this definition is plausible for the cotangent space of the origin of \( \mathbb{A}^n \). Functions on \( \mathbb{A}^n \) should restrict to a linear function on the tangent space. What (linear) function does \( x^2 + xy + x + y \) restrict to “near the origin”? You will naturally answer: \( x + y \). Thus we “pick off the linear terms”. Hence \( m/m^2 \) are the linear functionals on the tangent space, so \( m/m^2 \) is the cotangent space. In particular, you should picture functions vanishing at a point (i.e. lying in \( m \)) as giving functions on the tangent space in this obvious a way.

13.1.2. Old-fashioned example. Computing the Zariski-tangent space is actually quite hands-on, because you can compute it just as you did when you learned multivariable calculus. In \( \mathbb{A}^3 \), we have a curve cut out by \( x + y + z^2 + xyz = 0 \) and \( x - 2y + z + x^2y^2z^3 = 0 \). (You can use Krull’s Principal Ideal Theorem \( \S 22.3.3 \) to check that this is a curve, but it is not important to do so.) What is the tangent line near the origin? (Is it even smooth there?) Answer: the first surface looks like \( x + y = 0 \) and the second surface looks like \( x - 2y + z = 0 \). The curve has tangent line cut out by \( x + y = 0 \) and \( x - 2y + z = 0 \). It is smooth (in the traditional sense). In multivariable calculus, the students do a page of calculus to get the answer, because we aren’t allowed to tell them to just pick out the linear terms.

Let’s make explicit the fact that we are using. If \( A \) is a ring, \( m \) is a maximal ideal, and \( f \in m \) is a function vanishing at the point \([m] \in \text{Spec } A\), then the Zariski tangent space of \( \text{Spec } A/(f) \) at \( m \) is cut out in the Zariski tangent space of \( \text{Spec } A \) (at \( m \)) by the single linear equation \( f \, (\text{mod } m^2) \). The next exercise will force you think this through.

13.1.B. Important exercise (“Krull’s Principal Ideal Theorem for the Zariski tangent space” — but much easier than Krull’s Principal Ideal...
**Theorem 12.3.3.** Suppose $\mathcal{A}$ is a ring, and $m$ a maximal ideal. If $f \in m$, show that the Zariski tangent space of $\mathcal{A}/f$ is cut out in the Zariski tangent space of $\mathcal{A}$ by $f \pmod{m^2}$. (Note: we can quotient by $f$ and localize at $m$ in either order, as quotienting and localizing commute, (5.3.1).) Hence the dimension of the Zariski tangent space of $\text{Spec} \mathcal{A}$ at $[m]$ is the dimension of the Zariski tangent space of $\text{Spec} \mathcal{A}/(f)$ at $[m]$, or one less.

Here is another example to see this principle in action: $x + y + z^2 = 0$ and $x + y + x^2 + y^4 + z^5 = 0$ cuts out a curve, which obviously passes through the origin. If I asked my multivariable calculus students to calculate the tangent line to the curve at the origin, they would do a reams of calculations which would boil down to picking off the linear terms. They would end up with the equations $x + y = 0$ and $x + y = 0$, which cuts out a plane, not a line. They would be disturbed, and I would explain that this is because the curve isn’t smooth at a point, and their techniques don’t work. We on the other hand bravely declare that the cotangent space is cut out by $x + y = 0$, and (will soon) define this as a singular point. (Intuitively, the curve near the origin is very close to lying in the plane $x + y = 0$.) Notice: the cotangent space jumped up in dimension from what it was “supposed to be”, not down. We will see that this is not a coincidence soon, in Theorem 13.2.1.

Here is a nice consequence of the notion of Zariski tangent space.

**13.1.3. Problem.** Consider the ring $\mathcal{A} = k[x, y, z]/(xy - z^2)$. Show that $(x, z)$ is not a principal ideal.

As $\text{dim} \mathcal{A} = 2$ (by Krull’s Principal Ideal Theorem 12.3.3), and $\mathcal{A}/(x, z) \cong k[y]$ has dimension 1, we see that this ideal is codimension 1 (as codimension is the difference of dimensions for irreducible varieties, Exercise 12.2.D). Our geometric picture is that $\text{Spec} \mathcal{A}$ is a cone (we can diagonalize the quadric as $xy - z^2 = ((x + y)/2)^2 - ((x - y)/2)^2 - z^2$, at least if $\text{char} k \neq 2$ — see Exercise 6.4.I), and that $(x, z)$ is a ruling of the cone. (See Figure 13.1 for a sketch.) This suggests that we look at the cone point.

![Figure 13.1](image_url)

**Figure 13.1.** $V(x, z) \subset \text{Spec} k[x, y, z]/(xy - z^2)$ is a ruling on a cone

**Solution.** Let $m = (x, y, z)$ be the maximal ideal corresponding to the origin. Then $\text{Spec} \mathcal{A}$ has Zariski tangent space of dimension 3 at the origin, and $\text{Spec} \mathcal{A}/(x, z)$ has Zariski tangent space of dimension 1 at the origin. But $\text{Spec} \mathcal{A}/(f)$ must have Zariski tangent space of dimension at least 2 at the origin by Exercise 13.1.B.
13.1.C. Exercise. Show that \((x, z) \subset k[w, x, y, z]/(wz - xy)\) is a codimension 1 ideal that is not principal. (See Figure 13.2 for the projectivization of this situation.) This example was promised in Exercise 6.4.D. You might use it again in Exercise 13.1.D.

![Figure 13.2. The ruling \(V(x, z)\) on \(V(wz - xy) \subset \mathbb{P}^3\).](image)

13.1.D. Exercise. Show that \(A = k[w, x, y, z]/(wz - xy)\) is not a unique factorization domain. (One possibility is to do this “directly”. This might be hard to do rigorously — how do you know that \(x\) is irreducible? Another possibility, faster but less intuitive, is to use the intermediate result that in a unique factorization domain, any codimension 1 prime is principal, Lemma 12.2.2 and considering Exercise 13.1.C.) As \(A\) is integrally closed if \(k = \mathbb{F}\) and \(\text{char } k \neq 2\) (Exercise 6.4.I(c)), this yields an example of a scheme that is normal but not factorial, as promised in Exercise 6.4.F.

13.1.4. Morphisms and tangent spaces. Suppose \(f : X \to Y\), and \(f(p) = q\). Then if we were in the category of manifolds, we would expect a tangent map, from the tangent space of \(p\) to the tangent space at \(q\). Indeed that is the case; we have a map of stalks \(O_{Y,q} \to O_{X,p}\), which sends the maximal ideal of the former \(n\) to the maximal ideal of the latter \(m\) (we have checked that this is a “local morphism” when we briefly discussed locally ringed spaces). Thus \(n^2 \to m^2\), from which \(n/n^2 \to m/m^2\). If \((O_{X,p}, m)\) and \((O_{Y,q}, n)\) have the same residue field \(k\), so \(n/n^2 \to m/m^2\) is a linear map of \(k\)-vector spaces, we have a natural map \((m/m^2)^\vee \to (n/n^2)^\vee\). This is the map from the tangent space of \(p\) to the tangent space at \(q\) that we sought. (Aside: note that the cotangent map always exists, without requiring \(p\) and \(q\) to have the same residue field — a sign that cotangent spaces are more natural than tangent spaces in algebraic geometry.)

Here are some exercises to give you practice with the Zariski tangent space. If you have some differential geometric background, the first will further convince you that this definition correctly captures the idea of (co)tangent space.
13.1.E. Important Exercise (the Jacobian computes the Zariski tangent space). Suppose $X$ is a finite type $k$-scheme. Then locally it is of the form $\text{Spec } k[x_1, \ldots, x_n]/(f_1, \ldots, f_r)$. Show that the Zariski cotangent space at a closed point $p$ with residue field $k$ is given by the cokernel of the Jacobian map $k^r \to k^n$ given by the Jacobian matrix

$$J = \begin{pmatrix} \frac{\partial f_1}{\partial x_1}(p) & \cdots & \frac{\partial f_1}{\partial x_n}(p) \\ \vdots & \ddots & \vdots \\ \frac{\partial f_r}{\partial x_1}(p) & \cdots & \frac{\partial f_r}{\partial x_n}(p) \end{pmatrix}. \quad (13.1.4.1)$$

(This is makes precise our example of a curve in $\mathbb{A}^3$ cut out by a couple of equations, where we picked off the linear terms, see Example 13.1.2.) You might be alarmed: what does $\frac{\partial f}{\partial x}$ mean? Do you need deltas and epsilons? No! Just define derivatives formally, e.g.

$$\frac{\partial}{\partial x_1}(x_1^2 + x_1 x_2 + x_2^2) = 2x_1 + x_2.$$ 

Hint: Do this first when $p$ is the origin, and consider linear terms, just as in Example 13.1.2 and Exercise 13.1.B. For the general case, “translate $p$ to the origin”.

13.1.F. Less Important Exercise (“Higher-order data”). In Exercise 4.7.B you computed the equations cutting out the three coordinate axes of $\mathbb{A}^3_k$. (Call this scheme $X$.) Your ideal should have had three generators. Show that the ideal can’t be generated by fewer than three elements. (Hint: working modulo $m = (x, y, z)$ won’t give any useful information, so work modulo $m^2$.)

13.1.G. Exercise. Suppose $X$ is a $k$-scheme. Describe a natural bijection from $\text{Mor}_k(\text{Spec } k[\epsilon]/(\epsilon^2), X)$ to the data of a point $p$ with residue field $k$ (necessarily a closed point) and a tangent vector at $p$. (This turns out to be very important, for example in deformation theory.)

13.1.H. Exercise. Find the dimension of the Zariski tangent space at the point $[2, 2i]$ of $\mathbb{Z}[2i] \cong \mathbb{Z}[x]/(x^2 + 4)$. Find the dimension of the Zariski tangent space at the point $[(2, x)]$ of $\mathbb{Z}[\sqrt{-2}] \cong \mathbb{Z}[x]/(x^2 + 2)$. (If you prefer geometric versions of the same examples, replace $\mathbb{Z}$ by $\mathbb{R}$ or $\mathbb{C}$, and 2 by $y$: consider $\mathbb{C}[x, y]/(x^2 + y^2)$ and $\mathbb{C}[x, y]/(x^2 + y)$.)

### 13.2 The local dimension is at most the dimension of the tangent space, and nonsingularity

The key idea in the definition of nonsingularity is contained in the title of this section.

13.2.1. Theorem. — Suppose $(A, m, k)$ is a Noetherian local ring. Then $\dim A \leq \dim_k m/m^2$.

If equality holds, we say that $A$ is a regular local ring. (If a Noetherian ring $A$ is regular at all of its primes, $A$ is said to be a regular ring, but we won’t use this terminology.) A locally Noetherian scheme $X$ is regular or nonsingular at a point
p if the local ring \(O_{X,p}\) is regular. It is **singular** at the point otherwise. A scheme is **regular** or **nonsingular** if it is regular at all points. It is **singular** otherwise (i.e. if it is singular at at least one point).

You will hopefully become convinced that this is the right notion of “smoothness” of schemes. Remarkably, Krull introduced the notion of a regular local ring for purely algebraic reasons, some time before Zariski realized that it was a fundamental notion in geometry in 1947.

### 13.2.2. Proof of Theorem [13.2.7]

Note that \(m\) is finitely generated (as \(A\) is Noetherian), so \(m/m^2\) is a finitely generated \((A/m = k)\)-module, hence finite-dimensional. Say \(\text{dim}_k m/m^2 = n\). Choose a basis of \(m/m^2\), and lift them to elements \(f_1, \ldots, f_n\) of \(m\). Then by Nakayama’s lemma (version 4, Exercise 8.2.H), \((f_1, \ldots, f_n) = m\).

Recall Krull’s Theorem [12.3.7] any irreducible component of \(V(f_1, \ldots, f_n)\) has codimension at most \(n\). In this case, \(V((f_1, \ldots, f_n)) = V(m)\) is just the point \([m]\), so the codimension of \(m\) is at most \(n\). Thus the longest chain of prime ideals contained in \(m\) is at most \(n + 1\). But this is also the longest chain of prime ideals in \(A\) (as \(m\) is the unique maximal ideal), so \(n \geq \text{dim} A\).

### 13.2.A. Exercise

Show that Noetherian local rings have finite dimension. (Noetherian rings in general may have infinite dimension, see Exercise [12.1.E].)

### 13.2.3. The Jacobian criterion for nonsingularity, and \(k\)-smoothness.

A finite type \(k\)-scheme is locally of the form \(\text{Spec } k[x_1, \ldots, x_n]/(f_1, \ldots, f_r)\). The Jacobian criterion for nonsingularity (Exercise [13.1.B]) gives a hands-on method for checking for singularity at closed points, using the equations \(f_1, \ldots, f_r\), if \(k = \overline{k}\).


Suppose \(X = \text{Spec } k[x_1, \ldots, x_n]/(f_1, \ldots, f_r)\) has pure dimension \(d\). Show that a \(k\)-valued point \(p \in X\) is a smooth point of \(X\) if the corank of the Jacobian matrix [13.1.4.1] is \(d\) at \(p\).

### 13.2.C. Easy exercise.

Suppose \(k = \overline{k}\). Show that the singular closed points of the hypersurface \(f(x_1, \ldots, x_n) = 0\) in \(\mathbb{A}^n_k\) are given by the equations \(f = \frac{\partial f}{\partial x_1} = \cdots = \frac{\partial f}{\partial x_n} = 0\). (Translation: the singular points of \(f = 0\) are where the gradient of \(f\) vanishes. This is not shocking.)

### 13.2.4. Remark: The Jacobian criterion over fields in general.

If \(k = \overline{k}\), the Jacobian criterion tells you which closed points which are singular. We will see in [25.3.3] that this criterion works also when \(k\) is separably closed, and it is a sufficient (but not necessary) criterion for nonsingularity in general. The following example in characteristic \(p\) shows that the Jacobian criterion is not necessary in general: Let \(k = \mathbb{F}_p(u)\), and consider the hypersurface \(X = \text{Spec } k[x]/(x^p - u)\). Now \(k[x]/(x^p - u)\) is a field, hence nonsingular. But if \(f(x) = x^p - u\), then \(\frac{df}{dx}(u) = 0\), so the Jacobian criterion fails.

### 13.2.5. Smoothness over a field \(k\).

Before using the Jacobian criterion to get our hands dirty with some explicit varieties, I want to make some general philosophical comments. There seem to be two serious drawbacks with the Jacobian criterion.
For finite type schemes over $k$, the criterion gives a necessary condition for non-
singularity, but it is not obviously sufficient, as we need to check nonsingularity
at non-closed points as well. We can prove sufficiency by working hard to show
Fact 13.2.11 which shows that the non-closed points must be nonsingular as well.
A second failing is that the criterion requires $k$ to be algebraically closed. These
problems suggest that the old-fashioned ideas of using derivatives and Jacobians
are ill-suited to the correct modern notion of nonsingularity. But in fact the fault
is with nonsingularity. There is a better notion of smoothness over a field. Better
yet, this idea generalizes to the notion of a smooth morphism of schemes, which
behaves well in all possible ways (preserved by base change, composition, etc.).
This is another sign that some properties we think of as of objects ("absolute notions")
should really be thought of as properties of morphisms ("relative notions").
We know enough to define what it means for a scheme to be $k$-smooth, or smooth
over $k$: a $k$-scheme is smooth of dimension $d$ if it is reduced and locally of finite
type, pure dimension $d$, and for any local patch $\text{Spec} \ k[x_1, \ldots, x_n]/(f_1, \ldots, f_r)$, the
Jacobian has corank $d$ everywhere. We could then show that it suffices to check
this on any cover by affine open sets (and by any choice of generators of the ring
corresponding to such an open set), and also that it suffices to check at the closed
points (rank of a matrix of functions is an uppercontinuous function). But the cok-
ernel of the Jacobian matrix is secretly the space of differentials (which might not
be surprising if you have experience with differentials in differential geometry), so
we will hold off discussing this notion until §25.2.1.

So for now, let’s discuss some important examples over an algebraically closed
field.

13.2.D. Exercise. Suppose $k = \bar{k}$. Show that $A^1_k$ and $A^2_k$ are nonsingular. (Make
sure to check nonsingularity at the non-closed points! Fortunately you know what
all the points of $A^2_k$ are; this is trickier for $A^3_k$.) Show that $\mathbb{P}^1_k$ and $\mathbb{P}^2_k$ are nonsingular.
(This holds even if $k$ isn’t algebraically closed, and in higher dimension.)

13.2.E. Exercise (the Euler test for projective hypersurfaces). There is
an analogous Jacobian criterion for hypersurfaces $f = 0$ in $\mathbb{P}^n_k$. Suppose $k = k$.
Show that the singular closed points correspond to the locus

$$f = \frac{\partial f}{\partial x_1} = \cdots = \frac{\partial f}{\partial x_n} = 0.$$

If the degree of the hypersurface is not divisible by char $k$ (e.g. if char $k = 0$), show
that it suffices to check $\frac{\partial f}{\partial x_1} = \cdots = \frac{\partial f}{\partial x_n} = 0$. Hint: show that $(\deg f)f = \sum x_i \frac{\partial f}{\partial x_i}$.
(Fact: this will give the singular points in general, not just the closed points, cf.
[13.2.3] I don’t want to prove this, and I won’t use it.)

13.2.F. Exercise. Suppose that $k$ is algebraically closed. Show that $y^2z = x^3 - xz^2$
in $\mathbb{P}^2_k$ is an irreducible nonsingular curve. (Eisenstein’s criterion gives one way of
showing irreducibility. Warning: we didn’t specify char $k \neq 3$, so be careful when
using the Euler test.)

13.2.G. Exercise. Suppose $k = \bar{k}$ has characteristic 0. Show that there exists a
nonsingular plane curve of degree $d$. (Feel free to weaken the hypotheses.)
13.2.H. Exercise. Find all the singular closed points of the following plane curves. Here we work over $k = \overline{k}$ of characteristic 0 to avoid distractions.

(a) $y^2 = x^2 + x^3$. This is an example of a node.
(b) $y^2 = x^3$. This is called a cusp; we met it earlier in Exercise [10.6.F].
(c) $y^2 = x^4$. This is called a tacnode; we met it earlier in Exercise [10.6.G].

(A precise definition of a node etc. will be given in Definition [13.5.2].)

13.2.I. Exercise. Suppose $k = \overline{k}$. Use the Jacobian criterion to show that the twisted cubic $\text{Proj} \, k[w, x, y, z]/(wz - xy, wy - x^2, xz - y^2)$ is nonsingular. (You can do this, without any hypotheses on $k$, using the fact that it is isomorphic to $\mathbb{P}^1$. But do this with the explicit equations, for the sake of practice. The twisted cubic was defined in Exercise [9.2.A].)

13.2.J. Easy Exercise. Show that $\text{Spec} \, Z$ is a nonsingular curve.

13.2.K. Exercise. (This tricky exercise is for those who know about the primes of the Gaussian integers $\mathbb{Z}[i]$.) There are several ways of showing that $\mathbb{Z}[i]$ is dimension 1 (For example: (i) it is a principal ideal domain; (ii) it is the normalization of $\mathbb{Z}$ in the field extension $\mathbb{Q}(i)/\mathbb{Q}$; (iii) using Krull’s Principal Ideal Theorem [12.3.3] and the fact that $\dim \mathbb{Z}[x] = 2$ by Exercise [12.1.C].) Show that $\text{Spec} \, \mathbb{Z}[i]$ is a nonsingular curve. (There are several ways to proceed. You could use Exercise [13.1.B]. For example, consider the prime $(2, 1 + i)$, which is cut out by the equations 2 and $1 + x$ in $\text{Spec} \, \mathbb{Z}[x]/(x^2 + 1)$. We will later ($\S$ 13.3.10) have a simpler approach once we discuss discrete valuation rings.

13.2.L. Exercise. Show that $[(5, 5i)]$ is the unique singular point of $\text{Spec} \, \mathbb{Z}[5i]$. (Hint: $\mathbb{Z}[i]_5 \cong \mathbb{Z}[5i]_5$. Use the previous exercise.)

13.2.M. Two facts worth knowing about regular local rings.

Here are two pleasant facts. Because we won’t prove them in full generality, we will be careful when using them. In this section only, you may assume these facts in doing exercises. In some sense, the first fact connects regular local rings to algebra, and the second connects them to geometry.


Thus regular schemes are factorial, and hence normal by Exercise [5.4.F].

In particular, as you might expect, a scheme is “locally irreducible” at a “smooth” point: a (Noetherian) regular local ring is an integral domain. This can be shown more directly, [E] Cor. 10.14. (Of course, normality suffices to show that a Noetherian local ring is an integral domain — normal local rings are always integral domains.) Using “power series” ideas, we will prove the following case in [13.5] which will suffice for dealing with varieties.

13.2.O. Theorem. — Suppose $(A, m)$ is a regular local ring containing its residue field $k$ (i.e. $A$ is a $k$-algebra). Then $A$ is an integral domain.
13.2.M. Exercise. Suppose X is a variety over k, and p is a nonsingular k-valued point. Use Theorem 13.2.9 to show that only one irreducible component of X passes through p. (Your argument will apply without change to general Noetherian schemes using Fact 13.2.8.)

13.2.N. Easy Exercise. Show that a nonsingular Noetherian scheme is irreducible if and only if it is connected. (Hint: Exercise 6.2.I.)

13.2.10. Remark: factoriality is weaker than nonsingularity. There are local rings that are singular but still factorial, so the implication factorial implies nonsingular is strict. Here are two examples, that we will verify later.

(i) The ring \( A = k[x_1, \ldots, x_n]/(x_1^2 + \cdots + x_n^2) \) is a unique factorization domain when \( n \geq 5 \), so \( \text{Spec} A \) is factorial, but it is clearly singular at the origin. In this case where k is algebraically closed and characteristic not 2, the UFD fact will be shown in Exercise 15.2.1. More generally, it is a consequence of Grothendieck’s proof (of a conjecture of Samuel) that a local Noetherian ring that is a complete intersection — in particular a hypersurface — that is factorial in codimension at most 3 must be factorial, [SGA2, Exp. XI, Cor. 3.14]. The hypothesis \( n \geq 5 \) is necessary, because of our friend the nonsingular quadric, see Exercise 13.1.D.

(ii) If \( \text{char} k \neq 2 \), and k does not contain a square root of \(-1\), then \( k[x, y, z]/(x^2 + y^2 - z^2) \) is a unique factorization domain (see Exercise 15.2.K), but its spectrum is also clearly singular at the origin.

We come next to the second fact that will help us sleep well at night.

13.2.11. Fact [I, Cor. 19.14]. Suppose \((A, m)\) is a Noetherian regular local ring. Any localization of \( A \) at a prime is also a regular local ring.

Hence to check if Spec \( A \) (\( A \) a Noetherian ring) is nonsingular, then it suffices to check at closed points (at maximal ideals). This major theorem was an open problem in commutative algebra for a long time until settled by homological methods. The special case of local rings that are localizations of finite type \( k \)-algebras will be given in Exercise 22.7.E.

13.2.O. Exercise. Show (using Fact 13.2.11) that you can check nonsingularity of a Noetherian scheme by checking at closed points. (Caution: as mentioned in Exercise 6.1.E a scheme in general needn’t have any closed points!)

We will be able to prove two important cases of Exercise 13.2.O without invoking Fact 13.2.11. The first will be proved in 22.7.A.

13.2.12. Theorem. If \( X \) is a finite type \( \overline{k} \)-scheme that is nonsingular at all its closed points, then \( X \) is nonsingular.

13.2.P. Exercise. Suppose \( X \) is a Noetherian dimension 1 scheme that is nonsingular at its closed points. Show that \( X \) is reduced. Hence show (without invoking Fact 13.2.11) that \( X \) is nonsingular.

13.2.Q. Exercise (generalizing Exercise 13.2.G). Suppose \( k \) is an algebraically closed field of characteristic 0. Show that there exists a nonsingular hypersurface of degree \( d \) in \( \mathbb{P}^n \). (As in Exercise 13.2.G, feel free to weaken the hypotheses.)
Although we now know that $\mathbf{A}^n_k$ is nonsingular (modulo our later proof of Theorem 13.2.12), you may be surprised to find that we never use this fact (although we might make use the fact that it is nonsingular in dimension 0 and codimension 1, which we knew beforehand). Perhaps surprisingly, it is more important to us that $\mathbf{A}^n_k$ is factorial and hence normal, which we showed more simply. Similarly, geometers may be pleased to finally know that varieties are nonsingular if and only if they are nonsingular at closed points, but they likely cared only about the closed points anyway. In short, nonsingularity is less important than you might think, except in (co)dimension 1, which is the topic of the next section.

13.3 Discrete valuation rings: Dimension 1 Noetherian regular local rings

The case of (co)dimension 1 is important, because if you understand how primes behave that are separated by dimension 1, then you can use induction to prove facts in arbitrary dimension. This is one reason why Krull’s Principal Ideal Theorem 12.3.3 is so useful.

A dimension 1 Noetherian regular local ring can be thought of as a “germ of a smooth curve” (see Figure 13.3). Two examples to keep in mind are $k[x]_{(x)} = \{f(x)/g(x) : x \not| g(x)\}$ and $\mathbb{Z}_{(5)} = \{a/b : 5 \not| b\}$. The first example is “geometric” and the second is “arithmetical”, but hopefully it is clear that they are basically the same.

![Figure 13.3. A germ of a curve](image)

The purpose of this section is to give a long series of equivalent definitions of these rings. Before beginning, we quickly sketch these seven definitions. There are a number of ways a Noetherian local ring can be “nice”. It can be regular, or a principal domain, or a unique factorization domain, or normal. In dimension 1, these are the same. Also equivalent are nice properties of ideals: if $\mathfrak{m}$ is principal; or if all ideals are either powers of the maximal ideal, or 0. Finally, the ring can have a discrete valuation, a measure of “size” of elements that behaves particularly well.

13.3.1. Theorem. — Suppose $(A, \mathfrak{m})$ is a Noetherian local ring of dimension 1. Then the following are equivalent.

(a) $(A, \mathfrak{m})$ is regular.

(b) $\mathfrak{m}$ is principal.

Here is why (a) implies (b). If $A$ is regular, then $\mathfrak{m}/\mathfrak{m}^2$ is one-dimensional. Choose any element $t \in \mathfrak{m} - \mathfrak{m}^2$. Then $t$ generates $\mathfrak{m}/\mathfrak{m}^2$, so generates $\mathfrak{m}$ by Nakayama’s lemma 8.2.11. We call such an element a uniformizer.
Conversely, if \( m \) is generated by one element \( t \) over \( A \), then \( m/m^2 \) is generated by one element \( t \) over \( A/m = k \). Since \( \dim k \cdot m^2 / m \geq 1 \) by Theorem 13.2.1, we have \( \dim k \cdot m^2 / m = 1 \), and \((A, m)\) is regular.

We will soon use a useful fact, and we may as well prove it in much more generality than we need, because the proof is so short.

### 13.3.2. Proposition

If \((A, m)\) is a Noetherian local ring, then \( \cap_i m^i = 0 \).

The geometric intuition for this is that any function that is analytically zero at a point (vanishes to all orders) actually vanishes at that point. The geometric intuition also suggests an example showing that Noetherianness is necessary: consider the function \( e^{-1/x^2} \) in the germs of \( C^\infty \)-functions on \( \mathbb{R} \) at the origin.

It is tempting to argue that \( m(\cap_i m^i) = \cap_i m^i \), and then to use Nakayama’s Lemma 8.2.H to argue that \( \cap_i m^i = 0 \). Unfortunately, it is not obvious that this first equality is true: product does not commute with infinite intersections in general. But we will still make this work.

**Proof.** (A better proof, putting the result into a larger context, is via the Artin-Rees lemma [E Lem. 5.1], see [E Cor. 5.4].) Let \( I = \cap_i m^i \). We wish to show that \( I \subseteq mI \), and hence by Nakayama’s Lemma 8.2.H \( I = 0 \). Fix a primary decomposition of \( mI \). It suffices to show that \( q \) contains \( I \), for any \( q \) in this primary decomposition, as then \( I \) is contained in all the primary ideals in the decomposition of \( mI \), and hence \( mI \). Let \( p = \sqrt{q} \).

If \( p \neq m \), then choose \( x \in m \setminus p \). Now \( x \) is not nilpotent in \( A/q \), and hence is not a zero-divisor. (Recall that \( q \) is primary if and only if in \( A/q \), each zero-divisor is nilpotent.) But \( xI \subseteq mI \subseteq q \), so \( I \subseteq q \).

On the other hand, if \( p = m \), then as \( m \) is finitely generated, and each generator is in \( \sqrt{q} = m \), there is some \( a \) such that \( m^a \subseteq q \). But \( I \subseteq m^a \), so we are done. \( \square \)

### 13.3.3. Proposition

Suppose \((A, m)\) is a Noetherian regular local ring of dimension 1 (i.e. satisfying (a) above). Then \( A \) is an integral domain.

**Proof.** Suppose \( xy = 0 \), and \( x, y \neq 0 \). Then by Proposition 13.3.2 \( x \in m^i \setminus m^{i+1} \) for some \( i \geq 0 \), so \( x = at^i \) for some \( a \not\in m \). Similarly, \( y = bt^j \) for some \( j \geq 0 \) and \( b \not\in m \). As \( a, b \not\in m \), \( a \) and \( b \) are invertible. Hence \( xy = 0 \) implies \( t^{i+j} = 0 \). But as nilpotents don’t affect dimension,

\[
\dim A = \dim A/(t) = \dim A/m = \dim k = 0,
\]

contradicting \( \dim A = 1 \). \( \square \)

### 13.3.4. Theorem

Suppose \((A, m)\) is a Noetherian local ring of dimension 1. Then (a) and (b) are equivalent to:

(c) all ideals are of the form \( m^n \) or \( 0 \).

**Proof.** Assume (a): suppose \((A, m, k)\) is a Noetherian regular local ring of dimension 1. Then I claim that \( m^n \neq m^{n+1} \) for any \( n \). Otherwise, by Nakayama’s lemma, \( m^n = 0 \), from which \( t^n = 0 \). But \( A \) is an integral domain, so \( t = 0 \), from which \( A = A/m \) is a field, which can’t have dimension 1, contradiction.
I next claim that \(m^n/m^{n+1}\) is dimension 1. Reason: \(m^n = (t^n)\). So \(m^n\) is generated as a \(A\)-module by one element, and \(m^n/(mm^n)\) is generated as a \((A/m = k)\)-module by 1 element (non-zero by the previous paragraph), so it is a one-dimensional vector space.

So we have a chain of ideals \(A \supset m \supset m^2 \supset m^3 \supset \cdots\) with \(\cap m^i = \{0\}\) (Proposition 13.3.2). We want to say that there is no room for these, because “each pair is “separated by dimension 1”, and there is “no room at the end”. Proof: suppose \(1 \subset A\) is an ideal. If \(1 \neq \{0\}\), then there is some \(n\) such that \(I \subset m^n\) but \(I \not\subset m^{n+1}\). Choose some \(u \in I - m^{n+1}\). Then \((u) \subset I\). But \(u\) generates \(m^n/m^{n+1}\), hence by Nakayama it generates \(m^n\), so we have \(m^n \subset I\subset m^n\), so we are done. Conclusion: in a Noetherian local ring of dimension 1, regularity implies all ideals are of the form \(m^n\) or \(\{0\}\).

We now show that (c) implies (a). Assume (a) is false: suppose we have a dimension 1 Noetherian local integral domain that is not regular, so \(m/m^2\) has dimension at least 2. Choose any \(u \in m - m^2\). Then \((u, m^2)\) is an ideal, but \(m \subset (u, m^2) \subsetneq m^2\). \(\square\)

13.3.A. Easy Exercise. Suppose \((A, m)\) is a Noetherian dimension 1 local ring. Show that (a)--(c) above are equivalent to:

(d) \(A\) is a principal ideal domain.

13.3.5. Discrete valuation rings. We next define the notion of a discrete valuation ring. Suppose \(K\) is a field. A discrete valuation on \(K\) is a surjective homomorphism \(\nu : K^\times \to \mathbb{Z}\) (in particular, \(\nu(xy) = \nu(x) + \nu(y)\)) satisfying
\[
\nu(x + y) \geq \min(\nu(x), \nu(y))
\]
extcept if \(x + y = 0\) (in which case the left side is undefined). (Such a valuation is called non-archimedean, although we will not use that term.) It is often convenient to say \(\nu(0) = \infty\). More generally, a valuation is a surjective homomorphism \(\nu : K^\times \to G\) to a totally ordered group \(G\), although this isn’t so important to us.

Examples.

(i) \((the 5-adic valuation)\) \(K = \mathbb{Q}\), \(\nu(r)\) is the “power of 5 appearing in \(r\)”, e.g. \(\nu(35/2) = 1, \nu(27/125) = -3\).

(ii) \(K = k(x), \nu(f)\) is the “power of \(x\) appearing in \(f\).”

(iii) \(K = k(x), \nu(f)\) is the negative of the degree. This is really the same as (ii), with \(x\) replaced by \(1/x\).

Then \(0 \cup \{x \in K^\times : \nu(x) \geq 0\}\) is a ring, which we denote \(\mathcal{O}_\nu\). It is called the valuation ring of \(\nu\). (Not every valuation is discrete. Consider the ring of Puiseux series over a field \(k\), \(K = \bigcup_{n \geq 1} k\{x^{1/n}\}\), with \(\nu : K^\times \to \mathbb{Q}\) given by \(\nu(x^q) = q\).)

13.3.B. Exercise. Describe the valuation rings in the three examples above. (You will notice that they are familiar-looking dimension 1 Noetherian local rings. What a coincidence!)

13.3.C. Exercise. Show that \(\{0\} \cup \{x \in K^\times : \nu(x) \geq 1\}\) is the unique maximal ideal of the valuation ring. (Hint: show that everything in the complement is invertible.) Thus the valuation ring is a local ring.
An integral domain \( A \) is called a **discrete valuation ring** (or DVR) if there exists a discrete valuation \( v \) on its fraction field \( K = \mathbb{K}(A) \) for which \( \mathfrak{O}_v = A \).

Similarly, \( A \) is a **valuation ring** if there exists a valuation \( v \) on \( K \) for which \( \mathfrak{O}_v = A \).

Now if \( A \) is a Noetherian regular local ring of dimension 1, and \( t \) is a uniformizer (a generator of \( m \) as an ideal, or equivalently of \( m/m^2 \) as a \( k \)-vector space) then any non-zero element \( r \) of \( A \) lies in some \( m^{n} - m^{n+1} \), so \( r = t^n u \) where \( u \) is a unit (as \( t^n \) generates \( m^n \) by Nakayama, and so does \( r \), so \( K(A) = A_t = A[1/t] \)). So any element of \( K(A) \) can be written uniquely as \( ut^n \) where \( u \) is a unit and \( n \in \mathbb{Z} \).

Thus we can define a valuation \( v(ut^n) = n \).

13.3.D. Exercise. Show that \( v \) is a discrete valuation.

13.3.E. Exercise. Conversely, suppose \((A, m)\) is a discrete valuation ring. Show that \((A, m)\) is a Noetherian regular local ring of dimension 1. (Hint: Show that the ideals are all of the form \((0)\) or \(I_n = \{r \in A : v(r) \geq n\}\), and \((0)\) and \(I_1\) are the only primes. Thus we have Noetherianness, and dimension 1. Show that \(I_1/I_2\) is generated by the image of any element of \(I_1 - I_2\).)

Hence we have proved:

13.3.6. Theorem. — An integral domain \( A \) is a Noetherian local ring of dimension 1 satisfying (a)–(d) if and only if

(e) \( A \) is a discrete valuation ring.

13.3.F. Exercise. Show that there is only one discrete valuation on a discrete valuation ring.

13.3.7. Definition. Thus any Noetherian regular local ring of dimension 1 comes with a unique valuation on its fraction field. If the valuation of an element is \( n > 0 \), we say that the element has a **zero of order** \( n \). If the valuation is \(-n < 0\), we say that the element has a **pole of order** \( n \). We will come back to this shortly, after dealing with (f) and (g).

13.3.8. Theorem. — Suppose \((A, m)\) is a Noetherian local ring of dimension 1. Then (a)–(e) are equivalent to:

(f) \( A \) is a unique factorization domain,

(g) \( A \) is integrally closed in its fraction field \( K = \mathbb{K}(A) \).

Proof. (a)–(e) clearly imply (f), because we have the following stupid unique factorization: each non-zero element of \( r \) can be written uniquely as \( ut^n \) where \( n \in \mathbb{Z}^\geq0 \) and \( u \) is a unit.

Now (f) implies (g), because unique factorization domains are integrally closed in their fraction fields (Exercise 6.4.F).

It remains to check that (g) implies (a)–(e). We will show that (g) implies (b). Suppose \((A, m)\) is a Noetherian local integral domain of dimension 1, integrally closed in its fraction field \( K = \mathbb{K}(A) \). Choose any nonzero \( r \in m \). Then \( S = A/(r) \) is a Noetherian local ring of dimension 0 — its only prime is the image of \( m \), which we denote \( n \) to avoid confusion. Then \( n \) is finitely generated, and each generator is nilpotent (the intersection of all the prime ideals in any ring are the


13.3.9. Geometry of normal Noetherian schemes. We can finally make precise (and generalize) the fact that the function \((x - 2)^2 x/(x - 3)^4\) on \(A^1_3\) has a double zero at \(x = 2\) and a quadruple pole at \(x = 3\). Furthermore, we can say that 75/34 has a double zero at 5, and a single pole at 2. (What are the zeros and poles of \(x^3(x + y)/(x^2 + xy)^3\) on \(A^2_2\)?) Suppose \(X\) is a locally Noetherian scheme. Then for any regular codimension 1 points (i.e. any point \(p\) where \(O_{X,p}\) is a regular local ring of dimension 1), we have a discrete valuation \(v\). If \(f\) is any non-zero element of the fraction field of \(O_{X,p}\) (e.g. if \(X\) is integral, and \(f\) is a non-zero element of the function field of \(X\)), then if \(v(f) > 0\), we say that the element has a zero of order \(v(f)\), and if \(v(f) < 0\), we say that the element has a pole of order \(-v(f)\). (We aren’t yet allowed to discuss order of vanishing at a point that is not regular or codimension 1. One can make a definition, but it doesn’t behave as well as it does when have you have a discrete valuation.)

13.3.G. Exercise. Suppose \(X\) is an integral Noetherian scheme, and \(f \in K(X)^+\) is a non-zero element of its function field. Show that \(f\) has a finite number of zeros and poles. (Hint: reduce to \(X = \text{Spec} \, A\). If \(f = f_1/f_2\), where \(f_i \in A\), prove the result for \(f_1\).

Suppose \(A\) is an Noetherian integrally closed domain. Then it is regular in codimension 1 (translation: its points of codimension at most 1 are regular). If \(A\) is dimension 1, then obviously \(A\) is nonsingular.

13.3.H. Exercise. If \(f\) is a rational function on a Noetherian normal scheme with no poles, show that \(f\) is regular. (Hint: Algebraic Hartogs’ Lemma [12.3.10])

13.3.10. For example (cf. Exercise 13.2.K), \(\text{Spec} \, \mathbb{Z}[i]\) is nonsingular, because it is dimension 1, and \(\mathbb{Z}[i]\) is a unique factorization domain. Hence \(\mathbb{Z}[i]\) is normal, so all its closed (codimension 1) points are nonsingular. Its generic point is also nonsingular, as \(\mathbb{Z}[i]\) is an integral domain.

13.3.11. Remark. A (Noetherian) scheme can be singular in codimension 2 and still be normal. For example, you have shown that the cone \(x^2 + y^2 = z^2\) in \(A^3\) is normal (Exercise 6.4.1(b)), but it is singular at the origin (the Zariski tangent space is visibly three-dimensional).

But singularities of normal schemes are not so bad. For example, we have already seen Hartogs’ Theorem [12.3.10] for Noetherian normal schemes, which states that you could extend functions over codimension 2 sets.
13.3.12. Remark. We know that for Noetherian rings we have implications
unique factorization domain $\implies$ integrally closed $\implies$ regular in codimension 1.
Hence for locally Noetherian schemes, we have similar implications:

factorial $\implies$ normal $\implies$ regular in codimension 1.

Here are two examples to show you that these inclusions are strict.

13.3.13. Example. Suppose $k$ is algebraically closed of characteristic not 2. Then
$k[w, x, y, z]/(wz - xy)$ is integrally closed, but not a unique factorization domain, see Exercise [13.1.D].

13.3.14. Dedekind domains. A Dedekind domain is a Noetherian integral domain
of dimension at most one that is normal (integrally closed in its fraction field).
The localization of a Dedekind domain at any prime but $0$ (i.e. a codimension one prime) is hence a discrete valuation ring. This is an important notion, but we won’t use it much. Rings of integers of number fields are examples, see §10.6.1.

13.3.15. Remark: Serre’s criterion that “normal = $R1+S2$“. Suppose $A$ is a reduced
Noetherian integral domain. Serre’s criterion for normality states that $A$ is normal if and only if $A$ is regular in codimension 1, and every associated prime of a principal ideal generated by a non-zero-divisor is of codimension 1 (i.e. if $b$ is a non-zero-divisor, then $\text{Spec } A/(b)$ has no embedded points). The first hypothesis is sometimes called “$R1$”, and the second is called “Serre’s $S2$ criterion”. The $S2$ criterion says rather precisely what is needed for normality in addition to regularity in codimension 1. We won’t use this, so we won’t prove it here. (See [E] §11.2 for a proof.) Note that the necessity of $R1$ follows from the equivalence of (a) and (g) in Theorem [13.3.A]. An example of a variety satisfying $R1$ but not $S2$ is the knotted plane, Exercise [13.3.D].

13.3.16. Remark: Finitely generated modules over a discrete valuation ring. We record a
useful fact for future reference. Recall that finitely generated modules over a principal ideal domain are finite direct sums of cyclic modules (see for example [DE §12.1, Thm. 5]). Hence any finitely generated module over a discrete valuation ring $A$ with uniformizer $t$ is a finite direct sum of terms $A$ and $A/(t^r)$ (for various $r$). See Proposition [14.7.2] for an immediate consequence.
13.4 Valuative criteria for separatedness and properness

In reasonable circumstances, it is possible to verify separatedness by checking only maps from spectra of discrete valuations rings. There are three reasons you might like this (even if you never use it). First, it gives useful intuition for what separated morphisms look like. Second, given that we understand schemes by maps to them (the Yoneda philosophy), we might expect to understand morphisms by mapping certain maps of schemes to them, and this is how you can interpret the diagram appearing in the valuative criterion. And the third concrete reason is that one of the two directions in the statement is much easier (a special case of the Reduced-to-separated Theorem 11.2.1, see Exercise 13.4.A), and this is the direction we will repeatedly use.

We begin with a valuative criterion that applies in a case that will suffice for the interests of most people, that of finite type morphisms of Noetherian schemes. We will then give a more general version for more general readers.

13.4.1. Theorem (Valuative criterion for separatedness for morphisms of finite type of Noetherian schemes). Suppose \( f : X \to Y \) is a morphism of finite type of Noetherian schemes. Then \( f \) is separated if and only if the following condition holds. For any discrete valuation ring \( A \), and any diagram of the form

\[
\begin{array}{ccc}
\text{Spec } K(A) & \longrightarrow & X \\
& \downarrow f & \\
\text{Spec } A & \longrightarrow & Y
\end{array}
\]

(where the vertical morphism on the left corresponds to the inclusion \( A \hookrightarrow K(A) \)), there is at most one morphism \( \text{Spec } A \to X \) such that the diagram

\[
\begin{array}{ccc}
\text{Spec } K(A) & \longrightarrow & X \\
& \downarrow f & \\
\text{Spec } A & \longrightarrow & Y
\end{array}
\]

commutes.

13.4.A. Exercise (the easy direction). Use the Reduced-to-separated Theorem 11.2.A to prove one direction of the theorem: if \( f \) is separated, then the valuative criterion holds.

13.4.B. Exercise. Suppose \( X \) is an irreducible Noetherian separated curve. If \( p \in X \) is a nonsingular point, then \( O_{X,p} \) is a discrete valuation ring, so each nonsingular point yields a discrete valuation on \( K(X) \). Use the previous exercise to show that distinct points yield distinct valuations.

Here is the intuition behind the valuative criterion (see Figure 13.4). We think of Spec of a discrete valuation ring \( A \) as a “germ of a curve”, and Spec \( K(A) \) as the “germ minus the origin” (even though it is just a point!). Then the valuative criterion says that if we have a map from a germ of a curve to \( Y \), and have a lift of the map away from the origin to \( X \), then there is at most one way to lift the map.
from the entire germ. In the case where $Y$ is a field, you can think of this as saying that limits of one-parameter families are unique (if they exist).

**Figure 13.4.** The line with the doubled origin fails the valuative criterion for separatedness

For example, this captures the idea of what is wrong with the map of the line with the doubled origin over $k$ (Figure 13.5): we take $\text{Spec } A$ to be the germ of the affine line at the origin, and consider the map of the germ minus the origin to the line with doubled origin. Then we have two choices for how the map can extend over the origin.

**Figure 13.5.** The valuative criterion for separatedness

13.4.C. **Exercise.** Make this precise: show that map of the line with doubled origin over $k$ to $\text{Spec } k$ fails the valuative criterion for separatedness. (Earlier arguments were given in Exercises [11.1.D] and [11.1.L].)

13.4.2. **Remark for experts: moduli spaces and the valuative criterion of separatedness.** If $Y = \text{Spec } k$, and $X$ is a (fine) moduli space (a term I won’t define here) of some type of object, then the question of the separatedness of $X$ (over $\text{Spec } k$) has a natural interpretation: given a family of your objects parametrized by a “punctured discrete valuation ring”, is there always at most one way of extending it over the closed point?
13.4.3. Idea behind the proof. (One direction was done in Exercise [13.4.A]) If $f$ is not separated, our goal is to produce a diagram (13.4.1.1) that cannot be completed to (13.4.1.2). If $f$ is not separated, then $\delta : X \to X \times_\gamma X$ is a locally closed immersion that is not a closed immersion.

13.4.D. Exercise. Show that you can find points $p/\in X \times_\gamma X$ and $q/\in X \times_\gamma X$ such that $p/\in q$ and there are no points “between $p$ and $q$” (no points $r$ distinct from $p$ and $q$ with $p/\in r$ and $r/\in q$).

Let $Q$ be the scheme obtained by giving the induced reduced subscheme structure to $\overline{q}$. Let $B = O_{Q,p}$ be the local ring of $Q$ at $p$.

13.4.E. Exercise. Show that $B$ is a Noetherian local integral domain of dimension 1.

If $B$ were regular, then we would be done: composing the inclusion morphism $Q \to X \times_\gamma X$ with the two projections induces the same morphism $q \to X$ but different extensions to $Q$ precisely because $p$ is not in the diagonal. To complete the proof, one shows that the normalization of $B$ is Noetherian; then localizing at any prime above $p$ (there is one by the Lying Over Theorem [8.2.5]) yields the desired discrete valuation ring $A$.

With a more powerful invocation of commutative algebra, we can prove a valuative criterion with much less restrictive hypotheses.

13.4.4. Theorem (Valuative criterion of separatedness). — Suppose $f : X \to Y$ is a quasiseparated morphism. Then $f$ is separated if and only if the following condition holds. For any valuation ring $A$ with function field $K$, and any diagram of the form (13.4.1.1), there is at most one morphism $\text{Spec } A \to X$ such that the diagram (13.4.1.2) commutes.

Because I have already proved something useful that we will never use, I feel no urge to prove this harder fact. The proof of one direction, that separated implies that the criterion holds, follows from the identical argument as in Exercise [13.4.A]

13.4.5. Valuative criteria of properness.

There is a valuative criterion for properness too. It is philosophically useful, and sometimes directly useful, although we won’t need it.

13.4.6. Theorem (Valuative criterion for properness for morphisms of finite type of Noetherian schemes). — Suppose $f : X \to Y$ is a morphism of finite type of locally Noetherian schemes. Then $f$ is proper if and only if for any discrete valuation ring $A$ and any diagram (13.4.1.1), there is exactly one morphism $\text{Spec } A \to X$ such that the diagram (13.4.1.2) commutes.

Recall that the valuative criterion for separatedness was the same, except that $exact$ was replaced by $at$ $most$.

In the case where $Y$ is a field, you can think of this as saying that limits of one-parameter families always exist, and are unique. This is a useful intuition for the notion of properness.

13.4.F. Exercise. Use the valuative criterion of properness to prove that $P^n_A \to \text{Spec } A$ is proper if $A$ is Noetherian. (This is a difficult way to prove a fact that we already showed in Theorem [11.3.5].)
13.4.7. Remarks for experts. There is a moduli-theoretic interpretation similar to that for separatedness (Remark 13.4.2): $X$ is proper if and only if there is always precisely one way of filling in a family over a punctured discrete valuation ring.

Finally, here is a fancier version of the valuative criterion for properness.

13.4.8. Theorem (Valuative criterion of properness). — Suppose $f : X \rightarrow Y$ is a quasiseparated, finite type (hence quasicompact) morphism. Then $f$ is proper if and only if the following condition holds. For any valuation ring $A$ and any diagram of the form (13.4.1.1), there is exactly one morphism $\text{Spec} A \rightarrow X$ such that the diagram (13.4.1.2) commutes.

13.5. Completions

This section will briefly introduce the notion of completions of rings, which generalizes the notion of power series. Our short-term goal is to show that regular local rings appearing on $\overline{K}$-varieties are integral domains (Theorem 13.2.9), and a key fact (13.5.4) that will be used in the proof that nonsingularity for $\overline{K}$-varieties can be checked at closed points (Theorem 13.2.12). But we will also define some types of singularities such as nodes of curves.

13.5.1. Definition. Suppose that $I$ is an ideal of a ring $A$. Define $\hat{A}$ to be $\lim_{\leftarrow} A/I^i$, the completion of $A$ at $I$ (or along $I$).

13.5.A. Exercise. Suppose that $I$ is a maximal ideal $m$. Show that the completion construction factors through localization at $m$. More precisely, make sense of the following diagram, and show that it commutes.

\[
\begin{array}{ccc}
A & \longrightarrow & \hat{A} \\
\downarrow & & \downarrow \\
A_m & \longrightarrow & \hat{A}_m
\end{array}
\]

For this reason, one informally thinks of the information in the completion as coming from an even smaller shred of a scheme than the localization.

13.5.B. Exercise. If $J \subset A$ is an ideal, figure out how to define the completion $\hat{J} \subset \hat{A}$ (an ideal of $\hat{A}$) using $(I + I^m)/I^m \subset A/I^m$. With your definition, you will observe an isomorphism $A/J \cong \hat{A}/\hat{J}$, which is helpful for computing completions in practice.

13.5.2. Definition (cf. Exercise 13.2.H). If $X$ is a $\overline{K}$-variety of pure dimension 1, and $p$ is a closed point, where char $k \neq 2, 3$. We say that $X$ has a node (resp. cusp, tacnode, triple point) at $p$ if $\mathcal{O}_{X,p}$ is isomorphic to the completion of the curve $\text{Spec} \overline{K}[x, y]/(y^2 - x^2)$ (resp. $\text{Spec} \overline{K}[x, y]/(y^3 - x^3)$, $\text{Spec} \overline{K}[x, y]/(y^2 - x^4)$, $\text{Spec} \overline{K}[x, y]/(y^3 - x^3)$). One can define other singularities similarly. You may wish to extend these definitions to more general fields.
Suppose for the rest of this section that \((A, m)\) is Noetherian local ring containing its residue field \(k\) (i.e. it is a \(k\)-algebra), of dimension \(n\). Let \(x_1, \ldots, x_n\) be elements of \(A\) whose images are a basis for \(m/m^2\).

13.5.C. Exercise. Show that the natural map \(A \to \hat{A}\) is an injection. (Hint: Proposition 13.3.2)

13.5.D. Exercise. Show that the map of \(k\)-algebras \(k[[t_1, \ldots, t_n]] \to \hat{A}\) defined by \(t_i \mapsto x_i\) is a surjection. (First be clear why there is such a map!)

13.5.E. Exercise. Show that \(\hat{A}\) is a Noetherian local ring. (Hint: By Exercise 4.6.K, \(k[[t_1, \ldots, t_n]]\) is Noetherian.)

13.5.F. Exercise. Show that \(k[[t_1, \ldots, t_n]]\) is an integral domain. (Possible hint: if \(f, g \in k[[t_1, \ldots, t_n]]\) is nonzero, make sense of its “degree”, and its “leading term”.)

13.5.G. Exercise. Show that \(k[[t_1, \ldots, t_n]]\) is dimension \(n\). (Hint: find a chain of \(n + 1\) prime ideals to show that the dimension is at least \(n\). For the other inequality, use the multi-equation generalization of Krull, Theorem 12.3.7)

13.5.H. Exercise. If \(p \subset A\), show that \(\hat{p}\) is a prime ideal of \(\hat{A}\). (Hint: if \(f, g \notin p\), then let \(m_f, m_g\) be the first “level” where they are not in \(p\) (i.e. the smallest \(m\) such that \(f \notin p/m^m+1\)). Show that \(fg \notin p/m^m+1\).)

13.5.I. Exercise. Show that if \(I \subseteq J \subseteq A\) are nested ideals, then \(\hat{I} \subseteq \hat{J}\). Hence (applying this to prime ideals) show that \(\dim \hat{A} \geq \dim A\).

Suppose for the rest of this section that \((A, m)\) is a regular local ring.

13.5.J. Exercise. Show that \(\dim \hat{A} = \dim A\). (Hint: argue \(\dim \hat{A} \leq \dim m/m^2 = \dim A\).)

13.5.3. Theorem. — Suppose \((A, m)\) is a Noetherian regular local ring containing its residue field \(k\). Then \(k[[t_1, \ldots, t_n]] \to \hat{A}\) is an isomorphism.

(This is basically the Cohen Structure Theorem.) Thus you should think of the map \(A \to \hat{A} = k[[x_1, \ldots, x_n]]\) as sending an element of \(A\) to its power series expansion in the variables \(x_i\).

Proof. We wish to show that \(k[[t_1, \ldots, t_n]] \to \hat{A}\) is injective; we already know it is surjective (Exercise 13.5.D). Suppose \(f \in k[[t_1, \ldots, t_n]]\) maps to 0, so we get a surjection map \(k[[t_1, \ldots, t_n]]/f \to \hat{A}\). Now \(f\) is not a zero-divisor, so by Krull’s Principal Ideal Theorem 12.3.3 the left side has dimension \(n - 1\). But then any quotient of it has dimension at most \(n - 1\), yielding a contradiction.

13.5.K. Exercise. Prove Theorem 13.2.9 that regular local rings containing their residue field are integral domains.

13.5.4. Fact for later. We conclude by mentioning a fact we will use later. Suppose \((A, m)\) is a regular local ring of dimension \(n\), containing its residue field. Suppose
$x_1, \ldots, x_m$ are elements of $m$ such that their images in $m/m^2$ are linearly independent (over $k$). Let $I = (x_1, \ldots, x_m)$. Note that $(A/I, m)$ is a regular local ring: by Krull’s Principal Ideal Theorem \[12.3.3\] $\dim A/I \geq n - m$, and in $A/I$, $m/m^2$ is dimension $n - m$. Thus $I$ is a prime ideal, and $I/I^2$ is an $(A/I)$-module.

13.5.L. Exercise. Show that $\dim_k (I/I^2) \otimes_{A/I} k = n - m$. (Hint: reduce this to a calculation in the completion. It will be convenient to choose coordinates by extending $x_1, \ldots, x_m$ to $x_1, \ldots, x_n$.)
Part V

Quasicoherent sheaves
Quasicoherent and coherent sheaves

Quasicoherent and coherent sheaves generalize the notion of a vector bundle. To motivate them, we first discuss vector bundles, and their interpretation as locally free sheaves.

A free sheaf on \( X \) is an \( O_X \)-module isomorphic to \( \bigoplus I \) where the sum is over some index set \( I \). A locally free sheaf on a ringed space \( X \) is an \( O_X \)-module locally isomorphic to a free sheaf. This corresponds to the notion of a vector bundle (§14.1). Quasicoherent sheaves form a convenient abelian category containing the locally free sheaves that is much smaller than the full category of \( O \)-modules. Quasicoherent sheaves generalize free sheaves in much the way that modules generalize free modules. Coherent sheaves are roughly speaking a finite rank version of quasicoherent sheaves, which form a well-behaved abelian category containing finite rank locally free sheaves (or equivalently, finite rank vector bundles).

14.1 Vector bundles and locally free sheaves

We recall the notion of vector bundles on smooth manifolds. Nontrivial examples to keep in mind are the tangent bundle to a manifold, and the Möbius strip over a circle (interpreted as a line bundle). Arithmetically-minded readers shouldn’t tune out: for example, fractional ideals of the ring of integers in a number field (defined in §10.6.1) turn out to be an example of a “line bundle on a smooth curve” (Exercise 14.1.J).

A rank \( n \) vector bundle on a manifold \( M \) is a fibration \( \pi : V \to M \) with the structure of an \( n \)-dimensional real vector space on \( \pi^{-1}(x) \) for each point \( x \in M \), such that for every \( x \in M \), there is an open neighborhood \( U \) and a homeomorphism

\[
\phi : U \times \mathbb{R}^n \to \pi^{-1}(U)
\]

over \( U \) (so that the diagram

\[
\begin{array}{ccc}
\pi^{-1}(U) & \xleftarrow{=} & U \times \mathbb{R}^n \\
\downarrow{\pi|_{\pi^{-1}(U)}} & & \downarrow{\text{projection to first factor}} \\
U & \downarrow{y} & \\
& & \\
\end{array}
\]

commutes) that is an isomorphism of vector spaces over each \( y \in U \). An isomorphism (14.1.0.1) is called a trivialization over \( U \).

We call \( n \) the rank of the vector bundle. A rank 1 vector bundle is called a line bundle. (It can also be convenient to be agnostic about the rank of the vector
bundle, so it can have different ranks on different connected components. It is also sometimes convenient to consider infinite-rank vector bundles.)

14.1.1. Transition functions. Given trivializations over $U_1$ and $U_2$, over their intersection, the two trivializations must be related by an element $T_{12}$ of $\text{GL}(n)$ with entries consisting of functions on $U_1 \cap U_2$. If $\{U_i\}$ is a cover of $M$, and we are given trivializations over each $U_i$, then the $\{T_{ij}\}$ must satisfy the \textbf{cocycle condition}:

\[(14.1.1.1) \quad T_{ij}|_{U_i \cap U_j \cap U_k} \circ T_{jk}|_{U_j \cap U_i \cap U_k} = T_{ik}|_{U_i \cap U_j \cap U_k}.\]

(This implies $T_{ij} = T_{ji}^{-1}$.) The data of the $T_{ij}$ are called \textbf{transition functions} (or \textbf{transition matrices} for the trivialization.

Conversely, given the data of a cover $\{U_i\}$ and transition functions $T_{ij}$, we can recover the vector bundle (up to unique isomorphism) by “gluing together the various $U_i \times \mathbb{R}^n$ along $U_i \cap U_j$ using $T_{ij}$.”

14.1.2. The sheaf of sections. Fix a rank $n$ vector bundle $V \to M$. The sheaf of sections $F$ of $V$ (Exercise 3.2.G) is an $\mathcal{O}_M$-module — given any open set $U$, we can multiply a section over $U$ by a function on $U$ and get another section. Moreover, given a trivialization over $U$, the sections over $U$ are naturally identified with $n$-tuples of functions.

Thus given a trivialization, over each open set $U_i$, we have an isomorphism $F|_{U_i} \cong \mathcal{O}^\oplus_{U_i}$. We say that such an $F$ is a \textbf{locally free sheaf of rank} $n$. (A sheaf $F$ is \textbf{free of rank} $n$ if $F \cong \mathcal{O}^\oplus_{M}$.)

14.1.3. Transition functions for the sheaf of sections. Suppose we have a vector bundle on $M$, along with a trivialization over an open cover $U_i$. Suppose we have a section of the vector bundle over $M$. (This discussion will apply with $M$ replaced by any open subset.) Then over each $U_i$, the section corresponds to an $n$-tuple of functions over $U_i$.

\[
\begin{array}{c}
U \\
\downarrow \pi \downarrow \\
\pi|_U \downarrow \\
\uparrow \\
\pi|_{U_i} \downarrow \\
U_i \times \mathbb{R}^n
\end{array}
\]

Thus over $U_i$, we have an isomorphism $F|_{U_i} \cong \mathcal{O}^\oplus_{U_i}$. We say that such an $F$ is a \textbf{locally free sheaf of rank} $n$. (A sheaf $F$ is \textbf{free of rank} $n$ if $F \cong \mathcal{O}^\oplus_{M}$.)

14.1.A. Exercise. Show that over $U_i \cap U_j$, the vector-valued function $\vec{s}^i$ is related to $\vec{s}^j$ by the transition functions: $T_{ij} \vec{s}^i = \vec{s}^j$. (Don’t do this too quickly — make sure your $i$’s and $j$’s are on the correct side.)

Given a locally free sheaf $F$ with rank $n$, and a trivializing neighborhood of $F$ (an open cover $\{U_i\}$ such that over each $U_i$, $F|_{U_i} \cong \mathcal{O}^\oplus_{U_i}$ as $\mathcal{O}$-modules), we have transition functions $T_{ij} \in \text{GL}(n, \mathcal{O}(U_i \cap U_j))$ satisfying the cocycle condition $\{14.1.1.1\}$. Thus in conclusion the data of a locally free sheaf of rank $n$ is equivalent to the data of a vector bundle of rank $n$. This change of perspective is useful, and is similar to an earlier change of perspective when we introduced ringed spaces: understanding spaces is the same as understanding (sheaves of) functions on the spaces, and understanding vector bundles (a type of “space over $M$”) is the same as understanding functions.
14.1.4. Definition. A rank 1 locally free sheaf is called an invertible sheaf. (Unimportant aside: “invertible sheaf” is a heinous term for something that is essentially a line bundle. The motivation is that if $X$ is a locally ringed space, and $\mathcal{F}$ and $\mathcal{G}$ are $\mathcal{O}_X$-modules with $\mathcal{F} \otimes \mathcal{G} \cong \mathcal{O}_X$, then $\mathcal{F}$ and $\mathcal{G}$ are invertible sheaves $^{[MO]}$. Thus in the monoid of $\mathcal{O}_X$-modules under tensor product, invertible sheaves are the invertible elements. We will never use this fact.)

14.1.5. Locally free sheaves on schemes.

We can generalize the notion of locally free sheaves to schemes without change. A locally free sheaf of rank $n$ on a scheme $X$ is defined as an $\mathcal{O}_X$-module $\mathcal{F}$ that is locally a free sheaf of rank $n$. Precisely, there is an open cover $\{U_i\}$ of $X$ such that for each $U_i$, $\mathcal{F}|_{U_i} \cong \mathcal{O}_U^{\oplus n}$. This open cover determines transition functions — the data of a cover $\{U_i\}$ of $X$, and functions $T_{ij} \in \text{GL}(n, \mathcal{O}(U_i \cap U_j))$ satisfying the cocycle condition (14.1.1) — which in turn determine the locally free sheaf. As before, given this data, we can find the sections over any open set $U$. Informally, they are sections of the free sheaves over each $U \cap U_i$ that agree on overlaps. More formally, for each $i$, they are $s^i = \left( \begin{array}{c} s_1^i \\ \vdots \\ s_n^i \end{array} \right) \in \Gamma(U \cap U_i, \mathcal{O}_X^n)$, satisfying $T_{ij}s^i = s^j$ on $U \cap U_i \cap U_j$.

You should think of these as vector bundles, but just keep in mind that they are not the “same”, just equivalent notions. We will later (Definition [18.1.4]) define the “total space” of the vector bundle $V \to X$ (a scheme over $X$) in terms of the sheaf version of Spec (precisely, $\text{Spec Sym} V^\bullet$). But the locally free sheaf perspective will prove to be more useful. As one example: the definition of a locally free sheaf is much shorter than that of a vector bundle.

As in our motivating discussion, it is sometimes convenient to let the rank vary among connected components, or to consider infinite rank locally free sheaves.

14.1.6. Useful constructions, in the form of a series of important exercises.

We now give some useful constructions in the form of a series of exercises. Two hints: Exercises [14.1.B] [14.1.G] will apply for ringed spaces in general, so you shouldn’t use special properties of schemes. Furthermore, they are all local on $X$, so you can reduce to the case where the locally free sheaves in question are actually free.

14.1.B. Exercise. Suppose $\mathcal{F}$ and $\mathcal{G}$ are locally free sheaves on $X$ of rank $m$ and $n$ respectively. Show that $\text{Hom}_{\mathcal{O}_X}(\mathcal{F}, \mathcal{G})$ is a locally free sheaf of rank $mn$.

14.1.C. Exercise. If $\mathcal{E}$ is a (finite rank) locally free sheaf on $X$ of rank $n$, Exercise [14.1.B] implies that $\mathcal{E}^\vee := \text{Hom}(\mathcal{E}, \mathcal{O}_X)$ is also a locally free sheaf of rank $n$. This is called the dual of $\mathcal{E}$. Given transition functions for $\mathcal{E}$, describe transition functions for $\mathcal{E}^\vee$. (Note that if $\mathcal{E}$ is rank 1, i.e. invertible, the transition functions of the dual are the inverse of the transition functions of the original.) Show that $\mathcal{E} \cong \mathcal{E}^{\vee \vee}$. (Caution: your argument showing that there is a canonical isomorphism $(\mathcal{F}^{\vee})^{\vee} \cong \mathcal{F}$ better not also show that there is an isomorphism $\mathcal{F}^{\vee} \cong \mathcal{F}$! We will see an example in [15.1] of a locally free $\mathcal{F}$ that is not isomorphic to its dual: the invertible sheaf $\mathcal{O}(1)$ on $\mathbb{P}^n$.)
14.1.D. Exercise. If \( \mathcal{F} \) and \( \mathcal{G} \) are locally free sheaves, show that \( \mathcal{F} \otimes \mathcal{G} \) is a locally free sheaf. (Here \( \otimes \) is tensor product as \( \mathcal{O}_X \)-modules, defined in Exercise 3.5.H.) If \( \mathcal{F} \) is an invertible sheaf, show that \( \mathcal{F} \otimes \mathcal{F}^\vee \cong \mathcal{O}_X \).

14.1.E. Exercise. Recall that tensor products tend to be only right-exact in general. Show that tensoring by a locally free sheaf is exact. More precisely, if \( \mathcal{F} \) is a locally free sheaf, and \( \mathcal{G}' \to \mathcal{G} \to \mathcal{G}'' \) is an exact sequence of \( \mathcal{O}_X \)-modules, then so is \( \mathcal{G}' \otimes \mathcal{F} \to \mathcal{G} \otimes \mathcal{F} \to \mathcal{G}'' \otimes \mathcal{F} \). (Possible hint: it may help to check exactness by checking exactness at stalks. Recall that the tensor product of stalks can be identified with the stalk of the tensor product, so for example there is a “natural” isomorphism \( (\mathcal{G} \otimes_{\mathcal{O}_X} \mathcal{F})_x \cong \mathcal{G}_x \otimes_{\mathcal{O}_{X,x}} \mathcal{F}_x \), Exercise 3.5.H(b)).

14.1.F. Exercise. If \( \mathcal{E} \) is a locally free sheaf of finite rank, and \( \mathcal{F} \) and \( \mathcal{G} \) are \( \mathcal{O}_X \)-modules, show that \( \text{Hom}(\mathcal{F}, \mathcal{G} \otimes \mathcal{E}) \cong \text{Hom}(\mathcal{F} \otimes \mathcal{E}^\vee, \mathcal{G}) \). (Possible hint: first consider the case where \( \mathcal{E} \) is free.)

14.1.G. Exercise and Important Definition. Show that the invertible sheaves on \( X \), up to isomorphism, form an abelian group under tensor product. This is called the Picard group of \( X \), and is denoted \( \text{Pic} X \).

Unlike the previous exercises, the next one is specific to schemes.

14.1.H. Exercise. Suppose \( s \) is a section of a locally free sheaf \( \mathcal{F} \) on a scheme \( X \). Define the notion of the subscheme cut out by \( s = 0 \). (Hint: given a trivialization over an open set \( U \), \( s \) corresponds to a number of functions \( f_1, \ldots \) on \( U \); on \( U \), take the scheme cut out by these functions.)


We define rational (and regular) sections of a locally free sheaf on a scheme \( X \) just as we did rational (and regular) functions (see for example §6.5 and §7.5).

14.1.I. Less Important Exercise. Show that locally free sheaves on Noetherian normal schemes satisfy “Hartogs’ lemma”: sections defined away from a set of codimension at least 2 extend over that set. (Hartogs’ lemma for Noetherian normal schemes is Theorem 12.3.10.)

14.1.8. Remark. Based on your intuition for line bundles on manifolds, you might hope that every point has a “small” open neighborhood on which all invertible sheaves (or locally free sheaves) are trivial. Sadly, this is not the case. We will eventually see (21.10.1) that for the curve \( y^2 - x^3 - x = 0 \) in \( \mathbb{A}^2 \), every nonempty open set has nontrivial invertible sheaves. (This will use the fact that it is an open subset of an elliptic curve.)

14.1.J. Exercise (for those with sufficient arithmetic background; see also Proposition 15.2.6 and 15.2.9). Recall the definition of the ring of integers \( \mathcal{O}_K \) in a number field \( K \), Remark 10.6.1. A fractional ideal \( \mathfrak{a} \) of \( \mathcal{O}_K \) is an \( \mathcal{O}_K \)-submodule of \( K \) such that there is a nonzero \( \alpha \in \mathcal{O}_K \) such that \( \alpha \mathfrak{a} \subset \mathcal{O}_K \). Products of fractional ideals are defined analogously to products of ideals in a ring: \( \mathfrak{a} \mathfrak{b} \) consists of (finite) \( \mathcal{O}_K \)-linear combinations of products of elements of \( \mathfrak{a} \) and elements of \( \mathfrak{b} \). Thus fractional ideals form a semigroup under multiplication, with \( \mathcal{O}_K \) as the identity. In fact fractional ideals of \( \mathcal{O}_K \) form a group.
(a) Explain how a fractional ideal on a ring of integers in a number field yields an invertible sheaf.

(b) A fractional ideal is principal if it is of the form $r\mathcal{O}_K$ for some $r \in K$. Show that any two that differ by a principal ideal yield the same invertible sheaf.

(c) Show that two fractional ideals that yield the same invertible sheaf differ by a principal ideal.

The class group is defined to be the group of fractional ideals modulo the principal ideals (i.e. modulo $K^\times$). This exercise shows that the class group is (isomorphic to) the Picard group of $\mathcal{O}_K$. (This discussion applies to the ring of integers in any global field.)

14.1.9. The problem with locally free sheaves.

Recall that $\mathcal{O}_X$-modules form an abelian category: we can talk about kernels, cokernels, and so forth, and we can do homological algebra. Similarly, vector spaces form an abelian category. But locally free sheaves (i.e. vector bundles), along with reasonably natural maps between them (those that arise as maps of $\mathcal{O}_X$-modules), don’t form an abelian category. As a motivating example in the category of differentiable manifolds, consider the map of the trivial line bundle on $\mathbb{R}$ (with coordinate $t$) to itself, corresponding to multiplying by the coordinate $t$. Then this map jumps rank, and if you try to define a kernel or cokernel you will get confused.

This problem is resolved by enlarging our notion of nice $\mathcal{O}_X$-modules in a natural way, to quasicoherent sheaves. (You can turn this into two definitions of quasicoherent sheaves, equivalent to those we will give. We want a notion that is local on $X$ of course. So we ask for the smallest abelian subcategory of $\text{Mod}_{\mathcal{O}_X}$ that is “local” and includes vector bundles. It turns out that the main obstruction to vector bundles to be an abelian category is the failure of cokernels of maps of locally free sheaves — as $\mathcal{O}_X$-modules — to be locally free; we could define quasicoherent sheaves to be those $\mathcal{O}_X$-modules that are locally cokernels, yielding the definition at the start of the chapter. You may wish to later check that our future definitions are equivalent to these.)

$$\mathcal{O}_X\text{-modules} \supset \text{quasicoherent sheaves} \supset \text{locally free sheaves}$$

(abelian category) (abelian category) (not an abelian category)

Similarly, finite rank locally free sheaves will sit in a nice smaller abelian category, that of coherent sheaves.

$$\text{quasicoherent sheaves} \supset \text{coherent sheaves} \supset \text{finite rank locally free sheaves}$$

(abelian category) (abelian category) (not an abelian category)

14.1.10. Remark: Quasicoherent and coherent sheaves on ringed spaces in general. We will discuss quasicoherent and coherent sheaves on schemes, but they can be defined more generally on ringed spaces. Many of the results we state will hold in this greater generality, but because the proofs look slightly different, we restrict ourselves to schemes to avoid distraction.

14.2 Quasicoherent sheaves
We now define the notion of quasicoherent sheaf. In the same way that a scheme is defined by “gluing together rings”, a quasicoherent sheaf over that scheme is obtained by “gluing together modules over those rings”. Given an $A$-module $M$, we defined an $O$-module $\tilde{M}$ on $\text{Spec } A$ long ago (Exercise 5.1.D) — the sections over $D(f)$ were $M_f$.

14.2.1. Theorem. — Let $X$ be a scheme, and $F$ an $O_X$-module. Then let $P$ be the property of affine open sets that $F|_{\text{Spec } A} \cong \tilde{M}$ for an $A$-module $M$. Then $P$ satisfies the two hypotheses of the Affine Communication Lemma 6.3.2.

We prove this in a moment.

14.2.2. Definition. If $X$ is a scheme, then an $O_X$-module $F$ is quasicoherent if for every affine open subset $\text{Spec } A \subset X$, $F|_{\text{Spec } A} \cong \tilde{M}$ for some $A$-module $M$. By Theorem 14.2.1, it suffices to check this for a collection of affine open sets covering $X$. For example, $\tilde{M}$ is a quasicoherent sheaf on $X$, and all locally free sheaves on $X$ are quasicoherent.

14.2.A. IMPORTANT EXERCISE (NOT EVERY $O_X$-MODULE IS A QUASICOHERENT SHEAF). (a) Suppose $X = \text{Spec } k[t]$. Let $F$ be the skyscraper sheaf supported at the origin $[t]$, with group $k[t]$ and the usual $k[t]$-module structure. Show that this is an $O_X$-module that is not a quasicoherent sheaf. (More generally, if $X$ is an integral scheme, and $p \in X$ that is not the generic point, we could take the skyscraper sheaf at $p$ with group the function field of $X$. Except in a silly circumstances, this sheaf won’t be quasicoherent.) See Exercises 9.1.D and 14.3.F for more (pathological) examples of $O_X$-modules that are not quasicoherent.

(b) Suppose $X = \text{Spec } k[t]$. Let $F$ be the skyscraper sheaf supported at the generic point $[[0]]$, with group $k(t)$. Give this the structure of an $O_X$-module. Show that this is a quasicoherent sheaf. Describe the restriction maps in the distinguished topology of $X$. (Remark: your argument will apply more generally, for example when $X$ is an integral scheme with generic point $\eta$, and $F$ is the skyscraper sheaf $i_\eta,*,K(X)$.)

14.2.B. IMPORTANT EXERCISE (NOT EVERY QUASICOHERENT SHEAF IS LOCALLY FREE). Use the example of Exercise 14.2.A(b) to show that not every quasicoherent sheaf is locally free.

Proof of Theorem 14.2.1. Clearly if $\text{Spec } A$ has property $P$, then so does the distinguished open $\text{Spec } A_f$: if $M$ is an $A$-module, then $\tilde{M}|_{\text{Spec } A_f} \cong \tilde{M}_{if}$ as sheaves of $O_{\text{Spec } A_i}$-modules (both sides agree on the level of distinguished open sets and their restriction maps).

We next show the second hypothesis of the Affine Communication Lemma 6.3.2. Suppose we have modules $M_1, \ldots, M_n$, where $M_i$ is an $A_{i,-}$-module, along with isomorphisms $\phi_{ij} : (M_i)|_{f_i} \to (M_j)|_{f_i}$ of $A_{i,f_j}$-modules, satisfying the cocycle condition (14.1.1.1). We want to construct an $M$ such that $\tilde{M}$ gives us $\tilde{M}$ on $D(f_i) = \text{Spec } A_{i,}$, or equivalently, isomorphisms $\rho_i : \Gamma(D(f_i), \tilde{M}) \to M_i$, so that
the bottom triangle of

\[ (14.2.2.1) \]

\[
\begin{array}{c}
\otimes_{A_{f_1}} M \\
\phi_i \\
\otimes_{A_{f_i}} (M_i)_{f_i} \\
\otimes_{A_{f_j}} (M_j)_{f_j} \\
\otimes_{A_{f_k}} (M_k)_{f_k} \\
\otimes_{A_{f_l}} (M_l)_{f_l} \\
\otimes_{A_{f_m}} (M_m)_{f_m} \\
\end{array}
\]

commutes.

14.2.C. EXERCISE. Why does this suffice to prove the result? In other words, why does this imply that \( F_{\text{Spec } A} \cong \hat{M} \)?

We already know that \( M \) should be \( \Gamma(F, \text{Spec } A) \), as \( F \) is a sheaf. Consider elements of \( M_1 \times \cdots \times M_n \) that “agree on overlaps”; let this set be \( M \). In other words,

\[ (14.2.2.2) \quad 0 \longrightarrow M \longrightarrow M_1 \times \cdots \times M_n \longrightarrow M_{12} \times M_{13} \times \cdots \times M_{(n-1)n} \]

is an exact sequence (where \( M_{ij} = (M_i)_{f_i} \cong (M_j)_{f_j} \) and the map \( \gamma \) is the “difference” map. So \( M \) is a kernel of a morphism of \( A \)-modules, hence an \( A \)-module. We are left to show that \( M_i \cong M_{f_i} \) (and that this isomorphism satisfies (14.2.2.1)). (At this point, we may proceed in a number of ways, and the reader may wish to find their own route rather than reading on.)

For convenience assume \( i = 1 \). Localization is exact (Exercise 2.6.F(a)), so tensoring (14.2.2.2) by \( A_{f_1} \) yields

\[ (14.2.2.3) \quad 0 \longrightarrow M_{f_1} \longrightarrow (M_1)_{f_1} \times (M_2)_{f_1} \times \cdots \times (M_n)_{f_1} \]

\[ \longrightarrow M_{12} \times \cdots \times M_{1n} \times (M_{23})_{f_1} \times \cdots \times (M_{(n-1)n})_{f_1} \]

is an exact sequence of \( A_{f_1} \)-modules.

We now identify many of the modules appearing in (14.2.2.3) in terms of \( M_1 \). First of all, \( f_1 \) is invertible in \( A_{f_1} \), so \( (M_1)_{f_1} \) is canonically \( M_1 \). Also, \( (M_1)_{f_1} \cong (M_1)_{f_j} \) via \( \phi_{ij} \). Hence if \( i, j \neq 1 \), \( (M_i)_{f_1} \cong (M_1)_{f_i} \) via \( \phi_{i1} \) and \( \phi_{1j} \) (here the cocycle condition is implicitely used). Furthermore, \( (M_{11})_{f_i} \cong (M_1)_{f_i} \) via \( \phi_{i1} \). Thus we can write (14.2.2.3) as

\[ (14.2.2.4) \quad 0 \longrightarrow M_{f_1} \longrightarrow M_1 \times (M_1)_{f_2} \times \cdots \times (M_1)_{f_n} \]

\[ \longrightarrow \alpha (M_1)_{f_2} \times \cdots \times (M_1)_{f_n} \times (M_1)_{f_3} \times \cdots \times (M_1)_{f_{n-1}f_n} \]

By assumption, \( F_{\text{Spec } A_{f_1}} \cong \hat{M}_1 \) for some \( M_1 \), so by considering the cover

\[ \text{Spec } A_{f_1} = \text{Spec } A_{f_1} \cup \text{Spec } A_{f_1f_2} \cup \text{Spec } A_{f_1f_3} \cup \cdots \cup \text{Spec } A_{f_1f_n} \]
(notice the “redundant” first term), and identifying sections of $\mathcal{F}$ over $\text{Spec } \mathbb{A}_f$, in terms of sections over the open sets in the cover and their pairwise overlaps, we have an exact sequence of $\mathbb{A}_{f_1}$-modules

$$0 \longrightarrow M_1 \longrightarrow M_1 \times (M_1)_{f_2} \times \cdots \times (M_1)_{f_n} \longrightarrow (M_1)_{f_2} \times \cdots \times (M_1)_{f_n} \times (M_1)_{f_2 f_3} \times \cdots \times (M_1)_{f_{n-1} f_n} \longrightarrow \beta$$

which is very similar to (14.2.2.4). Indeed, the final map $\beta$ of the above sequence is the same as the map $\alpha$ of (14.2.2.4), so $\ker \alpha = \ker \beta$, i.e. we have an isomorphism $M_1 \cong M_{f_1}$.

Finally, the triangle of (14.2.2.1) is commutative, as each vertex of the triangle can be identified as the sections of $\mathcal{F}$ over $\text{Spec } \mathbb{A}_{f_1 f_2}$.

\section*{14.3 Characterizing quasicoherence using the distinguished affine base}

Because quasicoherent sheaves are locally of a very special form, in order to “know” a quasicoherent sheaf, we need only know what the sections are over every affine open set, and how to restrict sections from an affine open set $U$ to a distinguished affine open subset of $U$. We make this precise by defining what I will call the distinguished affine base of the Zariski topology — not a base in the usual sense. The point of this discussion is to give a useful characterization of quasicoherence, but you may wish to just jump to §14.3.3.

The open sets of the distinguished affine base are the affine open subsets of $X$. We have already observed that this forms a base. But forget that fact. We like distinguished open sets $\text{Spec } \mathbb{A}_f \hookrightarrow \text{Spec } \mathbb{A}$, and we don’t really understand open immersions of one random affine open subset in another. So we just remember the “nice” inclusions.

\section*{14.3.1. Definition.} The distinguished affine base of a scheme $X$ is the data of the affine open sets and the distinguished inclusions.

In other words, we remember only some of the open sets (the affine open sets), and only some of the morphisms between them (the distinguished morphisms). For experts: if you think of a topology as a category (the category of open sets), we have described a subcategory.

We can define a sheaf on the distinguished affine base in the obvious way: we have a set (or abelian group, or ring) for each affine open set, and we know how to restrict to distinguished open sets.

Given a sheaf $\mathcal{F}$ on $X$, we get a sheaf on the distinguished affine base. You can guess where we are going: we will show that all the information of the sheaf is contained in the information of the sheaf on the distinguished affine base.

As a warm-up, we can recover stalks as follows. (We will be implicitly using only the following fact. We have a collection of open subsets, and some subsets, such that if we have any $x \in U, V$ where $U$ and $V$ are in our collection of open sets, there is some $W$ containing $x$, and contained in $U$ and $V$ such that $W \hookrightarrow U$ and $W \hookrightarrow V$ are both in our collection of inclusions. In the case we are considering here,
this is the key Proposition 6.3.1 that given any two affine open sets Spec $A$, Spec $B$ in $X$, Spec $A \cap$ Spec $B$ could be covered by affine open sets that were simultaneously distinguished in Spec $A$ and Spec $B$. In fancy language: the category of affine open sets, and distinguished inclusions, forms a filtered set.)

The stalk $F_x$ is the colimit $\lim \limits_{\leftarrow} (f \in F(U))$ where the limit is over all open sets contained in $X$. We compare this to $\lim \limits_{\leftarrow} (f \in F(U))$ where the limit is over all affine open sets, and all distinguished inclusions. You can check that the elements of one correspond to elements of the other. (Think carefully about this!)

14.3.A. Exercise. Show that a section of a sheaf on the distinguished affine base is determined by the section's germs.

14.3.2. Theorem. —

(a) A sheaf on the distinguished affine base $F^b$ determines a unique sheaf $F$, which when restricted to the affine base is $F^b$. (Hence if you start with a sheaf, and take the sheaf on the distinguished affine base, and then take the induced sheaf, you get the sheaf you started with.)

(b) A morphism of sheaves on a distinguished affine base uniquely determines a morphism of sheaves.

c) An $O_X$-module “on the distinguished affine base” yields an $O_X$-module.

This proof is identical to our argument of §3.7 showing that sheaves are (essentially) the same as sheaves on a base, using the “sheaf of compatible germs” construction. The main reason for repeating it is to let you see that all that is needed is for the open sets to form a filtered set (or in the current case, that the category of open sets and distinguished inclusions is filtered).

For experts: (a) and (b) are describing an equivalence of categories between sheaves on the Zariski topology of $X$ and sheaves on the distinguished affine base of $X$.

Proof. (a) Suppose $F^b$ is a sheaf on the distinguished affine base. Then we can define stalks.

For any open set $U$ of $X$, define the sheaf of compatible germs

$$F(U) := \{ (f_x \in F^b_x)_{x \in U} : \text{for all } x \in U, \text{ there exists } U_x \text{ with } x \subset U_x \subset U, F^x \in F^b(U_x) \text{ such that } F^y_x = f_y \text{ for all } y \in U_x \}$$

where each $U_x$ is in our base, and $F^y_x$ means “the germ of $F^x$ at $y$”. (As usual, those who want to worry about the empty set are welcome to.)

This is a sheaf: convince yourself that we have restriction maps, identity, and gluability, really quite easily.

I next claim that if $U$ is in our base, that $F(U) = F^b(U)$. We clearly have a map $F^b(U) \to F(U)$. This is an isomorphism on stalks, and hence an isomorphism by Exercise 3.4.E


14.3.C. Exercise. Prove (c).
14.3.3. **A characterization of quasicoherent sheaves in terms of distinguished inclusions.** We use this perspective to give a useful characterization of quasicoherent sheaves. Suppose $\text{Spec } A_f \hookrightarrow \text{Spec } A \subset X$ is a distinguished open set. Let $\phi: \Gamma(\text{Spec } A, \mathcal{F}) \rightarrow \Gamma(\text{Spec } A_f, \mathcal{F})$ be the restriction map. The source of $\phi$ is an $A$-module, and the target is an $A_f$-module, so by the universal property of localization, $\phi$ naturally factors as:

$$
\begin{array}{ccc}
\Gamma(\text{Spec } A, \mathcal{F}) & \xrightarrow{\phi} & \Gamma(\text{Spec } A_f, \mathcal{F}) \\
\downarrow & & \uparrow \alpha \\
\Gamma(\text{Spec } A, \mathcal{F})_f & & \end{array}
$$

14.3.D. **Very important exercise.** Show that an $\mathcal{O}_X$-module $\mathcal{F}$ is quasicoherent if and only if for each such distinguished $\text{Spec } A_f \hookrightarrow \text{Spec } A$, $\alpha$ is an isomorphism.

Thus a quasicoherent sheaf is (equivalent to) the data of one module for each affine open subset (a module over the corresponding ring), such that the module over a distinguished open set $\text{Spec } A_f$ is given by localizing the module over $\text{Spec } A$. The next exercise shows that this will be an easy criterion to check.

14.3.E. **Important exercise (cf. Lemma [8.3.4]).** Suppose $X$ is a quasicompact and quasiseparated scheme (i.e. covered by a finite number of affine open sets, the pairwise intersection of which is also covered by a finite number of affine open sets). Suppose $\mathcal{F}$ is a quasicoherent sheaf on $X$, and let $f \in \Gamma(X, \mathcal{O}_X)$ be a function on $X$. Show that the restriction map $\text{res}_{X_f \subset X} : \Gamma(X, \mathcal{F}) \rightarrow \Gamma(X_f, \mathcal{F})$ (here $X_f$ is the open subset of $X$ where $f$ doesn’t vanish) is precisely localization. In other words show that there is an isomorphism $\Gamma(X, \mathcal{F})_f \rightarrow \Gamma(X_f, \mathcal{F})$ making the following diagram commute.

$$
\begin{array}{ccc}
\Gamma(X, \mathcal{F}) & \xrightarrow{\text{res}_{X_f \subset X}} & \Gamma(X_f, \mathcal{F}) \\
\downarrow & & \downarrow \sim \\
\otimes_A A_f & & \Gamma(X, \mathcal{F})_f \\
\end{array}
$$

All that you should need in your argument is that $X$ admits a cover by a finite number of open sets, and that their pairwise intersections are each quasicompact. (Hint: Apply the exact functor $\otimes_A A_f$ to the exact sequence

$$
0 \rightarrow \Gamma(X, \mathcal{F}) \rightarrow \oplus_i \Gamma(U_i, \mathcal{F}) \rightarrow \oplus \Gamma(U_{ijk}, \mathcal{F})
$$

where the $U_i$ form a finite affine cover of $X$ and $U_{ijk}$ form a finite affine cover of $U_i \cap U_j$.)

14.3.F. **Less important exercise.** Give a counterexample to show that the above statement need not hold if $X$ is not quasicompact. (Possible hint: take an infinite disjoint union of affine schemes. The key idea is that infinite direct products do not commute with localization.)
14.3.G. **IMPORTANT EXERCISE** (COROLLARY TO EXERCISE 14.3.E: **PUSHFORWARDS OF QUASICOHERENT SHEAVES ARE QUASICOHERENT IN NON-PATHOLOGICAL CIRCUMSTANCES**). Suppose $\pi : X \to Y$ is a quasicompact quasiseparated morphism, and $\mathcal{F}$ is a quasicoherent sheaf on $X$. Show that $\pi_*\mathcal{F}$ is a quasicoherent sheaf on $Y$.

14.3.4. **Grothendieck topologies.** The distinguished affine base isn’t a topology in the usual sense — the union of two affine sets isn’t necessarily affine, for example. It is however a first new example of a generalization of a topology — the notion of a **site** or a **Grothendieck topology**. We give the definition to satisfy the curious, but we certainly won’t use this notion. (For a clean statement, see [Stacks 00VH]; this is intended only as motivation.) The idea is that we should abstract away only those notions we need to define sheaves. We need the notion of open set, but it turns out that we won’t even need an underlying set, i.e. we won’t even need the notion of points! Let’s think through how little we need. For our discussion of sheaves to work, we needed to know what the open sets were, and what the (allowed) inclusions were, and these should “behave well”, and in particular the data of the open sets and inclusions should form a category. (For example, the composition of an allowed inclusion with another allowed inclusion should be an allowed inclusion — in the distinguished affine base, a distinguished open set of a distinguished open set is a distinguished open set.) So we just require the data of this category. At this point, we can already define presheaf (as just a contravariant functor from this category of “open sets”). We saw this idea earlier in Exercise 3.2.A.

In order to extend this definition to that of a sheaf, we need to know more information. We want two open subsets of an open set to intersect in an open set, so we want the category to be closed under fiber products (cf. Exercise 2.3.N). For the identity and gluability axioms, we need to know when some open sets cover another, so we also remember this as part of the data of a Grothendieck topology. This data of the coverings satisfy some obvious properties. Every open set covers itself (i.e. the identity map in the category of open sets is a covering). Coverings pull back: if we have a map $Y \to X$, then any cover of $X$ pulls back to a cover of $Y$. Finally, a cover of a cover should be a cover. Such data (satisfying these axioms) is called a **Grothendieck topology** or a **site**. We can define the notion of a sheaf on a Grothendieck topology in the usual way, with no change. A **topos** is a scary name for a category of sheaves on a Grothendieck topology.

Grothendieck topologies are used in a wide variety of contexts in and near algebraic geometry. Etale cohomology (using the etale topology), a generalization of Galois cohomology, is a central tool, as are more general flat topologies, such as the smooth topology. The definition of a Deligne-Mumford or Artin stack uses the etale and smooth topologies, respectively. Tate developed a good theory of non-archimedean analytic geometry over totally disconnected ground fields such as $\mathbb{Q}_p$ using a suitable Grothendieck topology. Work in K-theory (related for example to Voevodsky’s work) uses exotic topologies.

14.4 Quasicoherent sheaves form an abelian category
The category of $A$-modules is an abelian category. Indeed, this is our motivating example for the notion of abelian category. Similarly, quasicoherent sheaves on a scheme $X$ form an abelian category, which we call $\text{QCoh}_X$. Here is how.

When you show that something is an abelian category, you have to check many things, because the definition has many parts. However, if the objects you are considering lie in some ambient abelian category, then it is much easier. You have seen this idea before: there are several things you have to do to check that something is a group. But if you have a subset of group elements, it is much easier to check that it forms a subgroup.

You can look at back at the definition of an abelian category, and you’ll see that in order to check that a subcategory is an abelian subcategory, you need to check only the following things:

(i) $0$ is in the subcategory
(ii) the subcategory is closed under finite sums
(iii) the subcategory is closed under kernels and cokernels

In our case of $\text{QCoh}_X \subset \text{Mod}_{\mathcal{O}_X}$, the first two are cheap: $0$ is certainly quasicoherent, and the subcategory is closed under finite sums: if $\mathcal{F}$ and $\mathcal{G}$ are sheaves on $X$, and over $\text{Spec} A$, $\mathcal{F} \cong \mathcal{M}$ and $\mathcal{G} \cong \mathcal{N}$, then $\mathcal{F} \oplus \mathcal{G} = \mathcal{M} \oplus \mathcal{N}$ (do you see why?), so $\mathcal{F} \oplus \mathcal{G}$ is a quasicoherent sheaf.

We now check (iii), using the characterization of Important Exercise [14.3.3]. Suppose $\alpha: \mathcal{F} \rightarrow \mathcal{G}$ is a morphism of quasicoherent sheaves. Then on any affine open set $U$, where the morphism is given by $\beta: M \rightarrow N$, define $(\ker \alpha)(U) = \ker \beta$ and $(\coker \alpha)(U) = \coker \beta$. Then these behave well under inversion of a single element: if

$$0 \rightarrow K \rightarrow M \rightarrow N \rightarrow P \rightarrow 0$$

is exact, then so is

$$0 \rightarrow K_f \rightarrow M_f \rightarrow N_f \rightarrow P_f \rightarrow 0,$$

from which $(\ker \beta)_f \cong \ker(\beta_f)$ and $(\coker \beta)_f \cong \coker(\beta_f)$. Thus both of these define quasicoherent sheaves. Moreover, by checking stalks, they are indeed the kernel and cokernel of $\alpha$ (exactness can be checked stalk-locally). Thus the quasicoherent sheaves indeed form an abelian category.

14.4.A. Exercise. Show that a sequence of quasicoherent sheaves $\mathcal{F} \rightarrow \mathcal{G} \rightarrow \mathcal{H}$ on $X$ is exact if and only if it is exact on each open set in any given affine cover of $X$. (In particular, taking sections over an affine open Spec $A$ is an exact functor from the category of quasicoherent sheaves on $X$ to the category of $A$-modules. Recall that taking sections is only left-exact in general, see [3.5.E]). In particular, we may check injectivity or surjectivity of a morphism of quasicoherent sheaves by checking on an affine cover of our choice.

Warning: If $0 \rightarrow \mathcal{F} \rightarrow \mathcal{G} \rightarrow \mathcal{H} \rightarrow 0$ is an exact sequence of quasicoherent sheaves, then for any open set

$$0 \rightarrow \mathcal{F}(U) \rightarrow \mathcal{G}(U) \rightarrow \mathcal{H}(U)$$

is exact, and exactness on the right is guaranteed to hold only if $U$ is affine. (To set you up for cohomology: whenever you see left-exactness, you expect to eventually interpret this as a start of a long exact sequence. So we are expecting $H^1$’s on the right, and now we expect that $H^1(\text{Spec} A, \mathcal{F}) = 0$. This will indeed be the case.)
14.4.B. LESS IMPORTANT EXERCISE (CONNECTION TO ANOTHER DEFINITION, AND QUASICOHERENT SHEAVES ON RINGED SPACES IN GENERAL). Show that an \( \mathcal{O}_X \)-module \( \mathcal{F} \) on a scheme \( X \) is quasicoherent if and only if there exists an open cover by \( U_i \) such that on each \( U_i \), \( \mathcal{F}|_{U_i} \) is isomorphic to the cokernel of a map of two free sheaves:

\[
\mathcal{O}^{\oplus I}_{U_i} \rightarrow \mathcal{O}^{\oplus J}_{U_i} \rightarrow \mathcal{F}|_{U_i} \rightarrow 0
\]

is exact. We have thus connected our definitions to the definition given at the very start of the chapter. This is the definition of a quasicoherent sheaf on a ringed space in general. It is useful in many circumstances, for example in complex analytic geometry.

14.5 Module-like constructions

In a similar way, basically any nice construction involving modules extends to quasicoherent sheaves. (One exception: the Hom of two \( \Lambda \)-modules is an \( \Lambda \)-module, but the Hom of two quasicoherent sheaves is quasicoherent only in “reasonable” circumstances, see Exercise [14.7.A].)

14.5.1. Locally free sheaves from free modules.

14.5.A. EXERCISE (POSSIBLE HELP FOR LATER PROBLEMS). (a) Suppose

\[
0 \rightarrow \mathcal{F}' \rightarrow \mathcal{F} \rightarrow \mathcal{F}'' \rightarrow 0
\]

is a short exact sequence of locally free sheaves on \( X \). Suppose \( U = \text{Spec } \Lambda \) is an affine open set where \( \mathcal{F}', \mathcal{F}'' \) are free, say \( \mathcal{F}'|_{\text{Spec } \Lambda} = \Lambda^{\oplus a}, \mathcal{F}''|_{\text{Spec } \Lambda} = \Lambda^{\oplus b} \). (Here \( a \) and \( b \) are assumed to be finite for convenience, but this is not necessary, so feel free to generalize to the infinite rank case.) Show that \( \mathcal{F} \) is also free, and that \( 0 \rightarrow \mathcal{F}' \rightarrow \mathcal{F} \rightarrow \mathcal{F}'' \rightarrow 0 \) can be interpreted as coming from the tautological exact sequence \( 0 \rightarrow \Lambda^{\oplus a} \rightarrow \Lambda^{\oplus (a+b)} \rightarrow \Lambda^{\oplus b} \rightarrow 0 \). (As a consequence, given an exact sequence of quasicoherent sheaves where \( \mathcal{F}' \) and \( \mathcal{F}'' \) are locally free, \( \mathcal{F} \) must also be locally free.)

(b) In the finite rank case, show that given such an open cover, the transition functions (really, matrices) of \( \mathcal{F} \) may be interpreted as block upper-diagonal matrices, where the top \( a \times a \) block are transition functions for \( \mathcal{F}' \), and the bottom \( b \times b \) blocks are transition functions for \( \mathcal{F}'' \).

14.5.B. EXERCISE. Suppose \( 0 \rightarrow \mathcal{F}' \rightarrow \mathcal{F} \rightarrow \mathcal{F}'' \rightarrow 0 \) is an exact sequence of quasicoherent sheaves on \( X \). (a) If \( \mathcal{F}' \) and \( \mathcal{F}'' \) are locally free, show that \( \mathcal{F} \) is locally free. (Hint: Use the previous exercise.)

(b) If \( \mathcal{F} \) and \( \mathcal{F}' \) are locally free of finite rank, show that \( \mathcal{F}' \) is too. (Hint: Reduce to the case \( X = \text{Spec } \Lambda \) and \( \mathcal{F} \) and \( \mathcal{F}'' \) free. Interpret the map \( \phi : \mathcal{F} \rightarrow \mathcal{F}'' \) as an \( n \times m \) matrix \( M \) with values in \( \Lambda \), with \( m \) the rank of \( \mathcal{F} \) and \( n \) the rank of \( \mathcal{F}'' \). For each point \( p \) of \( X \), show that there exist \( n \) columns \( \{c_1, \ldots, c_n\} \) of \( M \) that are linearly independent at \( p \) and hence near \( p \) (as linear independence is given by nonvanishing of the appropriate \( n \times n \) determinant). Thus \( X \) can be covered by distinguished open subsets in bijection with the choices of \( n \) columns of \( M \). Restricting to one subset and renaming columns, reduce to the case where the determinant of the first \( n \) columns of \( M \) is invertible. Then change coordinates on \( \Lambda^{\oplus m} = \mathcal{F}(|\text{Spec } \Lambda|) \)
so that $M$ with respect to the new coordinates is the identity matrix in the first $n$ columns, and 0 thereafter. Finally, in this case interpret $F'$ as $A^{\oplus (m-n)}$.

(c) If $F'$ and $F$ are both locally free, show that $F/F'$ need not be. (Hint: over $k[t]$, consider $0 \to tk[t] \to k[t] \to k[t]/(t) \to 0$. We will soon interpret this as the closed subscheme exact sequence \[\begin{array}{c}
14.5.5.1
\end{array}\] for a point on $\mathbb{A}^1$.)

14.5.2. Tensor products. Another important example is tensor products.

14.5.C. Exercise. If $F$ and $G$ are quasicoherent sheaves, show that $F \otimes G$ is a quasicoherent sheaf described by the following information: If $\text{Spec } A$ is an affine open, and $\Gamma(\text{Spec } A, F) = M$ and $\Gamma(\text{Spec } A, G) = N$, then $\Gamma(\text{Spec } A, F \otimes G) = M \otimes N$, and the restriction map $\Gamma(\text{Spec } A, F \otimes G) \to \Gamma(\text{Spec } A_f, F \otimes G)$ is precisely the localization map $M \otimes_A N \to (M \otimes_A N)_f \cong M_f \otimes_A N_f$. (We are using the algebraic fact that $(M \otimes_R N)_f \cong M_f \otimes_R N_f$.) You can prove this by universal property if you want, or by using the explicit construction.)

Note that thanks to the machinery behind the distinguished affine base, sheafification is taken care of. This is a feature we will use often: constructions involving quasicoherent sheaves that involve sheafification for general sheaves don’t require sheafification when considered on the distinguished affine base. Along with the fact that injectivity, surjectivity, kernels and so on may be computed on affine opens, this is the reason that it is particularly convenient to think about quasicoherent sheaves in terms of affine open sets.

Given a section $s$ of $F$ and a section $t$ of $G$, we have a section $s \otimes t$ of $F \otimes G$. If $F$ is an invertible sheaf, this section is often denoted $st$.

14.5.3. Tensor algebra constructions.

For the next exercises, recall the following. If $M$ is an $A$-module, then the tensor algebra $T^*(M)$ is a non-commutative algebra, graded by $\mathbb{Z}^{\geq 0}$, defined as follows. $T^0(M) = A$, $T^n(M) = M \otimes_A \cdots \otimes_A M$ (where $n$ terms appear in the product), and multiplication is what you expect.

The symmetric algebra $\text{Sym}^* M$ is a symmetric algebra, graded by $\mathbb{Z}^{\geq 0}$, defined as the quotient of $T^*(M)$ by the (two-sided) ideal generated by all elements of the form $x \otimes y - y \otimes x$ for all $x, y \in M$. Thus $\text{Sym}^n M$ is the quotient of $M \otimes \cdots \otimes M$ by the relations of the form $m_1 \otimes \cdots \otimes m_n - m'_1 \otimes \cdots \otimes m'_n$, where $(m'_1, \ldots, m'_n)$ is a rearrangement of $(m_1, \ldots, m_n)$.

The exterior algebra $\wedge^* M$ is defined to be the quotient of $T^* M$ by the (two-sided) ideal generated by all elements of the form $x \otimes x$ for all $x \in M$. Expanding $(a+b) \otimes (a+b)$, we see that $a \otimes b = -b \otimes a$ in $\wedge^2 M$. This implies that if $2$ is invertible in $A$ (e.g. if $A$ is a field of characteristic not 2), $\wedge^n M$ is the quotient of $M \otimes \cdots \otimes M$ by the relations of the form $m_1 \otimes \cdots \otimes m_n - (-1)^{\text{sgn}(\sigma)} m_{\sigma(1)} \otimes \cdots \otimes m_{\sigma(n)}$, where $\sigma$ is a permutation of $\{1, \ldots, n\}$. The exterior algebra is a “skew-commutative” $A$-algebra.

It is most correct to write $T^*_A(M)$, $\text{Sym}^*_A(M)$, and $\wedge^*_A(M)$, but the “base ring” $A$ is usually omitted for convenience. (Better: both $\text{Sym}$ and $\wedge$ can be defined by universal properties. For example, $\text{Sym}^n_A(M)$ is universal among modules such that any map of $A$-modules $M^\otimes n \to N$ that is symmetric in the $n$ entries factors uniquely through $\text{Sym}^n_A(M)$.)
14.5.D. Exercise. Suppose $\mathcal{F}$ is a quasicoherent sheaf. Define the quasicoherent sheaves $\text{Sym}^n \mathcal{F}$ and $\Lambda^n \mathcal{F}$. (One possibility: describe them on each affine open set, and use the characterization of Important Exercise [14.3.3].) If $\mathcal{F}$ is locally free of rank $m$, show that $\text{Sym}^n \mathcal{F}$, $\Lambda^n \mathcal{F}$, and $\Lambda^n \mathcal{F}$ are locally free, and find their ranks.

You can also define the sheaf of non-commutative algebras $\mathcal{T} \mathcal{F}$, the sheaf of commutative algebras $\text{Sym}^\bullet \mathcal{F}$, and the sheaf of skew-commutative algebras $\Lambda^\bullet \mathcal{F}$.

14.5.E. Exercise. Suppose $0 \to \mathcal{F}' \to \mathcal{F} \to \mathcal{F}'' \to 0$ is an exact sequence of locally free sheaves. Show that for any $r$, there is a filtration of $\text{Sym}^r \mathcal{F}$

$$\text{Sym}^r \mathcal{F} = \mathcal{F}^0 \supset \mathcal{F}^1 \supset \cdots \supset \mathcal{F}^r \supset \mathcal{F}^{r+1} = 0$$

with subquotients

$$\mathcal{F}^p / \mathcal{F}^{p+1} \cong (\text{Sym}^p \mathcal{F}') \otimes (\text{Sym}^{r-p} \mathcal{F}'') .$$

(Here are two different possible hints for this and Exercise 14.5.D. (1) Interpret the transition matrices for $\mathcal{F}$ as block upper-diagonal, with two blocks, where one diagonal block gives the transition matrices for $\mathcal{F}'$, and the other gives the transition matrices for $\mathcal{F}''$. (cf. Exercise 14.5.1.1(b)). Then appropriately interpret the transition matrices for $\text{Sym}^r \mathcal{F}$ as block upper-diagonal as well, with $r+1$ blocks. (2) It suffices to consider a small enough affine open set $\text{Spec} A$, where $\mathcal{F}', \mathcal{F}, \mathcal{F}''$ are free, and to show that your construction behaves well with respect to localization at an element $f \in A$. In such an open set, the sequence is $0 \to A^\oplus p \to A^\oplus (p+q) \to A^\oplus q \to 0$ by the Exercise 14.5.A. Let $e_1, \ldots, e_n$ be the standard basis of $A^n$, and $f_1, \ldots, f_q$ be the the standard basis of $A^\oplus q$. Let $e'_1, \ldots, e'_p$ be denote the images of $e_1, \ldots, e_p$ in $A^\oplus q$. Let $f'_1, \ldots, f'_q$ be any lifts of $f_1, \ldots, f_q$ to $A^\oplus (p+q)$. Note that $f'_1$ is well-defined modulo $e'_1, \ldots, e'_p$.

Note that

$$\text{Sym}^s \mathcal{F}|_{\text{Spec} A} \cong \bigoplus_{i=0}^s \text{Sym}^i \mathcal{F}'|_{\text{Spec} A} \otimes \mathcal{O}_{\text{Spec} A} \otimes \text{Sym}^{s-i} \mathcal{F}''|_{\text{Spec} A} .$$

Show that $\mathcal{F}'^p := \bigoplus_{i=p}^s \text{Sym}^i \mathcal{F}'|_{\text{Spec} A} \otimes \mathcal{O}_{\text{Spec} A} \otimes \text{Sym}^{s-i} \mathcal{F}''|_{\text{Spec} A}$ gives a well-defined (locally free) subsheaf that is independent of the choices made, e.g. of the basis $e_1, \ldots, e_p$, $f_1, \ldots, f_q$, and the lifts $f'_1, \ldots, f'_q$.)

14.5.F. Exercise. Suppose $\mathcal{F}$ is locally free of rank $n$. Then $\Lambda^n \mathcal{F}$ is called the **determinant (line) bundle** or (both better and worse) the **determinant locally free sheaf**. Describe a map $\Lambda^r \mathcal{F} \times \Lambda^{n-r} \mathcal{F} \to \Lambda^n \mathcal{F}$ that induces an isomorphism of $\Lambda^r \mathcal{F} \to (\Lambda^{n-r} \mathcal{F}) \otimes \Lambda^n \mathcal{F}$. This is called a **perfect pairing of vector bundles**. (If you know about perfect pairings of vector spaces, do you see why this is a generalization?) You might use this later showing duality of Hodge numbers of nonsingular varieties over algebraically closed fields, Exercise 22.4.M.

14.5.G. Useful Exercise. Suppose $0 \to \mathcal{F}' \to \mathcal{F} \to \mathcal{F}'' \to 0$ is an exact sequence of locally free sheaves. Show that for any $r$, there is a filtration of $\Lambda^r \mathcal{F}$:

$$\Lambda^r \mathcal{F} = \mathcal{F}^0 \supset \mathcal{F}^1 \supset \cdots \supset \mathcal{F}^r \supset \mathcal{F}^{r+1} = 0$$

with subquotients

$$\mathcal{F}^p / \mathcal{F}^{p+1} \cong (\Lambda^p \mathcal{F}') \otimes (\Lambda^{r-p} \mathcal{F}'').$$
for each $p$. In particular, $\det F = (\det F') \otimes (\det F'')$. In fact we only need that $F''$ is locally free.

14.5.H. Exercise (Determinant Line Bundles Behave Well in Exact Sequences). Suppose $0 \to F_1 \to \cdots \to F_n \to 0$ is an exact sequence of finite rank locally free sheaves on $X$. Show that “the alternating product of determinant bundles is trivial”:

$$\det(F_1) \otimes \det(F_2)^\vee \otimes \det(F_3) \otimes \det(F_4)^\vee \otimes \cdots \otimes \det(F_n)^{(-1)^n} \equiv O_X.$$

(Hint: break the exact sequence into short exact sequences. Use Exercise 14.5.B(b) to show that they are short exact exact sequences of finite rank locally free sheaves. Then use the previous Exercise 14.5.G.)

14.5.4. Torsion-free sheaves (a stalk-local condition). Recall that an $A$-module $M$ is torsion-free if $rm = 0$ implies $r = 0$ or $m = 0$. An $O_X$-module $F$ is said to be torsion-free if $F_p$ is a torsion-free $O_{X,p}$-module for all $p$.

14.5.I. Exercise. Show that if $M$ is a torsion-free $A$-module, then so is any localization of $M$. Hence show that $\tilde{M}$ is a torsion free sheaf on $\text{Spec } A$.

14.5.J. Unimportant Exercise (Torsion-Freeness is Not Affine Local for Stupid Reasons). Find an example on a two-point space showing that $M := A$ might not be torsion-free on $\text{Spec } A$ even though $O_{\text{Spec } A} = \tilde{M}$ is torsion-free.

14.5.5. Quasicoherent sheaves of ideals correspond to closed subschemes. Recall that if $i : X \hookrightarrow Y$ is a closed immersion, then we have a surjection of sheaves on $Y$: $O_Y \longrightarrow i_* O_X$ (9.1). (The $i_*$ is often omitted, as we are considering the sheaf on $X$ as being a sheaf on $Y$.) The kernel $I_{X/Y}$ is a "sheaf of ideals" in $Y$: for each open subset $U$ of $Y$, the sections form an ideal in the ring of functions on $U$.

Compare (hard) Exercise 9.1.E and the characterization of quasicoherent sheaves given in (hard) Exercise 14.3.D. You will see that a sheaf is ideals is quasicoherent if and only if it comes from a closed subscheme. (An example of a non-quasicoherent sheaf of ideals was given in Exercise 9.1.D). We call

$$0 \to I_{X/Y} \to O_Y \to i_* O_X \to 0$$

the closed subscheme exact sequence corresponding to $X \hookrightarrow Y$.

14.6 Finiteness conditions on quasicoherent sheaves: finite type quasicoherent sheaves, and coherent sheaves

There are some natural finiteness conditions on an $A$-module $M$. I will tell you three. In the case when $A$ is a Noetherian ring, which is the case that almost all of you will ever care about, they are all the same.

The first is the most naive: a module could be finitely generated. In other words, there is a surjection $A^q \to M \to 0$.

The second is reasonable too. It could be finitely presented — it could have a finite number of generators with a finite number of relations: there exists a finite presentation

$$A^q \to A^p \to M \to 0.$$
14.6.A. Exercise. Suppose $M$ is a finitely presented $A$-module, and $\phi: A^p \rightarrow M$ is any surjection. Show that $\ker \phi$ is finitely generated. Hint: Write $M$ as the kernel of $A^p$ by a finitely generated module $K$. Figure out how to map the short exact sequence $0 \rightarrow K \rightarrow A^p \rightarrow M \rightarrow 0$, and use the Snake Lemma.

The third notion is frankly a bit surprising, and I will justify it soon. We say that an $A$-module $M$ is coherent if (i) it is finitely generated, and (ii) whenever we have a map $A^p \rightarrow M$ (not necessarily surjective!), the kernel is finitely generated.

Clearly coherent implies finitely presented, which in turn implies finitely generated.

14.6.1. Proposition. — If $A$ is Noetherian, then these three definitions are the same.

Before proving this, we take this as an excuse to develop some algebraic background.

14.6.2. Noetherian conditions for modules. If $A$ is any ring, not necessarily Noetherian, we say an $A$-module is Noetherian if it satisfies the ascending chain condition for submodules. Thus for example $A$ is a Noetherian ring if and only if it is a Noetherian $A$-module.

14.6.B. Exercise. Show that if $M$ is a Noetherian $A$-module, then any submodule of $M$ is a finitely generated $A$-module.

14.6.C. Exercise. If $0 \rightarrow M' \rightarrow M \rightarrow M'' \rightarrow 0$ is exact, show that $M'$ and $M''$ are Noetherian if and only if $M$ is Noetherian. (Hint: Given an ascending chain in $M$, we get two simultaneous ascending chains in $M'$ and $M''$. Possible further hint: prove that if $M' \longrightarrow M \overset{\phi}{\longrightarrow} M''$ is exact, and $N \subset N' \subset M$, and $N \cap M' = N' \cap M'$ and $\phi(N) = \phi(N')$, then $N = N'$.)

14.6.D. Exercise. Show that if $A$ is a Noetherian ring, then $A^{\oplus n}$ is a Noetherian $A$-module.

14.6.E. Exercise. Show that if $A$ is a Noetherian ring and $M$ is a finitely generated $A$-module, then $M$ is a Noetherian module. Hence by Exercise[14.6.B] any submodule of a finitely generated module over a Noetherian ring is finitely generated.

Proof of Proposition[14.6.1] As we observed earlier, coherent implies finitely presented implies finitely generated. So suppose $M$ is finitely generated. Take any $A^p \overset{\alpha}{\rightarrow} M$. Then $\ker \alpha$ is a submodule of a finitely generated module over $A$, and is thus finitely generated by Exercise[14.6.E] Thus $M$ is coherent.

Hence most people can think of these three notions as the same thing.


The proof in general is given in [14.8] in a series of short exercises. You should read this only if you are particularly curious.
Proof if \( A \) is Noetherian. Recall from our discussion at the start of \( \S 14.4 \) that we must check three things:

(i) The 0-sheaf is coherent.
(ii) The category of coherent modules is closed under finite sums.
(iii) The category of coherent modules is closed under kernels and cokernels.

The first two are clear. For (iii), suppose that \( f : M \to N \) is a map of finitely generated modules. Then \( \text{coker } f \) is finitely generated (it is the image of \( N \)), and \( \text{ker } f \) is too (it is a submodule of a finitely generated module over a Noetherian ring, Exercise \[14.6.E]\).

\[14.6.E. \] Easy Exercise (only important for non-Noetherian people). Show \( A \) is coherent as an \( A \)-module if and only if the notion of finitely presented agrees with the notion of coherent.

\[14.6.G. \] Exercise. If \( f \in A \), show that if \( M \) is a finitely generated (resp. finitely presented, coherent) \( A \)-module, then \( M_f \) is a finitely generated (resp. finitely presented, coherent) \( A_f \)-module. (The “coherent” case is the tricky one.)

\[14.6.H. \] Exercise. If \( (f_1, \ldots, f_n) = A \), and \( M_{f_i} \) is finitely generated (resp. finitely presented, coherent) \( A_{f_i} \)-module for all \( i \), then \( M \) is a finitely generated (resp. finitely presented, coherent) \( A \)-module. Hint for the finitely presented case: Exercise \[14.6.A\].

\[14.6.4. \] Definition. A quasicoherent sheaf \( \mathcal{F} \) is finite type (resp. finitely generated, coherent) if for every affine open \( \text{Spec } A, \Gamma(\text{Spec } A, \mathcal{F}) \) is a finitely generated (resp. finitely generated, coherent) \( A \)-module. Note that coherent sheaves are always finite type, and that on a locally Noetherian scheme, all three notions are the same (by Proposition \[14.6.1\]). By Proposition \[14.6.3\], implies that the coherent sheaves on \( X \) form an abelian category, which we denote \( \text{Coh}_X \). Coherence is basically only interesting if \( \mathcal{O}_X \) is coherent.

Thanks to the Affine Communication Lemma \[6.3.2\] and the two previous exercises \[14.6.G, 14.6.H\], it suffices to check this on the open sets in a single affine cover. Notice that locally free sheaves are always finite type, and if \( \mathcal{O}_X \) is coherent, locally free sheaves on \( X \) are coherent. (If \( \mathcal{O}_X \) is not coherent, then coherence is a pretty useless notion on \( X \).)

I want to say a few words on the notion of coherence. Proposition \[14.6.3\] is a good motivation for this definition: it gives a small (in a non-technical sense) abelian category in which we can think about vector bundles.

There are two sorts of people who should care about the details of this definition, rather than living in a Noetherian world where coherent means finite type. Complex geometers should care. They consider complex-analytic spaces with the classical topology. One can define the notion of coherent \( \mathcal{O}_X \)-module in a way analogous to this (see [S-FAC] Def. 2). Then Oka’s theorem states that the structure sheaf is coherent, and this is very hard.

The second sort of people who should care are the sort of arithmetic people who may need to work with non-Noetherian rings. For example, the ring of adeles is non-Noetherian.
Warning: it is common in the later literature to incorrectly define coherent as finitely generated. Please only use the correct definition, as the wrong definition causes confusion. Besides doing this for the reason of honesty, it will also help you see what hypotheses are actually necessary to prove things. And that always helps you remember what the proofs are — and hence why things are true.

14.7 Pleasant properties of finite type and coherent sheaves

We begin with the fact that $\text{Hom}$ behaves reasonably if the source is coherent.

14.7.A. EXERCISE. (a) Suppose $\mathcal{F}$ is a coherent sheaf on $X$, and $\mathcal{G}$ is a quasicoherent sheaf on $X$. Show that $\text{Hom}(\mathcal{F}, \mathcal{G})$ is a quasicoherent sheaf. (Hint: Describe it on affine open sets, and show that it behaves well with respect to localization with respect to $f$. To show that $\text{Hom}_A(M, N)_f \cong \text{Hom}_A(M, N)_f$, take a presentation $A^n \to A^m \to M \to 0$, and apply $\text{Hom}(\cdot, N)$ and localize. You will use the fact that $p$ and $q$ are finite. Up to here, you need only the fact that $\mathcal{F}$ is locally finitely presented.)

(b) If further $\mathcal{G}$ is coherent and $\mathcal{O}_X$ is coherent, show that $\text{Hom}(\mathcal{F}, \mathcal{G})$ is also coherent. Show that $\text{Hom}$ is a left-exact functor in both variables.

14.7.1. Duals of coherent sheaves. In particular, if $\mathcal{F}$ is coherent, its dual $\mathcal{F}^\vee := \text{Hom}(\mathcal{F}, \mathcal{O})$ is too. This generalizes the notion of duals of vector bundles in Exercise 14.1.C Your argument there generalizes to show that there is always a natural morphism $\mathcal{F} \to (\mathcal{F}^\vee)^\vee$. Unlike in the vector bundle case, this is not always an isomorphism. (For an example, let $\mathcal{F}$ be the coherent sheaf associated to $k[t]/(t)$ on $\mathbb{A}^1 = \text{Spec } k[t]$, and show that $\mathcal{F}^\vee = 0$.) Coherent sheaves for which the “double dual” map is an isomorphism are called reflexive sheaves, but we won’t use this notion. The canonical map $\mathcal{F} \otimes \mathcal{F}^\vee \to \mathcal{O}_X$ is called the trace map — can you see why?

We next come to a geometric interpretation of Nakayama’s lemma, which is why I consider Nakayama’s Lemma a geometric fact (with an algebraic proof).

14.7.B. USEFUL EXERCISE: GEOMETRIC NAKAYAMA (GENERATORS OF A FIBER GENERATE A FINITE TYPE SHEAF NEARBY). Suppose $X$ is a scheme, and $\mathcal{F}$ is a finite type quasicoherent sheaf. Show that if $U \subset X$ is a neighborhood of $x \in X$ and $a_1, \ldots, a_n \in \mathcal{F}(U)$ so that the images $\pi_1, \ldots, \pi_n \in \mathcal{F}_x$ generate $\mathcal{F}_x \otimes \kappa(x)$, then there is an affine neighborhood $x \in \text{Spec } A \subset U$ of $x$ such that “$a_1|_{\text{Spec } A}, \ldots, a_n|_{\text{Spec } A}$ generate $\mathcal{F}|_{\text{Spec } A}$” in the following senses:

(i) $a_1|_{\text{Spec } A}, \ldots, a_n|_{\text{Spec } A}$ generate $\mathcal{F}|_{\text{Spec } A}$ as an $A$-module;

(ii) for any $y \in \text{Spec } A$, $a_1, \ldots, a_n$ generate the stalk $\mathcal{F}|_{\text{Spec } A}$ as an $\mathcal{O}_{X, y}$-module (and hence for any $y \in \text{Spec } A$, the fibers $a_1|_y, \ldots, a_n|_y$ generate the fiber $\mathcal{F}|_y$ as a $\kappa(y)$-vector space).

In particular, if $\mathcal{F}_x \otimes \kappa(x) = 0$, then there exists a neighborhood $V$ of $x$ such that $\mathcal{F}|_V = 0$.

14.7.C. EXERCISE (THE SUPPORT OF A FINITE TYPE SHEAF IS CLOSED). This exercise is partially an excuse to discuss the useful notion of “support”. Suppose $s$ is a
section of a sheaf $\mathcal{F}$ of abelian groups. Define the **support** of $s$ by

$$\text{Supp } s = \{ p \in X : s_p \neq 0 \text{ in } \mathcal{F}_p \}.$$ 

Define the **support** of $\mathcal{F}$ by $\text{Supp } \mathcal{F} = \{ p \in X : \mathcal{F}_p \neq 0 \}$ (cf. Exercise 3.6.F(b)) — the union of “all the supports of sections on various open sets”. (Support is a stalk-local notion, and hence behaves well with respect to restriction to open sets, or to stalks. Warning: Support is where the germ(s) are nonzero, not where the value(s) are nonzero.) Show that the support of a finite type quasicoherent sheaf on a scheme $X$ is a closed subset. (Hint: Reduce to the case $X$ affine. Choose a finite set of generators of the corresponding module.) Show that the support of a quasicoherent sheaf need not be closed. (Hint: If $A = \mathbb{C}[t]$, then $\mathbb{C}[t]/(t - a)$ is an $A$-module supported at $a$. Consider $\bigoplus_{a \in \mathbb{C}} \mathbb{C}[t]/(t - a)$. Be careful: this example won’t work if $\bigoplus$ is replaced by $\prod$.)

**14.7.D. Useful Exercise (Local Freeness of a Coherent Sheaf Is a Stalk-local Property; And Locally Free Stalks Implies Local Freeness Nearby).** Suppose $\mathcal{F}$ is a coherent sheaf on scheme $X$. Show that if $\mathcal{F}_x$ is a free $\mathcal{O}_{X,x}$-module for some $x \in X$, then $\mathcal{F}$ is locally free in some open neighborhood of $X$. Hence $\mathcal{F}$ is locally free if and only if $\mathcal{F}_x$ is a free $\mathcal{O}_{X,x}$-module for all $x \in X$. Hint: Find an open neighborhood $U$ of $x$, and $n$ elements of $\mathcal{F}(U)$ that generate $\mathcal{F}_x$ and hence by Nakayama’s lemma they generate $\mathcal{F}_x$. Use Geometric Nakayama, Exercise 14.7.B show that the sections generate $\mathcal{F}_y$ for all $y$ in some neighborhood $Y$ of $x$ in $U$. Thus you have described a surjection $\mathcal{O}_Y^\oplus \to \mathcal{F}|_Y$. Show that the kernel this map is finite type, and hence has closed support (say $Z \subset Y$), which does not contain $x$. Thus $\mathcal{O}_Y^\oplus \to \mathcal{F}|_{Y\setminus Z}$.

This is enlightening in a number of ways. It shows that for coherent sheaves, local freeness is a stalk-local condition. Furthermore, on an integral scheme, any coherent sheaf $\mathcal{F}$ is automatically free over the generic point (do you see why?), so every coherent sheaf on an integral scheme is locally free over a dense open subset. And any coherent sheaf that is 0 at the generic point of an irreducible scheme is necessarily 0 on a dense open subset. The last two sentences show the utility of generic points; such statements would have been more mysterious in classical algebraic geometry.

**14.7.E. Exercise.** Show that torsion-free coherent sheaves on a nonsingular (hence implicitly locally Noetherian) curve are locally free. (Although “torsion sheaf” is has not yet been defined, you should also be able to make sense out of the statement: any coherent sheaf is a direct sum of a torsion-free sheaf and a torsion sheaf.)

To answer the previous exercise, use Useful Exercise 14.7.D (local freeness can be checked at stalks) to reduce to the discrete valuation ring case, and recall Remark 13.3.16, the structure theorem for finitely generated modules over a principal ideal domain $A$: any such module can be written as the direct sum of principal modules $A/(a)$. For discrete valuation rings, this means that the summands are of the form $A$ or $A/m^k$. Hence:

**14.7.2. Proposition.** — If $M$ is a finitely generated module over a discrete valuation ring, then $M$ is torsion-free if and only if $M$ is free.
14.7.3. Rank of a quasicoherent sheaf at a point.

Suppose $\mathcal{F}$ is a quasicoherent sheaf on a scheme $X$, and $p$ is a point of $X$. The vector space $\mathcal{F}_p/m\mathcal{F}_p = \mathcal{F}_p \otimes_{\kappa(p)} \kappa(p)$ can be interpreted as the fiber of the sheaf at the point, where $m$ is the maximal ideal corresponding to $p$, and $\kappa(p)$ is as usual the residue field at $p$. A section of $\mathcal{F}$ over an open set containing $p$ can be said to take on a value at that point, which is an element of this vector space. The rank of a quasicoherent sheaf $\mathcal{F}$ at a point $p$ is $\dim_{\kappa(p)} \mathcal{F}_p/m\mathcal{F}_p$ (possibly infinite). More explicitly, on any affine set $\text{Spec } A$ where $p = [p]$ and $\mathcal{F}(\text{Spec } A) = M$, then the rank is $\dim_{\kappa(A/p)} M_p/pM_p$. Note that this definition of rank is consistent with the notion of rank of a locally free sheaf. In the locally free case, the rank is a (locally) constant function of the point. The converse is sometimes true, see Exercise [14.7.J] below.

If $X$ is irreducible, and $\mathcal{F}$ is a quasicoherent (usually coherent) sheaf on $X$ on $X$, then $\text{rank } \mathcal{F}$ (with no mention of a point) by convention means at the generic point.

14.7.F. Exercise. Consider the coherent sheaf $\mathcal{F}$ on $\mathbb{A}^1_k = \text{Spec } k[t]$ corresponding to the module $k[t]/(t)$. Find the rank of $\mathcal{F}$ at every point of $\mathbb{A}^1_k$. Don’t forget the generic point!

14.7.G. Exercise. Show that at any point, $\text{rank}(\mathcal{F} \oplus \mathcal{G}) = \text{rank}(\mathcal{F}) + \text{rank}(\mathcal{G})$ and $\text{rank}(\mathcal{F} \otimes \mathcal{G}) = \text{rank } \mathcal{F} \text{ rank } \mathcal{G}$ at any point. (Hint: Show that direct sums and tensor products commute with ring quotients and localizations, i.e. $(M \otimes N) \otimes_R (R/I) \cong M/IM \oplus N/IN$, $(M \otimes_R N) \otimes_R (R/I) \cong (M \otimes_R R/I) \otimes_R (N \otimes_R R/I) \cong M/IM \otimes_{R/I} N/IM$, etc.)

If $\mathcal{F}$ is finite type, then the rank is finite, and by Nakayama’s lemma, the rank is the minimal number of generators of $M_p$ as an $A_p$-module.

14.7.H. Important Exercise. If $\mathcal{F}$ is finite type, show that $\text{rank}(\mathcal{F})$ is an upper semicontinuous function on $X$. Hint: generators at a point $p$ are generators nearby by Geometric Nakayama’s Lemma, Exercise [14.7.B]. (The example in Exercise [14.7.C] shows the necessity of the finite type hypothesis.)

14.7.I. Important Hard Exercise.

(a) If $X$ is reduced, $\mathcal{F}$ is finite type, and the rank is constant, show that $\mathcal{F}$ is locally free. Thus (as rank is uppersemicontinuous, Exercise [14.7.H]) finite type sheaves are locally free on a dense open set. Hint: Reduce to the case where $X$ is affine. Then show it in a neighborhood of a closed point $p$ as follows. (You will have to show that this suffices, using the affine assumption. But note that closed points aren’t necessarily dense in an affine scheme, see for example Exercise [4.4.J].)

Suppose $n = \text{rank } \mathcal{F}$. Choose $n$ generators of the fiber $\mathcal{F}_p$ (a basis as an $\kappa(p)$-vector space). By Geometric Nakayama’s Lemma [14.7.B] we can find a smaller
neighborhood \( p \in \text{Spec } A \subset X \), with \( \mathcal{F}|_{\text{Spec } A} = \mathcal{M} \), so that the chosen generators \( \mathcal{F}|_p \) lift to generators \( m_1, \ldots, m_n \) of \( M \). Let \( \phi : A^n \to M \) with \( (r_1, \ldots, r_n) \mapsto \sum r_i m_i \). If \( \ker \phi \neq 0 \), then suppose \( (r_1, \ldots, r_n) \) is in the kernel, with \( r_1 \neq 0 \). As \( r_1 \neq 0 \), there is some \( p \) where \( r_1 \notin p \) — here we use the reduced hypothesis. Then \( r_1 \) is invertible in \( A_p \), so \( M_p \) has fewer than \( n \) generators, contradicting the constancy of rank.

(b) Show that part (a) can be false without the condition of \( X \) being reduced. (Hint: \( \text{Spec } k[x]/x^2, M = k \).)

You can use the notion of rank to help visualize finite type sheaves, or even quasicoherent sheaves. For example, I think of a coherent sheaf as generalizing a finite rank vector bundle as follows: to each point there is an associated vector space, and although the ranks can jump, they fit together in families as well as one might hope. You might try to visualize the example of Example 14.7. You can use the notion of rank to help visualize finite type sheaves, or even quasicoherent sheaves. For example, I think of a coherent sheaf as generalizing a finite rank vector bundle as follows: to each point there is an associated vector space, and although the ranks can jump, they fit together in families as well as one might hope. You might try to visualize the example of Example 14.7. Nonreducedness can fit into the picture as well — how would you picture the coherent sheaf on \( \text{Spec } k[c]/(c^2) \) corresponding to \( k[c]/(c^2) \)?

14.7.4. Degree of a finite morphism at a point. Suppose \( \pi : X \to Y \) is a finite morphism. Then \( \pi_* \mathcal{O}_X \) is a finite type (quasicoherent) sheaf on \( Y \), and the rank of this sheaf at a point \( p \) is called the degree of the finite morphism at \( p \). By Exercise 14.7.1, the degree of \( \pi \) is a upper-semicontinuous function on \( Y \). The degree can jump: consider the closed immersion of a point into a line corresponding to \( k[t] \to k \) given by \( t \mapsto 0 \). It can also be constant in cases that you might initially find surprising — see Exercise 10.3.3, where the degree is always 2, but the 2 is obtained in a number of different ways.

14.7.J. Exercise. Suppose \( \pi : X \to Y \) is a finite morphism. By unwinding the definition, verify that the degree of \( \pi \) at \( p \) is the dimension of the space of functions of the scheme-theoretic preimage of \( p \), considered as a vector space over the residue field \( k(p) \). In particular, the degree is zero if and only if \( \pi^{-1}(p) \) is empty.

14.8 ** Coherent modules over non-Noetherian rings

This section is intended for people who might work with non-Noetherian rings, or who otherwise might want to understand coherent sheaves in a more general setting. Read this only if you really want to!

Suppose \( A \) is a ring. Recall the definition of when an \( A \)-module \( M \) is finitely generated, finitely presented, and coherent. The reason we like coherence is that coherent modules form an abelian category. Here are some accessible exercises working out why these notions behave well. Some repeat earlier discussion in order to keep this section self-contained.

The notion of coherence of a module is only interesting in the case that a ring is coherent over itself. Similarly, coherent sheaves on a scheme \( X \) will be interesting only when \( \mathcal{O}_X \) is coherent ("over itself"). In this case, coherence is clearly the same as finite presentation. An example where non-Noetherian coherence comes up is the ring \( R(x_1, \ldots, x_n) \) of "restricted power series" over a valuation ring \( R \) of a non-discretely valued \( K \) (for example, a completion of the algebraic closure of \( \mathbb{Q}_p \)). This is relevant to Tate’s theory of non-archimedean analytic geometry over \( K \).
The importance of the coherence of the structure sheaf underlines the importance of Oka’s theorem in complex geometry.

14.8.A. Exercise. Show that coherent implies finitely presented implies finitely generated. (This was discussed in the previous section.)

14.8.B. Exercise. Show that $0$ is coherent.


$$0 \rightarrow M \rightarrow N \rightarrow P \rightarrow 0$$

is an exact sequence of $A$-modules. In these series of problems, we will show that if two of $\{M, N, P\}$ are coherent, the third is as well, which will prove very useful.

14.8.1. Hint. The following hint applies to several of the problems: try to write

$$0 \rightarrow A^p \rightarrow A^{p+q} \rightarrow A^q \rightarrow 0$$

and possibly use the Snake Lemma 2.7.5.

14.8.C. Exercise. Show that $N$ finitely generated implies $P$ finitely generated. (You will only need right-exactness of (14.8.0.1).)

14.8.D. Exercise. Show that $M, P$ finitely generated implies $N$ finitely generated. (Possible hint: †.) (You will only need right-exactness of (14.8.0.1).)

14.8.E. Exercise. Show that $N, P$ finitely generated need not imply $M$ finitely generated. (Hint: if $I$ is an ideal, we have $0 \rightarrow I \rightarrow A \rightarrow A/I \rightarrow 0$.)

14.8.F. Exercise. Show that $N$ coherent, $M$ finitely generated implies $M$ coherent. (You will only need left-exactness of (14.8.0.1).)


$$0 \rightarrow M \rightarrow N \rightarrow P \rightarrow 0$$

(You will only need left-exactness of (14.8.0.1).)

14.8.H. Exercise. Show that $M$ finitely generated and $N$ coherent implies $P$ coherent. (Hint for (ii): †.)

14.8.I. Exercise. Show that $M, P$ coherent implies $N$ coherent. (Hint: †.)
14.8.J. Exercise. Show that a finite direct sum of coherent modules is coherent.

14.8.K. Exercise. Suppose $M$ is finitely generated, $N$ coherent. Then if $\phi : M \to N$ is any map, then show that $\text{Im} \phi$ is coherent.

14.8.L. Exercise. Show that the kernel and cokernel of maps of coherent modules are coherent.

At this point, we have verified that coherent $A$-modules form an abelian subcategory of the category of $A$-modules. (Things you have to check: $0$ should be in this set; it should be closed under finite sums; and it should be closed under taking kernels and cokernels.)

14.8.M. Exercise. Suppose $M$ and $N$ are coherent submodules of the coherent module $P$. Show that $M + N$ and $M \cap N$ are coherent. (Hint: consider the right map $M \oplus N \to P$.)

14.8.N. Exercise. Show that if $A$ is coherent (as an $A$-module) then finitely presented modules are coherent. (Of course, if finitely presented modules are coherent, then $A$ is coherent, as $A$ is finitely presented!)

14.8.O. Exercise. If $M$ is finitely presented and $N$ is coherent, show that $\text{Hom}(M, N)$ is coherent. (Hint: Hom is left-exact in its first argument.)

14.8.P. Exercise. If $M$ is finitely presented, and $N$ is coherent, show that $M \otimes N$ is coherent.

14.8.Q. Exercise. If $f \in A$, show that if $M$ is a finitely generated (resp. finitely presented, coherent) $A$-module, then $M_f$ is a finitely generated (resp. finitely presented, coherent) $A_f$-module. (Hint: localization is exact.) This exercise is repeated from Exercise [14.6.C] to make this section self-contained.

14.8.R. Exercise. Suppose $(f_1, \ldots, f_n) = A$. Show that if $M_{f_i}$ is finitely generated for all $i$, then $M$ is too. (Hint: Say $M_{f_i}$ is generated by $m_{ij} \in M$ as an $A_{f_i}$-module. Show that the $m_{ij}$ generate $M$. To check surjectivity $\oplus_{i,j} A \to M$, it suffices to check “on $D(f_i)$” for all $i$.)

14.8.S. Exercise. Suppose $(f_1, \ldots, f_n) = A$. Show that if $M_{f_i}$ is coherent for all $i$, then $M$ is too. (Hint: if $\phi : A^2 \to M$, then $(\ker \phi)_{f_i} = \ker(\phi_{f_i})$, which is finitely generated for all $i$. Then apply the previous exercise.)
CHAPTER 15

Line bundles: Invertible sheaves and divisors

We next describe convenient and powerful ways of working with and classifying line bundles (invertible sheaves). We begin with a fundamental example, the line bundles \( \mathcal{O}(n) \) on projective space, §15.1. We then introduce Weil divisors (formal sums of codimension 1 subsets), and use them to actually compute \( \text{Pic} \mathcal{X} \) in a number of circumstances, §15.2. We finally discuss sheaves of ideals that happen to be invertible (effective Cartier divisors), §15.3. A recurring theme is that invertible sheaves are closely related to “codimension 1 information”.

15.1 Some line bundles on projective space

We now describe an important family of invertible sheaves on projective space over a field \( k \).

As a warm-up, we begin with the invertible sheaf \( \mathcal{O}_{\mathbb{P}^1_k}(1) \) on \( \mathbb{P}^1_k = \text{Proj} \ k[x_0, x_1] \). The subscript \( \mathbb{P}^1_k \) refers to the space on which the sheaf lives, and is often omitted when it is clear from the context. We describe the invertible sheaf \( \mathcal{O}(1) \) using transition functions. It is trivial on the usual affine open sets \( U_0 = D(x_0) = \text{Spec} \ k[x_1/0] \) and \( U_1 = D(x_1) = \text{Spec} \ k[x_0/1] \). (We continue to use the convention \( x_i/j \) for describing coordinates on patches of projective space, see §5.4.9.) Thus the data of a section over \( U_0 \) is a polynomial in \( x_1/0 \). The transition function from \( U_0 \) to \( U_1 \) is multiplication by \( x_0/1 = x_1^{-1}/0 \). The transition function from \( U_1 \) to \( U_0 \) is hence multiplication by \( x_1/0 = x_0^{-1}/1 \).

This information is summarized below:

- open cover: \( U_0 = \text{Spec} k[x_1/0] \) \quad \( U_1 = \text{Spec} k[x_0/1] \)
- trivialization and transition functions:

\[
\begin{align*}
&k[x_1/0] \xrightarrow{x_0/1 = x_1^{-1}} k[x_0/1] \\
&k[x_1/0] \xrightarrow{x_1/0 = x_0^{-1}/1} k[x_0/1]
\end{align*}
\]

To test our understanding, let’s compute the global sections of \( \mathcal{O}(1) \). This will be analogous to our hands-on calculation that \( \Gamma(\mathbb{P}^1_k, \mathcal{O}_{\mathbb{P}^1_k}) \cong k \) (Example §5.4.6). A global section is a polynomial \( f(x_1/0) \in k[x_1/0] \) and a polynomial \( g(x_0/1) \in k[x_0/1] \) such that \( f(1/x_0/1)x_0/1 = g(x_0/1) \). A little thought will show that \( f \) must be linear: \( f(x_1/0) = ax_1/0 + b \), and hence \( f(x_0/1) = a + bx_0/1 \). Thus

\[
\dim \Gamma(\mathbb{P}^1_k, \mathcal{O}(1)) = 2 \neq 1 = \dim \Gamma(\mathbb{P}^1_k, \mathcal{O}).
\]
Thus $\mathcal{O}(1)$ is not isomorphic to $\mathcal{O}$, and we have constructed our first (proved) example of a nontrivial line bundle!

We next define more generally $\mathcal{O}_{\mathbb{P}^1}(n)$ on $\mathbb{P}^1_k$. It is defined in the same way, except that the transition functions are the $n$th powers of those for $\mathcal{O}(1)$.

Open cover

$$U_0 = \text{Spec } k[x_1/0] \quad U_1 = \text{Spec } k[x_0/1]$$

Trivialization and transition functions

$$k[x_1/0] \xrightarrow{x_0^n/1 = x_1^{-n}} k[x_0/1]$$

In particular, thanks to the explicit transition functions, we see that $\mathcal{O}(n) = \mathcal{O}(1)^{\otimes n}$ (with the obvious meaning if $n$ is negative: $(\mathcal{O}(1)^{\otimes (-n)})^\vee$). Clearly also $\mathcal{O}(m) \otimes \mathcal{O}(n) = \mathcal{O}(m+n)$.

**15.1.A. Important Exercise.** Show that $\dim \Gamma(\mathbb{P}^1, \mathcal{O}(n)) = n + 1$ if $n \geq 0$, and 0 otherwise.

**15.1.1. Example.** Long ago (§3.5.H), we warned that sheafification was necessary when tensoring $\mathcal{O}_X$-modules: if $\mathcal{F}$ and $\mathcal{G}$ are two $\mathcal{O}_X$-modules on a ringed space, then it is not necessarily true that $\mathcal{F}(X) \otimes_{\mathcal{O}_X(X)} \mathcal{G}(X) \cong (\mathcal{F} \otimes \mathcal{G})(X)$. We now have an example: let $X = \mathbb{P}^1_k$, $\mathcal{F} = \mathcal{O}(1)$, $\mathcal{G} = \mathcal{O}(-1)$.

**15.1.B. Exercise.** Show that if $m \neq n$, then $\mathcal{O}(m) \not\cong \mathcal{O}(n)$. Hence conclude that we have an injection of groups $\mathbb{Z} \hookrightarrow \text{Pic } \mathbb{P}^1_k$ given by $n \mapsto \mathcal{O}(n)$.

It is useful to identify the global sections of $\mathcal{O}(n)$ with the homogeneous polynomials of degree $n$ in $x_0$ and $x_1$, i.e. with the degree $n$ part of $k[x_0, x_1]$. Can you see this from your solution to Exercise 15.1.A? We will see that this identification is natural in many ways. For example, we will later see that the definition of $\mathcal{O}(n)$ doesn’t depend on a choice of affine cover, and this polynomial description is also independent of cover. As an immediate check of the usefulness of this point of view, ask yourself: where does the section $x_0^3 - x_0x_1^2$ of $\mathcal{O}(3)$ vanish? The section $x_0 + x_1$ of $\mathcal{O}(1)$ can be multiplied by the section $x_0^2$ of $\mathcal{O}(2)$ to get a section of $\mathcal{O}(3)$. Which one? Where does the rational section $x_0^2/(x_1 + x_0)/x_1$ of $\mathcal{O}(-2)$ have zeros and poles, and to what order? (We will rigorously define the meaning of zeros and poles shortly, but you should already be able to intuitively answer these questions.)

We now define the invertible sheaf $\mathcal{O}_{\mathbb{P}^n_k}(n)$ on the projective space $\mathbb{P}^n_k$. On the usual affine open set $U_i = \text{Spec } k[x_0/i, \ldots, x_m/i]/(x_i/i - 1) = \text{Spec } A_i$, it is trivial, so sections (as an $A_i$-module) are isomorphic to $A_i$. The transition function from $U_i$ to $U_j$ is multiplication by $x_j^n/i = x_i^{-n}/j$. Note that these transition functions clearly satisfy the cocycle condition.

$$U_i = \text{Spec } k[x_0/i, \ldots, x_m/i]/(x_i/i - 1) \quad U_j = \text{Spec } k[x_0/j, \ldots, x_m/j]/(x_j/j - 1)$$

$$k[x_0/i, \ldots, x_m/i]/(x_i/i - 1) \xrightarrow{x_i^n/j = x_j^{-n}/i} \text{Spec } k[x_0/j, \ldots, x_m/j]/(x_j/j - 1)$$
15.1.C. Essential exercise. Show that \( \dim_k \Gamma(\mathbb{P}^m_k, \mathcal{O}_{\mathbb{P}^m_k}(n)) = \binom{m+n}{n} \).

As in the case of \( \mathbb{P}^1 \), sections of \( \mathcal{O}(n) \) on \( \mathbb{P}^m \) are naturally identified with homogeneous degree \( n \) polynomials in our \( m + 1 \) variables. Thus \( x + y + 2z \) is a section of \( \mathcal{O}(1) \) on \( \mathbb{P}^2 \). It isn’t a function, but we know where this section vanishes — precisely where \( x + y + 2z = 0 \).

Also, notice that for fixed \( m \), \( \binom{m+n}{n} \) is a polynomial in \( n \) of degree \( m \), for \( n \geq 0 \) (or better: for \( n \geq -m - 1 \)). This should be telling you that this function “wants to be a polynomial,” but won’t succeed without some additional help. We will later define \( \mathcal{h}^0(\mathbb{P}^m_k, \mathcal{O}(n)) := \Gamma(\mathbb{P}^m_k, \mathcal{O}(n)) \), and later still we will define higher cohomology groups, and we will define the Euler characteristic \( \chi(\mathbb{P}^m_k, \mathcal{O}(n)) := \sum_{i=0}^{\infty} (-1)^i \mathcal{h}^i(\mathbb{P}^m_k, \mathcal{O}(n)) \) (cohomology will vanish in degree higher than \( n \)). We will discover the moral that the Euler characteristic is better-behaved than \( \mathcal{h}^0 \), and so we should now suspect (and later prove) that this polynomial is in fact the Euler characteristic, and the reason that it agrees with \( \mathcal{h}^0 \) for \( n \geq 0 \) because all the other cohomology groups should vanish.

We finally note that we can define \( \mathcal{O}(n) \) on \( \mathbb{P}^m \) for any ring \( A \): the above definition applies without change.

15.2 Invertible sheaves and Weil divisors

The notion of Weil divisors gives a great way of understanding and classifying line bundles, at least on Noetherian normal schemes. Some of what we discuss will apply in more general circumstances, and the expert is invited to consider generalizations by judiciously weakening hypotheses in various statements.

For the rest of this section, we consider only Noetherian schemes. We do this because we want to discuss codimension 1 subsets, and also have decomposition into irreducible components.

Define a Weil divisor as a formal sum of codimension 1 irreducible closed subsets of \( X \). In other words, a Weil divisor is defined to be an object of the form

\[
\sum_{Y \subset X \text{ codimension 1}} n_Y[Y]
\]

the \( n_Y \) are integers, all but a finite number of which are zero. Weil divisors obviously form an abelian group, denoted Weil \( X \).

For example, if \( X \) is a curve (such as the Spec of a Dedekind domain), the Weil divisors are linear combination of closed points.

We say that \( Y \) is an irreducible (Weil) divisor. A Weil divisor is said to be effective if \( n_Y \geq 0 \) for all \( Y \). In this case we say \( D \geq 0 \), and by \( D_1 \geq D_2 \) we mean \( D_1 - D_2 \geq 0 \). The support of a Weil divisor \( D \) is the subset \( \cup_{Y \neq 0} Y \). If \( U \subset X \) is an open set, there is a natural restriction map Weil \( X \to \text{Weil } U \), where \( \sum n_Y[Y] \mapsto \sum_{Y \cap U \neq \emptyset} n_Y[Y \cap U] \).

Suppose now that \( X \) is regular in codimension 1 (and Noetherian). We add this hypothesis because we will use properties of discrete valuation rings. Assume also that \( X \) is reduced. (This is only so we can talk about rational functions without worrying about them being defined at embedded points. Feel free to relax this hypothesis.) Suppose that \( \mathcal{L} \) is an invertible sheaf, and \( s \) a rational section not
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