# Abstract

These notes are based on Tony Pantev's "Algebra II" lectures at UPenn. Any mistake in what follows is my own.

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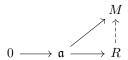
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# 1 Injective and flat modules

### 1.1 Lecture 1

**Proposition 1.1.1.** An R-module M is injective if and only if we can fill any injectivity diagram of ideal type, i.e.,



where  $\mathfrak{a}$  is an ideal in R.

Proof.

 $(\Longrightarrow)$ 

This is obvious.

 $(\Longleftrightarrow)$ 

Let

$$0 \longrightarrow X' \longrightarrow X$$

be an injectivity diagram of R-modules and define

$$S = \{(A, \xi) \mid X' \subset A \subset X, \ \xi : A \to M, \ \xi \upharpoonright_{X'} = \varphi\}.$$

By Zorn's lemma, there is some maximal element  $(N, \psi)$  of S. Suppose, toward a contradiction, that  $X \neq N$ . Pick any  $x \in X \setminus N$ . We have the ideal

$$\mathfrak{a} := \{a \in R : ax \in N\}$$

in R. Define the R-module morphism  $\theta : \mathfrak{a} \to M$  by  $a \mapsto \psi(ax)$ . By hypothesis, we get the following commutative diagram.

$$0 \longrightarrow \mathfrak{a} \stackrel{\theta}{\longleftrightarrow} \stackrel{\tilde{\theta}}{R}$$

Define the R-submodule  $\widetilde{N}=\langle N,x\rangle$ . We can write any  $z\in\widetilde{N}$  as z=y+ax for some  $y\in N$  and some  $a\in R$ . Define  $\widetilde{\psi}:\widetilde{N}\to M$  by  $y+ax\mapsto \psi(y)+\widetilde{\theta}(a)$ . To see that this is well-defined, let y+ax=y'+a'x. Then (y-y')=(a'-a)x, so that

$$\psi(y - y') = \psi((a' - a)x) = \tilde{\theta}(a' - a) = \tilde{\theta}(a') - \tilde{\theta}(a).$$

This implies that  $\tilde{\psi}$  is a well-defined homomorphism. But then  $\left(\tilde{N},\tilde{\psi}\right)>(N,\psi)$ , a contradiction.  $\Box$ 

Aside. The categorical dual P<sup>op</sup> of this recognition principle for injectivity expresses a recognition principle for projectivity, namely that for any R-module M, ideal  $I \subset R$ , and homomorphism  $\varphi : M \to R/I$ , we can fill the diagram

$$\begin{array}{ccc}
M \\
\downarrow & & \\
R & \longrightarrow & R/I & \longrightarrow & 0
\end{array}$$
(\*)

if and only if M is projective. This is equivalent to saying that M is projective if and only if the natural group map  $\operatorname{Hom}_R(M,R) \to \operatorname{Hom}_R\left(M,R/I\right)$  is surjective. But then  $\operatorname{P}^{\operatorname{op}}$  is precisely an affirmative answer to what is known as "Faith's problem on R-projectivity," which Trlifaj (2017) proved to be undecidable in ZFC + GCH. Therefore, both  $\operatorname{P}^{\operatorname{op}}$  and  $\operatorname{\neg}(\operatorname{P}^{\operatorname{op}})$  are consistent with ZFC + GCH.

### Corollary 1.1.2.

- 1. If R is an integral domain, then any injective R-module M is divisible.
- 2. If R is a PID, then M is injective if and only if it is divisible.

Proof.

1. Given any  $a \in R$ , we want to show that the homomorphism  $\operatorname{mult}_a : M \to M$  given by  $x \mapsto ax$  is surjective. The assumption that R is an integral domain entails that  $\operatorname{mult}_a : R \to R$  is injective. Note that  $\mathfrak{a} := \operatorname{mult}_a(R)$  is an ideal in R, giving the short exact sequence

$$0 \to \mathfrak{a} \to R \to R/\mathfrak{a} \to 0.$$

By assumption,  $\operatorname{Hom}_R(-,M)$  is exact, so that the sequence

$$0 \to \operatorname{Hom}\left(R/_{\mathfrak{A}}, M\right) \to \operatorname{Hom}_{R}(R, M) \to \operatorname{Hom}_{R}(\mathfrak{a}, M) \to 0$$

is exact. Since R and  $\mathfrak{a}$  are free R-modules of rank 1, it follows that  $\operatorname{Hom}_R(R,M)\cong M\cong \operatorname{Hom}_R(\mathfrak{a},M)$ . This means that the sequence

$$0 \to \operatorname{Hom}\left(R_{/\mathfrak{a}}, M\right) \to M \stackrel{\operatorname{mult}_a}{\longrightarrow} M \to 0$$

exact. In particular  $\operatorname{mult}_a$  is surjective.

2.  $(\Leftarrow)$  Suppose that M is divisible and R is a PID. We want to fill the injectivity diagram

$$0 \longrightarrow \mathfrak{a} \stackrel{\varphi}{\longrightarrow} \stackrel{M}{\underset{\psi|}{\uparrow}}.$$

where  $\mathfrak{a}$  is an ideal in R. We have that  $\mathfrak{a} = (a)$ . Therefore, the short exact sequence

$$0 \to \mathfrak{a} \to R \to R/\mathfrak{a} \to 0$$

is isomorphic to  $0 \to R \xrightarrow{\operatorname{mult}_a} R \to R/\mathfrak{a} \to 0$ . Since M is divisible, we know that  $M \xrightarrow{\operatorname{mult}_a} M \to 0$  is exact. Apply  $\operatorname{Hom}_R(-,M)$  to get the sequence

$$\operatorname{Hom}_R(R,M) \stackrel{(-) \circ \operatorname{mult}_a}{\longrightarrow} \operatorname{Hom}_R(R,M) \to 0,$$

which is isomorphic  $M \xrightarrow{\operatorname{mult}_a} M \to 0$ . This shows that  $\operatorname{Hom}_R(R,M) \xrightarrow{(-) \circ \operatorname{mult}_a} \operatorname{Hom}_R(R,M)$  is surjective. It follows that  $\varphi$  can be lifted to some  $\psi: R \to M$ .

# 1.2 Lecture 2

Corollary 1.2.1. Any abelian group is injective if and only if it's divisible.

Corollary 1.2.2. If R is a PID and M is an injective R-module, then every quotient of M is injective.

*Proof.* This follows from the fact that any quotient of a divisible group is divisible.  $\Box$ 

### Example 1.2.3.

- 1.  $\mathbb{Q}_{\mathbb{Z}}$  is injective.
- 2.  $S^1$  is injective.
- 3. Any non-trivial finitely generated abelian group G is never injective.

*Proof.* It suffices to show that G is never divisible. There exists a maximal proper subgroup  $H \leq G$ . Then G/H is a simple abelian group, so that  $G/H \cong C_p$  for some prime p. If G is divisible, then so must G/H. But  $C_p$  is not divisible, a contradiction.

**Theorem 1.2.4 (Baer embedding).** If R is a ring, then every module embeds into an injective module.

**Corollary 1.2.5.** For any R-module M, we can find an injective resolution

$$0 \to M \to I_0 \to I_1 \to \cdots \to I_k \to \cdots$$
.

*Proof.* We want to invent a duality operation that will convert  $R-\mathbf{Mod}^{\mathrm{op}}$  to  $R^{\mathrm{op}}-\mathbf{Mod}$  and then use projective objects in  $R^{\mathrm{op}}-\mathbf{Mod}$ . If T is an abelian group, then the functor

$$\mathbf{Ab} \xrightarrow{\mathrm{Hom}_{\mathbf{Ab}}(-,T)} \mathbf{Ab}^{\mathrm{op}}$$

will reverse arrows. The choice of T that ends up working is precisely  $\mathbb{Q}_{\mathbb{Z}}$ .

Claim. Let  $\operatorname{Hom}_{\mathbf{Ab}}(-,\mathbb{Q}_{\mathbb{Z}}) := (-)^D$ . Note that for any abelian group A, we have a canonical homomorphism  $\epsilon_A : A \to A^{DD}$  given by  $a \mapsto \left( \left[ \varphi : A \to \mathbb{Q}_{\mathbb{Z}} \right] \to \varphi(a) \right)$ . Then  $\epsilon_A$  is injective.

*Proof.* We need to show that if  $a \in A$  is nonzero, then we can find some homomorphism  $f: A \to \mathbb{Q}/\mathbb{Z}$  such that  $f(a) \neq 0$ .

<u>Case 1:</u> Suppose that  $|(a)| = n < \infty$ . Then define the homomorphism  $\varphi : (a) \to \mathbb{Q}_{\mathbb{Z}}$  by  $a \mapsto [\frac{1}{n}]$ . Since  $\mathbb{Q}_{\mathbb{Z}}$  is divisible in **Ab**, it is also injective. Thus, we may find some map  $\psi$  such that

$$0 \longrightarrow (a) \stackrel{\varphi}{\longleftrightarrow} A$$

commutes. This means that  $\psi(a) \neq 0$ , as required.

Case 2: If (a) has infinite order, then define  $\varphi:(a)\to\mathbb{Q}_{\mathbb{Z}}$  by  $a\mapsto\frac{1}{2}$  and apply a similar argument to Case 1.

The duality functor  $(-)^D$  extends to a functor  $(-)^D: R^{\mathrm{op}}-\mathbf{Mod} \to R-\mathbf{Mod}^{\mathrm{op}}$  that is compatible with forgetting the module structure. Indeed, if M is a left module over  $R^{\mathrm{op}}$ , then its module structure is given by a collection of maps  $\{\mathrm{mult}_a: M \to M \mid a \in R\}$ . Note that

$$\operatorname{mult}_a \circ \operatorname{mult}_b = \operatorname{mult}_{a \cdot_{R^{\operatorname{op}}} b} = \operatorname{mult}_{b \cdot_R a}.$$

For each  $a \in R$ , let  $\underline{\operatorname{mult}}_a(\varphi) = \varphi \circ \operatorname{mult}_a$ . Then the abelian group  $M^D$  has an R-module structure given by  $\underline{\operatorname{mult}}_a: M^D \to M^D$ , which clearly satisfies

$$\underline{\mathrm{mult}}_{ab} = \underline{\mathrm{mult}}_a \circ \underline{\mathrm{mult}}_b$$
.

**Lemma 1.2.6.** If M is a projective  $R^{op}$ -module, then  $M^D$  is an injective R-module.

*Proof.* Suppose that M is a projective  $R^{op}$ -module and consider the injectivity diagram

$$0 \longrightarrow X' \xrightarrow{\varphi} X$$

of R-modules. We want to lift  $\varphi: X' \to M^D$  to a map  $\psi: X \to M^D$ . Apply  $(-)^D$  to get a commutative diagram

$$0 \longleftarrow (X')^D \longleftarrow X^D$$

where the bottom row is exact because  $\mathbb{Q}_{\mathbb{Z}}$  is injective.

**Exercise 1.2.7.** Show that  $\epsilon_M: M \to M^{DD}$  is a map of  $R^{op}$ -modules.

We now have the following projectivity diagram of  $R^{\text{op}}$ -modules.

$$X^{D} \xrightarrow{\epsilon_{M} \circ \varphi^{D}} X^{D} \longrightarrow (X')^{D} \longrightarrow 0$$

By assumption, we may fill this diagram with some map  $\psi: M \to X^D$ . This induces the map  $\psi^D: X^{DD} \to M^D$ . Note that  $(\epsilon_M)^D \circ \varphi^{DD} = \psi^D \circ i^{DD}$  where  $i: X' \hookrightarrow X$ . But  $i^{DD} \upharpoonright_{X'} = i$  and  $\varphi^{DD} \upharpoonright_{X'} = \varphi$ , so that

$$\psi^D \circ i = (\epsilon_M)^D \circ \varphi = \varphi$$

on X'. It follows that

$$0 \longrightarrow X' \xrightarrow{\varphi} X^D \uparrow_{\psi^D \circ \epsilon_X}$$

commutes.

There is some surjection  $\bigoplus_{j\in J} R \to M^D$ . Therefore, we have a sequence of embeddings

$$M \hookrightarrow M^{DD} = \operatorname{Hom}_{\mathbb{Z}}\left(M^{D}, \mathbb{Q}_{\mathbb{Z}}\right) \hookrightarrow \operatorname{Hom}_{\mathbb{Z}}\left(\bigoplus_{j \in J} R, \mathbb{Q}_{\mathbb{Z}}\right) = \left(\bigoplus_{j \in J} R\right)^{D}$$
 injective by Lemma 1.2.6

**Definition 1.2.8.** Given two R-modules M and N, the additive invariants of M and N are the abelian groups

$$\operatorname{Ext}_R^i(M,N) := H^i(\operatorname{Hom}_R(P^{\bullet},N))$$

indexed by  $\mathbb{N}$  where  $P^{\bullet}$  is a chosen projective resolution of M.

# Proposition 1.2.9.

1.  $\operatorname{Ext}^i_R(M,N)$  is independent of our choice of projective resolution.

*Proof.* This follows from the fact that any two projective resolutions are chain homotopic.  $\Box$ 

2.  $\operatorname{Ext}_R^i(M,N) = H^i(\operatorname{Hom}_R(M,I_{\bullet}))$  for any injective resolution  $I_{\bullet}$  of N.

# Lemma 1.2.10.

- 1.  $\operatorname{Ext}_{R}^{0}(M, N) = \operatorname{Hom}_{R}(M, N)$
- 2.  $\operatorname{Ext}_{R}^{1}(M, N) = (\text{the group of isomorphism classes of extensions of } N \text{ by } M \text{ in } R-\operatorname{\mathbf{Mod}}).$

Proof. Let

$$\cdots \xrightarrow{\partial_1} P^1 \xrightarrow{\partial_0} P^0 \xrightarrow{\epsilon} M \to 0$$

be a projective resolution and let

$$(\xi): 0 \to N \xrightarrow{f} T \xrightarrow{g} M \to 0$$

be a short exact sequence of R-modules. Note that  $\operatorname{Hom}_R(P^k, -)$  is exact for each  $k \geq 0$ . Therefore, the sequence

$$0 \to \operatorname{Hom}_R(P^k, N) \xrightarrow{f_k} \operatorname{Hom}_R(P^k, T) \xrightarrow{g_k} \operatorname{Hom}_R(P^k, M) \to 0$$

is exact where  $f_k := f \circ (-)$  and  $g_k := g \circ (-)$ . Letting  $d_i := (-) \circ \partial_i$ , we get short exact sequences of complexes constituting the columns of

$$0 \longrightarrow \operatorname{Hom}_{R}(P^{0}, N) \xrightarrow{f_{0}} \operatorname{Hom}_{R}(P^{0}, T) \xrightarrow{g_{0}} \operatorname{Hom}_{R}(P^{0}, M) \longrightarrow 0$$

$$\downarrow^{d_{0}} \qquad \downarrow^{d_{0}} \qquad \downarrow^{d_{0}} \qquad \downarrow^{d_{0}}$$

$$0 \longrightarrow \operatorname{Hom}_{R}(P^{1}, N) \xrightarrow{f_{1}} \operatorname{Hom}_{R}(P^{1}, T) \xrightarrow{g_{1}} \operatorname{Hom}_{R}(P^{1}, M) \longrightarrow 0.$$

$$\downarrow^{d_{1}} \qquad \downarrow^{d_{1}} \qquad \downarrow^{d_{1}} \qquad \downarrow^{d_{1}}$$

$$\vdots \qquad \vdots \qquad \vdots \qquad \vdots$$

By definition,  $\operatorname{Ext}^i_R(M,N) = \underbrace{\ker d_i}_{\operatorname{im} d_{i-1}}$ . Since  $\operatorname{Hom}_R(-,M)$  is left-exact and  $P^1 \xrightarrow{\partial_1} P^0 \xrightarrow{\epsilon} M$ is exact, we also have the exact sequence

$$0 \to \operatorname{Hom}_R(M,M) \xrightarrow{(-) \circ \epsilon} \operatorname{Hom}_R(P^0,M) \xrightarrow{d_0} \operatorname{Hom}_R(P^1,M).$$

Let  $\psi \in \operatorname{Hom}_R(P^0, M)$  satisfy  $d_0(\psi) = 0$ . Then  $\psi = \varphi \circ \epsilon$  for some unique map  $\varphi : M \to M$ . Since  $g_0$ is surjective, there exists  $\alpha \in \operatorname{Hom}_R(P^0,T)$  such that  $g_0(\alpha) = \psi = \varphi \circ \epsilon$ . This implies that

$$g_1(d_0(\alpha)) = d_0(g_0(\alpha)) = d_0(\psi) = 0.$$

It follows that  $d_0(\alpha) \in \ker g_1 = \operatorname{im} f_1$ , so that  $d_0(\alpha) = f_1(\beta)$  for some  $\beta : P^1 \to N$ . Since  $f_2(d_1(\beta)) = f_1(\beta)$  $d_1(f_1(\beta)) = d_1(d_0(\alpha)) = 0$ , the fact that  $f_2$  is injective means that  $d_1(\beta) = 0$ . Hence  $\beta \in \ker d_1$ , and  $[\beta] \in \operatorname{Ext}^1_R(M,N)$ 

**Exercise 1.2.11.** Show that  $\psi \mapsto [\beta]$  is well-defined, i.e., that  $[\beta]$  is independent of  $\alpha$ .

This defines a map of abelian groups  $\delta_{\xi} : \operatorname{Hom}_{R}(M, M) \to \operatorname{Ext}_{R}^{1}(M, N)$  given by  $\varphi \mapsto [\beta]$ . Now, define the homomorphism

$$e: \operatorname{Ext}_R(M,N) \to \operatorname{Ext}^1_R(M,N), \quad (\xi) \mapsto \delta_{\xi}(\operatorname{id}_M).$$

Apply  $\operatorname{Hom}_R(M,-)$  to  $(\xi)$  to get the exact sequence

$$0 \to \operatorname{Hom}_R(M, N) \to \operatorname{Hom}_R(M, T) \to \operatorname{Hom}_R(M, M).$$

Claim. We can extend this sequence to a long exact sequence of abelian groups

$$0 \to \operatorname{Hom}_R(M,N) \to \operatorname{Hom}_R(M,T) \to \operatorname{Hom}_R(M,M) \xrightarrow{\delta_{\xi}} \operatorname{Ext}^1_R(M,N) \to \operatorname{Ext}^1_R(M,T) \to \operatorname{Ext}^1_R(M,M).$$

**Exercise 1.2.12.** Show that if  $(\xi)$  is split, then  $\delta_{\xi}(\mathrm{id}_M) = 0$ .

This implies that e is injective. We need to show that it is surjective as well. Suppose that How is it that eis injective?  $\gamma \in \operatorname{Ext}^1_R(M,N)$  and let  $I_{\bullet}$  be an injective resolution of N. Apply  $\operatorname{Hom}_R(M,-)$  to get

$$0 \to \operatorname{Hom}_R(M,N) \xrightarrow{\nu} \operatorname{Hom}_R(M,I_0) \xrightarrow{d_0} \operatorname{Hom}_R(M,I_1) \xrightarrow{d_1} \cdots$$

(where we have abused the notation  $d_i$ ). By Proposition 1.2.9(2), we have that  $\gamma = [f]$  for some  $f \in \ker d_1$ . Note that  $f: M \to \ker \partial_1 = \operatorname{im} \partial_0$ , giving

$$0 \longrightarrow N \longrightarrow I_0 \xrightarrow{\partial_0} \operatorname{im} \partial_0 \longrightarrow 0$$

$$f \uparrow \\ M$$

where the top row is exact. Take the pullback of  $\partial_0$  and f to obtain T such that

$$0 \longrightarrow N \longrightarrow I_0 \stackrel{\partial_0}{\longrightarrow} \operatorname{im} \partial_0 \longrightarrow 0$$

$$\downarrow \qquad \qquad \uparrow \qquad \qquad \uparrow \uparrow \qquad \qquad \cdot$$

$$0 \longrightarrow N \longrightarrow T \longrightarrow M \longrightarrow 0$$

#### Exercise 1.2.13.

- 1. Show that the map  $\rho : \operatorname{Ext}_R^1(M,N) \to \operatorname{Ext}_R(M,N)$  given by  $\gamma \mapsto \xi$  is independent of our choice of f.
- 2. Show that  $\rho$  is the inverse of e.

# 1.3 Lecture 3

Let N be a right R-module and M an R-module. Recall that  $N \otimes_R M \in \text{ob}(\mathbf{Ab})$  is precisely the object in  $\mathbb{Z}$ -Mod representing the functor  $B_{M,N} : \mathbf{Ab} \to \mathbf{Ab}$  given by

$$A \mapsto \{f : M \times N \to A \mid f(ax, y) = f(x, ay)\}.$$

Moreover, recall that N is flat if the functor

$$N \otimes_R (-) : R-\mathbf{Mod} \to \mathbf{Ab}$$

is exact.

**Definition 1.3.1.** Let N be a right R-module and M an R-module. Let  $x_1, \ldots, x_n \in M$ .

(1) A relation of the  $x_i$ 's with coefficients in R is a list of scalars  $a_1, \ldots, a_n \in R$  such that

$$\sum_{i=1}^{n} a_i x_i = 0.$$

(2) A relation of the  $x_i$ 's with coefficients in N is a list of elements  $y_1, \ldots, y_n \in N$  such that

$$\sum_{i=1}^{n} y_i \otimes x_i = 0.$$

Since  $R \otimes_R M \cong M$ , we see that (1) is a special case of (2). Let

$$a_1 \coloneqq (a_{11}, \dots, a_{1n})$$

$$a_2 \coloneqq (a_{21}, \dots, a_{2n})$$

$$\vdots$$

$$a_m \coloneqq (a_{21}, \dots, a_{mn}).$$

be relations of  $x_1, \ldots, x_n$  with coefficients in R. Let  $(z_1, \ldots, z_m) \in N^m$ . If A denotes the matrix  $(a_{ij})$ , then  $y = A^t z \in N^n$  is a relation with coefficients in N.

**Definition 1.3.2.** A relation y with coefficients in N follows from R-relations if y is of the form  $A^tz$  for some z and some matrix A of relations in R.

**Lemma 1.3.3.** A right R-module N is flat if and only if for any R-module M and any  $x_1, \ldots, x_n \in M$ , every N-relation among the  $x_i$  follows from R-relations.

Proof.

 $(\Longrightarrow)$ 

We have a module homomorphism  $\varphi: \mathbb{R}^n \to M$  given by  $(a_1, \ldots, a_n) \mapsto \sum_{i=1}^n a_i x_i$ . Then

$$\underbrace{\ker \varphi}_{K} = \{(r_1, \dots, r_n) \in \mathbb{R}^n \mid (r_1, \dots, r_n) \text{ is a relation of the } x_i\text{'s in } \mathbb{R}\}.$$

We have an exact sequence

$$0 \to K \xrightarrow{i} R^n \xrightarrow{\varphi} M.$$

If N is flat, then  $N \otimes_R (-)$  is exact, so that

$$0 \to N \otimes_R K \xrightarrow{\tilde{i}} N^n \xrightarrow{\tilde{\varphi}} N \otimes_R M$$

is exact. Thus,  $\ker \tilde{\varphi} = (N\text{-relations}) = N \otimes_R K$ .

(<del>\_\_\_</del>)

Let  $0 \to M' \xrightarrow{f} M \xrightarrow{g} M'' \to 0$  be a short exact sequence of R-modules. Since  $N \otimes_R (-)$  is right exact, it suffices to show that

$$N \times_R M' \stackrel{\mathrm{id}_N \otimes f}{\longrightarrow} N \otimes_R M$$

is injective. Let  $z \in \ker \operatorname{id}_N \otimes f$ . Then  $z = \sum_{i=1}^n y_i \otimes z_i$ . We know that

$$\sum_{i=1}^{n} y_i \otimes f(z_i) = \mathrm{id}_N \otimes f(z) = 0,$$

and thus  $(y_1, \ldots, y_n)$  is an N-relation among the  $f(z_i) \in M$ . This shows that there exist  $\left(a_i^j\right) \in R$  where  $i = 1, \ldots, n$  and  $j = 1, \ldots, m$  and elements  $v_1, \ldots, v_m \in N$  such that  $y_i = \sum_{i=1}^m v_j a_i^j$ . Therefore,  $\sum_{i=1}^n a_i^j f(z_i) = 0$  for each j. But

$$0 = \sum_{i=1}^{n} a_i^j f(z_i) = f\left(\sum_{i=1}^{n} a_i^j z_i\right).$$

As f is injective, it follows that  $\sum_{i=1}^{n} a_i^j z_i = 0$  for each j. Finally, we compute

$$\sum_{i=1}^{n} y_i \otimes z_i = \sum_{i=1}^{n} \sum_{j=1}^{m} (v_j a_i^j) \otimes z_i$$
$$= \sum_{j=1}^{n} v_j \otimes \left( \sum_{i=1}^{n} a_i^j z_i \right)$$
$$= \sum_{j=1}^{n} (v_j \otimes 0) = 0.$$

#### Corollary 1.3.4.

- 1. Any free module is flat.
- 2. Any colimit of flat modules is flat.
- 3. Any direct summand of a free module is flat, so that any projective module is flat.
- 4. Any colimit of projective modules is flat.

# 2 Localization

#### 2.1 Lecture 4

Let R be a commutative ring. Given  $x \in R$ , when can we make x multiplicatively invertible, perhaps in a new ring? This is a question of representability. We have a functor  $\Phi_x$ : CommRing  $\to$  Set given by

$$B \mapsto \{\varphi : R \to B \mid \varphi(x) \in B^{\times}\} \subset \operatorname{Hom}_{\mathbf{CommRing}}(R, B).$$

We are asking whether or not  $\Phi_x$  is representable. That is, we want to find some pair  $(R_x, h)$  where  $R_x$  is a commutative ring and  $h: R \to R_x$  is a morphism such that  $h(x) \in (R_x)^{\times}$  and if  $\varphi: R \to B$  with  $\varphi(x) \in B^x$ , then  $\varphi \circ h = \varphi$  for some map  $\varphi: R_x \to B$ .

In general, we can consider a set S of nonzero elements and ask for a universal way of making them invertible. But if we make S invertible, then we will also make the *multiplicative closure* cl(S) of S invertible.

**Definition 2.1.1.** Any  $S \subset R$  is called *multiplicatively closed* if  $0 \notin S$ ,  $1 \in S$ , and  $x, y \in S \implies xy \in S$ .

Given a multiplicatively closed subset  $S \subset R$ , we want to find a universal way of inverting every element of S. Equivalently, find a ring representing  $\Phi_S$ . Equivalently, we want to find a pair  $\left(S^{-1}R,h\right)$  where  $h:R\to S^{-1}R$  such that  $h(S)\subset (S^{-1}R)^{\times}$  and any  $\varphi:R\to B$  with  $\varphi(S)\subset B^{\times}$  has  $\varphi=\underline{\varphi}\circ h$  for some unique map  $\varphi:R\to B$ . We call the pair  $\left(S^{-1}R,h\right)$  the localization of R along S.

Formally adjoin to R fractions with numerator in R and denominator in S. Consider the set  $(R \times S, \sim)$  where  $(a, s) \sim (b, t)$  if u(at - bs) = 0 for some  $u \in S$ . Set  $S^{-1}R := R \times S / \sim$ . Let  $\frac{a}{s} := [(a, s)]$ . Define

$$\frac{a}{s} + \frac{b}{t} = \frac{at + bs}{st}$$

and

$$\frac{a}{s} \cdot \frac{b}{t} = \frac{ab}{st}.$$

Then  $S^{-1}R$  becomes a ring with unity  $\frac{1}{1}$ . Also, we see that  $h: R \to S^{-1}R$  given by  $a \mapsto \frac{a}{1}$  is a ring homomorphism. Given a map  $\varphi: R \to B$  such that im  $\varphi \subset B^{\times}$ , we have a well-defined map of rings  $\underline{\varphi}: S^{-1}R \to B$  given by  $\frac{a}{s} \mapsto \varphi(a)\varphi(s)^{-1}$ , which satisfies  $\underline{\varphi} \circ h = \varphi$ .

#### **Example 2.1.2.** Here are some natural choices for S.

- (a)  $\{1, x, x^2, \ldots\}$  with x not nilpotent.
- (b)  $R^{\times}$ .
- (c)  $\{r \in R : r \text{ is not a zero divisor}\}.$

If A is an integral domain and we take any multiplicatively closed subset  $S \subset A$ , then  $\operatorname{Frac}(A) \coloneqq (A \setminus \{0\})^{-1} A$  is a field and  $h: A \to (A \setminus \{0\})^{-1} A$  is injective. For now, let S denote the set of non zero-divisors. If  $\frac{a}{b} \in \operatorname{Frac}(A)$  is nonzero, then  $\frac{a}{b} \neq \frac{0}{1}$ , i.e.,  $a \cdot 1$  is not a zero divisor, so that  $a \neq 0$  and thus  $\frac{b}{a} \in \operatorname{Frac}(A)$ . This shows that  $\operatorname{Frac}(A)$  is a field. Moreover, if  $a \in A$  satisfies  $h(a) = \frac{a}{1} = 0 \in \operatorname{Frac}(A)$ , then  $\frac{a}{1} = \frac{0}{1} \implies a \cdot 1$  is a zero divisor. Hence a = 0, and h is injective.

If S is generic, then  $S^{-1}A \subset \operatorname{Frac}(A)$  since  $S^{-1}A$  equals the subring generated by  $A \cong h(A)$  and  $S^{-1} = \left\{ \frac{1}{s} \mid s \in S \right\}$ . In this case,  $\left( S^{-1}A, h \right)$  represents the functor  $\Phi : \mathbf{Field} \to \mathbf{Set}$  given by  $k \mapsto \{\varphi : A \to k \mid \varphi \text{ is injective.}\}$ . This means that for any ring map  $\varphi : A \to B$  with  $\varphi(S) \subset B^{\times}$ , there is some unique map  $\psi$  such that  $\psi \circ h = \varphi$ .

#### 2.2 Lecture 5

# Example 2.2.1.

- 1. If  $S = \{1, x, x^2, \ldots\}$  with x not nilpotent, then  $S^{-1}A = A_f \equiv \left\{\frac{a}{f^n} : n \geq 0, \ a \in A\right\}$ .
- 2. If  $S \subset A^{\times}$ , then  $h: A \to S^{-1}A$  is an isomorphism.
- 3. If A is any ring and  $S \subset A$  denotes the set of all non-zero divisors, then  $\operatorname{Frac}(A) = S^{-1}A$  is called the *fraction ring of* A. If A is an integral domain, then  $\operatorname{Frac}(A)$  is a field (called the *field of fractions of* A) and  $H: A \to \operatorname{Frac}(A)$  is injective. In this case,  $(\operatorname{Frac}(A), h)$  represents the functor  $F_A: \mathbf{Field} \to \mathbf{Set}$  given by  $K \mapsto \{\varphi: A \to K \mid \varphi \text{ monomorphism}\}.$

Let A be a commutative ring and  $S \subset A$  be multiplicatively closed. Let M be an A-module. Define the equivalence relation  $(M \times S, \sim)$  where  $(m, s) \sim (n, t)$  if u(tm - sn) = 0 for some  $u \in S$ .

Define the A-module  $S^{-1}M = M \times S / \infty$  where  $\frac{m}{s} + \frac{n}{t} := \frac{tm + sn}{st}$ . Define the module homomorphism  $h_M : M \to S^{-1}M$  by  $m \mapsto [(m, 1)]$ . Let  $\frac{m}{s}$  denote the equivalence class [(m, s)].

Moreover,  $S^{-1}M$  is naturally a module over  $S^{-1}A$  via the action  $\frac{a}{s} \cdot \frac{m}{t} := \frac{a \cdot m}{st}$ . This makes  $h_M$  a module over  $h: A \to S^{-1}A$  in that for any  $a \in A$  and  $m \in M$ , we have that  $h_M(a \cdot m) = h(a) \cdot h_M(m)$ .

We see that  $S^{-1}(-)$  is a functor which maps each homomorphism  $\varphi: M \to N$  to  $S^{-1}\varphi: S^{-1}M \to S^{-1}N$  given by  $\frac{m}{s} \mapsto \frac{\varphi(m)}{s}$ . It's easy to verify that  $S^{-1}(-)$  is left adjoint to the pullback functor  $h^{\bullet}$ .

If  $f: A \to B$  is a map of commutative rings, then there are natural functors  $f^{\bullet}: B-\mathbf{Mod} \to A-\mathbf{Mod}$  and  $f_{\bullet}: A-\mathbf{Mod} \to B-\mathbf{Mod}$ , called the *pullback* and *pushforward*, respectively.

On the one hand, the pullback functor is already familiar to us. On the other hand, the pushforward acts on objects by

$$f_{\bullet}(M) \equiv B \otimes_A M$$

where B is viewed as an A-module via f along with the action  $b \cdot (c \otimes m) \equiv (bc) \otimes m$ . It acts on morphisms by  $(\varphi : M \to N) \mapsto (\mathrm{id}_B \otimes \varphi : f_{\bullet}(M) \to f_{\bullet}(N))$ .

**Exercise 2.2.2.**  $(f_{\bullet}, f^{\bullet})$  is an adjoint pair.

Corollary 2.2.3.  $S^{-1}(-) \cong h_{\bullet}$ .

Naively, we could have tried to define fractions in A by  $(a,s) \sim_n (b,t)$  if (at-bs=0). But this is not in general an equivalence relation, for it is not transitive. Indeed, set  $A = \mathbb{C}[x,y]/(xy)$  and  $S = \{1, x, x^2, \ldots\}$ . Consider the localization  $A_x$ . Note that  $(y,1) \not\sim_n (0,1)$  but that  $(y,1) \sim_n (0,x)$  and  $(0,x) \sim_n (0,1)$ .

**Note 2.2.4.** We have that  $A_x = \mathbb{C}[x, x^{-1}]$ , which is a field, and that  $h: A \to A_x$  is given by

$$\underbrace{[f(x,y)]}_{[p(x)+uq(y)]} \mapsto p(x),$$

which is non-injective.

#### Proposition 2.2.5.

- 1. If  $h: A \to S^{-1}A$ , then  $\ker h = \{a \in A : (\exists s \in S) (sa = 0)\}$ .
- 2.  $S^{-1}A$  is flat as an A-module.

Corollary 2.2.6.  $S^{-1}(-)$  is an exact functor.

*Proof.* Let  $M \xrightarrow{f} T \xrightarrow{g} N$  be an exact sequence of A-modules. We want to show that

$$S^{-1}M \stackrel{S^{-1}f}{\longrightarrow} S^{-1}T \stackrel{S^{-1}g}{\longrightarrow} S^{-1}N$$

is exact as well. Let  $\frac{x}{s} \in S^{-1}T$  with  $\left(S^{-1}g\right)\left(\frac{x}{s}\right) = 0$ . This implies that  $\frac{g(x)}{s} = \frac{0}{1}$ , so that ug(x) = 0 for some  $u \in S$ . But since g is a morphism, we know that 0 = ug(x) = g(ux). This means that f(y) = ux for some  $y \in M$ . Then  $\frac{y}{us} \in S^{-1}M$  such that  $\left(S^{-1}f\right)\left(\frac{y}{us}\right) = \frac{f(y)}{us} = \frac{ux}{us} = \frac{x}{s}$ .

Suppose that  $f \in A$  is not nilpotent. We can compute  $A_f$  explicitly as follows. There is a natural map  $A_f[x] \to A_f$  given by  $x \mapsto \frac{1}{f}$ . This induces an isomorphism

$$A_f[x]/(x-\frac{1}{f}) \xrightarrow{\cong} A_f.$$

We also have a map  $A[x] \to A_f[x]$  from the map h on the coefficients. Define the map  $\alpha: A[x] \to A_f$  by  $a \mapsto h(a) = \frac{a}{1}$  and  $x \mapsto \frac{1}{f}$ . We must compute  $\ker \alpha$  as an ideal in A[x]. This is surjective since any element in  $A_f$  is of the form  $\frac{a}{f^n}$  for some  $a \in A$  and  $n \in \mathbb{N}$ , so that  $ax^n \mapsto \frac{a}{f^n}$ .

# 2.3 Lecture 6

Claim.  $\ker \alpha = (fx - 1)$ .

*Proof.* Note that  $xf - 1 \in \ker \alpha$ . Also, note that

$$\exists n \geq 0 \text{ s.t. } f^n g\left(\frac{1}{f}\right) = 0 \iff \alpha(g) = 0 \iff g\left(\frac{1}{f}\right) = 0 \text{ in } A_f.$$

Without loss of generality, we may assume that  $n \ge \deg g$ . Thus,  $f^n g(x)$  is a polynomial of fx with coefficients in A, so that there is some  $G(y) \in A[y]$  such that  $G(fx) = f^n g(x)$ . Then

$$g \in \ker \alpha \iff \exists G(y) \in A[y] \text{ s.t. } G(fx) = f^n g(x) \land G(1) = 0.$$

But then G(y) = (y-1) h(y) where  $h(y) \in A[y]$ . This implies that

$$g(x) \in \ker \alpha \iff \exists n \ge 0 \text{ s.t. } f^n g(x) \in (xf - 1).$$

But  $f, fx-1 \in A[x]$  are relatively prime since 1 = fx+(fx-1)(-1). Hence  $1^n = (fx+(fx-1)(-1))^n = f^nx^n + (fx-1)s(x)$  for some  $s(x) \in A[x]$ . Multiply by g(x) to get

$$g(x) = \underbrace{f^n g(x)}_{\bigcap} + \underbrace{(fx-1)s(x)g(x)}_{(xf-1)}.$$

Therefore,  $g(x) \in (xf - 1)$ , and  $(xf - 1) = \ker \alpha$ .

Suppose that  $\varphi:A\to B$  is a map of commutative rings. Then we can transport the ideals along  $\varphi$  as follows.

#### Definition 2.3.1.

- 1. Given an ideal  $\mathfrak{a} \subseteq A$ , the extension of  $\mathfrak{a}$  along  $\varphi$  is the ideal  $\mathfrak{a}^e \subseteq B$  that is generated by  $\varphi(\mathfrak{a})$ , i.e.,  $\mathfrak{a}^e = \varphi(\mathfrak{a}) \cdot B$ .
- 2. Given an ideal  $\mathfrak{b} \subseteq B$ , the contraction of  $\mathfrak{b}$  along  $\varphi$  is defined as the ideal  $\mathfrak{b}^c = \varphi^{-1}(\mathfrak{b})$ .

Suppose that A is a commutative ring and that  $S \subset A$  is multiplicatively closed. Recall the localization morphism  $h: A \to S^{-1}A$ . We want to study  $(-)^e$  and  $(-)^c$  along h.

# Proposition 2.3.2.

1. If  $\mathfrak{a} \leq A$ , then  $\mathfrak{a}^e = \left\{ \frac{a}{s} \mid a \in \mathfrak{a}, s \in S \right\}$ .

*Proof.* By definition,  $\mathfrak{a}^e = h(\mathfrak{a}) \cdot S^{-1}A = \left\{ \sum_i \frac{b_i}{t_i} \frac{a_i}{s_i} \mid a_i \in \mathfrak{a}, \ b_i \in A, \ s_i, t_i \in S \right\}$ . Since  $a_i b_i \in \mathfrak{a}$  and  $s_i t_i \in S$ , our proof is complete.

2. If  $\mathfrak{a} \subseteq A$ , then  $\mathfrak{a}^e = (1) \iff \mathfrak{a} \cap S \neq \emptyset$ .

*Proof.* Note that  $\left(S^{-1}A\right)^{\times}$  consists of every fraction  $\frac{a}{s}$  for which we can find some fraction  $\frac{b}{t}$  such that  $\frac{a}{s}\frac{b}{t}=1$ . Therefore, we must have some element  $u\in S$  such that  $u(ab-st)=0\iff \exists\beta\in A \text{ s.t. } \beta a\in S$ . Thus,  $\left(S^{-1}A\right)^{\times}=\left\{\frac{a}{s}\mid \exists\beta\in A \text{ s.t. } \beta a\in S\right\}$ . But then

$$\begin{split} \mathfrak{a}^e &= (1) \iff \mathfrak{a}^e \text{ contains some unit} \\ &\iff \left( \exists \frac{a}{s} \in \mathfrak{a}^e \right) \left( \exists \beta \in A \right) \left( \beta \cdot a \in S \right) \\ &\iff \beta \cdot a \in S \cap \mathfrak{a}. \end{split}$$

Suppose that  $I \subseteq S^{-1}A$  is an ideal . Then we can form  $I^{ce} \subseteq S^{-1}A$ . By definition,  $I \supset I^{ce}$ .

### Proposition 2.3.3.

1. In fact,  $I = I^{ce}$ .

*Proof.* If  $\frac{a}{s} \in I$ , then  $a \in h^{-1}(I)$  because  $h^{-1}(I) = \{r \in A \mid \frac{r}{1} \in I\}$ . But  $\frac{a}{1} = s \cdot \frac{a}{s}$  where  $s \in S^{-1}A$  and  $\frac{a}{s} \in I$ , so that  $a \in I$ . This implies that  $\frac{a}{s} \in I^{ce}$  for each  $s \in S$ , and thus  $I \subset I^{ce}$ .

2. If  $\mathfrak{a} \subseteq A$ , then  $\mathfrak{a}^{ec} = \{r \in A \mid \exists s \in S \text{ s.t. } sr \in \mathfrak{a}\}.$ 

Proof. Suppose that  $a \in \mathfrak{a}^{ec}$ . Then  $\frac{a}{1} = h(a) \in \mathfrak{a}^{e}$ , so that  $(\exists b \in \mathfrak{a}) (\exists s \in S) (\frac{a}{1} = \frac{b}{s})$ . This implies that  $\exists u \in S$  such that u(sa-b) = 0. Hence (us) a = b, and  $\mathfrak{a}^{ec} \subset \{r \in A \mid \exists s \in S \text{ s.t. } sr \in \mathfrak{a}\}$ . If  $r \in A$  satisfies  $rs \in \mathfrak{a}$  for some  $s \in S$ , then  $\frac{r}{1} = \frac{rs}{s} \in \mathfrak{a}^{e}$  and thus  $r \in \mathfrak{a}^{ec}$ .

- 3.  $\mathfrak{a} \subseteq A$  is contracted (i.e.,  $\mathfrak{a} = I^c$  for some  $I \subseteq S^{-1}A$ ) if and only if  $\mathfrak{a} = \mathfrak{a}^{ec}$  if and only if  $[s] \in A/\mathfrak{a}$  is not a zero divisor for any  $s \in S$ .
- 4. The map  $(-)^e$  induces a bijection

$$\left(-\right)^{e}:\left\{ \mathfrak{a}\trianglelefteq A\mid\mathfrak{a}\text{ is a contraction of some ideal}\right\} \rightarrow\left\{ I\mid I\trianglelefteq S^{-1}A\right\}$$

that preserves inclusions of ideals.

Suppose that M is an A-module.

**Definition 2.3.4.** A submodule  $N \subset M$  is S-saturated if  $N = \{x \in M \mid (\exists s \in S) (sx \in N)\}.$ 

If M = A and  $N = \mathfrak{a}$ , then N is S-saturated if and only if  $\mathfrak{a} = \mathfrak{a}^{ec}$ . The localization on modules induces an inclusion-preserving bijection

$$S^{-1}(-): \{N \subset M \mid N \text{ is } S\text{-saturated}\} \to \{M \mid M \subset S^{-1}M\}.$$

**Definition 2.3.5.** Let  $\mathfrak{a}$  be an ideal in A.

- 1. We say that  $\mathfrak{a}$  is a *maximal ideal* if it is properly contained in A and is maximal in the set of all properly contained ideals in A partially ordered by inclusion.
- 2. We say that  $\mathfrak{a}$  is a prime ideal if  $xy \in \mathfrak{a} \implies x \in \mathfrak{a} \vee y \in \mathfrak{a}$ .

**Exercise 2.3.6.** An ideal  $\mathfrak{b} \subseteq A$  is prime if and only if  $A \setminus \mathfrak{b}$  is multiplicatively closed.

### 2.4 Lecture 7

**Proposition 2.4.1.** If  $\mathfrak{p} \subseteq A$  is prime and  $S \subset A$  is multiplicatively closed, then  $\mathfrak{p}^e \subseteq S^{-1}A$  is prime if and only if  $S \cap \mathfrak{p} = \emptyset$ .

Proof. The forward direction is obvious. Conversely, suppose that  $S \cap \mathfrak{p} = \emptyset$ . Then  $\mathfrak{p}^{ec} = \mathfrak{p}$ . Indeed,  $\mathfrak{p}^{ec} = \{a \in A \mid \exists s \in S \text{ s.t. } sa \in \mathfrak{p}\}$ . But if  $sa \in \mathfrak{p}$ , then either  $s \in \mathfrak{p}$  or  $a \in \mathfrak{p}$ . Since  $S \cap \mathfrak{p} = \emptyset$ , we see that  $s \notin \mathfrak{p} \implies a \in \mathfrak{p}$ . Suppose that  $x \cdot y \in \mathfrak{p}^e$ . Then  $x = \frac{a}{s}$  for some  $a \in A$  and  $s \in S$ , and  $y = \frac{b}{t}$  for some  $b \in A$  and  $t \in B$ . Then  $\frac{ab}{st} \in \mathfrak{p}^e$ , so that  $\frac{ab}{t} \in \mathfrak{p}^e$  since  $\mathfrak{p}^e$  is an ideal. Hence  $ab \in \mathfrak{p}$ , which is prime by assumption. Say that  $a \in \mathfrak{p}$ . Then  $\frac{a}{s} \in \mathfrak{p}^e$ .

Corollary 2.4.2. If  $S \subset A$  is multiplicatively closed, then we get a bijection

$$\{\mathfrak{p} \in \operatorname{Spec}(A) \mid \mathfrak{p} \cap S = \emptyset\} \xrightarrow{(-)^e} \operatorname{Spec}(S^{-1}A).$$

*Proof.* This is because  $\mathfrak{p}^c$  is prime in A with  $\mathfrak{p}^c \cap S = \emptyset$  whenever  $\mathfrak{p}$  is prime in  $S^{-1}A$ .

Now, recall the property P that an ideal  $\mathfrak a$  in A is prime if and only if  $A \setminus \mathfrak a$  is multiplicatively closed.

**Proposition 2.4.3.**  $\mathfrak{a}$  is prime if and only if there is some multiplicatively closed  $S \subset \mathfrak{a}$  such that  $S \cap A = \emptyset$  and  $\mathfrak{a}$  is maximal among all ideals satisfying P.

*Proof.* If  $\mathfrak{a}$  is prime, then  $S = A \setminus \mathfrak{a}$  is multiplicatively closed and  $\mathfrak{a}$  is maximal. Conversely, let  $a, b \in A$  such that  $a, b \notin \mathfrak{a}$ . We must show that  $ab \notin \mathfrak{a}$ . Consider  $\mathfrak{a} + (a) \supsetneq \mathfrak{a}$  and  $\mathfrak{a} + (b) \supsetneq \mathfrak{a}$ . But we are given S such that  $\mathfrak{a} \cap S = \emptyset$ . Hence there are  $s \in S \cap (\mathfrak{a} + (a))$  and  $t \in S \cap (\mathfrak{a} + (b))$ . Then  $s = \alpha + x \cdot a$  and  $t = \beta + y \cdot b$  where  $\alpha, \beta \in \mathfrak{a}$  and  $x, y \in A$ . We compute

$$st = \alpha\beta + \alpha yb + \beta xa + xyab,$$

where  $st \in S$  and  $\alpha\beta$ ,  $\alpha yb$ ,  $\beta xa \in \mathfrak{a}$ . If we assume that  $ab \in \mathfrak{a}$ , then  $st \in S \cap \mathfrak{a}$ , a contradiction.

Note 2.4.4. If  $S \subset A$  is multiplicatively closed, then by Zorn's lemma there is some prime ideal  $\mathfrak{b}$  such that  $\mathfrak{b} = A \setminus S$ .

**Definition 2.4.5.** Let A be a ring. We call A a *local ring* if any of the following equivalent conditions holds.

- (a) A has a unique maximal ideal  $\mathfrak{m}$ .
- (b)  $A \setminus A^{\times}$  is an ideal.
- (c) If  $\mathfrak{m}$  is maximal and  $x \in \mathfrak{m}$ , then  $1 + x \in A^{\times}$ .

If A is a ring and  $\mathfrak{p}$  a prime ideal, we will denote the localization  $(A \setminus \mathfrak{p})^{-1} A$  by  $A_{\mathfrak{p}}$ .

**Proposition 2.4.6.** If  $\mathfrak{p}$  is prime, then  $A_{\mathfrak{p}}$  is a local ring with unique maximal ideal  $\mathfrak{p}^e$ .

Proof. Let  $S = A \setminus \mathfrak{p}$ . Then  $A_{\mathfrak{p}} = S^{-1}A$ . Suppose that  $I \leq A_{\mathfrak{p}} = S^{-1}A$  such that  $I \neq (1)$ . But any ideal in  $S^{-1}A$  is of the form  $I = \mathfrak{a}^e$  for some ideal  $\mathfrak{a}$  in A. Since  $(1) \neq I = \mathfrak{a}^e$ , it follows that  $\mathfrak{a} \cap S = \emptyset$ . Therefore,  $\mathfrak{a} = A \setminus S = \mathfrak{p}$ , so that  $I = \mathfrak{a}^e \subset \mathfrak{p}^e$ . Hence every nontrivial ideal in  $A_{\mathfrak{p}}$  is contained in  $\mathfrak{p}^e$ , implying that  $\mathfrak{p}^e$  is the unique maximal ideal.

Corollary 2.4.7. In particular, the map

(prime ideal of 
$$A \mid A \subset \mathfrak{p}$$
)  $\stackrel{(-)^e}{\longrightarrow} \operatorname{Spec}(A_{\mathfrak{p}})$ 

is a bijection that preserves inclusions of ideals.

**Definition 2.4.8.** Let A be a commutative ring, For every  $\mathfrak{p} \subseteq A$  prime, the height of  $\mathfrak{p}$  is

$$ht(\mathfrak{p}) \equiv \sup\{k \mid \mathfrak{p} = \mathfrak{p}_0 \supseteq \mathfrak{p}_1 \supseteq \mathfrak{p}_2 \supseteq \cdots \supseteq \mathfrak{p}_k, \ \mathfrak{p}_i \trianglelefteq A\}.$$

Note that  $ht(\mathfrak{p}) = ht(\mathfrak{p}^e \text{ in } A_{\mathfrak{p}}).$ 

**Definition 2.4.9.** The Krull dimension of A is

$$\dim A \equiv \sup \{ \mathsf{ht}(\mathfrak{m}) \mid \text{maximal } \mathfrak{m} \not\subseteq A \}.$$

Note that dim  $A_{\mathfrak{p}} = \mathsf{ht}(\mathfrak{p})$  and that dim  $A = \sup \{ \dim A_{\mathfrak{m}} \mid \text{maximal } \mathfrak{m} \not\subseteq A \}$ .

# Example 2.4.10.

- 1. If k is a field, then dim k = 0. (The converse is also true.)
- 2. If A is a PID, then dim A = 1. For example,  $\mathbb{Z}$ ,  $\mathbb{Q}[x]$ , and  $\mathbb{Z}[i]$  have dimension 1.

#### Exercise 2.4.11.

- 1. Show that  $\mathbb{Z}\left[-\sqrt{5}\right]$  is not a PID but has dimension 1.
- 2. Sow that dim  $\mathbb{C}[x_1,\ldots,x_n]=n$ .

# 3 Basic algebraic geometry

Any information about a commutative ring A, a prime ideal in A, a localization in A, and the relations between them can be packaged into a geometrical object, specifically, a topological space along with a distinguished class of maps.

Let  $X = \mathsf{Spec}(A)$ , the set of all prime ideals in A, or spectrum of A. For any  $f \in A$ , define the principal open subset associated with f as

$$X_f \equiv \{ \mathfrak{p} \in X \mid f \notin \mathfrak{p} \} .$$

Such subsets satisfy

- (a)  $X_f \cap X_g = X_{fg}$ .
- (b)  $X_{f^n} = X_f, X_f = X \iff f \notin \mathfrak{p} \ \forall \mathfrak{p} \ \text{prime} \iff f \in A^{\times}.$
- (c)  $X_f = \emptyset \iff f \in \bigcap_{\mathfrak{p} \subset X} \mathfrak{p}$ .

**Definition 3.0.1.** The minimal topology on X generated by  $\{X_f\}_{f\in A}$  is called the *Zariski topology* on X.

The subset  $U \subset X$  is open if and only if there is some  $T \subset A$  such that  $U = \bigcup_{f \in T} X_f$ . Also,  $Y \subset X$  is closed if  $Y = \bigcap_{f \in T} (X \setminus X_f)$  for some  $T \subset A$ . Hence  $Y \subset X$  is closed if there is some  $T \subset A$  such that

$$Y = \{ \mathfrak{p} \mid \mathfrak{p} \supset \langle T \rangle \}.$$

In particular, for any ideal  $\mathfrak{a} \subseteq A$ , we can define a Zariski-closed subset  $V(\mathfrak{a}) = \{\mathfrak{p} \in X \mid \mathfrak{p} \supset \mathfrak{a}\}$ . (Note that replacing  $\mathfrak{a}$  with a set  $S \subset A$  determines an equivalent topology.) Every closed subset is of this form.

**Exercise 3.0.2.** Write arbitrary intersections of closed sets, finite unions of closed sets, X, and  $\emptyset$  in this form.

Any  $f \in A$  can be viewed as a function on X in two ways. First, view f as a mapping  $X \to \coprod_{\mathfrak{p} \in X} A_{\mathfrak{p}}$  given by  $\mathfrak{p} \mapsto \frac{f}{1} \in A_{\mathfrak{p}}$ . Then for any  $\mathfrak{p}$ , the value of f on  $\mathfrak{p}$  is in  $A_{\mathfrak{p}}$ . Second, view f as a mapping  $X \to \coprod_{\mathfrak{p} \in X} k_{\mathfrak{p}}$  given by

$$f \mapsto \frac{f}{1} + \mathfrak{p}^e \in k_{\mathfrak{p}} = A_{\mathfrak{p}/\mathfrak{p}^e}.$$

We call  $k_{\mathfrak{p}}$  the residue field of  $A_{\mathfrak{p}}$ .

**Example 3.0.3.** Suppose that k is a field and  $A = k[x_1, \ldots, x_n]$  such that for any  $\mathfrak{m}$ ,  $k_{\mathfrak{m}} = k$ . Then  $f \in A$  induces a function (prime ideals in A)  $\to k$  given by  $(x_1 - a_1, \ldots, x_n - a_n) \mapsto f(a_1, \ldots, a_n)$ .

**Lemma 3.0.4.** X is quasi-compact, meaning that for any Zariski-open  $U \subset X$  and any open cover  $\{U_{\alpha}\}$  of X, there is some finite subcover and  $U = \bigcup_{\alpha} U_{\alpha}$ .

**Note 3.0.5.** X is *not* Hausdorff in general.

**Exercise 3.0.6.** Let  $A = \mathbb{C}[x]$ . Show that  $X = \operatorname{Spec}(A)$  is not Hausdorff.

# 3.1 Lecture 8

#### Note 3.1.1.

- 1. We have that  $V(S) = V(\mathfrak{a})$  whenever  $\mathfrak{a} = \langle S \rangle$ .
- 2. The Zariski topology is generated by the collection of principal open subsets on X, i.e., subsets of the form  $X_f = \{ \mathfrak{b} \in X \mid f \notin \mathfrak{b} \}$  where  $f \in A$ . The elements in the ring A may be viewed as kinds of functions on X. View  $f \in A$  as a function  $X \to \coprod_{\mathfrak{b} \in X} A_{\mathfrak{b}} \to \coprod_{\mathfrak{b} \in X} A + \mathfrak{b}/\mathfrak{b}^c$  defined by  $\mathfrak{b} \mapsto \frac{f}{1} \in A_{\mathfrak{b}}$ .

If k is a field and  $A = k[x_1, \ldots, x_n]$ , then  $V(a_1, \ldots, a_n) \in A^n$ . We get a maximal ideal

$$\langle x_1 - a_1, x_2 - a_2, \dots, x_n - a_n \rangle$$
.

Thus, if  $f(x) \in A$  and we restrict this function, then we get the evaluation of f on points  $a \in A^n$ .

$$a \xrightarrow{\in} \langle x_1 - a_1, x_2 - a_2, \dots, x_n - a_n \rangle X \xrightarrow{} \coprod_{\mathfrak{b}} k_{\mathfrak{b}}$$

$$\uparrow \qquad \qquad \uparrow \qquad \qquad \uparrow$$

$$A^n \xrightarrow{} \coprod_a k \cong A^n \times k$$

This geometric way of packaging information about A and all of its prime ideals is compatible with all natural rings homomorphisms between the  $A_p$ 's. If we have a principal open, then for every  $f \in A$ , we get a ring  $A_f$ , provided that f is not nilpotent, and a functor (poset of principal open sets in X)<sup>op</sup>  $\to$  **CommRing** given by  $X_f \mapsto A_f$ . This extends to another functor (opens in X)<sup>op</sup>  $\to$  **CommRing**. Given  $f \in A$ , viewing f as a kind of function on X thus induces a compatible system of elements of all rings A(U) where  $U \in X$  is open.

**Lemma 3.1.2.** X is quasi-compact, meaning that any open cover  $\{U_{\alpha}\}$  of X admits some finite subcover.

*Proof.* Let  $X = \bigcup_{\alpha} U_{\alpha}$ . The principal opens generate the Zariski topology, so that for any  $\alpha$ , we can find a cover  $U_{\alpha} = \bigcup_{\beta} X_{f_{\alpha}^{\beta}}$  where  $f_{\alpha}^{\beta} \in A$ . Then  $X = \bigcup_{\alpha,\beta} X_{f_{\alpha}^{\beta}}$ , so that

$$\emptyset = \bigcap_{\alpha,\beta} (\underbrace{X - X_{f_{\alpha}^{\beta}}}_{V(f_{\alpha}^{\beta})}).$$

But  $\emptyset = \bigcap_{\alpha,\beta} V(f_{\alpha}^{\beta}) = V(\{f_{\alpha}^{\beta}\}_{\alpha,\beta})$ . Hence  $\langle (f_{\alpha}^{\beta})_{\alpha,\beta} \rangle$  is not contained in a ny prime ideal, so that  $\langle (f_{\alpha}^{\beta})_{\alpha,\beta} \rangle = A$ , hence  $1 \in \langle (f_{\alpha}^{\beta})_{\alpha,\beta} \rangle$ . We can find a collection of elements  $\{a_{\alpha}^{\beta}\}_{\alpha,\beta}$  where  $a_{\alpha}^{\beta} \in A$  such that  $1 = \sum_{\alpha,\beta} a_{\alpha}^{\beta} f_{\alpha}^{\beta}$  and at most finitely many  $a_{\alpha}^{\beta}$  are nonzero. Thus, there is sequence  $(\alpha_{1},\beta_{1}),\ldots,(\alpha_{k},\beta_{k})$  with

$$\langle f_{\alpha_1}^{\beta_1}, \dots, f_{\alpha_k}^{\beta_k} \rangle = \langle 1 \rangle = A.$$

Hence  $V(f_{\alpha_1}^{\beta_1}) \cap \cdots \cap V(f_{\alpha_k}^{\beta_k}) = \emptyset$ , and  $X = X_{f_{\alpha_1}^{\beta_1}} \cup \cdots \cup X_{f_{\alpha_k}^{\beta_k}}$ . But  $X_{f_{\alpha_i}^{\beta_i}} \in U_{\alpha_i}$  for each  $i = 1, \dots, k$ . Therefore,  $X = U_{\alpha_1} \cup \cdots \cup U_{\alpha_k}$ .

**Example 3.1.3.** Suppose that X is a compact and Hausdorff space. Let A = C(X), the ring of complex-valued continuous functions on X. Consider Spec(A) equipped with the Zariski topology and its subset  $Max(A) := \{\mathfrak{a} \in C(X) \mid \mathfrak{a} \text{ maximal}\}$  equipped with the subspace topology.

**Claim.** The natural map  $X \xrightarrow{\varphi} \operatorname{Max}(A)$  given by  $x \mapsto \{f \in C(X) \mid f(x) = 0\}$  is a homomorphism.

Proof. Let  $\mathfrak{a}_x \coloneqq \ker(\operatorname{ev}_x : A \to \mathbb{C})$ . By Urysohn's lemma, for any two distinct points  $x,y \in X$ , there is some  $f \in A$  such that f(x) = 0 and f(y) = 1. But  $f \in \mathfrak{a}_x$  and  $f \notin \mathfrak{a}_y$ , making  $\mathfrak{a}_x \neq \mathfrak{a}_y$ . Now, suppose  $\mathfrak{a} \in \operatorname{Max}(A)$  and  $\mathfrak{a} \neq \mathfrak{a}_x$  for any  $x \in X$ . This means that for any  $x \in X$ , there is some  $f_x \in \mathfrak{a}$  such that  $f_x(x) \neq 0$ . Let  $U_x \subset U$  be an open neighborhood of  $x \in X$  such that  $f_x \mid_{U_x} \neq 0$ . Then  $X = \bigcup_{x \in X} U_x$ , so that there is some finite subcover  $U_{x_1}, \ldots, U_{x_k}$  of X. Let  $f = \sum_{i=1}^k |f_{x_i}|^2$ , which does not vanish at any point of X. Note that  $f = \sum_{i=1}^k f_{x_i} \cdot \bar{f}_{x_i}$ , so that  $f \in \mathfrak{a}$ . But f is nowhere vanishing, so that  $\frac{1}{f}$  is a well defined continuous function on X. Thus,  $\frac{1}{f} \in A$ , and  $1 \in \mathfrak{a}$ , contrary to the fact that  $\mathfrak{a}$  is maximal.

**Exercise 3.1.4.** Check that  $\varphi$  is continuous, hence a homeomorphism.

Let A be a commutative ring and M an A-module. Then M defines a subset of  $X := \operatorname{\mathsf{Spec}}(A)$ , namely

$$\operatorname{supp}(M) \equiv \{ \mathfrak{b} \in X \mid M_{\mathfrak{b}} \neq 0 \},\$$

called the support of M.

# Proposition 3.1.5.

1.  $\operatorname{supp}(M) \subset V(\operatorname{ann}(M))$  where  $\operatorname{ann}(M) \equiv \{a \in A \mid a \cdot m = 0 \text{ for each } m \in M\}$ .

Proof. Let  $\mathfrak{b} \in \operatorname{supp}(M)$ . Then  $M_{\mathfrak{b}} \neq (0)$ . We need to show that  $\operatorname{ann}(M) \subset \mathfrak{b}$ . Suppose that there is some  $a \in \operatorname{ann}(M)$  with  $a \notin \mathfrak{b}$ . Let  $x \in M_{\mathfrak{b}}$ . Then  $x = \frac{m}{s}$  where  $m \in M$  and  $s \notin \mathfrak{b}$ . We compute  $\frac{a}{1} \cdot \frac{m}{s} = \frac{am}{s} = 0$  in  $M_{\mathfrak{b}}$ . Since  $a \notin \mathfrak{b}$ , it follows that  $\frac{a}{1}$  is invertible in  $A_{\mathfrak{b}}$ , i.e,  $\frac{1}{a} \in A_{\mathfrak{b}}$ . Hence  $\frac{m}{s} = \frac{1}{a} \left( \frac{a}{1} \frac{m}{s} \right) = 0$  in  $M_{\mathfrak{b}}$ , so that  $M_{\mathfrak{b}} = (0)$ , a contradiction.

2. If M is finitely generated, then  $supp(M) \supset V(ann(M))$ .

Proof. Let  $\mathfrak{b} \in V(\operatorname{ann}(M))$  and  $\mathfrak{b} \supset \operatorname{ann}(M)$ . We want to show that  $M_{\mathfrak{b}} \neq (0)$ . Suppose to the contrary. Then for any  $m \in M_{\mathfrak{f}}$  we have that  $\frac{m}{1} = 0$  in  $M_{\mathfrak{b}}$ . This shows that there exists  $s \notin \mathfrak{b}$  such that  $s \cdot m = 0$  in M. But M is finitely generated. Let  $m_1, \ldots, m_k \in M$  be generators of  $M \setminus A$ . Then there are  $s_1, \ldots, s_k \in A \setminus \mathfrak{b}$  such that  $s_i m_i = 0$  in M for each i. Let  $s = s_1 \cdots s_k \in A \setminus \mathfrak{b}$ . Then for any  $m \in M$ , we have that  $s \cdot m = 0$ . Hence  $s \in A \setminus \mathfrak{b}$ , and  $s \in \operatorname{ann}(M)$ , a contradiction.

# 3.2 Lecture 9

**Proposition 3.2.1.**  $M = (0) \iff \operatorname{supp}(M) = \emptyset \iff \operatorname{supp}(M) \cap \operatorname{Max}(A) = \emptyset.$ 

*Proof.* It's clear that  $M = \emptyset \implies \operatorname{supp}(M) = (0) \implies \operatorname{supp}(M) \cap \operatorname{Max}(A) = \emptyset$ . Hence it suffices to show that

$$\operatorname{supp}(M) \cap \operatorname{Max}(A) = \emptyset \implies M = (0).$$

On the one hand, if M is finitely generated, then  $\operatorname{supp}(M) = V(\operatorname{ann}(M))$ , so that  $\operatorname{supp}(M)$  must contain any maximal idea that contains  $\operatorname{ann}(M) \leq A$ . Thus, the assumption that  $\operatorname{supp}(M) \cap \operatorname{Max}(A) = \emptyset$  implies that  $\operatorname{ann}(A)$  is not contained in any maximal ideal, meaning that  $\operatorname{ann}(M) = A$ . This means that M = (0).

On the other hand, if M is arbitrary, then  $M = \operatorname{colim}_{\alpha} N_{\alpha}$  with each  $N_{\alpha} \subset M$  finitely generated. But then  $M_{\mathfrak{a}} = \operatorname{colim}_{\alpha}(N_{\alpha})_{\mathfrak{a}}$  because localization is exact. Since each  $N_{\alpha} = (0)$ , it follows that  $\operatorname{colim}_{\alpha} N_{\alpha} = 0$  as well.

Corollary 3.2.2. If we have a sequence of modules

$$\eta: M \stackrel{f}{\longrightarrow} T \stackrel{g}{\longrightarrow} N,$$

then  $\eta$  is exact at  $T \iff \eta_{\mathfrak{p}}$  is exact at  $T_{\mathfrak{p}}$  for each  $\mathfrak{p} \in \operatorname{Spec}(A) \iff \eta_{\mathfrak{a}}$  is exact at  $T_{\mathfrak{a}}$  for each  $\mathfrak{a} \in \operatorname{Max}(A)_{\dot{\mathcal{C}}}$ 

*Proof.* All of the forward directions are clear. Conversely, if  $\eta_{\mathfrak{a}}$  is for every  $\mathfrak{a}$ , then  $M_{\mathfrak{a}} \xrightarrow{f_{\mathfrak{a}}} T_{\mathfrak{a}} \xrightarrow{g_{\mathfrak{a}}} N_{\mathfrak{a}}$  is exact. If  $H = \ker g_{\text{im } f}$ , then  $H_{\mathfrak{a}} = \ker g_{\mathfrak{a}}/_{\text{im } f_{\mathfrak{a}}} = 0$ . Thus,  $\operatorname{supp}(H) \cap \operatorname{Max}(A) = \emptyset$ , so that H = 0.

**Definition 3.2.3.** Suppose that  $\Pi$  is a property of A-modules or of morphisms of A-modules. We say that  $\Pi$  holds locally for A if  $\Pi_{\mathfrak{a}}$  holds for every  $\mathfrak{a} \in \mathsf{Spec}(A)$ .

#### Example 3.2.4.

- 1. M = (0) holds locally if and only if it holds globally.
- 2.  $M \to T \to N$  is exact locally if and only if it's exact globally.

# **Lemma 3.2.5.** *TFAE*.

- (a) M is flat over A.
- (b) M is locally flat over A.
- (c)  $M_{\mathfrak{a}}$  is flat over  $A_{\mathfrak{a}}$  for every  $\mathfrak{a} \in \operatorname{Max}(A)$ .
- (d)  $M_{\mathfrak{a}}$  is flat over A for every  $\mathfrak{a} \in \operatorname{Max}(A)$ .

*Proof.* The fact that (a)  $\Longrightarrow$  (b)  $\Longrightarrow$  (c)  $\Longrightarrow$  (d) is obvious. To see that (c)  $\Longrightarrow$  (a), suppose that M is an A-module such that  $M_{\mathfrak{a}}$  is flat as an  $A_{\mathfrak{a}}$ -module. Suppose that  $0 \to X \to Y$  is an exact sequence of A-modules. Let  $K = \ker(X \otimes_A M \to Y \otimes_A M)$ . We want to show that K = 0.

Localizing  $0 \to K \to X \otimes_A M \to Y \otimes_A M$  along  $\mathfrak a$  gives an exact sequence  $0 \to K_{\mathfrak a} \to X_{\mathfrak a} \otimes_{A_{\mathfrak a}} M_{\mathfrak a} \to Y_{\mathfrak a} \otimes_{A_{\mathfrak a}} M_{\mathfrak a}$ , where we have used the fact that  $(X \otimes_A M)_{\mathfrak a} = X \otimes_A M_{\mathfrak a} = X_{\mathfrak a} \otimes_{A_{\mathfrak a}} M_{\mathfrak a}$ . But  $M_{\mathfrak a}$  is flat over  $A_{\mathfrak a}$ . Hence if we tensor the exact sequence  $0 \to X_{\mathfrak a} \to Y_{\mathfrak a}$  with  $M_{\mathfrak a}$  over  $A_{\mathfrak a}$ , then it will remain exact. This implies that  $\ker(X_{\mathfrak a} \otimes_{A_{\mathfrak a}} M_{\mathfrak a} \to Y_{\mathfrak a} \otimes_{A_{\mathfrak a}} M_{\mathfrak a}) = 0$ , so that  $K_{\mathfrak a} = 0$  for each  $\mathfrak a$ . It follows that  $\sup(K) = \emptyset$ , which implies that K = (0).

**Definition 3.2.6.** If A is commutative ring, then the Jacobson radical of A is the ideal

$$\operatorname{Jac}(A) \equiv \bigcap_{\mathfrak{a} \in \operatorname{Max}(A)} \mathfrak{a}.$$

**Lemma 3.2.7 (Nakayama).** If A is a commutative ring and M is a finitely generated A-module with  $Jac(A) \cdot M = M$ , then M = (0).

*Proof.* Let M be finitely generated over A. Choose some finite set of generators  $m_1, \ldots, m_t$  of M of minimal cardinality. If  $M \neq (0)$ , then t > 0. Then  $m_t \in M = \operatorname{Jac}(A) \cdot M$ . Thus there are  $a_1, \ldots, a_t \in \operatorname{Jac}(A)$  such that  $m_t = \sum_{i=1}^t a_i m_i$ . Then

$$(1 - a_t) m_t = \sum_{i=1}^{t-1} a_i m_i.$$

But  $a_t \in \operatorname{Jac}(A)$ , meaning that  $m_t$  belongs to every maximal ideal. Then  $1 - a_t$  cannot be in any maximal ideal. Hence  $1 - a_t$  is a unit in A. Let  $u \in A$  such that  $u(1 - a_t) = 1$ . Then  $m_t = \sum_{i=1}^{t-1} a_i u m_i$ . This contradicts that t is minimal.

Corollary 3.2.8 (Classical Nakayama). Suppose A is a local ring with maximal ideal  $\mathfrak{a}_A$ . Let M be a finitely generated A-module such that  $\mathfrak{a}_A M = M$ . Then M = (0).

#### Proposition 3.2.9.

- 1. If A is a commutative ring, then the functor  $(-) \otimes {}^{A}/_{Jac(A)} : A-\mathbf{Mod^{fg}} \to {}^{A}/_{Jac(A)} \mathbf{Mod^{fg}}$  is faithful.
- 2. If M is a finitely generated A-module and  $m_1, \ldots, m_t \in M$  are such that their images  $\bar{m}_1, \ldots, \bar{m}_t \in M$ / $\operatorname{Jac}(A) \cdot M$  generate the module M/ $\operatorname{Jac}(A) \cdot M$ , then they generate M.

*Proof.* If  $N = \langle m_1, \dots, m_t \rangle \subset M$ , then  $\overline{M/N} = (0)$  since  $\overline{M/N} = \overline{M}/\overline{N}$ . But then  $M \setminus N = 0$  by Lemma 3.2.7.

**Proposition 3.2.10.** If A is a local ring and t is the minimal number of generators of a finitely generated A-module M, then every generating set for M contains a generating set of t elements.

Proof. Let  $m_1, \ldots, m_k$  be a generating set for M. Then  $\bar{m}_1, \ldots, \bar{m}_k$  generate  $M_{\mathfrak{q}_A M} =$  (finite dimensional vector space over  $k_A = A_{\mathfrak{q}_A}$ ). This must have dimension t since every spacing subset in  $M_{\mathfrak{q}_A M}$  lifts to a spanning subset of M. Choose a linearly independent subset in  $\{\bar{m}_1, \ldots, \bar{m}_k\}$  and lift this to M.

**Theorem 3.2.11.** Let A be a local ring and M an A-module. Assume that one of the following conditions holds.

- (a) A is Noetherian with M finitely generated.
- (b) M is finitely presentable.

Then M is free  $\iff$  M is projective  $\iff$  M is flat.

*Proof.* We only need to show that if M is flat, then M is free. Suppose that M is flat and finitely presentable. We want to show that M is free. Let  $0 \to K \to A^t \to M \to 0$  be a finite presentation where K is finitely generated. Since M being flat implies that  $(-) \otimes k_A$  is exact, we have that

$$\eta: 0 \to K \otimes_A k_A \to k_A^t \to M \otimes_A k_A \to 0$$

is exact. Indeed, if  $0 \to N' \to N \to N'' \to 0$  is a short exact sequence of A-modules and  $N^n$  is flat, then for every A-module, the sequence

$$0 \to N' \otimes M \to N \otimes M \to N'' \otimes M \to 0$$

is exact. To see this, choose a presentation  $0 \to K \to F \to M \to 0$ , where F is free. Then we get a commutative diagram

$$0 \longrightarrow K \longrightarrow N''$$

$$\downarrow \qquad \qquad \downarrow \qquad \qquad \downarrow$$

$$K \otimes N' \longrightarrow K \otimes N \longrightarrow K \otimes N''$$

$$\downarrow \delta_1 \qquad \qquad \downarrow \delta_2 \qquad \qquad \downarrow \delta_3$$

$$0 \longrightarrow F \otimes N' \longrightarrow F \otimes N \longrightarrow F \otimes N''$$

$$\downarrow \qquad \qquad \downarrow \qquad \qquad \downarrow$$

$$M \otimes N' \stackrel{\theta}{\longrightarrow} M \otimes N \longrightarrow M \otimes N''$$

$$\downarrow \qquad \qquad \downarrow \qquad \qquad \downarrow$$

$$0 \longrightarrow 0 \longrightarrow 0$$

Apply the snake lemma (Lemma 11.2.6 below) to the first two rows.

Returning to  $\eta$ , note that  $k_A^t$  and  $M \otimes_A k_A$  are t-dimensional vector spaces over  $k_A$ . Hence  $K \otimes_A k_A = 0$ . But K is a finitely generated K module. Therefore, Lemma 3.2.7 implies that K = 0.

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# 4 Algebraic extensions

### 4.1 Lecture 10

**Definition 4.1.1.** Suppose that  $A \subset B$  where A and B are commutative rings.

- 1. We say that  $u \in B$  is algebraic over A if there is some  $f(x) \in A[x]$  such that f(u) = 0 in B and  $f \neq 0$ . We say that u is transcendental over A if it is not algebraic over A.
- 2. In general, we say that a collection of elements  $u_1, \ldots, u_k \in B$  are algebraically independent over A if there is some  $f(x_1, \ldots, x_k) \in A[x_1, \ldots, x_k]$  such that  $f \neq 0$  and  $f(u_1, \ldots, u_k) = 0$  in B. We say that  $u_1, \ldots, u_k \in B$  are independent transcendentals over A if they are not algebraically independent over A.
- 3. We say that  $B \supset A$  is algebraic if each  $u \in B$  is algebraic over A.

Our goal is to understand any algebraic extension of a ring. If A and B are domains, then we have a Cartesian diagram

$$A \hookrightarrow \operatorname{Frac}(A) 
\downarrow \qquad \qquad \downarrow 
B \hookrightarrow \operatorname{Frac}(B)$$

We have that B is an algebraic extension of A if and only if Frac(B) is an algebraic extension of Frac(A). This motivates the study of algebraic extensions of fields.

**Definition 4.1.2.** If  $L \supset K$  is a field extension, we say that L is a *finite extension* if L is finite dimensional as a vector space over K. We call  $[L:K] := \dim_K L$  the degree of the extension.

Remark 4.1.3. Finite field extensions arise naturally from polynomials.

**Definition 4.1.4.** If K is a field, then  $f(x) \in K[x]$  is called *irreducible* if deg f > 0 and f cannot be written as f = gh where  $g, h \in K[x]$  not units.

**Theorem 4.1.5.** If  $h(x) \in K[x]$  is irreducible, then the ring K[x]/(h) is a field and the inclusion  $K \subset K[x]/(h)$  is a finite field extension of degree  $\deg h$ .

*Proof.* Recall that K[x] is a Euclidean domain, in particular, a PID.

**Lemma 4.1.6.** Let A be a PID and  $u \in A$  be nonzero. TFAE.

- (a)  $A_{(u)}$  is a field.
- (b) (u) is prime.
- (c) u is simple.

*Proof.* The fact that (b) and (c) are equivalent is obvious.

Suppose that u is not simple, so that u = vw with  $v, w \in A$  not units. Then in  $A_{(u)}$  we have two elements [v] and [w] such that  $[v] \cdot [w] = [u] = [0]$ . But both [v] and [w] are nonzero since A has cancellations as a PID. Thus,  $A_{(u)}$  is not a field.

Conversely, if  $u \in A$  is simple, then for any  $x \in A \setminus (u)$  we have that (x, u) = (1) since x and u are coprime. This means that we can find  $a, b \in A$  such that ax + bu = 1. Then  $[x] \cdot [a] = [1]$ . Hence [x] is a unit, so that A/(u) is a field.

From this our theorem follows immediately.

#### Note 4.1.7.

- 1. If  $h(x) \in K[x]$  is irreducible and L = K[x]/(h), then h(x) has a natural root in L, namely, t + (h). Moreover, every element in L can be written in the form  $g(\alpha)$  for some  $g(x) \in K[x]$ .
- 2. If  $B \supset A$  is a ring extension and  $\alpha_1, \ldots, \alpha_k \in B$ , we get an intermediate ring  $A \subset A[\alpha_1, \ldots, \alpha_k] \subset B$  where  $A[\alpha_1, \ldots, \alpha_k]$  is the image of the evaluation map  $\operatorname{ev}_\alpha : f(x_1, \ldots, x_k) \mapsto f(\alpha_1, \ldots, \alpha_k)$ . Thus, if K is a field and  $h(x) \in K[x]$  is irreducible and  $\alpha = t + +(h)$ , then  $L := K[x]/(h) = K[\alpha]$ . Observe that  $\alpha$  is algebraic over K, meaning that L is generated by a single algebraic element  $\alpha$ .

**Definition 4.1.8.** We say that field extension  $L \supset K$  is *simple* if it is isomorphic to  $K[x]_{(h)}$  for some irreducible h.

#### Example 4.1.9.

- 1.  $\mathbb{C} = \mathbb{R}[i] = \mathbb{R}[x]/(x^2 + 1)$ .
- 2. If K is any field and  $a \in K$  is not a square, then  $x^2 a$  is irreducible and we get a simple field extension  $K[\sqrt{a}] := K[t]/(t^2 a)$ .

Let  $L \supset K$  be any field extension and  $u \in L$  be algebraic over K. Consider

$$\operatorname{ann}(u) \equiv \{g(x) \in K[x] \mid g(u) = 0\},\$$

which is an ideal in K[x]. Since K[x] is a PID, we see that this ideal is generated by a single element s(x). If we require that s(x) be monic, then it is uniquely determined. We call this the *minimal* polynomial of u, denoted by  $\min_{u}(x)$ .

**Lemma 4.1.10.** If  $L \supset K$  is a field extension and  $u \in L$  is algebraic over K, then  $\min_u(x)$  is irreducible and K[u] is isomorphic to the simple field extension  $K[t]/\min_u(x)$ .

Proof. If  $\min_u(x) = f(x)g(x)$ , then  $0 = \min_u(u) = f(u)g(u)$ , so that either f(u) = 0 or g(u) = 0. But  $f, g \mid \min_u$ , so that  $\deg f, \deg g \leq \deg \min_u$ . By the minimality of  $\min_u$ , this implies that  $\deg f = \deg \min_u$  or  $\deg g = \deg \min_u$ . Then either  $\deg f = 0$  or  $\deg g = 0$ . **Theorem 4.1.11.** Let  $L \supset K$  be a field extension and  $u \in L$ .

- (a) u is algebraic over K if and only if K[u] is a finite dimensional vector space over K.
- (b) If u is algebraic, then  $[K(u):K] = \deg \min_{u}$ .

Proof. We have proven (b) in Lemma 4.1.10. For (a), suppose that the ring K[u] is finite dimensional as a vector space over K. Then there exist nonnegative integers  $k_1, \ldots, k_s$  such that  $k(u) = \operatorname{span}_K(u^{k_1}, \ldots, u^{k_2})$ . Thus, if  $m > \max(k_1, \ldots, k_s)$ , then  $u^m$  is a K-linear combination of  $u^{k_1}, \ldots, u^{k_s}$ . Write  $u^m = a_1 u^{k_1} + \cdots + a_s u^{k_s}$ . Then  $f(x) = x^m - \sum_{i=1}^s a_i x^{k_i}$  satisfies f(u) = 0. Conversely, if  $u \in L$  is algebraic over K, then there is some n > 0 such that  $u^n = \operatorname{span}_K(1, u, \ldots, u^{n-1})$ . Then  $u^m \in \operatorname{span}_K(1, u, \ldots, u^{n-1})$  for any m. This implies that K[u] is finite dimensional over K.  $\square$ 

**Corollary 4.1.12.** If  $L \supset K$  is a finite field extension, then L is algebraic over K.

### 4.2 Lecture 11

**Definition 4.2.1.** A finite field extension of  $\mathbb{Q}$  is called a *number field*.

Fix a prime p > 0. Let  $\epsilon_p = e^{\frac{2\pi i}{p}} \in \mathbb{C}$ . Then  $\mathbb{Q}(\epsilon_p) \supset \mathbb{Q}$  is a finite extension because  $\epsilon_p$  is annihilated by the polynomial  $x^p - 1$ . It is called the *p-th cyclotomic field*. Note that  $x^p - 1$  is not minimal since we can factor out (x - 1). We claim that  $\frac{x^p - 1}{x - 1}$  is the minimal polynomial, so that  $[\mathbb{Q}(\epsilon_p) : \mathbb{Q}] = p - 1$ . This will hold if we can prove that  $\frac{x^p - 1}{x - 1}$  is irreducible in  $\mathbb{Q}[x]$ .

**Lemma 4.2.2 (Gauss).** If  $f(x) \in \mathbb{Z}[x]$  is irreducible, then it is irreducible in  $\mathbb{Q}[x]$ .

Proof. Note that if  $p(x) \in \mathbb{Q}[x]$ , then there exists  $N \in \mathbb{Z}_{>0}$  such that  $Np(x) \in \mathbb{Z}[x]$  and the coefficients of Np are pairwise coprime. Suppose that  $f(x) \in \mathbb{Z}[x]$  is irreducible in  $\mathbb{Z}[x]$ . Suppose, towards a contradiction, that there are  $g(x), h(x) \in \mathbb{Q}[x]$  non-units such that f(x) = g(x)h(x). Then g(x) and h(x) are  $\mathbb{Q}$ -proportional to some  $\tilde{g}(x)$  and  $\tilde{h}(x)$ , respectively, over  $\mathbb{Z}$  with each having pairwise coprime coefficients. Thus,  $f(x) = \lambda \tilde{g}(x)\tilde{h}(x)$  for some  $\lambda \in \mathbb{Q}^{\times}$ . Let  $\lambda = \frac{a}{b}$  with (a,b) = 1. If  $b \neq \pm 1$ , then there is some p > 0 where  $p \mid b$  and  $pf = a\tilde{g}\tilde{h}$ . We have that  $bf, a\tilde{g}, \tilde{h} \in \mathbb{Z}[x]$ . We can reduce p to get  $p \in [bf]_p = [a]_p[\tilde{g}]_p[\tilde{h}]_p$  But  $p \in \mathbb{Z}[x]_p$ , so that  $p \in [bf]_p = [b]_p[f]_p = 0$ . Hence  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p = 0$  in  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p$ , and  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p = 0$ . Since each of  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p = 0$ . Since each of  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p = 0$ . Since each of  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p = 0$ . Since each of  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{g}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{h}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{h}]_p[\tilde{h}]_p[\tilde{h}]_p$  and  $p \in [a]_p[\tilde{h}]_p[\tilde{h$ 

Thus, it suffices to show that  $\frac{x^p-1}{x-1}$  is irreducible in  $\mathbb{Z}[x]$ . Let  $f(x)=x^{p-1}+x^{p-2}+\cdots+x+1$ . Then  $f(x)(x-1)=x^p-1$ . By the binomial formula, we see that  $[(x-1)^p]_p=[x^p-1]_p$ . Thus,  $[f]_p[x-1]_p=[(x-1)^p]_p$ , so that  $[f]_p[(x-1)]_p=\left([x-1]_p\right)^p$  and  $[f]_p=[(x-1)]^{p-1}$  If f=gh for some non-units g and h, then  $[g]_p[h]_p=([(x-1)]_p)^{p-1}$ , which implies that  $[g]_p=[(x-1)^r]_p$  and  $[h]_p=[(x-1)^s]_p$  for some r and s. Thus,  $[g(1)]_p=[g]_p(1)=0=[h]_p(1)=[h(1)]_p$ , meaning that  $p\mid g(1)$  and  $p\mid h(1)$ . Since f=gh, it follows that  $p^2\mid f(1)=p$ , a contradiction.

**Theorem 4.2.3.** Suppose that  $M \supset L \supset K$  is a chain of finite field extensions. Then  $M \supset K$  is also finite with [M:K] = [M:L][L:K].

*Proof.* Let  $e_1, \ldots, e_n$  be a basis of L over K and  $f_1, \ldots, f_m$  be a basis of M over K. Then  $\{e_i \cdot f_j\}_{i,j}$  forms a basis of M over K.

Note 4.2.4. Suppose that  $L \supset K$  is a field extension with  $u_1, \ldots, u_n \in L$ . We get a ring  $K[u_1, \ldots, u_n] = \operatorname{im} \operatorname{ev}_u$ , which is a domain since it's contained in L. Let  $K(u_1, \ldots, u_n) := \operatorname{Frac}(K[u_1, \ldots, u_n])$ . Then we fave that  $K \subset K[u_1, \ldots, u_n] \subset K(u_1, \ldots, u_n) \subset L$ . Note that if  $u \in L$  is algebraic over K, then  $K \subset K[u] = K(u) \subset L$ .

**Theorem 4.2.5.** Suppose that  $L \supset K$  is a field extension and let  $u_1, \ldots, u_n \in L$  be algebraic over K. Then  $\dim_K K(u_1, \ldots, u_n) < \infty$ . In particular,  $K(u_1, \ldots, u_n) \supset K$  is an algebraic extension.

Proof. Note that

$$K \subset K(u_1) \subset K(u_1, u_2) \subset \cdots \subset K(u_1, \dots, u_n)$$
  
 $K(u_1, \dots, u_k) = K(u_1, \dots, u_{k-1})(u_k).$ 

Since each  $u_k$  is algebraic over K, we see that  $u_k$  is algebraic over any field containing K. Thus,  $u_k$  is algebraic over  $K(u_1, \ldots, u_{k-1})$ . Hence  $\dim_{K(u_1, \ldots, u_{k-1})} K(u_1, \ldots, u_k) < \infty$ . By Theorem 4.2.3,  $\dim_K K(u_1, \ldots, u_n) < \infty$ .

#### Definition 4.2.6.

- 1. A field K is algebraically closed if for every  $L \supset K$  and every  $u \in L$  algebraic over K, we have that  $u \in K$ .
- 2. We say that  $K \subset L$  is algebraically closed in L if any  $u \in L$  that is algebraic over K belongs to K.

**Theorem 4.2.7.** If  $L \supset K$  is a field extension, then

$$\overline{K} := \{u \in L \mid u \text{ is algebraic over } K\}$$

is a field that is algebraically closed in L.

*Proof.* Let  $u, v \in \overline{K}$ . Then both are algebraic over K. If  $K \subset K(u, v) \subset L$ , then Theorem 4.2.5 shows that  $K(u, v) \supset K$  is an algebraic extension. Since  $K(u, v) \subset \overline{K}$ , it follows that  $\overline{K}$  is a field.  $\square$ 

Suppose  $u \in L$  is algebraic over  $\overline{K}$ . Then we can find  $f(x) = \sum_{i=1}^n a_i x^i \in \overline{K}[x]$  such that  $\deg f > 0$  and f(u) = 0. Hence  $f \in K(a_1, \ldots, a_n)[x]$ , so that u is algebraic over  $K(a_1, \ldots, a_n)$ . Hence  $K(a_1, \ldots, a_n, u)$  is finite dimensional over  $K(a_1, \ldots, a_n)$ . But  $a_1, \ldots, a_n \in K$  are algebraic over K, so that  $K(a_1, \ldots, a_n)$  is algebraic over K. This means that u is algebraic over K.

Let h is an irreducible polynomial over K. Write  $\widetilde{K} = {}^{K[x]}/_{(h)}$  and let  $\alpha$  denote the marked root x + (h) of h viewed as a polynomial in  $\widetilde{K}[t]$ .

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Lemma 4.2.8 (Main lemma of Galois theory). For any  $\varphi: K \to F$  field homomorphism, the natural map

$$\left\{\psi: \widetilde{K} \to F \mid \psi \upharpoonright_K = \varphi\right\} \to (\textit{distinct roots of } h^\varphi \in F[x])$$

given by  $\psi \mapsto \psi(\alpha)$  is a bijection, where  $h^{\varphi}$  denotes the polynomial obtained by applying  $\varphi$  to the coefficients of h.

*Proof.* Let  $\psi: \widetilde{K} \to F$  be a homomorphism with  $\psi \upharpoonright_K = \varphi$ . Then

$$h^{\varphi}(\psi(\alpha)) = \varphi(a_n)\psi(\alpha)^n + \varphi(a_{n-1})\psi(\alpha)^{n-1} + \dots + \varphi(a_1)\psi(\alpha) + \varphi(a_0)$$
$$= \psi(a_n)\psi(\alpha)^n + \psi(a_{n-1})\psi(\alpha)^{n-1} + \dots + \psi(a_1)\psi(\alpha) + \psi(a_0) = \psi(h(\alpha))$$
$$= 0.$$

Now, let  $\xi \in F$  be a root of  $h^{\varphi}$ . Define a homomorphism  $K[x] \to F$  by  $f(x) \mapsto f^{\varphi}(\xi)$ . Then  $h(x) \mapsto h^{\varphi}(\xi) = 0$ . Thus, this homomorphism descends to a homomorphism  $\psi : K[x]/(h) \to F$  such that  $\psi(\alpha) = \xi$ . This implies that the assignment  $\psi \mapsto \psi(\alpha)$  is surjective.

Finally, suppose that  $\tilde{\varphi}: \widetilde{K} \to F$  is any homormophism such that  $\tilde{\varphi} \upharpoonright_K = \varphi$ . Then  $\tilde{\varphi}(\alpha)$  is a root of  $h^{\varphi}$ . Let  $\psi_{\tilde{\varphi}(\alpha)}: \widetilde{K} \to F$  be the extension that we constructed. Then  $\tilde{\varphi} \upharpoonright_K = \varphi$ , and  $\psi_{\tilde{\varphi}(\alpha)\upharpoonright_K} = \varphi$ . Also, we have that  $\tilde{\varphi}(\alpha) = \xi$  and  $\psi_{\tilde{\varphi}(\alpha)}(\alpha) = \xi$ . This shows that  $\tilde{\varphi} \upharpoonright_{K(\alpha)} = \psi_{\tilde{\varphi}(\alpha)} \upharpoonright_{K(\alpha)}$ . But  $K(\alpha) = \widetilde{K}$ .  $\square$ 

# 5 Splitting fields

### 5.1 Lecture 12

**Definition 5.1.1.** If K is a field and  $f(x) \in K[x]$ , then a field extension  $L \supset K$  is a *splitting field for* f if

- (a)  $f(x) = a \prod_{i=1}^{n} (x c_i)$  with  $a, c_i \in L$  and
- (b)  $L = K(c_1, \ldots, c_n)$ .

**Theorem 5.1.2.** For every  $f(x) \in K[x]$ , a splitting field for f exists and is unique up to an isomorphism over K.

Proof. Consider the tower of fields  $K = K_0 \subset K_1 \subset K_2 \subset \cdots$  where  $K_i = K_{i-1}[\alpha_i]$  and  $\alpha_i$  is a root of an irreducible factor  $f_i$  of f over  $K_{i-1}$  with deg  $f_i > 0$ . The degree of f is fixed, but the number of irreducible factors of f strictly increases after each step. Hence this sequence of fields will stabilize at some  $K_s$ , which is thus a splitting field for f.

To prove uniqueness, suppose that  $L \supset K$  is another splitting field for f. We have  $\varphi_0 : K_0 = K \hookrightarrow L$ . By Lemma 4.2.8, we can extend  $\varphi_0$  to a homomorphism  $\varphi_1 : K_1 \to L$  provided that  $f_1^{\varphi_0}$  has a root in L. But by assumption,  $f^{\varphi_0}$  has each of its roots in L. Since  $f_1 \mid f$ , it follows that  $f_1^{\varphi_0}$  has each of its roots in L as well. This implies that  $\varphi_1 : K_1 \to L$  will extend to a map provided that  $\varphi_2 : K_2 \to L$   $f_2^{\varphi_1} = f_2^{\varphi_0}$  has some root in L. But this holds since  $f_2 \mid f$ . Continuing in this way, we

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get  $\varphi_s: K_s \to L$  such that  $f_1^{\varphi_{s-1}}$  has all of its roots in L. Thus,  $f^{\varphi_s} = f^{\varphi_0}$  has all of its roots in L. But  $\varphi_s \upharpoonright_K = \varphi_0$ , so that  $\varphi_s$  is injective. But L = K(all roots of f). By construction, all roots of f belong to im  $\varphi_s$ . Also,  $K \subset \operatorname{im} \varphi_s$ . Hence  $\varphi_s$  is surjective and thus an isomorphism.

Exercise 5.1.3. Describe all splitting fields of polynomials of degree 2.

**Example 5.1.4.** Suppose that K is a field of characteristic  $\neq 2$ . Let  $f(x) = x^3 + a_2x^2 + a_1x + a_0 \in K[x]$ . Let L be a splitting field for f. What can L be? This depends on the splitting behavior of f over K

- (A) Suppose that f has all of its roots in K. Then L = K, and [L : K] = 1.
- (B) Suppose that f has exactly one root in K Then  $f(x) = (x \alpha) g(x)$  with  $\alpha \in K$  and g(x) a quadratic irreducible in K[x]. Consider  $L = {K[x] / (g)}$ . Then [L:K] = 2, and g has a root in L. This implies that g has all of its roots in L. Hence L is the splitting field for f.
- (C) Suppose that f has no roots in K. Then f is irreducible in K[x]. Let  $K_1 = K[x]/(f)$ , which is a simple extension of degree 3. Note that f has a root  $\alpha_1$  in  $K_1$ . Thus,  $K_1 = K[\alpha_1]$ . Consider  $f(x) = (x \alpha_1) g(x)$  with  $g \in K_1[x]$  and  $\deg g = 2$ . There are two sub-cases to consider.
  - (a) Suppose that g has two roots in  $K_1$ . Then  $L = K_1$ , so that [L : K] = 3.
  - (b) Suppose that f is irreducible in  $K_1$ . Then  $L = K_2 = K_1[x]/(g)$ , so that [L:K] = 6.

We conclude that if L is the splitting field for f, then  $[L:K] \in \{1,2,3,5\}$ .

How can we compute [L:K] from the coefficients of f? We have that  $f(x) = (x - \alpha_1)(x - \alpha_2)(x - \alpha_3)$  in L[x]. Look at  $Discr(f) := (\alpha_1 - \alpha_2)^2(\alpha_1 - \alpha_3)^2(\alpha_2 - \alpha_3)^2 \in L$ . This is a symmetric function in  $\alpha_1, \alpha_2, \alpha_3$ . Hence it is expressible in terms of  $a_2, a_1, a_0$ . Note that

$$Discr(f) = a_1^2 a_2^2 - 4a_2^2 a_0 - 4a_1^3 + 18a_0 a_1 a_2 - 27a_0^2.$$

**Proposition 5.1.5.** Suppose that f has no roots in K. Then  $[L:K] = 3 \iff \operatorname{Discr} f \in K^2$ .

Proof. We know that f is irreducible over K. Hence  $K_1 = K[x]/(f)$  is an extension of degree 3 in which f has a root  $\alpha_1$ . Note that  $\operatorname{Discr}(f) \notin K^2 \iff \operatorname{Discr}(f) \notin K_1^2$ . The  $(\iff)$  direction is obvious. For the reverse direction, suppose, towards a contradiction, that  $\operatorname{Discr}(f) \notin K^2$  but  $\operatorname{Discr}(f) \in K_1^2$ . This means that  $\left[K\left[\sqrt{\operatorname{Discr}(f)}\right]:K\right] = 2$  and  $K \subset K\left[\sqrt{\operatorname{Discr}(f)}\right] \subset K_1$ . Thus,  $3 = [K_1:K] = \left[K\left[\sqrt{\operatorname{Discr}(f)}\right]:K\right] \cdot \left[K_1:K\left[\sqrt{\operatorname{Discr}(f)}\right]\right] = 2 \cdot 1$ , a contradiction.

Now,  $\operatorname{Discr}(f) \in K_1^2 \iff (\alpha_1 - \alpha_2)(\alpha_1 - \alpha_3)(\alpha_2 - \alpha_3) \in K_1$ . This implies that  $\alpha_2 - \alpha_3 \in K_1$ . Indeed,  $f \in K_1[x]$  satisfies  $f(x) = (x - \alpha_1) g(x)$ , and  $\alpha_2, \alpha_3 \in L$  are roots of g. Therefore, we have in L that  $g(\alpha_1) = (\alpha_1 - \alpha_2) (\alpha_1 - \alpha_3) \in L$ . But  $g \in K_1[x]$  and  $\alpha_1 \in K_1$ , so that  $g(\alpha_1) \in K_1$ . It follows that  $\alpha_2 - \alpha_3 \in K_1$ . Hence  $\operatorname{Discr}(f) \in K_1^2 \iff \alpha_2 - \alpha_1 \in K_1$ . But  $-\alpha_2 - \alpha_3$  is a coefficient of g in  $K_1$ . Therefore,  $\alpha_2, \alpha_3 \in K_1$ .

#### Note 5.1.6.

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- 1. Suppose that K is a finite field. Then  $\operatorname{char} K > 0$ .
- 2. Suppose that K is any field such that  $\operatorname{char} K = p > 0$ . Then the natural map  $\varphi : K \to K$  given by  $x \mapsto x^p$  respects addition due to the binomial theorem. Hence it is a field homomorphism, called the *Frobenius morphism*. If K is finite, then this map is an automorphism. In general, im  $\varphi = K^p \subset K$  is a subfield.

3. If K has characteristic p, then the natural map  $\mathbb{F}_p \to K$  given by  $[n] \mapsto \underbrace{1+1+\cdots+1}_{n \text{ times}}$  is a field extension. Therefore, if K is finite, then  $K \supset \mathbb{F}_p$  is a finite field extension. In this case, if K has degree n, then  $K \cong \mathbb{F}_p^{\oplus n}$  is a vector space over  $\mathbb{F}_p$ . Hence  $|K| = |\mathbb{F}_p|^n = p^n$ .

# 6 Finite fields

**Theorem 6.0.1.** For every prime p and integer n > 0, there is some finite field K consisting of  $p^n$  elements that is unique up to an isomorphism over  $\mathbb{F}_p$ .

Proof. We first prove uniqueness. If F is a finite field with  $q := p^n$  elements, then  $|F^{\times}| = q - 1$ . It follows that for any  $a \in F^{\times}$ ,  $a^{q-1} = 1$ . But then for any  $a \in F$ ,  $a^q = a$ , so that each element of F is a root of  $x^q - x \in \mathbb{F}_p[x]$ . Then  $\prod_{a \in F} (x - a) \mid x^q - x$  in F[x]. This implies that  $x^q - x = \prod_{a \in F} (x - a)$  in F[x]. This means that F is a splitting field for  $x^q - x$  over  $\mathbb{F}_p$ , which must be unique up to isomorphism.

To prove existence, consider F the splitting field for  $x^q - x$  over  $\mathbb{F}_p$ . We want to show that |F| = q.

Note 6.0.2. If A is any commutative ring, then A[x] has a natural derivation. There exists a unique map  $\frac{d}{dx}:A[x]\to A[x]$  such that  $\frac{d}{dx}(a)=0$  for any  $a\in A$ ,  $\frac{d}{dx}(x)=1$ , and  $\frac{d}{dx}$  satisfies the Leibniz rule, i.e.,  $\frac{d}{dx}(fg)=\frac{df}{dx}g+f\frac{dg}{dx}$ . Note that  $\frac{d}{dx}$  is given by  $\frac{d}{dx}(a_nx^n+a_{n-1}x^{n-1}+\cdots+a_0)=na_n+(n-1)a^{n-1}x^{n-2}+\cdots+0$ . Then  $\frac{d}{dx}$  is an A-module homomorphism. If  $A\subset B$  is a subring, then we get compatible derivations  $\frac{d}{dx}\hookrightarrow A[x]\subset B[x]\hookrightarrow \frac{d}{dx}$ .

Returning to our proof, consider  $f(x) = x^q - x$ . Then since  $F \supset \mathbb{F}_p$  is the splitting field for f(x), it follows that  $f(x) = \prod_{i=1}^q (x-c_i)$  where  $c_i \in F$ . How many distinct roots does f(x) have in F? If f(x) has a repeated root, then we can write  $f(x) = (x-c)^2 g(x)$  in F[x]. This implies that  $\frac{df}{dx}(x) = 2(x-c)g + (x-c)^2 \frac{dg}{dx}$  will also have c has a root. But  $\frac{df}{dx} = qx^{q-1} - 1 = -1$  in  $\mathbb{F}_p[x] \subset F[x]$ . But in this case  $\frac{df}{dx}$  has no roots. Thus, f(x) has no repeated roots in F, so that  $|F| \geq q$ .

Now consider  $R_f := \{c \in F \mid f(c) = 0\}$ . Note that  $\mathbb{F}_p \subset R_f \subset F$  and that  $R_f = \{c \in F \mid \varphi^n(c) = c\}$  where  $\varphi$  denotes the Frobenius map. But since  $\varphi$  is a field automorphism of F, so is  $\varphi^n$ . Hence the fixed points of  $\varphi^n$  form a subfield. This means that  $R_f$  is a subfield, hence a splitting field for f. Thus,  $R_f \cong F$ .

#### 6.1 Lecture 13

We write  $\mathbb{F}_q$  for the splitting field for  $x^q - x \in \mathbb{F}_p[x]$ .

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**Proposition 6.1.1.** The group  $\mathbb{F}_q^{\times}$  is a cyclic group of order q-1.

*Proof.* By the structure theorem for finite abelian groups, we get

$$\mathbb{F}_q^{\times} \cong \mathbb{Z}/p_1^{m_{11}} \times \mathbb{Z}/p_1^{m_{12}} \times \cdots \times \mathbb{Z}/p_1^{m_{1k}} \times \mathbb{Z}/p_2^{m_{21}} \times \mathbb{Z}/p_2^{m_{22}} \times \cdots \times \mathbb{Z}/p_s^{m_{sk}} \times \mathbb{Z}/p_s^{m_{sk}} \times \mathbb{Z}/p_s^{m_{sk}} \times \mathbb{Z}/p_s^{m_{sk}} \times \mathbb{Z}/p_s^{m_{sk}}$$

Let  $\alpha_i = p_1^{m_{1i}} p_2^{m_{2i}} \cdots p_s^{m_{si}}$  for each i = 1, ..., k. Hence  $|\mathbb{F}_q^{\times}| |d_1 d_2 \cdots d_k$  where  $d_1 |d_2 | \cdots |d_k$ . Hence every element in  $\mathbb{F}_q^{\times}$  has order dividing  $d_k$ . For any  $a \in \mathbb{F}_q$ ,  $a^{d_k+1} = a$ , so that  $|\mathbb{F}_q| = \deg x^{d_k+1} - x = d_k + 1$ . Then  $q \leq d_k + 1$ , so that  $q - 1 \leq d_k$ . Since  $d_k | q - 1$ , we have that  $d_k = q - 1$ , and thus  $d_1 = d_2 = \cdots = d_{k-1} = 1$ . Hence

$$\mathbb{F}_q^{\times} \cong \mathbb{Z}/p_1^{m_{1k}} \times \mathbb{Z}/p_2^{m_{2k}} \times \cdots \times \mathbb{Z}/p_s^{m_{sk}}.$$

Since the  $p_i^{m_{ik}}$  are pairwise coprime, it follows that  $\mathbb{F}_q^{\times} \cong \mathbb{Z}/p_1^{m_{1k}}p_2^{m_{2k}}\cdots p_s^{m_{sk}}$ .

Corollary 6.1.2.  $\mathbb{F}_q = \mathbb{F}_p(\sigma)$ .

*Proof.* Since  $\mathbb{F}_q^{\times}$  is cyclic, we know that  $\mathbb{F}_q^{\times} = \langle \sigma \rangle$ .

**Proposition 6.1.3.** Aut( $\mathbb{F}_q$ ) is a cyclic group of order n. In fact, Aut( $\mathbb{F}_q$ )  $\cong \langle \varphi \rangle$  where  $\varphi$  denotes the Frobenius map.

*Proof.* We have arranged it so that  $\mathbb{F}_q$  is unique up to isomorphism over  $\mathbb{F}_p$ , so that each  $\psi \in \operatorname{Aut}(\mathbb{F}_q)$  restricts to the identity on  $\mathbb{F}_p \subset \mathbb{F}_q$ . This implies that  $\operatorname{Aut}(\mathbb{F}_q) \cong \operatorname{Aut}(\mathbb{F}_q/\mathbb{F}_p)$ , which is the subgroup of all  $\psi : \mathbb{F}_q \xrightarrow{\cong} \mathbb{F}_q$  such that  $\psi \upharpoonright_{\mathbb{F}_p} = \operatorname{id}_{\mathbb{F}_p}$ .

**Lemma 6.1.4.** Let  $L \supset K$  be a finite field extension of degree n. Then we have that  $|\operatorname{Aut}(L/K)| \leq n$ .

Proof. Since  $[L:K]=n<\infty$ , we can construct L as a tower  $K=K_0\subset K_1\subset\cdots\subset K_s=L$  where  $K_{i+1}=K_i[\alpha_i]$  and  $\alpha_i$  is a root of an irreducible  $f_i(x)\in K_i[x]$ . Consider  $\varphi_0:K\hookleftarrow L$  the natural inclusion. Applying Lemma 4.2.8, we see that  $\varphi_0$  extends to  $\varphi_1:K_1\to L$  in finitely many ways such that the number of such  $\varphi_1$ 's equals the number of distinct roots of  $f_0^{\varphi_0}$  in L. This quantity is  $\leq \deg f_0=[K_1:K_0]$  Each  $\varphi_1$  extends to a map  $\varphi_2:K_2\to L$  in at most  $\deg f_1=[K_2:K_1]$  ways. Therefore,  $\varphi_0$  will extend to a map  $\varphi_s:L\to L$  in  $[K_1:K_0][K_2:K_1]\cdots [K_s:K_{s-1}]$  many ways. It follows that

$$|\operatorname{Aut}(L/K)| \le \prod_{i=0}^{s-1} [K_{i+1} : K_i] = [L : K] = n.$$

**Corollary 6.1.5.** If  $f(x) \in K[x]$  and L is a splitting field for f and f has distinct roots in L, then  $|\operatorname{Aut}(L/K)| = [L:K]$ .

We have that  $|\operatorname{Aut}(\mathbb{F}_q/\mathbb{F}_p)| = [\mathbb{F}_q : \mathbb{F}_p] = n$ . If  $\varphi \in \operatorname{Aut}(\mathbb{F}_q/\mathbb{F}_p)$ , then  $\varphi^n = \operatorname{id}$ . Thus, it suffices to show that  $\varphi^m \neq \operatorname{id}$  for any m < n. Suppose that m has  $\varphi^m = \operatorname{id}$ . Then  $\varphi^m(a) = a$  for every  $a \in \mathbb{F}_q$ . Therefore,  $q^{p^m} = a$  for each  $a \in \mathbb{F}_q$ , so that  $p^n = q = |\mathbb{F}_q| \leq p^m$ . Then  $m \geq n$ .

This completes our main proof.

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**Proposition 6.1.6.** There is a bijection (subfields of  $\mathbb{F}_q$ )  $\cong_{\mathbf{Set}}$  (subgroups of  $\mathrm{Aut}(\mathbb{F}_q)$ ).

Proof. Let  $F \subset \mathbb{F}_q$  be a subfield, so that  $\mathbb{F}_p \subset F \subset \mathbb{F}_q$ . We have that  $|F| \mid |\mathbb{F}_q| = p^n$ , so that  $|F| \leq p^d$  for some  $d \leq n$ . Since  $\mathbb{F}_q \supset F$ , we have that  $\mathbb{F}_q$  is a vector space over F. If  $[\mathbb{F}_q : F] = r$ , then  $\mathbb{F}_q \cong F^{\oplus r}$  as F-vector spaces. Note that

$$p^{n} = |\mathbb{F}_{q}| = |F|^{r} = (p^{d})^{r} = p^{dr},$$

which implies that  $d \mid n$ .

Since F is a finite field, it follows that  $F^{\times}$  is cyclic of order  $p^d-1$ . Hence any  $a \in F \subset \mathbb{F}_q$  satisfies  $a^{p^d}=a$ . But if  $d\mid n$ , then  $x^q-x=x^{p^n}-x=\left(x^{p^d}-x\right)g(x)$  because  $p^n-1=p^{dr}-1=\left(p^d\right)^r-1=\left(p^d-1\right)m$  so that  $x^{q-1}-1=\left(x^{p^d-1}-1\right)g(x)$ . But  $\mathbb{F}_q$  is the splitting field for  $x^q-x$ , and all roots of this are distinct. Thus, there are exactly  $p^d$  roots of  $x^q-x$  that are the distinct roots of  $x^{p^d}-x$ . Therefore,

$$F = \mathbb{F}_{p^d} = \left( \text{subfield of } \mathbb{F}_q \text{ that is the splitting field for } x^{p^d} - x \right) = \left( \text{fixed subfield of } \varphi^d \right).$$

Hence F is the fixed point subgroup of  $\langle \varphi^d \rangle \subseteq \operatorname{Aut}(\mathbb{F}_q)$ .

Let  $\psi \in \operatorname{Aut}(\mathbb{F}_q)$  with  $\psi \notin \langle \varphi^d \rangle$ . Then  $\psi = \varphi^e$  for some  $e \geq 0$  such that  $d \nmid e$ . If  $\xi$  generates  $F^{\times}$  and  $\xi^{p^e} = \psi(\xi) = \xi$ , then  $p^d - 1 \mid p^e - 1$  since  $|F^{\times}| = p^d - 1$ . But this is impossible, which implies that  $\psi \upharpoonright_F \neq \operatorname{id}_F$ . Therefore,  $\langle \varphi^d \rangle = \operatorname{Aut}(\mathbb{F}_q/F)$ , and we have a bijection

(subfields of 
$$\mathbb{F}_q$$
)  $\cong_{\mathbf{Set}}$  (subgroups of  $\mathrm{Aut}(\mathbb{F}_q)$ )  
 $F \mapsto \mathrm{Aut}(\mathbb{F}_q/F)$   
 $\mathbb{F}_q^G \longleftrightarrow G.$ 

# 6.2 Lecture 14

#### Proposition 6.2.1.

1. Let  $\mathbb{F}_q^{\times} = \langle \theta \rangle$ . Then  $\mathbb{F}_q = \mathbb{F}_p(\theta)$ , meaning that  $\theta$  is a primitive element for the extension  $\mathbb{F}_q \supset \mathbb{F}_p$ . Further, if h denotes the minimal polynomial of  $\theta$  over  $\mathbb{F}_p$ , then  $\mathbb{F}_q$  is the splitting field for h.

Proof. Every nonzero element of  $\mathbb{F}_q$  is a power of  $\theta$ . Hence  $\mathbb{F}_q = \mathbb{F}_p(\theta)$ . Now, note that  $\deg h = n$  because  $[\mathbb{F}_q : \mathbb{F}_p] = n$ . Write  $h(x) = x^n + a_{n-1}x^{n-1} + \cdots + a_0$  where each  $a_i \in \mathbb{F}_p$ . If we view h over  $\mathbb{F}_q$ , then  $\varphi(a_i) = a_i$  due to Fermat's little theorem. Hence  $\varphi(h(x)) = h(\varphi(x))$  for any  $x \in \mathbb{F}_q$ , meaning that  $\varphi(c)$  is a root of h whenever c is a root. Thus, we get n roots of h.

$$\theta, \ \theta^p, \ \theta^{p^2}, \dots, \theta^{p^{n-1}}$$

If K is the splitting field for h, then  $\mathbb{F}_p \subset K \subset \mathbb{F}_q$ . But  $[K : \mathbb{F}_p] = n = [\mathbb{F}_q : \mathbb{F}_p]$ , so that  $K = \mathbb{F}_q$ .

2. Let  $m \geq 0$  be any integer and  $q = p^n$ . Then there is some irreducible polynomial over  $\mathbb{F}_q$  of degree m

*Proof.* Let  $\mathbb{F}_{p^{mn}}^{\times} = \langle \theta \rangle$ . Then the minimal polynomial p(x) of  $\theta$  over  $\mathbb{F}_q$  has degree m, and p(x) is irreducible since it is minimal.

# 7 Cyclotomic fields

Let  $q = p^n$  and d > 0 be any integer. Among the finitely many polynomials over  $\mathbb{F}_q$  of degree d, how many of these are irreducible? We have just shown that at least one is irreducible.

Define the Möbius function  $\mu: \mathbb{Z}_{>0} \to \{-1, 0, 1\}$  by

$$n \mapsto \begin{cases} -1 & n = 1 \\ (-1)^k & n = p_1 \cdots p_k \text{ where the } p_i \text{ are pairwise distinct } \cdot \\ 0 & n \text{ is divisible by a square} \end{cases}$$

### Proposition 7.0.1.

- (i)  $\mu(k) \neq 0$  for some k.
- (ii)  $\mu(nm) = \mu(n)\mu(m)$  when (n, m) = 1.

(iii) 
$$\sum_{d|n} \mu(d) = \begin{cases} 1 & n=1\\ 0 & n \neq 1 \end{cases}.$$

*Proof.* Let n > 0 be an integer and write  $p_1^{k_1} \cdots p_k^{r_k}$  where the prime  $p_i$  are pairwise distinct. Let  $n_0 = p_1 \cdots p_k$ . Then  $\sum_{d|n} \mu(d) = \sum_{d|n_0} \mu(d)$ . If  $d \mid n_0$ , then  $d = p_{i_1} \cdots p_{i_s}$ , so that  $\mu(d) = (-1)^s$ . By the binomial theorem, it follows that

$$\sum_{d|n_0} \mu(d) = \sum_{s=0}^k \binom{k}{s} (-1)^2$$
$$= (1-1)^k$$
$$= \begin{cases} 1 & k=0\\ 0 & k>0 \end{cases}.$$

Therefore, 
$$\sum_{d|n} \mu(d) = \begin{cases} 1 & n=1 \\ 0 & n \neq 1 \end{cases}$$
.

**Corollary 7.0.2.** For any  $m, d \in \mathbb{Z}_{>0}$  such that  $d \mid m$ , we have that

$$\sum_{d|n|m} \mu\left(\frac{m}{n}\right) = \begin{cases} 1 & d=m\\ 0 & d \neq m \end{cases}.$$

Remark 7.0.3. Proposition 7.0.1 completely characterizes the Möbius function.

Lemma 7.0.4 (Möbius inversion formula). Let A be an abelian group and  $f, g : \mathbb{Z}_{>0} \to A$  be functions such that  $f(n) = \sum_{d|n} g(d)$  for every n. Then

$$g(n) = \sum_{d|n} \mu\left(\frac{n}{d}\right) f(d). \tag{\dagger}$$

If A is written multiplicatively, then this becomes

$$g(n) = \prod_{d|n} f(d)^{\mu(\frac{n}{d})}.$$

*Proof.* We compute

$$\sum_{d|n} \mu\left(\frac{n}{d}\right) f(d) = \sum_{d|n} \mu\left(\frac{n}{d}\right) \sum_{k|d} g(k)$$

$$= \sum_{d|n} \sum_{k|d} \mu\left(\frac{n}{d}\right) g(k)$$

$$= \sum_{k|n} g(k) \sum_{d: k|d|n} \mu\left(\frac{n}{d}\right)$$

$$= \sum_{k|n} g(k) \delta(k, n)$$

$$= g(n).$$

**Definition 7.0.5.** Define the Euler (totient) function  $\varphi : \mathbb{Z}_{>0} \to \mathbb{Z}$  by

$$\varphi(n) = \# \{ m \in \mathbb{Z}_{>0} : m \le n, (m, n) = 1 \}.$$

If  $n \in \mathbb{Z}_{>0}$ , then  $n = \sum_{d|n} \varphi(d)$ . Therefore, if  $f : \mathbb{Z}_{>0} \to \mathbb{Z}$  is given by f(n) = n and  $g := \varphi$ , then we can apply  $(\dagger)$  to get

$$\begin{split} \varphi(n) &= g(n) = \sum_{d \mid n} \mu\left(\frac{n}{d}\right) f(d) \\ &= \sum_{d \mid n} \mu\left(\frac{n}{d}\right) d \\ &= \sum_{m \mid n} \mu(m) \frac{n}{m} = \left(\sum_{m \mid n} \frac{\mu(m)}{m}\right) n. \end{split}$$

**Lemma 7.0.6.** If  $n = p_1^{r_1} \cdots p_k^{r_k}$ , then  $\varphi(n) = n \prod_{i=1}^k \left(1 - \frac{1}{p_i}\right)$ .

*Proof.* Let  $n_0 = p_1 \cdots p_k$ . Then

$$\sum_{m|n} \frac{\mu(m)}{m} = \sum_{m|n_0} \frac{\mu(m)}{m}$$

$$= \underbrace{1}_{m=1} - \sum_{i=1}^k \frac{1}{p_i}$$

$$+ \sum_{i < j} \frac{1}{p_i p_j} + \dots + (-1)^s \sum_{i_1 < \dots < i_s} \frac{1}{p_{i_1} \dots p_{i_s}}$$

$$+ \dots + (-1)^k \frac{1}{p_1 \dots p_k}$$

$$= \prod_{i=1}^k \left(1 - \frac{1}{p_i}\right).$$

Let  $p(x) = x^n - 1 \in \mathbb{Q}[x]$  with n > 0. Let  $\Gamma_n$  be the splitting field for p(x). We know that  $\Gamma_n = \mathbb{Q}(\zeta_n)$ , where  $\zeta_n$  denotes a primitive n-th root of unity in  $\mathbb{C}$ . Let the set  $\operatorname{Prim}_n$  consist of all the primitive n-th roots of unity.

We have that  $\mu_n = \{1, \zeta_n, \zeta_n^2, \dots, \zeta_n^{n-1}\} = \coprod_{d|n} \operatorname{Prim}_d$ . Define the d-th cyclotomic polynomial as

$$\Phi_d(x) = \prod_{\alpha \in \text{Prim}_d} (x - \alpha).$$

For example,

$$\begin{split} &\Phi_1 = x - 1 \\ &\Phi_2 = x + 1 \\ &\Phi_3 = x^2 + x + 1 \\ &\Phi_4 = x^2 + 1 \\ &\vdots \\ &\Phi_p = x^{p-1} + x^{p-2} + \dots + x + 1 \text{ with } p \text{ prime.} \end{split}$$

Note that

$$x^n - 1 = \prod_{d|n} \Phi_d(x).$$

Consider the functions  $\Phi_{(-)}: \mathbb{Z}_{>0} \to \mathbb{C}(x)^{\times}$  and  $f: \mathbb{Z}_{>0} \to \mathbb{C}(x)^{\times}$  where  $f(n) = x^n - 1$ . We can apply  $(\dagger)$  to get

$$\Phi_n(x) = \prod_{d|n} (x^d - 1)^{\mu(\frac{n}{d})},$$

which is a rational function over  $\mathbb{Z}$ . We can write  $\Phi_n(x) = \frac{a(x)}{b(x)}$  with b(x) monic. Write  $\Phi_n(x) = x^m + p_{m-1}X^{m-1} + \dots + p_0$  and set  $M = \text{lcm}\{c_i \mid p_i = \frac{t_i}{c_i}, i = 1, \dots, m-1\}$ . Let  $P_i = Mp_i$  for each  $i = 1, \dots, m-1$  and  $P_m = M$ . Since  $M\Phi_n(x)b(x) = Ma(x) \in \mathbb{Z}[x]$ , we see that M divides

each coefficient of  $M\Phi_n(x)b(x)$ . Suppose, towards a contradiction, that M>1. Then there exists a prime divisor p of M. By our choice of M, there exists a maximal  $0 \le i_0 \le m$  such that  $p \nmid P_{i_0}$ . If deg b(x) = s, then the coefficient of  $X^{m+s}$  in  $M\Phi_n(x)b(x)$  has the form  $M + p \cdot t$  for some  $t \in \mathbb{Z}$ . But this is not divisible by p and thus not divisible by M, a contradiction. Thus, M = 1, so that  $\Phi_n(x) \in \mathbb{Z}[x]$ .

Moreover, since deg  $\Phi_n = \varphi(n)$ , it follows that  $[\Gamma_n : \mathbb{Q}] = \varphi(n)$ .

# 7.1 Lecture 15

Let  $q = p^n$ . Let  $\psi_d(q) = \#\{p(x) \text{ irreducible over } \mathbb{F}_q \mid \deg p(x) = d\}$ . If  $f(x) \in \mathbb{F}_q[x]$  is irreducible, then  $F = \mathbb{F}_q[x]/(f)$  is a finite field. Thus,  $\#F = g^d = p^{nd}$ , so that F is the splitting field for  $x^{p^{nd}} - x$  over  $\mathbb{F}_p$ . Also, F is just the set of roots of  $x^{p^{nd}} - x$ . By construction, the polynomial  $f(x) \in \mathbb{F}_q[x]$  has a root over F, and  $x^{p^{nd}} - x \in \mathbb{F}_q[x]$  has a root in F.

Since f(x) is irreducible in  $\mathbb{F}_q[x]$ , we see that  $\left(f, x^{p^{nd}} - x\right) \in \{1, f\}$  in  $\mathbb{F}_q[x]$ . But if  $\left(f, x^{p^{nd}} - x\right) = 1$ , then  $1 = a(x)f(x) + b(x)(x^{p^{nd}} - x)$ . with  $a, b \in \mathbb{F}_q[x]$ . If we write this as an equation in F[x], then evaluating on  $\alpha \in F$  a common root of f(x) and  $x^{p^{nd}} - x$  will give us a contradiction. Hence  $f(x) \mid x^{p^{nd}} - x$  in  $\mathbb{F}_q[x]$ . Since all roots of  $x^{p^{nd}} - x$  are pairwise distinct, we see that any irreducible monic polynomial of degree d over  $\mathbb{F}_q$  appears exactly once in the decomposition of  $x^{p^{nd}} - x$  into irreducibles. Note that if m = dr, then

$$x^{q^d} - x \mid \underbrace{x^{q^m} - x}_{\text{distinct roots}},$$

and thus every irreducible monic polynomial over  $\mathbb{F}_q$  of degree dividing m appears exactly once in the irreducible decomposition of  $x^{q^m} - x$ .

For each  $d \geq 1$ , let  $f_{d,1}, f_{d,2}, \ldots, f_{d,\psi_d(q)}$  be irreducible monic polynomials over  $\mathbb{F}_q$  of degree d. Then for any  $m \geq 1$ , we get

$$x^{q^m-x} - x = \prod_{d|m} \prod_{k=1}^{\psi_d(q)} f_{d,k}(x),$$

so that  $q^m = \sum_{d|m} d\psi(q)$ . Then

$$\psi_d(q) = \frac{1}{m} \sum_{d|m} \mu\left(\frac{m}{d}\right) q^d.$$

# Example 7.1.1.

$$\psi_2(2) = \frac{1}{2}(2^2 - 2) = 1.$$
  
$$\psi_3(2) = \frac{1}{2}(2^3 - 2) = 2.$$

Remark 7.1.2. A randomly chosen polynomial over  $\mathbb{F}_q$  of degree d will be irreducible with probability  $\approx \frac{1}{d}$ . Given a polynomial p(x) over  $\mathbb{F}_q$  of degree d, there is no known algorithm that is polynomial in d for deciding whether p(x) is irreducible.

# 8 Galois theory

**Definition 8.0.1.** If  $L \supset K$  is a field extension, then the *Galois group* is

$$Gal(L/K) \equiv \{ \varphi \in Aut(L) : \varphi \upharpoonright_K = id_K \}.$$

**Theorem 8.0.2.** Let  $L \supset K$  be a field extension of degree  $n < \infty$ . Let  $G \leq \operatorname{Gal}(L/K)$ .

- (a)  $L^G = K \iff |G| = n$ .
- (b) If  $L^G = K$  and  $K \subset P \subset Q \subset L$  is a chain of field extensions, then every homomorphism  $\varphi: P \to L$  over K extends to a homomorphism  $Q \to L$  in exactly [Q: P] many ways.

Proof.

(a) For the  $(\Leftarrow)$  direction, note that if  $G \leq \operatorname{Aut}(L)$ , then tautologically  $G \leq \operatorname{Gal}(L/L^G)$ . Hence  $|G| \leq |\operatorname{Gal}(L/L^G)| = [L:L^G]$ . If  $G \leq \operatorname{Gal}(L/K)$ , then  $L \subset L^G \subset L$ , so that  $[L:L^G] \leq [L:K] = n$ . This means that  $|G| \leq n$ .

Conversely, let  $L^G = K$ . Take  $\alpha \in L$  and let  $\mathrm{Orb}_G(\alpha) = \{\lambda_1, \dots, \lambda_m\} \subset L$ . Consider

$$f(x) = \prod_{i=1}^{m} (x - \lambda_i) \in L[x].$$

But the coefficients are symmetric polynomials in  $\lambda_i$ , and any  $g \in G$  permutes the  $\lambda_i$ . In this case, g permutes the coefficients of f(x). Hence  $f(x) \in L^G[x] = K[x]$ . By construction,  $\alpha$  is a root of f(x), the minimal polynomial of  $\alpha$ . We can decompose f(x) into linear factors in L.

Apply part (b) to P = K and Q = L. In this case, (b) implies that if  $L^G = K$ , then |Gal(L/K)| = [L:K] = n. Thus, we must show that G = Gal(L/K).

Let  $\varphi \in \operatorname{Gal}(L/K)$ . Recall that f(x) is the minimal polynomial of  $\alpha$  over K. Note that  $\varphi(x)$  is a root of  $f^{\varphi}(x)$ . Indeed, since  $\varphi_K = \operatorname{id}_K$ , we have that  $f^{\varphi}(x) = f(x)$ . Hence  $\varphi(\alpha) \in \operatorname{Orb}_G(\alpha)$ , so that there exists  $g \in G$  such that  $\varphi(\alpha) = g(\alpha)$ . If L is a finite field, then we can take  $\alpha$  to be the generator of  $L^{\times}$ , in which case  $\varphi(\alpha) = g(\alpha) \Longrightarrow \varphi(\alpha^k) = g(\alpha^k)$  for each  $k \Longrightarrow \varphi = g$ . If L is infinite, then K is infinite and for any  $g \in G$ , we consider  $L_g = \{a \in L \mid \varphi(a) = g(a)\} \subset L$ . By definition,  $L_g = L^{g^{-1} \circ \varphi}$  is a subfield in L. This contains K because  $g^{-1} \circ \varphi \in \operatorname{Gal}(L/K)$ . Therefore,  $K \subset L_g \subset L$  is a field extension, meaning that  $L_g$  is a K-vector subspace in L.

We have shown that  $L = \bigcup_{g \in G} L_g$ .

**Claim.** If K is an infinite field and V is a finite-dimensional K-vector space and  $V_1, V_2, \ldots, V_g \subset V$  are subspaces, then  $V = \bigcup_{i=1}^g V_i \implies V = V_k$  for some k.

Proof. Suppose that each  $V_i \subsetneq V$  and that  $V = \bigcup_{i=1}^g V_i$ . Then there exists a linear map  $f_i: V \to K$  such that  $f_i \upharpoonright_{V_i} = 0$  and  $f_i \neq 0$ . Then  $f: V \to K$  given by  $f = \prod_{i=1}^k f_i$  is the function associates with a nonzero polynomial in  $V = K^n$  of degree s. But f is the zero function since  $V = \bigcup_{i=1}^k V_i$ , a contradiction.

(b) Suppose that  $K \subset P \subset L$ . Let  $K \subset P \subset Q \subset L$  where  $Q = P(\alpha) = P[\alpha]$  and  $\alpha$  is a root of some irreducible  $h(x) \in P[x]$ . Let  $f(x) = \prod_{i=1}^m (x - \lambda_i)$  where  $\operatorname{Orb}_G(\alpha) = \{\lambda_1, \dots, \lambda_m\} \subset L$ . Then  $f(x), g(x) \in P[x]$  have a common root, and h(x) is irreducible. Hence  $h \mid f$  in P[x]. Let  $\varphi : P \to L$  be any field homomorphism over K. Then  $h^{\varphi} \mid \underbrace{f^{\varphi}}_{f}$  in L[x]. But f decomposes into distinct linear factors in L[x]. Hence  $h^{\varphi}$  equals a product of a subcollection of these factors. It follows that  $h^{\varphi}$  has  $\deg h^{\varphi} = \deg h$  distinct roots in L. By Lemma 4.2.8, since  $Q \cong P[x]/(h)$ , we see that  $\varphi$  extends in exactly  $\deg h - [Q:P]$  many ways.

This proves our result for simple field extensions. Since every finite extension is a tower of simple extensions, we are done by induction on the length of the tower.

**Definition 8.0.3.** A finite field extension  $L \supset K$  is a Galois extension if |Gal(L/K)| = [L:K].

**Corollary 8.0.4.** If  $L \supset K$  is a Galois extension, then  $K \subset P \subset L \implies L \supset P$  is Galois as well.

*Proof.* Take Q = L and apply (b) then (a).

**Definition 8.0.5.** If K is a field and  $f(x) \in K[x]$ , then we say that f is separable over K if f has no repeated roots in any finite extension of K. Equivalently, f has no repeated roots in its splitting field.

## 8.1 Lecture 16

**Proposition 8.1.1.** A polynomial  $f(x) \in K[x]$  is separable over K if and only if (f, f') = 1.

Proof. If  $f,g \in K[x]$ , then  $(f,g) \in K[x]$ . Suppose there exists  $L \supset K$  such that f has a multiple root in L. Then there exists an irreducible polynomial  $h(x) \in L[x]$  such that  $h^2 \mid f$ . This implies that  $f = h^2q$ , so that  $f' = 2hh'q + h^2q' = h(2h'q + hq')$ . Hence  $h \mid f'$  in L[x]. Then  $h \mid (f, f')$  in L[x], making  $(f, f') \neq 1$ .

Conversely, suppose that  $(f, f') \neq 1$ . Then there exists h irreducible in K[x] such that  $h \mid f'$  and  $h \mid f$  in K[x]. We can write f = hg, so that f' = h'g + hg'. Either  $h \mid g$  or h' = 0. In the former case, we have that  $h \mid g \implies h^2 \mid f \implies f$  has a double root in L = K[x]/(h). In the latter case, we see that  $\operatorname{char} K = p > 0$  and  $h(x) = a_0 + a_1 x^p + a_2 x^{2p} + \dots + a_s x^{sp}$  with  $a_s \neq 0$ . Let  $L \supset K$  be a finite field extension such that for any  $i = 0, \dots, s$ , we have  $b_i$  such that  $b_i^p = q_i$ . Then viewing  $h(x) \in L[x]$ , we get  $h(x) = (b_0 + b_1 x + b_2 x^2 + \dots + b_s x^s)^p$  since  $b_s \neq 0$ . Hence if  $\varphi(x)$  is an irreducible factor of  $b_0 + b_1 x + \dots + b_s x^s$  in L[x], then if  $F = L[x]/(\varphi)$ , then  $\varphi$  has a root in F and h will have a root with multiplicity p in F. In this case, f has a root with multiplicity p in F.

**Corollary 8.1.2.** If K has char K = 0, then every irreducible  $f(x) \in K[x]$  is separable.

Proof. If char K = 0, then  $f \neq 0$  and f is irreducible. Since  $\deg f > 0$ , it follows that  $f' \neq 0$ . Hence (f, f') is a polynomial of degree  $\geq 0$ . Since f is irreducible, this means that  $(f, f') \in \{1, f\}$ . But  $\deg f' < \deg f$ , so that (f, f') = 1.

**Corollary 8.1.3.** If  $f(x) \in K[x]$  is irreducible and char  $K \nmid \deg f$ , then f is separable.

Corollary 8.1.4. Every irreducible polynomial f over a finite field F is separable.

Proof. If f is irreducible and  $f' \neq 0$ , then apply a similar argument to the proof of Corollary 8.0.4. Suppose f' = 0. Then  $f(x) = a_0 + a_1 x^p + \cdots + a_s x^{sp}$  with  $p = \operatorname{char} F$ . But as F is finite, we know that the Frobenius map  $\varphi$  is an automorphism. Thus, any element in F has a p-th root in F. Hence there exists  $b_i \in F$  such that  $b_i^p = a_i$ . This shows that  $f(x) = (b_0 + b_1 x + \cdots + b_s x^s)^p$ , which contradicts that f is irreducible over F.

**Example 8.1.5.** There are irreducible polynomials over fields of characteristic > 0 that are not separable. For example, let  $K = \mathbb{F}_p(t)$  and  $f(x) = x^p - t$ . This is irreducible in K[x] but not separable over K.

Indeed, if  $L \supset K$  is such that f has a root  $\alpha$  in L, then f splits in L[x]. We can write  $f(x) = (x - \alpha)^p$ . But if 0 < k < p, then  $\alpha^k \notin K$ . This shows that f is irreducible in F[x] but has a root of multiplicity p.

**Theorem 8.1.6.** If  $f(x) \in K[x]$  and every irreducible factor of f is separable over K, then the splitting field L of f is Galois over K.

*Proof.* We constructed L as a tower

$$K = K_0 \subset K_1 \subset \cdots \subset K_s = L$$

where  $K_{i+1} = K_i(\alpha_{i+1})$  and  $\alpha_{i+1}$  is a root of some irreducible factor  $f_{i+1}(x)$  of  $f(x) \in K_i[x]$ . Since  $f_{i+1}$  is irreducible in  $K_i[x]$  and  $f_{i+1} \mid f$  in  $K_i[x]$ , it follows that  $f_{i+1}$  must divide one of the irreducible factors of f(x) in K[x]. But these are separable, which implies that  $f_{i+1}$  is separable for each i. By Lemma 4.2.8, a field homomorphism  $\varphi: K \to L$  extends to an isomorphism  $\varphi: L \to L$  in

(# of distinct roots in  $f_1$ ) · (# of distinct roots in  $f_2$ ) · · · (# of distinct roots in  $f_{s-1}$ )

many ways. Note that

(# of distinct roots in 
$$f_1$$
) · (# of distinct roots in  $f_2$ ) · · · (# of distinct roots in  $f_{s-1}$ )
$$= \deg f_1 \cdot \deg f_2 \cdots \deg f_{s-1}.$$

Hence

$$|\operatorname{Gal}(L/K)| = \operatorname{deg} f_1 \cdot \operatorname{deg} f_2 \cdots \operatorname{deg} f_{s-1} = [K_1 : K_0] [K_2 : K_1] \cdots [K_s : K_{s-1}] = [L : K].$$

If  $f(x) \in K[x]$  and  $L \supset K$  is the splitting field for f, then let  $\alpha_1, \ldots, \alpha_m$  denote the distinct roots of f in L. We have that  $L = K(\alpha_1, \ldots, \alpha_m)$  and any  $\varphi \in \operatorname{Gal}(L/K)$  sends  $\{\alpha_1, \ldots, \alpha_m\}$  to itself. This gives us a homomorphism  $\operatorname{Gal}(L/K) \to S_m$  that is injective by Lemma 4.2.8. Therefore,  $\operatorname{Gal}(L/K) \subset S_m$ .

### Example 8.1.7.

1. Let K be a field and let  $f(x) \in K[x]$  be irreducible of degree 2. Let L denote the splitting field for f(x). Then  $K\left[\sqrt{D}\right]$  where  $D = \operatorname{Discr}(f) \in K$ . In this case, [L:K] = 2, and  $\operatorname{Gal}(L/K) \subset S_2$  since  $D \neq 0$ . Thus, f must have distinct roots in L. Note that  $\operatorname{Gal}(L/K) \neq \{\operatorname{id}\}$ , since these roots are not in K. This shows that  $\operatorname{Gal}(L/K) = \langle \sigma \rangle = S_2$  where  $\sigma: L \to L$  is given by  $a + b\sqrt{D} \mapsto a - b\sqrt{D}$ .

- 2. Let  $q = p^n$ . Consider the extension  $\mathbb{F}_q \supset \mathbb{F}_p$ . Then  $\operatorname{Gal}(\mathbb{F}_q/\mathbb{F}_p) = \langle \varphi \rangle \cong \mathbb{Z}/n$ .
- 3. Recall that the cyclotomic field  $\Gamma_n \supset \mathbb{Q}$  is the splitting field for

$$x^{n} - 1 = (x - 1)(x^{n-1} + x^{n-2} + \dots + x + 1).$$

This polynomial has irreducible factors each of which is separable. Thus,  $\Gamma_n \supset \mathbb{Q}$  is a Galois extension such that  $\operatorname{Gal}(\Gamma_n/\mathbb{Q}) \subset S_{\mu_n}$ . Since any  $g \in \operatorname{Gal}(\Gamma_n/\mathbb{Q})$  respects multiplication in  $\Gamma_n$ , we see that  $g \upharpoonright_{\mu_n} : \mu_n \to \mu_n$  is a group automorphism. It follows that  $\operatorname{Gal}(\Gamma_n/\mathbb{Q}) \subset \operatorname{Aut}_{\mathbf{Grp}}(\mu_n) \cong (\mathbb{Z}/n)^{\times}$ , which has order  $\phi(n)$ . We have shown that the minimal polynomial of a root of 1 over  $\mathbb{Q}$  is precisely  $\Phi_n(x)$ , where  $\operatorname{deg} \Phi_n(x) = \phi(n)$ . Hence  $[\Gamma_n : \mathbb{Q}] = \phi(n)$ , so that  $\operatorname{Gal}(\Gamma_n/\mathbb{Q}) = (\mathbb{Z}/n)^{\times}$ .

4. Suppose that char  $K \notin \{2,3\}$ . Let  $f(x) \in K[x]$  be irreducible and monic of degree 3. Let  $D \in K$  denote the discriminant of f. Let  $L \supset K$  be the splitting field for f, so that  $L \supset K$  is Galois. Then

$$|\operatorname{Gal}(L/K)| = \begin{cases} 6 & D \notin K^2 \\ 3 & D \in K^2 \end{cases}.$$

But  $Gal(L/K) \subset S_3$ . This shows that

$$\operatorname{Gal}(L/K) = \begin{cases} S_3 & D \notin K^2 \\ A_3 & D \in K^2 \end{cases}.$$

#### 8.2 Lecture 17

**Definition 8.2.1.** Let k be a field and A be a finitely generated k-algebra. A collection  $u_1, \ldots, u_n \in A$  is a transcendence basis of A/k if

- (i) the  $u_i$  are independent transcendentals over k and
- (ii) every  $a \in A$  is algebraically dependent with  $k[u_1, \ldots, u_n]$ .

If A is a domain and  $u_1, \ldots, u_n$  forms a transcendence basis of A/k, then they also form a transcendence basis of  $\operatorname{Frac}(A)$  over k. Observe that  $x \in \operatorname{Frac}(A)$  is algebraic over  $k[u_1, \ldots, u_n]$  if and only if it is algebraic over  $k(u_1, \ldots, u_n)$ . Then

$$S := \{x \in \operatorname{Frac}(A) \mid x \text{ is algebraic over } k[u_1, \dots, u_n]\}$$

is a subfield. But  $A \subset S \subset \operatorname{Frac}(A)$ , so that, by the universal property,  $S = \operatorname{Frac}(A)$ . Hence  $\operatorname{Frac}(A)$  is algebraic over  $k(u_1, \ldots, u_n)$ .

Let  $A = k[u_1, \ldots, u_n]$  and suppose that  $\{u_1, \ldots, u_d\}$  is a maximal subset of algebraically independent elements over k in  $\{u_1, \ldots, u_n\}$ . Then  $u_1, \ldots, u_d$  form a transcendence basis of A/k. Indeed, K equals the algebraic closure of  $k(u_1, \ldots, u_d)$  in  $\operatorname{Frac}(A)$ . Thus,  $u_1, \ldots, u_n \in K$ , so that  $K = \operatorname{Frac}(A)$ . It follows that  $K \supset A$ .

As a result, if A is a finitely generated algebra without zero divisions, then A has a transcendence basis over k. Indeed, choose any system of generators of A/k and then choose a maximal subset of algebraically independent elements.

**Lemma 8.2.2.** Suppose that  $\{u_1, \ldots, u_n\}$  is a transcendence basis of A/k and that v is transcendental over  $k[u_1, \ldots, u_n]$ . Then  $\{v, u_2, u_3, \ldots, u_n\}$  is also a transcendence basis of A/k.

Proof. Note that  $v, u_2, \ldots, u_n$  are algebraically independent over k whereas  $v, u_1, u_2, \ldots, u_n$  are algebraically dependent. A nontrivial algebraic relation among these will be given by a polynomial p(x) over k such that p(x) includes a monomial involving  $u_1$  with a nonzero coefficient. Then p(x) can be viewed as a nonzero polynomial in  $(k[v, u_1, \ldots, u_n])[u_1]$  with  $\deg \geq 1$  on u. We have that  $u_1$  is algebraic over  $k[v, u_1, \ldots, u_n]$ . Thus, the algebraic closure of  $k[v, u_2, \ldots, u_n]$  in  $\operatorname{Frac}(A)$  contains  $u_1$ , hence contains A. It follows that the algebraic closure of  $k[v, u_2, \ldots, u_n]$  equals  $\operatorname{Frac}(A)$ .

This shows that any transcendence basis of A/k has the same cardinality. Indeed, let  $u_1, \ldots, u_n$  and  $v_1, \ldots, v_m$  be transcendence bases of A/k. Then at least one of the  $v_i$ 's must be transcendental over  $k[u_2, \ldots, u_n]$ . This is because if each  $v_i$  is algebraic over  $k[u_2, \ldots, u_n]$ , then  $A \supset k[u_2, \ldots, u_n]$  will be algebraic, in which case  $u_2, \ldots, v_m$  is also a basis, a contradiction.

Say that  $v_1$  is transcendental over  $k[u_2,\ldots,u_n]$ . Then  $A\supset k[v_1,u_2,\ldots,u_n]$  is algebraic  $\ldots$  One of  $v_1,\ldots,v_m$  must be transcendental over  $k[v_1,\ldots,u_2,\ldots,u_n]$ . Hence  $A\supset k[v_1,v_2,u_3,\ldots,u_n]$  is algebraic. If  $m\leq n$ , then  $A\supset k[v_1,v_2,\ldots,v_m,u_{m+1},\ldots,u_n]$  is algebraic and  $v_1,v_2,\ldots,v_m,u_{m+1},\ldots,u_n$  are dependent. This is a contradiction unless n=m.

If F is a field and  $\widetilde{F} \supset F$  is a field extension, then we can measure how far  $\widetilde{F}$  is from being an algebraic extension of F by its transcendence degree over F

 $\operatorname{trdeg}\left(\widetilde{F}/F\right) \equiv \operatorname{card}(\operatorname{independent} \ \operatorname{transcendentals} \ \operatorname{we} \ \operatorname{need} \ \operatorname{to} \ \operatorname{add} \ \operatorname{to} \ F \ \operatorname{to} \ \operatorname{generate} \ \widetilde{F}).$ 

Corollary 8.2.3. trdeg is an invariant of the extension  $\widetilde{F}$ .

**Example 8.2.4.** Let k be a field and  $a_1, a_2, \ldots, a_n$  be indeterminates. Let  $K := k(a_1, \ldots, a_n)$ . Consider  $f(x) = x^n + a_1 x^{n-1} + a_2 x^{n-2} + \cdots + a_n \in K[x]$ . Then  $Gal(L/K) = S_n$  where L denotes the splitting field for f.

Proof. Let  $x_1, x_2, \ldots, x_n \in L$  denote the roots of f. Then  $a_i = (-1)^i \sigma_i(x_1, \ldots, x_n)$  where  $\sigma_i$  denotes the i-th elementary symmetric function. Hence  $L = K(x_1, \ldots, x_n) = k(x_1, \ldots, x_n)$ . Consider the chain of field extensions  $L \supset K \supset k$ . Note that  $L \supset K$  is an algebraic extension and that  $K \supset k$  is a transcendental extension because K is obtained from adding n independent transcendentals to K.

Since  $\operatorname{trdeg}(L/k) = \operatorname{trdeg}(K/k) = n$  and  $L = K(x_1, \dots, x_n)$ , we see that  $x_1, \dots, x_n$  are algebraically independent over k. Therefore, there are pairwise distinct. This shows that  $f(x) \in K[x]$  has distinct roots, so that  $L \supset K$  is separable and thus a Galois extension. It follows that  $\operatorname{Gal}(L/K) = S_n$  and

$$L^{S_n} = (k(x_1, \dots, x_n))^{S_n} = K = k(\sigma_1, \dots, \sigma_n).$$

Theorem 8.2.5 (Main theorem of Galois theory). Let  $L \supset K$  be a Galois extension. Then the mappings

$$(K \subset P \subset L : P \text{ field}) \mapsto (G \leq \operatorname{Gal}(L/K))$$

$$L^G \leftrightarrow G$$

are inverse to each other.

Furthermore, if  $L \supset P \supset K$ , then  $P \supset K$  is a Galois extension of K if and only if  $Gal(L/P) \subseteq Gal(L/K)$ .

*Proof.* Consider  $K \subset P \subset L$  and  $K \subset L^{\operatorname{Gal}(L/P)} \subset L$ . Then  $L^{\operatorname{Gal}(L/P)} \supset P$ . From a theorem from two lectures ago, we have the following two results.

- (a)  $[L:P] = |\operatorname{Gal}(L/P)|$  for any  $K \subset P \subset L$ .
- (b)  $[L:L^G] = |G|$  for any  $G \leq \operatorname{Gal}(L/K)$ .

Therefore,  $\left[K:L^{\operatorname{Gal}(L/P)}\right]\cdot [L:P]=[L:P]$ , so that  $\left[L^{\operatorname{Gal}(L/P)}:L\right]=1$ . Hence  $L^{\operatorname{Gal}(L/P)}=L$ . Similarly,  $\operatorname{Gal}(L/L^G)\leq G$  satisfies  $\left|\operatorname{Gal}(L/L^G)\right|=|G|$ , so that  $\operatorname{Gal}(L/L^G)=G$ .

For the second part our theorem, note that any automorphism of P/K will extend to an automorphism of L/K. This shows that the map  $\{\varphi \in \operatorname{Gal}(L/K) \mid \varphi(P) \subset P\} \to \operatorname{Gal}(P/K)$  given by  $\varphi \mapsto \varphi \upharpoonright_P$  is surjective. Then  $P \supset K$  will be Galois if and only if the elements of  $\{\varphi \in \operatorname{Gal}(L/K) \mid \varphi(P) \subset P\}$  induce [P : K] distinct elements of  $\operatorname{Gal}(P/K)$ .

We compute

$$\begin{split} |\mathrm{Gal}(L/P)| &= [L:P] \\ [P:K] &= \frac{[L:K]}{[L:P]} = \frac{|\mathrm{Gal}(L/K)|}{|\mathrm{Gal}(L/K)|} \\ [P:K] &= [\mathrm{Gal}(L/K):\mathrm{Gal}(L/P)] \,. \end{split}$$

Thus,  $P \supset K$  is a Galois extension if and only if any element of Gal(L/K) leaves P invariant. But  $P = L^{Gal(L/P)}$ , and  $P = P^{Gal(L/P)}$ . Hence any  $g \in Gal(L/P)$  satisfies

$$g(P) = g(L^{Gal(L/P)}) = L^{g Gal(L/P)g^{-1}}.$$

It follows that  $g(P) = P \iff g \operatorname{Gal}(L/P)g^{-1} = \operatorname{Gal}(L/P)$ .

# 8.3 Lecture 18

**Example 8.3.1.** Let K be a field with char  $K \notin \{2,3\}$ . Let f be an irreducible, monic, cubic polynomial over K. Let L be the splitting field for f. Let  $D := \operatorname{Discr} f \in K \setminus K^2$ . Then  $\operatorname{Gal}(L/K) = S_3$ . We get

$$L \supset L^{A_3} \supset K$$
$$Gal(L/K) \trianglerighteq A_3 \trianglerighteq \{e\}.$$

It follows that  $L^{A_3} \supset K$  is Galois with  $\operatorname{Gal}(L^{A_3}/K) \cong \operatorname{Gal}(L/K)/A_3 \cong C_2$ . In fact,  $L^{A_3} \cong K \left\lceil \sqrt{D} \right\rceil$ .

Now, let p > 2 be prime. Consider the cyclotomic field  $\Gamma_p \supset \mathbb{Q}$ . We have that

$$\operatorname{Gal}(\Gamma_p/\mathbb{Q}) \cong (\mathbb{Z}/p)^{\times} \cong \mathbb{Z}/(p-1).$$

Let  $H \subseteq \operatorname{Gal}(\Gamma_p/\mathbb{Q})$  be the unique subgroup of index 2. Then  $\left[\Gamma_p^H:\mathbb{Q}\right]=2$ .

Let  $\langle \varphi \rangle = \operatorname{Gal}(\Gamma_p/\mathbb{Q})$ . Then  $\varphi \upharpoonright_{\mu_p} : \mu_p \to \mu_p$  is a group automorphism and uniquely determines  $\varphi$ , which in turn is uniquely determined by the image of  $\zeta$  the positive p-th root of 1. Write  $\varphi(\zeta) = 1\zeta^r$  for some  $r \in \mathbb{Z}_{>0}$  so that  $[r]_p \in \mathbb{Z}/p$  is a generator of  $(\mathbb{Z}/p)^{\times}$ .

**Definition 8.3.2.** Given  $k \in \mathbb{Z}_{>0}$  and prime p > 2, define the Legendre symbol

$$\left(\frac{k}{p}\right) = \begin{cases} 1 & [k]_p \in \left(\left(\mathbb{Z}/p\right)^{\times}\right)^2 \\ -1 & [k]_p \notin \left(\left(\mathbb{Z}/p\right)^{\times}\right)^2 \end{cases}.$$

Consider  $\alpha \in \Gamma_p$  given by

$$\alpha = \sum_{k=1}^{p-1} \left(\frac{k}{p}\right) \zeta^{r^{k-1}} = \zeta - \zeta^r + \zeta^{r^2} - \dots \zeta^{r^{p-2}}.$$

If  $g \in \operatorname{Gal}(\Gamma_p/\mathbb{Q})$ , then

$$g(\alpha) = \begin{cases} \alpha & g \in H \\ -\alpha & g \notin H \end{cases}.$$

Then  $\alpha \in \Gamma_p^H$ . Also,  $\alpha^2$  is fixed by every element of  $\operatorname{Gal}(\Gamma_p/\mathbb{Q})$  and thus is rational. This implies that  $\Gamma_p^H = \mathbb{Q}[\alpha]$ .

**Lemma 8.3.3.**  $\alpha^2 = (-1)^{\frac{p-1}{2}} p$ , so that

$$\Gamma_p^H = \begin{cases} \mathbb{Q} \left[ \sqrt{p} \right] & p \equiv 1 \mod 4 \\ \mathbb{Q} \left[ \sqrt{-p} \right] & p \not\equiv 1 \mod 4 \end{cases}.$$

*Proof.* Let  $L \supset K$  be a finite extension of fields. Then for any  $u \in L$ , we get a map  $\operatorname{mult}_u : L \to L$ , which is linear over K. Applying trace determines a K-linear map  $L \to K$  given by  $u \mapsto \operatorname{tr}(\operatorname{mult}_u)$ . This induces a symmetric bilinear map  $\langle \cdot, \cdot \rangle : L \otimes_K L \to K$  given by  $u \otimes v \mapsto \operatorname{tr}(\operatorname{mult}_u \circ \operatorname{mult}_v)$ . Note that if  $u \neq 0$ , then

$$\langle u, u^{-1} \rangle = \operatorname{tr}(\operatorname{mult}_{uu^{-1}}) = \operatorname{tr}(\operatorname{id}_L) = [L:K]$$

since  $(\operatorname{char} K, [L:K]) = 1$ . Now, the vector space  $\Gamma_p$  has a  $\mathbb{Q}$ -basis  $\{1, \zeta, \zeta^2, \ldots, \zeta^{p-2}\}$ . Hence  $\operatorname{mult}_{\zeta^2}$  is a cyclic operator, and  $\operatorname{tr}(\operatorname{mult}_1) = p-1$  and  $\operatorname{tr}(\operatorname{mult}_{\zeta^k}) = -1$  for each  $k = 1, \ldots, p-2$ . It follows that

$$\langle \zeta^k, \zeta^l \rangle = \begin{cases} p-1 & k+l \equiv 0 \mod p \\ 1 & \text{otherwise} \end{cases}.$$

If  $x = \sum_{i=0}^{p-1} x_i \zeta^i$  and  $y = \sum_{i=0}^{p-1} y_i \zeta^i$  are two elements of  $\Gamma_p$ , then we can choose  $x_i$  and  $y_i$  such that  $\sum x_i = 0$  and  $\sum y_i = 0$ . Thus,

$$\langle x, y \rangle = p(x_0 y_0 + \sum k = 1^{p-1} x_k y_{p-k}).$$

But  $\alpha = \sum_{k=1}^{p-1} \left(\frac{k}{p}\right) \zeta^{r^{k-1}}$ , so that

$$\alpha^{2} = \frac{1}{p-1} \langle \alpha, \alpha \rangle = \frac{1}{p-1} \sum_{k=1}^{p-1} p\left(\frac{k}{p}\right) \left(\frac{-k}{p}\right)$$
$$= \frac{p}{p-1} \sum_{k=1}^{p} \left(\frac{k}{p}\right) \left(\frac{-k}{p}\right)$$
$$= p\left(\frac{-1}{p}\right) = (-1)^{\frac{p-1}{2}} p.$$

**Definition 8.3.4.** Let  $L \supset K$  be a field extension and  $\alpha \in L_{\dot{\iota}}$ . We sat that  $\alpha$  can be expressed in radicals over K if it can be obtained from elements in K by applying +,  $\cdot$ , and  $\sqrt[r]{\cdot}$ , i.e., there exists a tower of subfields

$$K = K_0 \subset K_1 \subset \cdots \subset K_s \subset L$$

such that  $K_{i+1} = K_i(\alpha_{i+1})$  where  $\alpha_{i+1}^{r_{i+1}} \in K_i$  and  $\alpha \in K_s$ .

**Proposition 8.3.5.** If  $f(x) \in K[x]$  is irreducible,  $L \supset K$  is an extension, and  $\alpha$  is a root of f(x), then  $\alpha$  can be expressed in radicals if and only if any root of f can be expressed in radicals in the splitting field for f.

*Proof.* If  $L_1, L_2 \supset K$  are field extensions and  $\alpha_1 \in L_1$  and  $\alpha_2 \in L_2$  are roots of f, then by Lemma 4.2.8 there is a unique map  $\varphi : K(\alpha_1) \xrightarrow{\cong} K(\alpha_2)$  such that  $\varphi(\alpha_1) = \alpha_2$ . Now transport all suitable expressions by  $\varphi$  or  $\varphi^{-1}$ .

# 9 Solvability in radicals

**Definition 9.0.1.** We say that  $f(x) \in K[x]$  is solvable in radicals if every root of f can be expressed in radicals over K.

This is equivalent to saying that L is a splitting field for f, then there is a tower of subfields  $K = K_0 \subset K_1 \subset \cdots \subset K_s = L$  such that  $K_{i+1} = K_i(\alpha_{i+1})$  where  $\alpha_{i+1}^{r_i+1} \in K_i$ .

**Theorem 9.0.2.** If K is a field with characteristic 0,  $f(x) \in K[x]$  is irreducible, and  $L \supset K$  is the splitting field for f, then f is solvable in radicals over K if and only if Gal(L/K) is solvable.

#### Note 9.0.3.

- 1. A generic polynomial equation over K of deg  $\geq 5$  will not be solvable in radicals, since Gal  $\cong S_n$ .
- 2. If  $f(x) \in \mathbb{Q}[x]$  is irreducible of degree 5, then f will not be solvable in radicals as soon as  $\operatorname{Gal}(L/\mathbb{Q}) \in \{S_5, A_5\}$ . Suppose  $f \in \mathbb{Q}[x]$  is such a polynomial and let  $\alpha_1, \ldots, \alpha_5$  be the roots of f. Note that  $\operatorname{Gal}(L/\mathbb{Q}) \subset S_5$ . Since f is irreducible, it must be separable, which means that the  $\alpha_i$  are pairwise distinct. Hence  $5 \mid |\operatorname{Gal}(L/\mathbb{Q})|$ . Therefore,  $\operatorname{Gal}(L/\mathbb{Q})$  must contain an element of order 5, so that  $\operatorname{Gal}(L/\mathbb{Q})$  contains a 5-cycle. If we can choose f so that  $\operatorname{Gal}(L/\mathbb{Q})$  contains a transposition, then  $\operatorname{Gal}(L/\mathbb{Q}) = S_5$ .

Choose f so that it has exactly three real roots. In this case, complex conjugation will belong to  $Gal(L/\mathbb{Q})$ , so that  $Gal(L/\mathbb{Q}) = S_5$ . Start with  $x^5 - 16x = x(x-2)(x+2)(x^2+4)$ , which has exact three real roots. To make this irreducible, shift its graph to obtain the polynomial  $f(x) = x^5 - 16x + 2$ .

#### 9.1 Lecture 19

**Theorem 9.1.1.** If K is a field with characteristic 0,  $f(x) \in K[x]$  is irreducible, and  $L \supset K$  is the splitting field for f, then f is solvable in radicals over K if and only if Gal(L/K) is solvable.

Proof.

 $(\Longleftrightarrow)$ 

We have a series

$$G = G^{(0)} \trianglerighteq G^{(1)} \trianglerighteq \cdots \trianglerighteq G^{(s)} \trianglerighteq \{e\},\$$

which we can refine to get a normal series

$$G = G_0 \triangleright G_1 \triangleright \cdots \triangleright G_r = \{e\}$$

such that  $G_{i+1}/G_i \cong \mathbb{Z}/n_i$ . Letting  $K_i = L^{n_i}$ , we have a tower

$$K = K_0 \subset K_1 \subset \cdots \subset K_r = L.$$

Let F be the cyclotomic field that contains all roots of 1 of order  $n = n_1 n_2 \cdots n_r$ . Consider the tower

$$KF \subset K_1F \subset \cdots \subset K_rF = LF.$$

Then  $K_{n_i}F \supset K_1F$  is a cyclic extension of degree dividing  $n_i$ .

**Lemma 9.1.2.** Let K be a field and let n have char  $K \nmid n$ . Suppose that  $K \supset \mu_n$ .

- (a) For any  $\alpha \in K$ , the extension  $K(\sqrt[n]{\alpha}) \supset K$  is cyclic of order dividing n.
- (b) For every  $\widetilde{K} \supset K$  Galois and cyclic of order n, there exists  $\alpha \in K$  such that  $\widetilde{K} \cong K(\sqrt[n]{\alpha})$ .

Before proving this, note that (b) implies that  $K_{i+1}F \supset K_iF$  for every i.

Proof.

(a) By definition,  $K(\sqrt[n]{\alpha})$  contains some root of  $x^n - \alpha$ . But K contains  $\mu_n$ , so that  $K(\sqrt[n]{\alpha})$  contains every root of  $x^n - \alpha$ . Thus,  $K(\sqrt[n]{\alpha})$  is the splitting field for  $x^n - \alpha$ . Hence  $K(\sqrt[n]{\alpha}) \supset K$  is Galois. If  $\sigma \in \text{Gal}(K(\sqrt[n]{\alpha})/K)$ , then  $\sigma(\sqrt[n]{\alpha}) = \zeta_{\sigma}\sqrt[n]{\alpha}$  where  $\zeta_{\sigma}$  is some n-th root of 1 depending on  $\sigma$ . Then we get a map

$$\zeta: \operatorname{Gal}(K(\sqrt[n]{\alpha})/K) \to \mu_n$$

given by  $\sigma \mapsto \zeta_{\sigma}$ . But since  $K \supset \mu_n$ , if  $\sigma, \tau \in \operatorname{Gal}(K(\sqrt[n]{\alpha})/K)$ , then

$$\sigma(\tau(\sqrt[n]{\alpha})) = \sigma(\zeta_{\tau} \sqrt[n]{\alpha}) = \sigma(\zeta_{\tau})\sigma(\sqrt[n]{\alpha}).$$

As  $\zeta_{\tau} \in \mu_n \subset K$ , we see that  $\sigma(\zeta_{\tau}) = \zeta_{\tau}$ . This implies that

$$\sigma(\tau(\sqrt[n]{\alpha})) = \zeta_{\tau}\zeta_{\sigma}\sqrt[n]{\alpha}.$$

But  $\sigma(\tau(\sqrt[n]{\alpha})) = \zeta_{\tau\sigma} \sqrt[n]{\alpha}$  as well, so that  $\zeta_{\tau\sigma} = \zeta_{\tau}\zeta_{\sigma}$ . This shows that  $\zeta$  is a homomorphism.

Moreover, if  $\sigma \in \ker \zeta$ , i.e.,  $\zeta_{\sigma} = 1$ , then  $\sigma(\sqrt[n]{\alpha}) = 1 \cdot \sqrt[n]{\alpha} = \sqrt[n]{\alpha}$ . Since any  $\sigma \in \operatorname{Gal}(K(\sqrt[n]{\alpha})/K)$  preserving  $\sqrt[n]{\alpha}$  must be the identity, it follows that  $\zeta$  is injective. As a result, we get an embedding  $\operatorname{Gal}(K(\sqrt[n]{\alpha})/K) \leq \mu_n$ .

(b) Suppose that  $\operatorname{Gal}\left(\widetilde{K}/K\right)$  is cyclic of order  $d\mid n$ . We want to show that there exists  $\alpha\in K$  such that  $\widetilde{K}\cong K(\sqrt[d]{\alpha})$ .

Let  $\alpha \in \widetilde{K}$  and  $\xi \in \mu_d \subset \mu_n \subset K$ . The Lagrange resolvent of  $(\alpha, \xi)$  is the element

$$\ell(\alpha,\xi) = \alpha + \xi \sigma(\alpha) + \xi^2 \sigma^2(\alpha) + \dots + \xi^{d-1} \sigma^{d-1}(\alpha)$$

of  $\widetilde{K}$  where  $\sigma \in \operatorname{Gal}\left(\widetilde{K}/K\right)$  is a generator.

Note that  $\sigma(\ell(\alpha,\xi)) = \xi^{-1}\ell(\alpha,\xi)$ , so that  $\sigma(\ell(\alpha,\xi)^2) = \xi^{-k}\ell(\alpha,\xi)^k$ .

Suppose that  $\xi$  is a primitive d-th root of unity. We see that  $\operatorname{id} + \xi \sigma + \xi^2 \sigma^2 + \dots + \xi^{d-1} \sigma^{d-1}$  is a linear combination of operators  $L \to L$  viewed as a K-vector space. But in  $\operatorname{End}_K\left(\widetilde{K}\right)$  the generators are linearly independent. Therefore, module this statement, we have that

$$\sum_{k=0}^{d-1} \xi^k \sigma^k \neq 0$$

in End<sub>K</sub>  $(\widetilde{K})$ . Hence there exists  $\alpha \in \widetilde{K}$  such that  $\ell(\alpha, \xi) = \sum_{k=0}^{d-1} \xi^k \sigma^k(\alpha) \neq 0$ . But for each  $i = 0, \ldots, d-1$ , we see that  $\sigma^i(\ell(\alpha, \xi)) = \xi^{-i}\ell(\alpha, \xi)$ . This implies that

$$\ell(\alpha,\xi), \ \sigma(\ell(\alpha,\xi)), \ \sigma^2(\ell(\alpha,\xi)), \ \dots, \ \sigma^{d-1}(\ell(\alpha,\xi))$$

are pairwise distinct in  $\widetilde{K}$ . Therefore,  $\ell(\alpha,\xi) \in \widetilde{K}$  does not belong to any proper subfield of  $\widetilde{K}$ . It follows that  $\widetilde{K} = K(\ell(\alpha,\xi))$ . But  $\sigma(\ell(\alpha,\xi)^d) = \underbrace{\xi^{-1}}_{} \ell(\alpha,\xi)^d = \ell(\alpha,\xi)^d$ . Hence

$$\ell(\alpha, \xi)^d = \widetilde{K}^{\operatorname{Gal}(\widetilde{K}/K)} = K.$$

This proves our lemma modulo the statement that id,  $\sigma$ ,  $\sigma^2$ , ...,  $\sigma^{d-1}$  are linearly independent linear operators.

Note 9.1.3. The  $\sigma^i$  belong to  $\operatorname{End}_K\left(\widetilde{K}\right)$  and commute with each other. They can be simultaneously diagonalized over  $L\supset K$  the splitting field for  $f(x)=\det(\sigma-x\cdot\operatorname{id})$ . Writing a linear combination of the  $\sigma^i$  and evaluating it on a basis of eigenvectors will produce a homogenous linear system with a Vandermonde coefficient matrix. Then one needs to show that  $\sigma$  has distinct eigenvalues.

 $(\Longrightarrow)$ 

See Section 9.2.

#### 9.2 Lecture 20

**Definition 9.2.1.** If G is a group and k a field, then a k-character of G is a group homomorphism  $\chi: G \to \operatorname{GL}_1(k) = k^{\times}$ .

Each k-character  $\chi$  of G can be viewed as a function with values in k.

**Lemma 9.2.2 (Dedekind).** If  $\chi_1, \ldots, \chi_s$  are pairwise distinct k-characters of G, then they are linearly independent in  $\operatorname{Fun}(G,k)$ .

*Proof.* We induct on s. If s=1, then  $\chi_1$  must be linearly independent since  $\chi_1 \neq 0$ . Suppose, inductively, that any collection  $\sigma_1, \ldots, \sigma_t$  of characters with  $t \leq s$  is linearly independent. Suppose that  $\chi_1, \ldots, \chi_s$  are linearly dependent. Then there are  $a_1, \ldots, a_s \in k$  such that  $a_1\chi_1 + \cdots + a_s\chi_s$  is the zero function. By our IH, each  $a_i$  must be nonzero, say,  $a_s$ . Let  $b_i = -\frac{a_i}{a_s}$ . Then

$$\sum_{i=1}^{s-1} b_i \chi_i = \chi_s.$$

If  $g, h \in G$ , then

$$\chi_s(h)\chi_s(g) = \sum_{i=1}^{s-1} b_i \chi_i(h)\chi(g),$$

in which case  $\chi_s(g) = \sum_{i=1}^{s-1} (b_i \frac{\chi_i(h)}{\chi_s(h)}) \chi_i(g)$ . Fix  $h \in G$ , so that

$$\chi_s = \chi_s = \sum_{i=1}^{s-1} (b_i \frac{\chi_i(h)}{\chi_s(h)}) \chi_i$$

and  $\chi_s = \sum_{i=1}^{s-1} b_i \chi_i$ . It follows that  $0 = \sum_{i=1}^{s-1} (b_i \frac{\chi_i(h)}{\chi_s(h)} - b_i) \chi_i$ . By our IH, we see that  $b_i \frac{\chi_i(h)}{\chi_s(h)} - b_i = 0$ . But  $b_i \neq 0$  for any i. We have that  $\chi_i(h) = \chi_s(h)$  for any  $i = 1, \ldots, s-1$ . This proves that  $\chi_i = \chi_s$  for any  $i = 1, \ldots, s-1$ . This contradicts the assumption that the  $\chi_1, \ldots, \chi_s$  are pairwise distinct.  $\square$ 

**Definition 9.2.3.** If  $K_1, K_2 \subset L$ , then the *composite of*  $K_1$  *and*  $K_2$  *in* L is the field

$$K_1K_2 = \bigcap \left\{ P \mid P \subset L, K_1, K_2 \subset P \right\}.$$

Let  $K_1$  and  $K_2$  be finite extensions of k, so that  $K_1 = k(a_1, \ldots, a_s)$  and  $K_2 = k(b_1, \ldots, b_t)$ . Then the field  $k(a_1, \ldots, a_s, b_1, \ldots, b_t)$  both contains  $K_1K_2$  and is contained in some L such that  $K_1, K_2 \subset L$ . Hence

$$K_1K_2 = k(a_1, \ldots, a_s, b_1, \ldots, b_t).$$

**Lemma 9.2.4.** Suppose that K and F are two finite field extensions of k. Then

- (a) If  $K \supset k$  is Galois, then so is  $KF \supset F$ .
- (b)  $Gal(KF/F) = Gal(K/K \cap F)$ .

Proof.

- (a) If  $K \supset k$  is Galois, then K is the splitting field of some separable polynomial  $f(x) \in k[x]$ . Thus, KF is the splitting field of f(x) viewed over F. But if f is separable over k, then it is separable over F. Therefore,  $KF \supset F$  is Galois.
- (b) Consider the tower of extensions  $k \subset K \subset KF$ . The main theorem of Galois theory says that  $\operatorname{Gal}(K/k) \leq \operatorname{Gal}(KF/k)$  since  $K \supset k$  is assumed to be Galois. Thus, if  $\sigma \in \operatorname{Gal}(KF/k)$ , then  $\sigma(K) \subset K$ . Indeed,  $\sigma(K) = K$  as a subfield in KF if and only if  $\sigma(K) = (KF)^{\operatorname{Gal}(KF/K)}$ . Let  $g \in \operatorname{Gal}(KF/K) \subset \operatorname{Gal}(KF/k)$ . Then g(x) = x for any  $x \in K$ .

Let  $x \in K$ . Consider  $\sigma(x) \in KF$ . We must show that  $g(\sigma(x)) = \sigma(x)$  for any  $g \in \operatorname{Gal}(KF/K)$ , i.e.,  $(\sigma^{-1}g\sigma)(x) = x$  for any g. But since  $\sigma^{-1}g\sigma \in \sigma^{-1}\operatorname{Gal}(KF/K)\sigma$ , we see that  $\sigma^{-1}g\sigma(x) = x$  for any  $x \in K$ . Hence we get a natural homomorphism  $\rho : \operatorname{Gal}(KF/F) \to \operatorname{Gal}(K/K)$  given by  $\sigma \mapsto \sigma \upharpoonright_K$ . Note that

$$\ker \rho = \{ \sigma \in \operatorname{Gal}(KF/F) \mid \sigma \upharpoonright_K = \operatorname{id}_K \}$$
$$= \{ \sigma \in \operatorname{Gal}(KF/k) \mid \sigma \upharpoonright_K = \operatorname{id}_K, \ \sigma \upharpoonright_F = \operatorname{id}_F \} .$$

But KF is generated by K and F, so that  $\sigma = \mathrm{id}_{KF}$ . This shows that  $\ker \rho = \{\mathrm{id}_{KF}\}$ . We see that  $\mathrm{Gal}(KF/F) \subset \mathrm{Gal}(K/k)$ .

Let  $H := \operatorname{im} \rho \subset \operatorname{Gal}(K/k)$  and consider its fixed subfield  $K^H$ . Note that (b) is equivalent to saying that  $K^H - K \cap F$ . We have that  $K^H \supset K \cap F$  because  $K^H = \{x \in K \mid (\forall \sigma \in \operatorname{Gal}(KF/F)) \ (\sigma(x) = x)\}$ . Moreover, if we view  $K^H$  as subfield of KF, then  $k \subset K^H \subset KF$  and  $k \subset F \subset KF$ . Since  $\operatorname{Gal}(KF/F)$  fixes  $K^H$  and F (pointwise), it follows that  $\operatorname{Gal}(KF/F)$  fixes  $K^HF$ . Therefore,  $K^HF \subset KF^{\operatorname{Gal}(KF/F)} = F$ , so that  $K^H \subset H$ . This proves that  $K^H \subset F \cap K$ .

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Corollary 9.2.5. If both K and F are Galois field extensions of k, then



is a lattice of Galois field extensions.

**Theorem 9.2.6.** If K is a field with characteristic 0,  $f(x) \in K[x]$  is irreducible, and  $L \supset K$  is the splitting field for f, then f is solvable in radicals over K if and only if Gal(L/K) is solvable.

Proof.

 $(\Longleftrightarrow)$ 

This was proven in Section 9.1.

 $(\Longrightarrow)$ 

For any root  $\alpha$  of f, we can find an extension  $K_{\alpha} \supset K$  such that  $\alpha \in L_{\alpha} \subset L$  and there exists a tower of radical extensions

$$K = K_0 \subset K_1 \subset \cdots \subset K_s = K_\alpha \subset L$$

with  $K_{i+1} = K_i(\alpha_{i+1})$  and  $\alpha_{i+1}^{n_{i+1}} \in L_i$ .

Claim. Without loss of generality, we may assume that  $K_{\alpha}$  satisfies the following properties.

- $\alpha \in K_{\alpha}$ .
- $K_{\alpha} \supset K$  is Galois.
- Each step of  $K_{\alpha}$  (viewed as our tower of radical extensions) is Galois and cyclic.

Proof. Since  $K_{\alpha} \supset K$  is a finite extension, we can find a K-basis  $e_1, \ldots, e_n$  of  $K_{\alpha}$ . Let  $f_i \in K[x]$  denote the minimal polynomial of  $e_i$ . Let  $S_i$  denote the splitting field of  $f_i$ . Then  $S_i \supset K$  is a Galois extension and contains  $e_i$ . Note that the composite of the  $S_i$  contains each  $e_i$ . Let  $L_{\alpha} = S_1 S_2 \cdots S_n$ . Then  $K \subset K_{\alpha} \subset L_{\alpha}$ . (We call  $L_{\alpha}$  the Galois closure of  $K_{\alpha}$ .) Consider the tower  $K = K_0 \subset K_1 \subset \cdots \subset K_s = K_{\alpha}$  where  $K_{i+1} \supset K_i$  is a radical extension of degree  $n_i$ . If  $\sigma \in \operatorname{Gal}(L_{\alpha}/K)$ , then  $K = \sigma K \subset \sigma K_1 \subset \cdots \subset \sigma K_{\alpha}$  is still a tower of radical extensions.

By taking the composites  $K_1 \sigma K_1 \subset \cdots \subset K_1 \sigma K_s$  and  $K_2 K_1 \sigma K_1 \cdots$ , we get a composite of all  $\{\sigma K\alpha\}_{\sigma \in \operatorname{Gal}(L_{\alpha}/K)}$ , which will be a tower of radical extensions. But  $K \subset \prod_{\sigma} \sigma K_{\alpha} \subset L_{\alpha}$ , and  $L_{\alpha}$  is generated by all  $\sigma K_{\alpha}$ . Hence  $L_{\alpha} = \prod_{\sigma} \sigma K_{\alpha} = L$ .

We still must prove that each step in our radical tower is Galois and cyclic. Let  $n = n_1 n_2 \cdots n_k$ . Let  $F = K[\mu_n]$ . If the tower  $K = K_0 \subset K_1 \subset \cdots \subset K_t = L_\alpha$  has  $K_i = K_{i-1} \begin{bmatrix} \frac{n_i}{\sqrt[n]} a_i \end{bmatrix}$ , then we can pass to composites

$$K \subset K_0F \subset K_1F \subset \cdots \subset K_tF = L_\alpha F.$$

We see that  $LF \supset K$  is radical and Galois as the splitting field for  $x^n - 1$  and that  $K_iF \supset K_{i+1}F$  is radical of degree  $n_i$  and contains  $\mu_{n_i}$ . Thus,  $K_iF \supset K_{i+1}F$  is Galois and cyclic of degree dividing  $n_i$  by Lemma 8.3.3(a).

We have constructed an extension  $LF \supset K$  such that

- $\alpha \in LF$ ,
- $LF \supset L$  is Galois, and
- LF is a tower of radical, cyclic, Galois extensions.

It follows that  $\operatorname{Gal}(LF/K)$  is solvable. But  $LF \supset L \supset K$  where  $L \supset K$  is Galois. Hence  $\sigma(L) \subset L$  for any  $\sigma \in \operatorname{Gal}(LF/K)$ , so that  $\operatorname{Gal}(L/K) < \operatorname{Gal}(LF/K)$ . This proves that  $\operatorname{Gal}(L/K)$  is solvable.  $\square$ 

#### 9.3 Lecture 21

**Definition 9.3.1.** Let K be a field and  $f(x) \in K[x]$ . We say that f is solvable in quadratic radicals if the splitting field L for f is a tower

$$K = K_0 \subset K_1 \subset \cdots \subset K_s = L$$

such that  $K_i = K_{i-1} \left[ \sqrt{a_i} \right]$  for some  $a_i \in K_{i-1}$ .

**Theorem 9.3.2.** Let K be a field with char  $K \neq 2$  and  $f(x) \in K[x]$  be irreducible. Then f is solvable in quadratic radicals if and only if  $[L:K] = 2^n$  for some n where L denotes the splitting field for f.

Proof.

 $(\Longrightarrow)$ 

We have that  $L \supset K$  is a tower of quadratic extensions. Hence  $[L:K] = 2^n$  for some n.

 $(\longleftarrow)$ 

We have that  $[L:K]=2^n$  for some  $n \geq 0$  and  $\deg f = [K(\alpha):K] \mid [L:K]$  where  $\alpha$  is a root of f(x). Thus,  $[K(\alpha):K]$  equals a power of 2, so that f is separable. This shows that  $L \supset K$  is Galois and thus that  $G := \operatorname{Gal}(L/K)$  has order  $2^n$ . It follows that there is some normal series

$$G = G^0 \triangleright G^1 \triangleright \cdots \triangleright G^s = \{e\}$$

such that  $G^i/_{G^{i+1}} \cong \mathbb{Z}/2$ . This induces a tower of field extensions

$$K = L^{G^0} \subset L^{G_1} \subset \dots \subset L^{G_s} = L$$

such that  $[L^{G^{i+1}}:L^{G_i}]=2.$ 

Note 9.3.3 (The construction problem). Given a unit measure and segments of lengths  $a_1, \ldots, a_k$ , we want to construct a segment of length  $\alpha$  using ruler and compass. Elementary geometry shows that such a construction is possible if and only if  $\alpha$  can be expressed in quadratic radicals over  $\mathbb{Q}(a_1, \ldots, a_k)$ . If  $\alpha$  is transcendental over  $\mathbb{Q}(a_1, \ldots, a_k)$ , then our construction is impossible.

**Example 9.3.4.** We see that  $\pi$  cannot be constructed over  $\mathbb{Q}$ , i.e., we cannot square the circle.

Moreover, if  $\alpha$  is algebraic over  $\mathbb{Q}(a_1,\ldots,a_k)$ , then  $\alpha$  can be constructed by Theorem 9.3.2 if and only if the minimal polynomial of  $\alpha$  has degree power of 2.

#### Example 9.3.5.

- (a) <u>Doubling the cube.</u> Given a segment of length one, construct a segment of length  $\sqrt[3]{2}$ . Since the minimal polynomial of  $\sqrt[3]{2}$  is  $x^3 2$ , such a construction is impossible.
- (b) Trisecting an angle  $\varphi$ . Given a segment of length  $\cos \varphi$ , construct a segment of length  $\cos \left(\frac{\varphi}{3}\right)$ . The minimal polynomial of  $\cos \left(\frac{\varphi}{3}\right)$  over  $\mathbb{Q}(\cos \varphi)$  is  $4x^3 3x \cos \varphi$ . In general, this is irreducible, in which case our construction is impossible.
- (c) Constructing regular n-gons. Given a segment of length i, construct a segment of length  $\cos\left(\frac{2\pi}{n}\right)$ . This is possible if and only if  $e^{\frac{2\pi i}{n}}$  is expressible in quadratic radicals over  $\mathbb{Q}$ . In turn, this happens if and only if

$$\underbrace{\left[\Gamma_n:\mathbb{Q}\right]}_{\varphi(n)}=2^s.$$

For example, if p is prime, then we can construct a regular p-gon if and only if  $1 + 2^k$  for some k. Currently, the largest known such p is 65,537.

# 10 Further applications of Galois theory

#### 10.1 Lecture 22

To begin, note that the following statements are true.

- If  $f(x) \in \mathbb{R}[x]$  has odd degree, then it has a real root.
- Every  $\alpha \in \mathbb{C}$  has a square root in  $\mathbb{C}$ .

Now, suppose that  $K \supseteq \mathbb{R}$  is a finite field extension. If  $[K : \mathbb{R}]$  is odd and  $\alpha \in K \setminus \mathbb{R}$ , then  $K \supset \mathbb{R}(\alpha) \supset \mathbb{R}$ , in which case  $\deg f \mid [K : \mathbb{R}]$  where f denotes the minimal polynomial of  $\alpha$  over  $\mathbb{R}$ . In this case, f has odd degree and thus has a root in  $\mathbb{R}$ , so that  $\mathbb{R}(\alpha) = \mathbb{R}$ , a contradiction. This proves that  $[K : \mathbb{R}]$  is odd.

We want to prove the fundamental theorem of algebra: that any  $f(x) \in \mathbb{C}[x]$  has a root in  $\mathbb{C}$ . Note that if c is a complex root of  $f(x)\overline{f(x)} \in \mathbb{R}[x]$ , then either c or  $\overline{c}$  is a root of f(x). Thus, it suffices to show that any polynomial over  $\mathbb{R}$  has a root in  $\mathbb{C}$ .

To this end, let  $g(x) \in \mathbb{R}[x]$  be non-constant and irreducible. Let L denote the splitting field for g. Then  $[L:K] = |\operatorname{Gal}(L/\mathbb{R})|$  is even, so that there is some nontrivial 2-Sylow subgroup  $H \leq \operatorname{Gal}(L/\mathbb{R})$ . This means that the intermediate extension  $L \supset L^H \supset \mathbb{R}$  has odd degree. But then  $L^H = \mathbb{R}$ . This means that  $L \supset L^H$  is Galois, so that

$$[L:\mathbb{R}] = [L:L^H] = |Gal(L/L^H)| = |H| = 2^n$$

for some n. By Theorem 9.3.2, it follows that g(x) is solvable in quadratic radicals. Therefore,  $L = \mathbb{C}$  since  $[\mathbb{C} : \mathbb{R}] = 2$ .

**Theorem 10.1.1 (Primitive element theorem).** Suppose that  $L \supset K$  is a finite field extension. This has a primitive element, i.e.,  $L = K(\theta)$  for some  $\theta \in L_j$ , if and only if there are at most finitely many intermediate fields  $K \subset F \subset L$ .

*Proof.* If K is finite, then L is a finite group with cyclic multiplicative group  $\langle \theta \rangle$ . In this case, we have shown that  $L = K(\theta)$ .

 $(\Longleftrightarrow)$ 

For any  $\alpha, \beta \in L$ , consider the collection of intermediate fields

$$K \subset K(\alpha + c\beta) \subset L$$

where  $c \in K$ . Thus,  $\exists c, c' \in K$  such that  $E := K(\alpha + c\beta) = K(\alpha + c'\beta)$ . Hence  $(c - c')\beta \in E$ , and  $c - c' \in K \setminus \{0\}$ . Then  $\beta \in E$ , so that  $\alpha \in E$ . This shows that  $E \supset K(\alpha, \beta)$ . It's clear that  $E \subset K(\alpha, \beta)$ . Hence  $E = K(\alpha, \beta)$ . But  $E \supset K$  is a finite extension, which implies that  $E \supset K(\alpha, \beta)$  for some  $\alpha_1, \ldots, \alpha_n$ . By indiction on  $E \subset K$ , we can find elements  $E \subset K$  such that

$$K(\alpha_1, \dots, \alpha_n) = K(\alpha_1 + c_2\alpha_2 + c_3\alpha_3 + \dots + c_n\alpha_n).$$

 $(\Longrightarrow)$ 

We have that  $L = K(\theta)$ . Let  $f(x) \in K[x]$  denote the minimal polynomial of  $\theta$  over K. Let  $K \subset F \subset L$  be an intermediate field extension. Let  $g_F(x) \in F[x]$  denote the minimal polynomial over F. This proves that  $g_F(x) \mid f(x)$  in F[x]. We get a map

(intermediate field extensions 
$$K \subset F \subset L$$
)  $\rightarrow$  (divisors of  $f(x)$ )

given by  $F \mapsto g_F(x)$ . Since there are at most finitely many divisors of f(x), it suffices to check that this map is injective.

Suppose that  $K \subset F \subset L$ . Let  $F_0 \subset F$  be the subfield obtained from K by adjoining the coefficients of  $g_F(x)$ . It is enough to show that  $F_0 = F$ . Note that  $g_F(x)$  is irreducible in F[x], so that  $g_F(x)$  is irreducible in  $F_0[x]$  Therefore,  $g_F(x) \in F_0[x]$ , which means that  $g_F(x)$  is the minimal polynomial of  $\theta$  over  $F_0$ . Then  $[L:F_0] = \deg g_F = [L:F]$ , so that  $F_0 = F$ .

**Corollary 10.1.2.** If  $L \supset K$  is a (finite) separable extension, then L has a primitive element.

*Proof.* It suffices to show that if  $\alpha, \beta \in L$  are separable over K, then  $K(\alpha, \beta) = K(\theta)$  for some  $\theta$ . If K is finite, then we're done. Suppose that K is infinite. Let  $\varphi_1, \ldots, \varphi_n$  denote the distinct embeddings of  $K(\alpha, \beta)$  in  $\overline{K}$  over K. Consider

$$f(x) = \prod_{i \neq j} (\varphi_i(\alpha) + x\varphi_i(\beta) - \varphi_j(\alpha) - x\varphi_j(\beta)).$$

Since this is not the zero polynomial, there is some  $c \in K$  such that  $f(c) \neq 0$ . It follows that the  $\varphi_i(\alpha + c\beta)$  are pairwise distinct in  $\overline{K}$ . Then  $[K(\alpha + c\beta) : K] \geq n$ . But  $[K(\alpha, \beta) : K] = n$ , so that  $K(\alpha, \beta) = K(\alpha + c\beta)$ .

Let K be a field and  $f(x) \in K[x]$  be a monic separable polynomial. Let L denote the splitting field of f, so that  $L \supset K$  is Galois. Let  $G_f := \operatorname{Gal}(L/K) \subset S_n$  where  $n = \deg f$ . Let  $\operatorname{char} K \neq 2$ .

**Theorem 10.1.3.**  $L^{G_f \cap A_n} = K(\Delta(f))$  where  $\Delta(f) = \prod_{i < j} (\lambda_i - \lambda_j)$  and  $\lambda_1, \ldots, \lambda_n$  denote the distinct roots of f.

Before proving this, note that  $\Delta(f)$  is a square root of  $\mathrm{Discr}(f) \in K$ .

*Proof.* Consider  $x_1, \ldots, x_n$  purely transcendental elements over K. Let  $K(x_1, \ldots, x_n) \supset K$  be the corresponding extension. There is a group homomorphism  $\Phi: S_n \to \operatorname{Gal}(K(x_1, \ldots, x_n)/K)$  given by  $\sigma \mapsto \Phi_{\sigma}$  where

$$\Phi_{\sigma}(f)(x_1,\ldots,x_n) = f(x_{\sigma^{-1}(1)},\ldots,x_{\sigma^{-1}(n)}).$$

This is injective, and  $K(x_1,\ldots,x_n)^{S_n}=K(\sigma_1,\ldots,\sigma_n)$  where  $\sigma_1,\ldots,\sigma_n\in K[x_1,\ldots,x_n]$  are the alternating symmetric polynomials. Further,  $\mathrm{Gal}(K(x_1,\ldots,x_n)/K(\sigma_1,\ldots,\sigma_n))=S_n$ . Let  $\Delta_n=\prod_{i< j}(x_i-x_j)\in K(x_1,\ldots,x_n)$ . Then  $\Phi_{\sigma}(\Delta_n)=\mathrm{sgn}(\sigma)\Delta_n$ , and  $\Delta_n\notin K(\sigma_1,\ldots,\sigma_n)$ .

Define ev:  $K(x_1, \ldots, x_n) \to L$  by  $x_i \mapsto \lambda_i$ . Then ev  $\circ \Phi_{\sigma} = \sigma^{-1} \circ \text{ev}$ . Thus,

$$\operatorname{ev}(\Delta(f)) = \operatorname{ev}(\Phi_{\sigma}(\Delta_n)) = \sigma^{-1}(\Delta(f)).$$

This shows that the subgroup in  $G_f$  fixing  $\Delta(f)$  is precisely  $G_f \cap A_n$ .

Corollary 10.1.4. If char  $K \neq 2$  and f(x) is monic and separable over K, then  $G_f \subset A_n$  if and only if  $\operatorname{Discr}(f) \in K^2$ .

## 10.2 Lecture 23

**Theorem 10.2.1.** Suppose that K is a field and  $f(x) \in K[x]$  is separable. Then f is irreducible if and only if the Galois group  $G_f$  acts transitively on the set of roots of f.

Proof.

 $(\Longrightarrow)$ 

For any two roots  $\lambda_i, \lambda_j$  of f, we have that  $K(\lambda_I) \cong K(\lambda_j)$  as fields over K because both  $\operatorname{ev}_{\lambda_i} : K[x] \to K(\lambda_i)$  and  $\operatorname{ev}_{\lambda_j} : K[x] \to K(\lambda_j)$  induces isomorphisms with  $K[x]_{f}$ . By Lemma 4.2.8, we

can extend this isomorphisms to an automorphism  $\sigma: L \to L$  of the splitting field L for f. Thus,  $\sigma \in \operatorname{Gal}(L/K)$  with  $\sigma(\lambda_i) = \lambda_j$ .

 $(\Longleftrightarrow)$ 

Let  $\{\lambda_1, \ldots, \lambda_n\}$  denote the set of roots of f. Let f(x) = g(x)h(x) where  $\deg g \geq 1$  and g is irreducible. We must show that h is constant. Let  $\lambda$  be any root of g. Then there exists  $\sigma_i \in G_f$  such that  $\sigma_i(\lambda) = \lambda_i$  for each  $i = 1, \ldots, n$ . Note that

$$g(\lambda_i) = g(\sigma_i(\lambda)) = \sigma_i(g(\lambda)) = 0,$$

so that each  $\lambda_i$  is a root of g. Hence  $f \mid g$ , which implies that h is constant.

**Theorem 10.2.2.** Suppose that p is prime and that  $f(x) \in \mathbb{Q}[x]$  is monic and irreducible with  $\deg f = p$ . Suppose that f has exactly two non-real roots in  $\mathbb{C}$ . Then  $G_f = S_p$ .

*Proof.* Let L be the splitting field for f(x). Write  $f(x) = \prod_{i=1}^{p} (x - x_i)$  with each  $\lambda_i \in \mathbb{C}$ . Then  $\mathbb{Q}(\lambda_1, \ldots, \lambda_p) \subset \mathbb{C}$ . We see that

$$\mathbb{Q} \subset \mathbb{Q}(\lambda_i) \subset \mathbb{Q}(\lambda_1, \dots, \lambda_p) \subset \mathbb{C},$$

so that  $[\mathbb{Q}(\lambda_i):\mathbb{Q}] \mid [L:\mathbb{Q}]$ . Since  $p \mid [L:\mathbb{Q}] = |G_f| \subset S_p$ , it follows from Sylow that  $G_p$  contains an element of order p, i.e., that  $G_f$  contains a p-cycle. Also, the element in  $G_f$  that switches the roots is the complex conjugate pair of a transposition.

**Theorem 10.2.3 (Brouwer).** For any prime  $p \geq 5$ , there are infinitely many polynomials in  $\mathbb{Q}[x]$  of degree p with Galois group  $S_p$ .

*Proof.* Let k be an odd integer and let  $0 \le m, n_1 \le n_2 < \cdots < n_{k-2}$  be even integers. Consider

$$g(x) = (x^2 + m) (x - n_1) (x - n_2) \cdots (x - n_{k-2}).$$

This polynomial has  $\frac{k-3}{2}$  local maxima. Also, for each odd  $h \in \mathbb{Z}$ , |g(h)| > 2. Hence if c denotes a local maximum of g, then g(c) > 2. This shows that if f(x) = g(x) - 2, then there are

- $\frac{k-3}{2}$  positive local maxima in  $[n_1, n_{k-2}]$  and
- $\frac{k-3}{2}$  negative local maxima in  $[n_1, n_{k-2}]$ .

It follows that f(x) has k-3 real roots in  $[n_1, n_{k-2}]$  with  $f(n_{k-2}) = -2$  and  $\lim_{x\to\infty} f(x) > 0$ . Therefore, we have another real roots  $> n_{k-2}$ . Hence f(x) has at least k-2 real roots. Let  $\lambda, \ldots, \lambda_n \in \mathbb{C}$  denote the distinct roots of f. Then

$$\prod_{i=1}^{k} (x - \lambda_i) = f(x) = (x^2 + m) (x - n_1) (x - n_2) \cdots (x - n_{k-2}) - 2$$

, and  $-\sum_{i=1}^k \lambda_i = -\sum_{i=1}^{k-2} n_i$ . From this, we compute

$$\begin{split} \sum_{i < j} \lambda_i \lambda_j &= m + \sum_{a < b} n_a n_b \\ \sum_{i = 1}^k \lambda_i^2 &= \left(\sum_{i = 1}^k \lambda_i\right)^2 = \sum_{i < j} \lambda_i \lambda_j \\ &= \left(\sum_{i = 1}^{k - 2} n_i\right)^2 - 2m - 2\left(\sum_{a < b} n_a n_b\right) \\ &= \sum_{i = 1}^{k - 2} n_i^2 - 2m. \end{split}$$

Choose  $m \gg \sum n_i^2$  so that  $\sum_{i=1}^k \lambda_i^2 < 0$ . This implies that there exists a non-real root. Hence we must have exactly two real roots. Further, we can write  $f(x) = x^k + a_1 x^{k-1} + \cdots + a_{k-1} x + a_k$  with each  $a_i \in 2\mathbb{Z}$ . Since  $a_k = f(0) = g(0) - 2$ , we see that  $2 \mid a_{k-1}$  but  $4 \nmid a_{k-1}$ . By Eisenstein's criterion, f must be irreducible. We thus get infinitely many f's such that  $G_f = S_p$ .

# 11 Chain complexes and chain maps

The originators of homological algebra include Betti, Poincaré, and Riemann. The main goal of this subject is to extract invariants from topological spaces. Decompose X into contractible pieces (such as cells or simplices) to reduce X to combinatorial data. Specifically, reduce X to a collection of pieces of various dimensions where the boundary of a piece of dimension n is glued to a sub-collection of pieces of dimension n-1.

Emmy Noether introduced groups of chains  $C_i(X)$ , a free abelian group generated by the collection of *i*-dimensional pieces, equipped with boundary relations  $\partial_i: C_i(X) \to C_{i-1}(X)$ . From this, we obtain abelian groups  $H_i(X) \equiv \ker \partial_i /_{\text{im } \partial_{i+1}}$ , which are algebraic invariants of X.

Hilbert wanted to extract numerical invariants from a module. Specifically, if k is a field and  $K := k[x_1, \ldots, x_n]$ , then he wanted to understand the complexity of a module over K (or, more generally, any graded module over k).

A typical graded module over R will be a module of the form MR/I where  $I \subseteq R$  is a homogeneous ideal. By the Hilbert basis theorem,  $I \subseteq R$  is generated by finitely many homogeneous polynomials  $f_1, f_2, \ldots, f_{r_0}$ . Thus, we have surjective map  $\psi: R^{\oplus r_o} \to I$  given by  $(a_1, \ldots, a_{r_0}) \mapsto \sum a_i f_i$ . But, there generators are not, in general, independent. Therefore, we consider the module of relations  $Z_0(I) \equiv \ker \psi$  among the  $f_i$ . Note that  $Z_0(I)$  is finitely generated. We can choose generators and get a map  $\psi': R^{\oplus r_1} \to Z_0(I)$ . Then

$$R^{\oplus r_1} \to R^{\oplus r_0} \to I \to 0$$

is an exact sequence of graded R-modules. If  $Z_1(I) \equiv \ker \psi'$  is not zero, then choose generators again to get a map  $\psi'': R^{\oplus r_2} \to Z_1(I)$ . Continuing in this way, we get an exact sequence

$$\cdots \to R^{\oplus r_2} \to R^{\oplus r_1} \to R^{\oplus r_0} \to I \to 0.$$

The length of this sequence is defined to be  $\max\{i \mid r_i \neq 0\}$ . This is an invariant of I and of R/I.

Theorem 11.0.1 (Hilbert's syzygy theorem). Hilbert's syzygy theorem states that  $Z_{n-1}(I)$  is free, i.e., that there is an exact sequence of graded R-modules

$$0 \to R^{\oplus r_n} \to R^{\oplus r_{n-1}} \to \cdots \to R^{\oplus r_0} \to I \to 0.$$

#### 11.1 Lecture 24

#### Definition 11.1.1.

1. A chain complex (in **Ab**) is a pair  $(M_{\bullet}, \partial_{\bullet})$  where  $M_{\bullet} = \{M_i\}_{i \in \mathbb{Z}}$  is a set of abelian groups and  $\partial_{\bullet} = \{\partial_i\}_{i \in \mathbb{Z}}$  is a set of morphisms in **Ab** such that the *i-th differential*  $\partial_i : M_i \to M_{i-1}$  satisfies  $\partial_{i-1} \circ \partial_i = 0$ .

We call  $Z_n \equiv \ker \partial_n$  the group of degree n cycles and  $B_n \equiv \operatorname{im} \partial_{n+1}$  the group of degree n boundaries. Finally, we call  $H_n \equiv Z_n/B_n$  the degree n homology group.

2. A ((co)chain) complex (in  $\mathbf{Ab}$ ) is a pair  $(M^{\bullet}, d^{\bullet})$  where  $M^{\bullet} = \{M^{i}\}_{i \in \mathbb{Z}}$  is a set of abelian groups and  $d^{\bullet} = \{d^{i}\}_{i \in \mathbb{Z}}$  is a set of morphisms in  $\mathbf{Ab}$  such that the *i-th differential*  $d^{i}: M^{i} \to M^{i+1}$  satisfies  $d^{i+1} \circ d^{i} = 0$ .

We call  $Z^n \equiv \ker d^n$  the group of degree n cocycles and  $B^n \equiv \operatorname{im} d^{n-1}$  the group of degree n coboundaries.

Finally, we call  $H^n \equiv Z^n/B^n$  the degree n cohomology group.

**Definition 11.1.2.** Let  $(A^{\bullet}, d_A^i)$  and  $(B^{\bullet}, d_B^i)$  be complexes. A chain map  $f^{\bullet}: (A^{\bullet}, d_A^{\bullet}) \to (B^{\bullet}, d_B^{\bullet})$  consists of group homomorphisms  $f^i: A^i \to B^i$  for each  $i \in \mathbb{Z}$  such that  $d_B^i \circ f^i = f^{i+1} \circ d_A^i$ .

#### Note 11.1.3.

- 1. Any chain map  $f^{\bullet}: (A^{\bullet}, d_A^{\bullet}) \to (B^{\bullet}, d_B^{\bullet})$  restricts term-wise to maps  $f^i: Z^i(A^{\bullet}) \to Z^i(B^{\bullet})$  and maps  $f^i: B^i(A^{\bullet}) \to B^i(B^{\bullet})$ . Thus, it induces a map  $f^*: H^i(A^{\bullet}) \to H^i(B^{\bullet})$ .
- 2. We have a natural isomorphism  $\mathbf{Ch}(\mathbf{Ab}) \to \mathbf{CoCh}(\mathbf{Ab})$  given by  $N_i \mapsto M^{-i}$  and  $\partial_i \mapsto d^{-i}$ .

**Definition 11.1.4.** We say that  $(A^{\bullet}, d^{\bullet})$  is bounded above if there is some N such that  $A^n = 0$  for any  $n \geq N$ . We define bounded below similarly. We say that  $(A^{\bullet}, d^{\bullet})$  is bounded if it is both bounded above and bounded below.

As a result, we have the subcategories  $CoCh^{-}(Ab)$ ,  $CoCh^{+}(Ab)$ , and  $CoCh^{b}(Ab)$ , respectively.

If  $C^{\bullet} = \bigoplus_{i \in \mathbb{Z}} C^i$  is a graded abelian group, then it induces a natural complex  $(\underline{C}^{\bullet}, 0)$  where  $\underline{C}^i \equiv C^i$ . In particular, any abelian group may be viewed as a complex.

Conversely, given a complex  $(M^{\bullet}, d^{\bullet})$ , we can form the graded abelian group  $M^{\bullet} \equiv \bigoplus_{i \in \mathbb{Z}} M^i$  and package the differential  $d^i$  into a single group map  $D: M^{\bullet} \to M^{\bullet}$  such that  $D \upharpoonright_{M^i} = d^i$  and  $D^2 = 0$ .

We can write D as the block diagonal matrix

$$\begin{bmatrix} 0 & & & & & \\ d^i & 0 & & & & \\ & d^{i+1} & 0 & & & \\ & & d^{i+2} & 0 & & \\ & & \ddots & \ddots \end{bmatrix}.$$

As a result, we obtain the *cochain functor* given by  $(A^{\bullet}, d^{\bullet}) \to \bigoplus_{i \in \mathbb{Z}} A^i$  and  $f^{\bullet} \mapsto (f^i)_{i \in \mathbb{Z}}$ .

**Definition 11.1.5.** We say that  $(A^{\bullet}, d^{\bullet})$  is acyclic or exact if  $H^{\bullet}(A^{\bullet}, d^{\bullet}) = 0$ .

**Theorem 11.1.6.** Let  $K^{\bullet}$  be an exact complex of R-modules and  $I^{\bullet}$  a bounded below complex of injective R-modules. Any chain map  $f: K^{\bullet} \to I^{\bullet}$  is homotopic to zero.

Proof. By hypothesis, there is some  $r \in \mathbb{Z}$  such that  $I^k = 0$  for any k < r. Then  $f^k = 0$  for any k < r. Define  $h^k : K^k \to I^{k-1}$  by  $h^k = 0$  for each  $k \le r$ . Then  $f^k = 0 = d_I h^k + h^{k+1} d_K$  for any k < r. Let s > r and assume, for induction, that, for each k < s, we have constructed a map  $h^k : K^k \to I^{k-1}$  such that  $f^{k-1} = d_I h^{k-1} + h^k d_K$ . We must construct a map  $h^s : K^s \to I^{s-1}$  such that  $f^{s-1} = d_I h^{s-1} + h^s d_K$ .

Let  $g^{s-1} = f^{s-1} - d_I h^{s-1}$ . Note that

$$g^{s-1}d_K = (f^{s-1} - d_I h^{s-1}) d_K$$

$$= f^{s-1}d_K - d_I h^{s-1}d_K$$

$$= d_I f^{s-2} - d_I (f^{s-2} - d_I h^{s-2})$$

$$= 0.$$

Therefore,  $g^{s-1}$  descends to a map  $g^{s-1}: K^{s-1}/_{\operatorname{im} d_K} \to I^{s-1}$ . Since  $K^{\bullet}$  is exact, we have

$$g^{s-1}: \overset{K^{s-1}}{/_{\ker d_K}} \to I^{s-1}.$$

Moreover, since  $I^{s-1}$  is injective, we can find some map  $h^s: K^s \to I^{s-1}$  such that

$$I^{s-1}$$

$$g^{s-1} \uparrow \qquad \qquad h^{s}$$

$$K^{s-1} / \ker d_{K} \xrightarrow{\cong} \operatorname{im} d_{K} \stackrel{\longleftarrow}{\longleftrightarrow} K^{s}$$

commutes. Hence  $h^s d_K = g^{s-1}$ . It follows that

$$\begin{split} d_I h^{s-1} + h^s d_K &= d_I h^{s-1} + g^{s-1} \\ &= d_I h^{s-1} + f^{s-1} - d_I h^{s-1} \\ &= f^{s-1}, \end{split}$$

as desired  $\Box$ 

**Definition 11.1.7.** If A is an abelian group, then a *left resolution of* A is an exact complex  $(C^{\bullet}, d^{\bullet}) \in$  ob  $\mathbf{CoCh}^{\leq 0}(\mathbf{Ab})$  of the form

$$\cdots \to C^{i-1} \to C^i \to \cdots \to C^0 \to A \to 0.$$

**Example 11.1.8.** If  $I \subseteq k[x_1, \ldots, x_n]$  is a homogenous ideal, then Hilbert's syzygy theorem says that I has a left resolution of length n+1 with n+1 terms free finitely generated R-modules.

Let  $a \in \mathbb{Z}$ . Define the *shift functor* 

$$-[a]: \mathbf{CoCh}(\mathbf{Ab}) \to \mathbf{CoCh}(\mathbf{Ab})$$

as follows. Let  $(M^{\bullet}, d_{M}^{\bullet})$  be a complex. Form the pair  $\left(M^{\bullet}[a], d_{M[a]}^{\bullet}\right)$  where  $(M^{\bullet}[a])^{n} \equiv M^{a+n}$  and  $\left(d_{M[a]}\right)^{n} \equiv (-1)^{a} d_{M}^{a+n}$ . If  $f^{\bullet}$  is a chain map, then let  $(f^{\bullet}[a])^{n} = f^{a+n}$ .

Proposition 11.1.9. The shift functor is an equivalence that preserves  $CoCh^{-}(Ab)$ ,  $CoCh^{+}(Ab)$ , and  $CoCh^{b}(Ab)$ .

**Definition 11.1.10.** Let  $f: M \to N$  be a chain map. Form  $\operatorname{cone}(f)$  the *cone of* f as a new complex where  $\operatorname{cone}(f)^{\bullet} \equiv N \oplus M[1]$  and  $d_{\operatorname{cone}(f)}^{\bullet} \equiv \begin{bmatrix} d_N & f \\ 0 & d_{M[1]} \end{bmatrix}$ .

We see that

$$cone(f)^n = N^n \oplus M^{n+1}$$

and  $d_{\text{cone}(f)}^n: N^n \oplus M^{n+1} \to N^{n+1} \oplus M^{n+2}$  with

$$d_{\operatorname{cone}(f)}^{n} = \begin{bmatrix} d_{N}^{n} & f^{n+1} \\ 0 & -d_{M}^{n+1} \end{bmatrix}.$$

Exercise 11.1.11. Show that  $d_{\text{cone}(f)}^{i+1} \circ d_{\text{cone}(f)}^{i} = 0$ .

## Definition 11.1.12.

1. A double complex is a triple  $(A^{\bullet,\bullet}, d^{\bullet}, \delta^{\bullet})$  where  $A^{i,j} = \{A^{i,j}\}_{(i,j)\in\mathbb{Z}^2}$  and both  $d: A^{\bullet,\bullet} \to A^{\bullet+1,\bullet}$  and  $\delta: A^{\bullet,\bullet} \to A^{\bullet,\bullet+1}$  are homomorphisms such that  $d\delta = \delta d$  and  $d^2 = \delta^2 = 0$ . As a commutative diagram, this has the form

2. The total complex of  $(A^{\bullet,\bullet}, d^{\bullet}, \delta^{\bullet})$  is the complex Tot(A) where  $\text{Tot}(A)^n \equiv \bigoplus_{p+q=n} A^{p,q}$  and  $d_{\text{Tot}(A)} \upharpoonright_{A^{p,q}} \equiv d + (-1)^p \delta$ .

**Proposition 11.1.13.** Any chain map  $f: M \to N$  induces a double complex

$$M^{i-1,0} \xrightarrow{d_M} M^{i,0} \xrightarrow{d_M} M^{i+1,0} \xrightarrow{d_M} M^{i+2,0} \xrightarrow{d_M} \cdots$$

$$\downarrow f \qquad \qquad \downarrow f \qquad \qquad \downarrow f \qquad \qquad \downarrow f$$

$$N^{i-1,1} \xrightarrow{d_N} N^{i,1} \xrightarrow{d_N} N^{i+1,1} \xrightarrow{d_N} N^{i+2,1} \xrightarrow{d_N} \cdots$$

The total complex of this is precisely cone(f).

Let N and C be complexes. Suppose that  $C \xrightarrow{\iota} N$  is a chain map where each  $\iota^n : N^n \to C^n$  is injective. Let  $s^n : C^n \to N^n$  be a group homomorphism such that  $s^n \circ \iota^n = \mathrm{id}_{N^n}$ . Then  $M := \left(C/N, d_{C/N}\right)$  is a complex. Our choice of  $s^n$  produces a splitting  $C^{\bullet} \cong N^{\bullet} \oplus M^{\bullet}[1]$  in the category of graded abelian groups. Thus, we have the map  $d_C = \begin{bmatrix} d_N & f \\ 0 & d_{M[1]} \end{bmatrix}$  where  $f: M \to N$  is a map of graded abelian groups.

**Exercise 11.1.14.** Show that f is a chain map and  $C \cong \text{cone}(f)$ .

#### 11.2 Lecture 25

**Definition 11.2.1.** Let  $f, g: A^{\bullet} \to B^{\bullet}$  be two chain maps. A homotopy between f and g is a map of graded abelian groups  $h: A^{\bullet} \to B^{\bullet - 1}$  such the

$$d_B h + h d_A = f - g.$$

We say that f and g are homotopy equivalent (written as  $f \sim g$ ) if there is a homotopy between them.

#### Proposition 11.2.2.

- 1. Homotopy is an equivalence relation.
- 2. The class mor  $^{\circ}$  CoCh(Ab) of all chain maps homotopic to 0 is a two-sided ideal in mor CoCh(Ab).
- 3. If  $f \simeq g : A^{\bullet} \to B^{\bullet}$ , then  $H^{\bullet}(f) = H({}^{\bullet}g)$ .
- 4. If  $f \simeq g$  and c is a cocycle, then  $f(c) g(c) = d_B h(c)$ , which is a coboundary.

Notation. Let  $C(\mathbf{Ab})$  denote the category with complexes as objects and homotopy classes of chain maps as morphisms.

### Note 11.2.3.

1. We have that  $\operatorname{Hom}_{\mathcal{C}(\mathbf{Ab})}(A, B) = \frac{\operatorname{Hom}_{\mathbf{CoCh}(\mathbf{Ab})}(A, B)}{\operatorname{Hom}_{\mathbf{CoCh}(\mathbf{Ab})}^{\sim 0}(A, B)}$ .

2.  $H^{\bullet}$  descends to a well-defined functor in the sense that the diagram

$$\begin{array}{c} \mathbf{CoCh}(\mathbf{Ab}) \longrightarrow \mathcal{C}(\mathbf{Ab}) \\ \\ H^{\bullet} \downarrow \\ \mathbf{grAb} \end{array}.$$

commutes.

**Definition 11.2.4.** A short exact sequence of complex is a sequence  $0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0$  of complexes such that each sequence

$$0 \to A^n \xrightarrow{f^n} B^n \xrightarrow{g^n} C^n \to 0$$

is exact in Ab.

Let

$$0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0$$

be a short exact sequence of complexes. Consider the commutative diagram

$$0 \longrightarrow A^{n-1} \xrightarrow{f^{n-1}} B^{n-1} \xrightarrow{g^{n-1}} C^{n-1} \longrightarrow 0$$

$$\downarrow d_A^{n-1} \downarrow \qquad \downarrow d_B^{n-1} \qquad \downarrow d_C^{n-1} \downarrow$$

$$0 \longrightarrow A^n \xrightarrow{f^n} B^n \xrightarrow{g^n} C^n \longrightarrow 0$$

$$\downarrow d_A^n \downarrow \qquad \downarrow d_B^n \qquad \downarrow d_C^n$$

$$0 \longrightarrow A^{n+1} \xrightarrow{f^{n+1}} B^{n+1} \xrightarrow{g^{n+1}} C^{n+1} \longrightarrow 0$$

Define a collection of edge homomorphisms  $\left\{\delta^n: H^n(C) \to H^{n+1}(A)\right\}_{n \in \mathbb{Z}}$  as follows. Let  $c \in C^n$  with  $d^n_C(c) = 0$ . By exactness, there is some  $b \in B^n$  such that  $g^n(b) = c$ . But then

$$d_B^n(b) \in \ker g^{n+1} = \operatorname{im} f^{n+1}.$$

Since  $f^{n+1}$  is injective, this means that there is a unique  $a \in A^{n+1}$  such that  $f^n(a) = d_B^n(b)$ . Let  $\delta^n([c]) = [a]$ .

**Exercise 11.2.5.** Check that  $\delta^n$  is a homomorphism and that it is independent both of our choice of c and of our choice of b.

Lemma 11.2.6 (Snake). Any short exact sequence

$$0 \to A \xrightarrow{f} B \xrightarrow{g} C \to 0$$

complexes induces a long exact sequence in cohomology

$$H^{n}(A) \xrightarrow{\delta^{n-1}} H^{n}(C)$$

$$H^{n}(A) \xrightarrow{f^{*}} H^{n}(B) \xrightarrow{g^{*}} H^{n}(C)$$

$$H^{n+1}(A) \xrightarrow{f^{*}} H^{n+1}(B) \xrightarrow{\delta^{n}} \cdots$$

Proof.

Exactness at  $H^n(B)$ : We have that  $0_{H^n(C)} = H^n(0) = H^n(g \circ f) = H^n(g) \circ H^n(f)$ . Hence  $\lim_{h \to \infty} H^n(f) \subset \ker_{h} H^n(g)$ .

For the reverse inclusion, let  $[b] \in \ker H^n(g) \subset H^n(B)$ . Then  $g(b) \in C^n$  must be a coboundary, so that there is some  $c \in C^{n-1}$  such that  $g(b) = d_C c$ . Choose a lift  $b_1 \in B^{n-1}$  of c, meaning that  $g(b_1) = c$ . Then  $b - d_B b_1 \in Z^n(B)$ , and  $[b] = [b - d_B b_1]$ . But

$$g(b - d_B b_1) = g(b) - g(d_B b_1) = g(b) - d_C g(b_1) = g(b) - d_C c = 0.$$

Hence  $b-d_Bb_1 \in \ker g \subset B^n$ . This implies that there exists a unique  $a \in A^n$  such that  $b-d_Bb_1 = f(a)$ . Also,

$$f(d_A a) = d_B(f(a)) = d_B(b - d_B b_1) = 0.$$

Since f is injective, we see that  $d_A a = 0$ , i.e.,  $a \in Z^n(A)$ . Thus,  $H^n(f)([a]) = [f(a)] = [b - d_B b_1] = [b]$ . This proves that  $[b] \in \operatorname{im} H^n(f)$ .

Exactness at  $H^n(C)$ : Let  $[b] \in H^n(B)$ . Note that  $\delta^n(H^n(g)([b])) = [a]$  where  $a \in A^{n+1}$  denotes the unique element such that  $f(a) = d_B b$ . Since  $d_B b = 0$  and f is injective, it follows that a = 0. Hence im  $H^n(g) \subset \ker \delta^n$ .

Conversely, let  $[c] \in \ker \delta^n$ . Choose  $b \in B^n$  such that g(b) = c and then the unique  $a \in A^{n+1}$  such that  $f(a) = d_B b$ . Thus,  $\delta^n([c]) = [a] = 0$ , so that  $a \in B^{n+1}(A)$ , i.e.,  $d_A a_1 = a$  for some  $a_1 \in A^n$ . Note that  $g(b - f(a_1)) = g(b) - g(f(a_1)) = c - 0 = c$ . Further,

$$d_B(b - f(a_1)) = d_B(b) - d_B(f(a_1))$$

$$= f(a) - f(d_A a_1)$$

$$= f(a) - f(a)$$

$$= 0.$$

This shows that  $b - f(a_1)$  is a cocycle. Thus,  $H^n(g)([b - f(a_1)]) = [g(b - f(a_1))] = [c]$ , so that  $[c] \in \operatorname{im} H^n(g)$ .

Exactness at  $H^{n+1}(A)$ : Let  $[c] \in H^n(C)$  and find  $[a] = \delta^n([c])$ , where

$$\begin{array}{c}
b \xrightarrow{g} c \\
\downarrow \\
a \xrightarrow{f} d_B b
\end{array}$$

Then  $H^{n+1}(f)([a]) = [f(a)] = [d_B b] = 0$ . It follows that im  $\delta^n \subset \ker H^{n+1}(f)$ .

Conversely, let  $[a] \in \ker H^{n+1}(f)$ , so that  $H^{n+1}(f)([a]) = [f(a)] = 0$ . This means that  $f(a) = d_B b$  for some  $b \in B^n$ . Then  $\delta^n([g(b)]) = [a]$ . This shows that im  $\delta^n \supset \ker H^{n+1}(f)$ .

# 12 Additive categories

Definition 12.0.1.

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1. A category  $\mathscr C$  is enhanced over  $\mathbf A \mathbf b$  if  $\mathrm{Hom}_{\mathscr C}(a,b)$  is an abelian group for any  $a,b\in\mathrm{ob}\,\mathscr C$  and

$$\operatorname{Hom}_{\mathscr{C}}(y,z) \times \operatorname{Hom}_{\mathscr{C}}(x,y) \stackrel{\circ}{\longrightarrow} \operatorname{Hom}_{\mathscr{C}}(x,z)$$

is bilinear for any  $x, y, z \in \text{ob } \mathscr{C}$ .

2. A category  $\mathscr{C}$  is called *additive* if it is enhanced over  $\mathbf{Ab}$  and has finite products.

**Example 12.0.2.** The following are additive categories.

- 1. **Ab**.
- 2. *R*-**Mod**.

#### Note 12.0.3.

- 1. Let  $\mathscr{C}$  be category with finite products. The product of the empty diagram is the terminal object in  $\mathscr{C}$  since it is the initial object in **Set**.
- 2. If  $\mathscr{C}$  is additive and \* is the terminal object in  $\mathscr{C}$ , then  $\operatorname{Hom}_{\mathscr{C}}(*,*)$  consists of a single element, which must equal the group identity element.

Exercise 12.0.4. Verify the following statements.

- 1. If C is a additive, then its terminal object is also initial and thus is a zero object in C.
- 2. A zero object  $0_{\mathscr{C}}$  satisfies  $\operatorname{Hom}_{\mathscr{C}}(x,0_{\mathscr{C}})=0$  and  $\operatorname{Hom}_{\mathscr{C}}(0_{\mathscr{C}},x)=0$  for any  $x\in\operatorname{ob}\mathscr{C}$ .
- 3. Any additive category has finite coproducts that are equal to finite products.

### 12.1 Lecture 26

**Definition 12.1.1.** Let  $\mathscr C$  be an additive category. Let  $f: x \to y$  be a morphism in  $\mathscr C$ .

1. A kernel (object) for f is a pair (k,q) where  $k \in \text{ob}\,\mathscr{C}$  and  $q: k \to x$  such that for any  $z \in \text{ob}\,\mathscr{C}$ , the natural sequence

$$\operatorname{Hom}(z,k) \xrightarrow{q \circ -} \operatorname{Hom}(z,x) \xrightarrow{f \circ -} \operatorname{Hom}(z,y)$$

is exact.

2. A cokernel (object) for f is a pair (c, p) where  $c \in \text{ob } \mathscr{C}$  and  $p : y \to c$  such that for any  $z \in \text{ob } \mathscr{C}$ , the natural sequence

$$\operatorname{Hom}(c,z) \xrightarrow{-\circ p} \operatorname{Hom}(y,z) \xrightarrow{-\circ f} \operatorname{Hom}(x,z)$$

is exact.

**Definition 12.1.2.** We say that a category  $\mathscr{A}$  is abelian if

1.  $\mathscr{A}$  is additive and

- 2. for any morphism  $f: x \to y$  in  $\mathscr{A}$ , there exists a sequence  $k \xrightarrow{q} x \xrightarrow{a} i \xrightarrow{b} y \xrightarrow{p} c$  in  $\mathscr{A}$  such that
  - (a) (k,q) is a kernel for f,
  - (b) (c, p) is a cokernel for f,
  - (c) (c, a) is a cokernel for q, and
  - (d) (i, b) is a kernel for p.

We call i the image of f.

**Definition 12.1.3.** If  $\mathscr{A}$  is a abelian, then a sequence  $x \xrightarrow{f} y \xrightarrow{g} z$  in  $\mathscr{A}$  is exact if im  $f = \ker g$ .

## Example 12.1.4.

- 1. **Ab**.
- 2. *R*-**Mod**.
- 3. **PreShAb**X where X is a space.

*Remark* 12.1.5. Our notion of and results for cohomology for complexes of abelian groups hold for complexes of objects in an abelian category.

**Theorem 12.1.6 (Freyd-Mitchell).** Every abelian category admits a fully faithful embedding into R-Mod for some ring R.

Remark 12.1.7. It is not, in general, possible to complete an additive category  $\mathscr{C}$  to an abelian one. Still, we can always add enough images to  $\mathscr{C}$  to get cones of maps of complexes.

Let  $\mathscr{C}$  be additive. A map  $e: x \to x$  in  $\mathscr{C}$  is an *idempotent* if  $e^2 = e$ . Let  $\mathscr{C} = \mathbf{Vect}_k$ . Then an idempotent map  $e: x \to x$  is a projection map, i.e.,  $x = x_1 \oplus x_2$  such that  $e = i_1 \circ p_1$ .

If  $\mathscr C$  is additive and  $e: x \to x$  is idempotent in  $\mathscr C$ , then we say that e has an image in  $\mathscr C$  if there exists a decomposition  $x = x_1 \oplus x_2$  such that

$$e = \begin{bmatrix} i\mathbf{d}_{x_1} & 0\\ 0 & 0 \end{bmatrix}$$

with respect to this decomposition. We say that  $x_1$  is the *image of e*.

Let  $e: x \to x$  be an idempotent. Then  $\mathrm{id}_x e: x \to x$  is also an idempotent. Indeed,

$$(\mathrm{id}_x - e)^2 = \mathrm{id}_x^2 - \mathrm{id}_x e - e\mathrm{id}_x + e^2 = \mathrm{id}_x = e.$$

If 
$$x = x_1 \oplus x_2$$
 has  $e = \begin{bmatrix} id_{x_1} & 0 \\ 0 & 0 \end{bmatrix}$ , then  $id_x - e = \begin{bmatrix} 0 & 0 \\ 0 & id_{x_2} \end{bmatrix}$ , so that  $id_x - e$  has  $x_2$  as an image.

**Definition 12.1.8.** A category  $\mathscr{C}$  is *idempotent complete* or *Karoubian* if  $\mathscr{C}$  is additive and any idempotent in  $\mathscr{C}$  has an image in  $\mathscr{C}$ .

**Exercise 12.1.9.** Show that for any additive category  $\mathscr{C}$ , there exists a unique (up to unique isomorphism) category  $\mathscr{C}^{\text{Kor}}$  together with a functor  $F:\mathscr{C}\to\mathscr{C}^{\text{Kor}}$  such that

- 1.  $\mathscr{C}^{\mathrm{Kor}}$  is idempotent complete,
- 2. F is fully faithful, and
- 3. every object in  $\mathscr{C}^{\mathrm{Kor}}$  is an image of an idempotent in  $\mathscr{C}$ .

**Definition 12.1.10.** A graded additive category is an additive category  $\mathscr C$  such that for any  $x,y\in \operatorname{ob}\mathscr C$ ,  $\operatorname{Hom}(x,y)$  is a graded abelian group, i.e.,  $\operatorname{Hom}(x,y)\cong\bigoplus_{n\in\mathbb Z}\operatorname{Hom}^n(x,y)$  and  $\operatorname{Hom}(x,y)\times \operatorname{Hom}(y,z)\stackrel{\circ}{\longrightarrow}\operatorname{Hom}(x,z)$  has the form  $\operatorname{Hom}^n(x,y)\times\operatorname{Hom}^m(y,z)\stackrel{\circ}{\longrightarrow}\operatorname{Hom}^{n+m}(x,z)$  where  $\circ$  is bilinear.

**Definition 12.1.11.** A graded additive category  $\mathscr C$  is a differential graded category if for any  $x,y\in \operatorname{ob}\mathscr C$ , the graded group  $\operatorname{Hom}(x,y)$  is equipped with with a homomorphism  $d:\operatorname{Hom}(x,y)\to\operatorname{Hom}(x,y)$  such that

- (a)  $d: \operatorname{Hom}^n(x, y) \to \operatorname{Hom}^{n+1}(x, y),$
- (b)  $d^2 = 0$ , and
- (c) d satisfies the graded Leibniz rule, i.e., if  $f \in \operatorname{Hom}^n(x,y)$  and  $g \in \operatorname{Hom}(a,x)$ , then

$$d(f \circ g) = df \circ g + (-1)^n f \circ dg.$$

Proposition 12.1.12. Let  $\mathscr{C}$  be a category.

- 1. If  $\mathscr{C}$  is additive, then for any  $x \in \operatorname{ob} \mathscr{C}$ ,  $\operatorname{Hom}(x,x)$  is a ring (in fact, a  $\mathbb{Z}$ -algebra).
- 2. If  $\mathscr{C}$  is a graded additive category, then for any  $x \in \operatorname{ob}\mathscr{C}$ ,  $\operatorname{End}(x)$  is a graded ring.
- 3. If  $\mathscr{C}$  is differential graded category, then for any  $x \in \text{ob}\mathscr{C}$ , End(x) is a differential graded algebra.

**Definition 12.1.13.** If  $\mathscr{C}$  is a differential graded category, then the *homotopy category of*  $\mathscr{C}$  is the category  $\operatorname{Ho}(\mathscr{C})$  (or  $[\mathscr{C}]$ ) given by

$$\begin{split} \operatorname{ob} \operatorname{Ho}(\mathscr{C}) &\equiv \operatorname{ob} \mathscr{C} \\ \operatorname{Hom}_{\operatorname{Ho}(\mathscr{C})}(x,y) &\equiv H^0(\operatorname{Hom}_{\mathscr{C}}(x,y),d) \\ &= \frac{\ker(\operatorname{Hom}_{\mathscr{C}}^0(x,y) \overset{d}{\longrightarrow} \operatorname{Hom}_{\mathscr{C}}^1(x,y))}{\operatorname{im}(\operatorname{Hom}_{\mathscr{C}}^{-1}(x,y) \overset{d}{\longrightarrow} \operatorname{Hom}_{\mathscr{C}}^0(x,y))}. \end{split}$$

Let  $\mathscr{B}$  be an additive category. Define the category  $\mathbf{Compl}(\mathscr{B})$  of complexes in  $\mathscr{B}$  by

ob 
$$\mathbf{Compl}(\mathscr{B}) = (\text{complexes of objects in } \mathscr{B})$$
  
mor  $\mathbf{Compl}(\mathscr{B}) = (\text{morphisms of complexes})$ .

This is an additive category. We can also refine this definition by incorporating degree-shifting maps to get a differential graded category of complexes in  $\mathscr{B}$ . Define the category  $\mathbf{Compl}^{\bullet}(\mathscr{B})$  by

$$\operatorname{ob}\mathbf{Compl}^{\bullet}(\mathscr{B}) = (\text{complexes of objects in }\mathscr{B})$$
 
$$\operatorname{Hom}_{\mathbf{Compl}^{\bullet}(\mathscr{B})}(M,N) = \bigoplus_{n \in \mathbb{Z}} \operatorname{Hom}^n(M,N)$$

where

$$\operatorname{Hom}^n(M,N) \equiv \prod_{a \in \mathbb{Z}} \operatorname{Hom}_{\mathscr{B}}(M^a,N^{a+n}).$$

The composition is obtained component-wise from the composition in  $\mathscr{B}$ . Define  $d: \operatorname{Hom}^n(M,N) \to \operatorname{Hom}^{n+1}(M,N)$  by

$$(f_a)_{a\in\mathbb{Z}}\to (d_N\circ f_a+(-1)^nf_{a+1}\circ d_M)_{a\in\mathbb{Z}}.$$

This makes  $Compl^{\bullet}(\mathcal{B})$  a differential graded category.

Let  $M, N \in \text{ob } \mathbf{Compl}^{\bullet}(\mathscr{B})$ . Then

$$\begin{split} Z^0(\operatorname{Hom}_{\mathbf{Compl}^{\bullet}(\mathscr{B})}(M,N)) &= \ker(\operatorname{Hom}^0 \stackrel{d}{\longrightarrow} \operatorname{Hom}^1) \\ &= \operatorname{Hom}_{\mathbf{Compl}(\mathscr{B})}(M,N) \\ B^0(\operatorname{Hom}_{\mathbf{Compl}^{\bullet}(\mathscr{B})}(M,N)) &= \operatorname{im}(\operatorname{Hom}^{-1} \stackrel{d}{\longrightarrow} \operatorname{Hom}^0) \\ &= (\operatorname{homotopies of 0-maps of complexes}) \,. \end{split}$$

Also, we have that

$$H^0(\operatorname{Hom}_{\mathbf{Compl}^{\bullet}(\mathscr{B})}(M,N)) = \text{(maps of complexes)/(homotopies)}.$$

Example 12.1.14.  $\operatorname{Ho}(\operatorname{Comp}^{\bullet}(\operatorname{Ab})) = \mathcal{C}(\operatorname{Ab}), \text{ and } Z^{0}(\operatorname{Comp}^{\bullet}(\operatorname{Ab})) = \operatorname{CoCh}(\operatorname{Ab}).$ 

# 13 Triangulated categories

#### 13.1 Lecture 27

Let  $\mathscr{C}$  be a category. For any  $x \in \text{ob}\,\mathscr{C}$ , define x[n] as the object, if it exists, in  $\mathscr{C}$  that represents the shift functor on morphisms  $\text{Hom}_{\mathscr{C}}(-,x)[n]:\mathscr{C}^{\text{op}}\to\mathbf{Compl}(\mathbf{Ab})$ . If  $f:x\to y$  is a morphism in  $\mathscr{C}$ , then define the cone cone(f) of f to be the object, it it exists, in  $\mathscr{C}$  that represents the functor  $\mathscr{C}^{\text{op}}\to\mathbf{Comp}(\mathbf{Ab})$  given by  $z\mapsto \text{cone}(\text{Hom}_{\mathscr{C}}(z,x)\xrightarrow{f\circ -}\text{Hom}_{\mathscr{C}}(z,y))$ .

**Definition 13.1.1.** A category  $\mathscr{C}$  is called *strongly pre-triangulated* if every object in  $\mathscr{C}$  has shifts in  $\mathscr{C}$  and every morphism in  $\mathscr{C}$  has cones in  $\mathscr{C}$ . We call  $\mathscr{C}$  *pre-triangulated* if every object in  $\mathscr{C}$  has shifts in  $Ho(\mathscr{C})$  and every morphism in  $\mathscr{C}$  has cones in  $Ho(\mathscr{C})$ .

**Note 13.1.2.** Both the assignment  $x \mapsto x[n]$  and the assignment  $f \mapsto \text{cone}(f)$  are functorial.

**Definition 13.1.3.** Given a differential graded category  $\mathscr{C}$ , we define  $\operatorname{Ho}^{\bullet}(\mathscr{C})$  as the graded additive category such that

$$\mathrm{ob}\,\mathrm{Ho}^\bullet(\mathscr{C})=\mathrm{ob}\,\mathscr{C}$$
 
$$\mathrm{Hom}_{\mathrm{Ho}^\bullet(\mathscr{C})}(x,y)=H^\bullet(\mathrm{Hom}_\mathscr{C}(x,y)).$$

If  $\mathscr{C}$  is strongly pre-triangulated, then  $\operatorname{Ho}^{\bullet}(\mathscr{C})$  and  $\operatorname{Ho}(\mathscr{C})$  contain the same information. Indeed,  $\operatorname{Ho}(\mathscr{C})$  is precisely the degree zero piece of  $\operatorname{Ho}^{\bullet}(\mathscr{C})$ . Conversely,  $ifx,y \in \operatorname{ob}\mathscr{C}$ , then

$$\operatorname{Hom}_{\operatorname{Ho}^{\bullet}(\mathscr{C})}(x,y) = \bigoplus_{a \in \mathbb{Z}} \operatorname{Hom}_{\operatorname{Ho}^{\bullet}(\mathscr{C})}^{a}(x,y)$$

where  $\operatorname{Hom}_{\operatorname{Ho}^{\bullet}(\mathscr{C})}^{a}(x,y) = \operatorname{Hom}_{\operatorname{Ho}(\mathscr{C})}(x,y[a]) = \operatorname{Hom}_{\operatorname{Ho}(\mathscr{C})}(x,y)[a].$ 

*Notation.* From now on, if  $\mathscr{C}$  is strongly pre-triangulated, then we write  $Ho(\mathscr{C})$  for the graded homotopy category.

**Definition 13.1.4.** If  $\mathscr{C}$  is strongly pre-triangulated, then a *triangle*  $\triangle$  *in* Ho( $\mathscr{C}$ ) is a sequence of degree zero maps  $x \xrightarrow{u} y \xrightarrow{v} z \xrightarrow{w} x[1]$ . We represent this as

$$x \xrightarrow{u} y \\ \downarrow_{v} \\ \downarrow_{v}.$$

Let  $\mathscr{C}$  be strongly pre-triangulated. Given a triangle

$$\begin{array}{c}
x \xrightarrow{u} y \\
\downarrow v, \\
z
\end{array}$$

we have a long sequence of maps

$$x[-1] \xrightarrow{u[-1]} y[-1] \xrightarrow{v[-1]} z[-1]$$

$$x \xrightarrow{w} y \xrightarrow{v} z$$

$$x[1] \xrightarrow{u[1]} y[1] \xrightarrow{v[1]} z[1] \xrightarrow{w[1]} \cdots$$

in  $\mathscr{C}$ .

**Definition 13.1.5.** Let  $\mathscr{C}$  be strongly pre-triangulated. We say that a triangle in  $\operatorname{Ho}(\mathscr{C})$  is *exact* if it is isomorphic to the triangle

$$x \xrightarrow{u} y \xrightarrow{\text{``inclusion''}} \text{cone}(u) \xrightarrow{\text{``projection''}} x[1]$$
.

**Definition 13.1.6.** A graded additive category  $\mathscr{D}$  is *triangulated* if  $\mathscr{D}$  is equipped with a shift functor  $[1]: \mathscr{D} \to \mathscr{D}$  and a collection of *distinguished triangles* such that the following axioms hold.

- (0) Every triangle that is isomorphic to a distinguished triangle is distinguished.
- (1) For any object x in  $\mathscr{D}$ , the triangle  $x \xrightarrow{\mathrm{id}_x} x \to 0 \to x[1]$  is distinguished.
- (2) (rotation invariance) The shift rotation of a triangle  $\triangle$  is distinguished if and only if  $\triangle$  is, i.e., the triangle  $x \xrightarrow{u} y \xrightarrow{v} z \xrightarrow{w} x[1]$  is distinguished if and only if the triangle

$$y \xrightarrow{v} z \xrightarrow{w} x[1] \xrightarrow{-u[1]} y[1]$$

is distinguished.

(3) Every morphism  $u: x \to y$  can be included in a distinguished triangle  $x \xrightarrow{u} y \xrightarrow{v} z \xrightarrow{w} x[1]$ , and every commutative square

$$\begin{array}{ccc}
x & \xrightarrow{u} & y \\
f \downarrow & & \downarrow g \\
x' & \xrightarrow{u'} & y'
\end{array}$$

can be completed to a commutative diagram of distinguished triangles, i.e.,

(4) (octahedron axiom) Given any two distinguished triangles  $x \xrightarrow{u} y \xrightarrow{v} z \xrightarrow{w} x[1]$  and  $y \xrightarrow{f} y' \xrightarrow{g} q \xrightarrow{h} y[1]$ , we can complete them to a commutative diagram

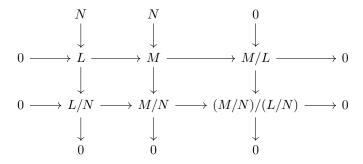
$$\begin{array}{ccccc}
x & \xrightarrow{u} & y & \xrightarrow{v} & z & \xrightarrow{w} & x[1] \\
\parallel & & \downarrow^f & & \downarrow^a & & \parallel \\
x & \longrightarrow & y' & \longrightarrow & z' & \longrightarrow & x[1] \\
\downarrow^g & & \downarrow^b & & \downarrow^{u[1]} \\
\downarrow^q & & & \downarrow^c & & \downarrow^c \\
y[1] & \longrightarrow & z[1]
\end{array}$$

where each new triangle is distinguished.

The octahedron axiom is the formal transplant of the second isomorphism theorem for  $\mathbf{Comp}(\mathbf{Ab})$ , which states that given two complexes L and M, an inclusion  $f: L \hookrightarrow M$ , and a subcomplex N of L

 $<sup>^{1}\</sup>mathrm{The}$  second isomorphism theorem holds in some form for any abelian category.

and of M, we have that  $M/L \cong (M/N)/(L/N)$ , i.e., if



has exact rows and exact left two columns, then the third column is also exact.

Now, suppose that  $\mathscr C$  is strongly pre-triangulated and let  $\alpha:M\to N$  be a morphism in  $\mathscr C$  such that  $\alpha$  is injective (i.e., ker  $\alpha$  exists and is trivial) with  $d\alpha=0$  and  $\alpha$  is split (i.e., there exists  $\beta:N\to M$  with  $p\circ\alpha=\mathrm{id}_M$ ). We call such an  $\alpha$  a split monomorphism in  $\mathscr C$ .

#### Lemma 13.1.7.

- (i) The map  $cone(\alpha) \to N/M$  is a homotopy equivalence.
- (ii) Any morphism in  $\mathscr C$  is homotopy equivalent to a split mono, i.e., given  $f:M\to L$ , we can construct a natural diagram

$$M \xrightarrow{\alpha} N \\ \downarrow^g \\ L$$

in  $\mathscr{C}$  such that  $\alpha$  is a split mono and q is an iso in  $Ho(\mathscr{C})$ .

Partial proof. For (ii), take  $N = L \oplus \text{cone}(\text{id}_M)$ .

**Theorem 13.1.8.** If  $\mathscr{C}$  is a strongly pre-triangulated differential graded category and  $\mathscr{D} = \operatorname{Ho}(\mathscr{C})$ , then  $\mathscr{D}$  is triangulated with exact triangles as the distinguished triangles.

Proof.

Verifying axioms (0) and (1) is trivial.

For axiom (2), if  $x \to y \to z \to x[1]$  is a triangle, then we can use Lemma 13.1.7 to rewrite it as a homotopy equivalent triangle  $M \to N \to L \to M[1]$  where  $M \stackrel{\alpha}{\longrightarrow} N$  is a split mono. In this case, we can check that  $N \to L \to M[1] \to N[1]$  is exact by using the splitting.

For axiom (3), note that any  $u: x \to y$  is included in  $x \to y \to \text{cone}(u) \to x[1]$ . Moreover, if

$$\begin{array}{ccc}
x & \xrightarrow{u} & y \\
f \downarrow & & \downarrow^{g} \\
x' & \xrightarrow{u'} & y'
\end{array}$$

is commutative in  $Ho(\mathscr{C})$  and we lift f, g, u, and u' to maps  $\tilde{\cdot}$  in  $\mathscr{C}$ , then we get a diagram

in 
$$\mathscr C$$
 where  $M \equiv \begin{bmatrix} \delta & 0 \\ 0 & 0 \end{bmatrix}$ ,  $\delta \in \operatorname{Hom}^{-1}(x,y')$ , and  $\tilde g \circ \tilde u - \underbrace{\tilde u' \circ \tilde f}_{d(\delta)} \sim 0$ .

For axiom (4), given a distinguished triangle  $M \to N \to L \to M[1]$ , we apply Lemma 13.1.7 twice to get a homotopy equivalent distinguished triangle  $M \to N' \to L'' \to M[1]$  where each map in this is a split mono. We are done after an application of the second isomorphism theorem.

*Remark* 13.1.9. Such reasoning can be applied to complete any differential graded category to a triangulated one.

#### 13.2 Lecture 28

**Definition 13.2.1.** If  $\mathscr{A}$  and  $\mathscr{B}$  are differential graded categories, then a differential graded functor  $F:A\to B$  has the following properties.

- (i) F is additive, i.e., F:  $\operatorname{Hom}_{\mathscr{A}}(x,y) \to \operatorname{Hom}_{\mathscr{B}}(F(x),F(y))$  is a group homomorphism for any  $x,y \in \operatorname{ob} \mathscr{A}$ .
- (ii) F respects differentials, i.e., if  $x, y \in \text{ob } \mathscr{A}$ , then  $F : \text{Hom}_{\mathscr{A}}(x, y) \to \text{Hom}_{\mathscr{B}}(F(x), F(y))$  is a map of complexes.

If  $F, G : \mathcal{A} \to \mathcal{B}$  are two differential graded functors between differential graded categories, then define, for each  $n \in \mathbb{Z}$ , the group

$$\operatorname{Hom}^n(F,G) \equiv \{\varphi_x \mid \varphi_x : F(x) \to G(x) \text{ in } \operatorname{Hom}^n_{\mathscr{B}}(F(x),G(x)), \ x \in \operatorname{ob} \mathscr{A}\}.$$

A map  $F \to G$  is defined as a natural transformation  $F \to G$  such that each component  $\varphi_x : F(x) \to G(x)$  belongs to  $\operatorname{Hom}_{\mathscr{B}}^n(F(x), G(x))$ . The differential on  $\prod_{x \in \operatorname{ob} \mathscr{A}} \operatorname{Hom}^{\bullet}(F(x), G(x))$  induces a differential on

$$\operatorname{Hom}^{\bullet}(F,G) \equiv \bigoplus_{n \in \mathbb{Z}} \operatorname{Hom}^{n}(F,G).$$

This produces a complex of maps between F and G, and we get a differential graded category  $\mathbf{dgFun}(\mathscr{A},\mathscr{B})$ .

Exercise 13.2.2. Prove the following assertions.

1. If  $F: \mathscr{A} \to \mathscr{B}$  is a differential graded functor, then  $H^0(F): H^0(\mathscr{A}) \to H^0(\mathscr{B})$  is an additive functor.

2. If  $F, G : \mathcal{A} \to \mathcal{B}$  are differential graded functors, then there is an embedding  $H^0(\text{Hom}(F,G)) \subset \text{Hom}(H^0(F), H^0(G))$ .

**Definition 13.2.3.** If  $\mathscr{A}$  is a differential graded category, then a *left*  $\mathscr{A}$ -module is a differential graded functor  $\mathscr{A} \to \mathbf{Compl}(\mathbf{Ab})$  and a *right*  $\mathscr{A}$ -module is a differential graded functor  $\mathscr{A}^{\mathrm{op}} \to \mathbf{Compl}(\mathbf{Ab})$ .

If  $\mathscr{A}$  is a differential graded category with a single object \*, then  $\mathscr{A} \leftrightarrow R := \operatorname{Hom}_{\mathscr{A}}(*,*)$ , which is precisely the complex of abelian groups equipped with a multiplication-like operation  $\cdot$  such that  $\lambda$  satisfies the graded Leibniz rule for  $\cdot$ .

**Exercise 13.2.4.** Show that a module over  $\mathscr A$  is precisely the data of a complex x of abelian groups together with a differential graded algebra homomorphism  $R \to \operatorname{Hom}_{\mathbf{Compl}(\mathbf{Ab})}(x,x)$ .

Given a differential graded category  $\mathscr{A}$ , we have respective categories of left and right modules over  $\mathscr{A}$  that are linear over a field k, namely

$$\begin{split} \mathscr{A}-\mathbf{dgmod}_k &\equiv \mathbf{dgFun}(\mathcal{A},\mathbf{Compl}(k\mathbf{-Vect})) \\ \mathbf{dgmod}_k - \mathscr{A} &\equiv \mathbf{dgFun}(\mathcal{A}^\mathrm{op},\mathbf{Compl}(k\mathbf{-Vect})). \end{split}$$

Exercise 13.2.5. Show that the functors

$$\begin{split} h^{\bullet} : \mathscr{A}^{\mathrm{op}} &\to \mathscr{A}\mathbf{-dgmod}_k \\ x &\mapsto h^{\times} \equiv \mathrm{Hom}_{\mathscr{A}}(x,-) \\ h_{\bullet} : \mathscr{A} &\to \mathscr{A}^{\mathrm{op}}\mathbf{-dgmod}_k \\ h_{\times} &\equiv \mathrm{Hom}_{\mathscr{A}^{\mathrm{op}}}(x,-) = \mathrm{Hom}_{\mathscr{A}}(-,x) \end{split}$$

are fully faithful differential graded functors.

#### Proposition 13.2.6.

- 1. If  $\mathscr{A}$  is a small differential graded category, then  $H^0(\mathscr{A}^{\mathrm{op}}-\mathbf{dgmod}_k)$  is triangulated.
- 2. If  $\mathscr A$  is a pre-triangulated differential graded category, then the fully faithful functor

$$H^0(h_{\bullet}): H^0(\mathscr{A}) \to H^0(\mathscr{A}^{\mathrm{op}} - \mathbf{dgmod}_k)$$

gives a triangulated structure on  $H^0(\mathscr{A})$ .

**Definition 13.2.7.** We say that an object F in  $\mathscr{A}^{\mathrm{op}}$ - $\mathbf{dgmod}_k$  is compact or perfect if  $F: \mathcal{A}^{\mathrm{op}} \to \mathbf{Compl}(k-\mathbf{Vect})$  commutes with arbitrary coproducts.

Note 13.2.8.  $h^{\times}$  is compact for any  $x \in \text{ob } \mathscr{A}$ .

**Definition 13.2.9.** We say that a k-linear differential graded category  $\mathscr{A}$  is *triangulated* if every compact object in  $\mathscr{A}^{\mathrm{op}}$ - $\mathbf{dgmod}_k$  is representable.

Note 13.2.10. A triangulated differential graded category is automatically strongly pre-triangulated, and  $H^0(\mathscr{A})$  is triangulated.

#### Exercise 13.2.11.

1. Suppose that  $\mathscr{D}$  is a triangulated additive category. Let  $M \to N \to C \to M[1]$  be a distinguished triangle. Show that for every  $L \in \text{ob } \mathscr{D}$ , the sequence

$$\cdots \longrightarrow \operatorname{Hom}(L,M) \longrightarrow \operatorname{Hom}(L,N) \longrightarrow \operatorname{Hom}(L,C)$$
 
$$\operatorname{Hom}(L,M[1]) \stackrel{\longleftarrow}{\longrightarrow} \operatorname{Hom}(L,N[1]) \longrightarrow \operatorname{Hom}(L,C[1]) \longrightarrow \cdots$$

is a long exact sequence of abelian groups.

2. Suppose that  $\mathscr{D}$  is triangulated. Show that the sum  $\triangle_1 \oplus \triangle_2$  of two triangles in  $\mathscr{D}$  is distinguished if and only if both  $\triangle_1$  and  $\triangle_2$  are distinguished.

**Definition 13.2.12.** If  $\mathscr{D}_1$  and  $\mathscr{D}_2$  are triangulated additive categories, then a *triangulated* (or *exact*) functor  $F: \mathscr{D}_1 \to \mathscr{D}_2$  is an additive functor such that

- (i) F is equipped with an isomorphism  $\sigma: F \circ [1] \to [1] \circ F$  and
- (ii) F sends distinguished triangles to distinguished triangles.

A morphism of two triangulated functors  $(F, \theta_F)$  and  $(G, \theta_G)$  is a morphism  $f : F \to G$  of additive functors such that f intertwines  $\theta_F$  and  $\theta_G$ . As a result, we have a category of triangulated functors  $\mathcal{D}_1 \to \mathcal{D}_2$ .

If  $\mathscr{A}$  and  $\mathscr{B}$  are differential graded categories and  $F:\mathscr{A}\to\mathscr{B}$  is a differential graded functor, then we have a natural differential graded functor  $\mathscr{A}^{\mathrm{op}}-\mathbf{dgmod}_k \stackrel{F}{\longrightarrow} \mathscr{B}^{\mathrm{op}}-\mathbf{dgmod}_k$  so that  $H^0(F)$  is triangulated.

**Definition 13.2.13.** If  $\mathscr{D}$  is a triangulated category and  $\mathscr{A}$  is an abelian category, then a *cohomological functor* is a functor  $H: \mathscr{D} \to \mathscr{A}$  such that

- (i) H is additive and
- (ii) H sends distinguished  $\triangle$ 's in  $\mathcal{D}$  into long exact sequences in  $\mathcal{A}$ .

#### Example 13.2.14.

- 1. If  $\mathcal{C}(\mathbf{Ab})$  denotes the triangulated category of homotopy classes of complexes of abelian groups, then  $H^{\bullet}: \mathcal{C}(\mathbf{Ab}) \to \mathbf{grAb}$  is a cohomological functor.
- 2. If  $\mathscr{D}$  is a triangulated category and  $L \in \text{ob } \mathscr{D}$ , then  $h^L : \mathscr{D} \to \mathbf{Ab}$  given by  $M \mapsto Z^0(\text{Hom}_{\mathscr{B}}(L, M))$  is a cohomological functor.

### 13.3 Lecture 29

Let  $\mathscr{D}$  be a triangulated category and  $\mathscr{V} \subset \mathscr{D}$  a triangulated subcategory (i.e., the inclusion functor is triangulated). We wish to construct a quotient category  $\mathscr{D}/\mathscr{V}$ , i.e., a triangulated category  $\mathscr{D}/\mathscr{V}$  together with a triangulated functor  $q: \mathscr{D} \to \mathscr{D}/\mathscr{V}$  such that

- q(x) = 0 for any  $x \in \text{ob } \mathcal{V}$  and
- for any triangulated functor  $f: \mathcal{D} \to \mathcal{D}'$  satisfying  $x \in \text{ob } \mathcal{V} \implies f(x) = 0$ , we have  $g \circ q = f$ .

#### Note 13.3.1.

- 1. In the triangulated category of triangulated categories with exact functors, the triangle  $\mathscr{V} \to \mathscr{D}/\mathscr{V} \to \mathscr{V}[1]$  is exact.
- 2. If  $\mathscr{D}$  is triangulated and  $u: x \to y$  is a morphism in  $\mathscr{D}$ , then there exists an object cone(u) in  $\mathscr{D}$  that is unique up to a non-unique isomorphism. This is the third term in a distinguished  $\triangle$  completing u.

**Exercise 13.3.2.** Show that if  $u: x \to y$  is a map in  $\mathscr{D}$ , then it is an isomorphism in  $\mathscr{D}$  if and only if cone(u) = 0.

**Definition 13.3.3.** If  $\mathscr{D}$  is a triangulated category and  $\mathscr{V} \subset \mathscr{D}$  a triangulated subcategory, then a morphism  $u: x \to y$  in  $\mathscr{D}$  is a  $\mathscr{V}$ -quasi-isomorphism if  $cone(u) \in ob(\mathscr{V})$ .

**Exercise 13.3.4.** Let  $\mathcal{V} \subset \mathcal{D}$  be a pair of triangulated categories. Use the octahedron axiom to show that if f and g are compassable morphisms in  $\mathcal{D}$ , then every morphism in  $\{f, g, g \circ f\}$  is a  $\mathcal{V}$ -quasi-isomorphism if and only if at least two morphisms in it are  $\mathcal{V}$ -quasi-isomorphisms.

Remark 13.3.5. One may define  $\mathscr{D}/\mathscr{V}$  as the localization of  $\mathscr{D}$  in the set of all  $\mathscr{V}$ -quasi-isomorphisms. But doing so requires a lot of work.

**Definition 13.3.6.** Suppose that  $\mathscr{I}$  is a small category. We say that  $\mathscr{I}$  is a *directed category* if it satisfies the following properties.

(1) If  $x_1, x_2 \in \text{ob } \mathscr{I}$ , then there is some diagram

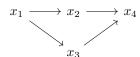


of maps in  $\mathscr{I}$ .

(2) If



is a diagram of maps in  $\mathscr{I}$ , then there exist maps  $x_2 \to x_4$ ;  $x_3 \to x_4$  in  $\mathscr{I}$  such that



commutes in  $\mathscr{I}$ .

(3) For any two parallel maps  $f, g: x \to y$ , there exists a map  $h: y \to z$  such that  $h \circ f = h \circ g$ .

Let  $\mathscr{I}$  be small. There is a well-defined functor colim:  $\operatorname{Fun}(\mathscr{I}, \mathbf{Ab}) \to \mathbf{Ab}$ , but this need *not* be exact even though both  $\mathbf{Ab}$  and  $\operatorname{Fun}(\mathscr{I}, \mathbf{Ab})$  are abelian.

**Exercise 13.3.7.** Show, however, that if  $\mathscr I$  is directed, then colim is an exact functor.

Now, let  $\mathscr{V} \subset \mathscr{D}$  be a pair of triangulated categories. Let  $x \in \text{ob}\,\mathscr{D}$  and let  $\mathscr{Q}/x$  be the full subcategory of  $\mathscr{D}/x$  consisting of morphisms  $y \to x$  that are  $\mathscr{V}$ -quasi-isomorphisms. Similarly, let  $x/\mathbb{Q}$  be the full subcategory of  $x/\mathscr{D}$  consisting morphisms  $x \to z$  that are  $\mathscr{V}$ -quasi-isomorphisms.

Exercise 13.3.8. Prove the following assertions.

- 1. Both x/2 and  $(2/x)^{op}$  are directed categories.
- 2. Any map in x/2 or 2/x is automatically a  $\mathcal{V}$ -quasi-isomorphism.

**Definition 13.3.9.** Define the Verdier quotient of  $\mathscr{D}$  by  $\mathscr{V}$  as the category  $\mathscr{D}/\mathscr{V}$  with ob  $\mathscr{D}/\mathscr{V} \equiv \text{ob }\mathscr{D}$  and  $\text{Hom}_{\mathscr{D}/\mathscr{V}}(a,b) \equiv \text{colim}_{a' \in (\mathscr{D}/a)^{\text{op}}} \text{Hom}_{\mathscr{D}}(a',b)$ .

There exists a canonical isomorphism

$$\operatorname*{colim}_{a' \in (\mathscr{Q}/a)^{\operatorname{op}}} \operatorname{Hom}_{\mathscr{D}}(a',b) \cong \operatorname*{colim}_{b' \in (b/\mathscr{Q})} \operatorname{Hom}_{\mathscr{D}}(a,b').$$

For this, we must check that given a top triangle



we can form a commutative double triangle

As a result, we get  $q: \mathcal{D} \to \mathcal{D}/\mathcal{V}$ .

**Lemma 13.3.10.** If  $x \in \text{ob } \mathscr{D}$  has q(x) = 0 in  $\mathscr{D}/\mathscr{V}$ , then x is a direct summand of an object in  $\mathscr{V}$ .

*Proof.* We have that 
$$q(x) = 0 \iff$$
 there is some  $y \in \mathscr{D}$  such that  $\varphi : y \to x$  is a  $\mathscr{V}$ -quasi-isomorphism. In this case,  $\underline{\mathrm{cone}(\varphi)} \in \mathscr{V}$ .

**Definition 13.3.11.** A triangulated subcategory  $\mathcal{V} \subset \mathcal{D}$  is *thick* if any object in  $\mathcal{D}$  that is isomorphic to a direct summand of an object in  $\mathcal{V}$  is an object in  $\mathcal{V}$ .

**Note 13.3.12.** If  $\mathscr{V}$  is a strict full thick triangulated subcategory of  $\mathscr{D}$ , then  $q: \mathscr{D} \to \mathscr{D}/\mathscr{V}$  kills all and only objects in  $\mathscr{V}$ .

**Definition 13.3.13.** If  $\mathscr{D}$  is triangulated and  $\mathscr{U}, \mathscr{V} \subset \mathscr{D}$  are strict full triangulated subcategories, then  $(\mathscr{U}, \mathscr{V})$  is an *admissible pair of subcategories* if

- (a)  $\operatorname{Hom}_{\mathscr{D}}(x,y) = 0$  for any  $x \in \operatorname{ob} u$  and  $y \in \operatorname{ob} \mathscr{V}$  and
- (b) any object  $z \in \text{ob } \mathscr{D}$  fits in a distinguished triangle  $x \to z \to y \to x[1]$  with  $x \in \text{ob } \mathscr{U}$  and  $y \in \text{ob } \mathscr{V}$ .

Exercise 13.3.14. Prove the following assertions.

- 1. The  $\triangle$  in condition (b) is unique up to a unique isomorphism and is functorial in z.
- 2. The functor  $\mathscr{D} \to \mathscr{U}$  given by  $z \mapsto x(z)$  is triangulated and is right adjoint to  $\mathscr{U} \hookrightarrow \mathscr{D}$ .

  Dually, the functor  $\mathscr{D} \to \mathscr{V}$  given by  $z \mapsto y(z)$  is triangulated and is left adjoint to  $\mathscr{V} \hookrightarrow \mathscr{D}$ .
- 3. Each of  $\mathcal{U}$  and  $\mathcal{V}$  determines the other. Specifically,

$$\mathcal{V} = \mathcal{U}^{\perp} \equiv \underbrace{\left\{ y \in \text{ob} \, \mathscr{D} \mid \text{Hom}_{\mathscr{D}}(x,y) = 0, \ x \in \text{ob} \, \mathscr{U} \right\}}_{full \ subcategory}$$

$$\mathcal{U} = {}^{\perp} \mathcal{V} \equiv \underbrace{\left\{ x \in \text{ob} \, \mathscr{D} \mid \text{Hom}_{\mathscr{D}}(x,y) = 0, \ y \in \text{ob} \, \mathscr{V} \right\}}_{full \ subcategory}.$$

In particular, both  $\mathcal{U}$  and  $\mathcal{V}$  are thick subcategories.

4. The natural compositions  $\mathscr{U} \hookrightarrow \mathscr{D} \to \mathscr{D}/\mathscr{V}$  and  $\mathscr{V} \hookrightarrow \mathscr{D} \to \mathscr{D}/\mathscr{U}$  are triangulated equivalences.

**Definition 13.3.15.** An additive pair  $(\mathcal{U}, \mathcal{V})$  is called a *semiorthogonal decomposition of*  $\mathcal{D}$  *into*  $\mathcal{U}$  and  $\mathcal{V}$ .

**Proposition 13.3.16.** If  $\mathcal{U} \subset \mathcal{D}$  is a strict full triangulated thick subcategory, then TFAE.

- 1. The inclusion  $\mathscr{U} \hookrightarrow \mathscr{D}$  has a left adjoint.
- 2. The quotient  $\mathcal{D} \to \mathcal{D}/\mathcal{U}$  has a right adjoint.
- 3.  $(\mathcal{U}, \mathcal{U}^{\perp})$  is admissible.

**Definition 13.3.17.** If  $\mathscr A$  is an abelian category, then the *derived category of*  $\mathscr A$  is the triangulated category

$$\mathcal{D}(\mathscr{A}) \equiv \frac{\mathcal{C}(\mathscr{A})}{\mathcal{C}(\mathscr{A})^{\mathrm{acyclic}}},$$

where  $\mathcal{C}(\mathscr{A})^{\text{acyclic}}$  is the full subcategory of  $\mathcal{C}(a)$  consisting of those x with zero cohomology.

To do computations in  $\mathcal{D}(\mathscr{A})$ , we must understand when  $\mathcal{D}(a)$  can be embedded in  $\mathcal{C}(\mathscr{A})$  so that  $(\mathcal{C}(\mathscr{A})^{\text{acyclic}}, \mathcal{D}(\mathscr{A}))$  is an adjoint pair. This requires  $(\mathcal{C}(\mathscr{A})^{\text{acyclic}})^{\perp}$  to be large.

Define  ${}^{\perp}\mathcal{C}(\mathscr{A})^{\text{acyclic}}$  as the category of homotopically projective objects in  $\mathcal{C}(\mathscr{A})$  and  $(\mathcal{C}(\mathscr{A})^{\text{acyclic}})^{\perp}$  as the category of homotopically injective objects in  $\mathcal{C}(\mathscr{A})$ .

#### Proposition 13.3.18.

- 1. Every bounded-above complex of projectives is a homotopically projective object in  $\mathcal{C}(\mathscr{A})$ .
- 2. Any bounded-below complex of injectives is a homotopically injective object in  $\mathcal{C}(\mathscr{A})$ .