

# Arithmetic geometry notes

Soham Chowdhury

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# **Part I**

## **Basic notions**

# Chapter 1

## Algebraic integers

Pellentesque condimentum,  
magna ut suscipit hendrerit,  
ipsum augue ornare nulla, non  
luctus diam neque sit amet urna.

---

The Dude

Fix a domain  $A$  integrally closed in  $K := K(A)$ . Let  $L|K$  be a finite extension, and  $B$  the integral closure of  $A$  in  $L$ . This is the  $AKLB$  *diagram*:

$$\begin{array}{ccc} K & \hookrightarrow & L \\ \uparrow & & \uparrow \\ A & \longrightarrow & B \end{array}$$

### 1.1 Properties of integrality

Integrality is stable under the ring operations: one would like the following to hold, and they do:

This is a corollary of the following:

**Theorem 1.1.1** (Module-theoretic characterization of integrality). *A finite number of  $b_i$  are integral over  $A \iff$  the ring  $A[b_1, \dots, b_n]$  is finitely generated as an  $A$ -module.*

*Proof.* TODO. ■

**Corollary 1.1.2.** *If  $a$  and  $b$  are integral over  $A$ , so are  $a + b$  and  $ab$ .*

**Theorem 1.1.3** (Integrality is transitive). *Consider ring extensions  $A \subseteq B \subseteq C$ .  $A \subseteq B$  integral and  $B \subseteq C$  integral  $\iff A \subseteq C$  integral.*

*Proof.*  $A \subseteq C$  integral implies  $A \subseteq B$  integral. (Why?) ■

**Theorem 1.1.4.** *Any element  $l \in L$  is equal to  $b/a$  for  $b \in B$  and  $a \in A$ .*

*Proof.* Consider an element  $l \in L$ . The minimal polynomial  $m_l$  of  $l$  over  $K$  gives rise to a polynomial over  $A$

$$a_n l^n + a_{n-1} l^{n-1} + \cdots + a_0 = 0$$

by clearing denominators. Now observe that  $\ell := a_n l$  is integral over  $A$ : multiplying by  $a_n^{n-1}$  gives an equation of the form

$$\ell^n + a'_{n-1} \ell^{n-1} + \cdots + a'_0 = 0.$$

This shows that taking  $b/a = \ell/a_n$  works. ■

*Remark 1.1.5.* Notice that  $K(B) = L$ . Indeed,  $B \subset L$  so  $K(B) \subset L$ , and the result above shows that  $L \subset K(B)$  (set-theoretically,  $L \subset B \times A \subset B \times B$ ).

**Theorem 1.1.6.**  $l \in L$  is integral over  $A$  iff its minimal polynomial  $\mu_l$  over  $K$  has coefficients in  $A$ .

*Proof.* If  $\mu := \mu_l \in A[x]$  then we have integrality of  $l$  over  $A$  by definition. Consider now the case of an integral  $l$  with minimal polynomial  $\mu \in K[x]$ . From integrality over  $A$  we know that  $l$  is a root of some  $g \in A[x]$ . Then  $\mu|g$  in  $K[x]$ , so all zeros of  $\mu$  are zeros of  $g$  and hence integral over  $A$ .

By Vieta, the coefficients  $a_i$  are given by elementary symmetric polynomials in the roots and are hence, by Corollary 1.1.2, integral over  $A$  themselves. The  $a_i$  are elements of  $K$ , so, in this case, integrality over  $A$  means that  $a_i \in K$ , and hence  $\mu \in A[x]$ . ■

## 1.2 The trace and the norm

Given  $x \in L$ , multiplication by  $x$  determines an endomorphism

$$T_x : \alpha \mapsto x\alpha$$

of the  $K$ -vector space  $L$ . We define the trace and norm maps

$$\begin{aligned} \text{Tr}_{L|K}(z) &= \text{tr } T_z \\ \text{Nm}_{L|K}(z) &= \det T_z \end{aligned}$$

Let  $n = [L : K]$ . The characteristic polynomial

$$\begin{aligned} \chi_z(t) &= \det(tI - T_z) \\ &= t^n - a_1 t^{n-1} + \cdots + (-1)^n a_n \in K[t] \end{aligned}$$

contains coefficients  $a_1 = \text{Tr}_{L|K}(z)$  and  $a_n = \text{Nm}_{L|K}(z)$ .

*Remark 1.2.7.* If this isn't immediately clear, think Vieta. (This will be one of the recurring themes throughout this chapter.)

### 1.3 Galois-theoretic interpretations

Fix an algebraic closure  $\bar{K} = K^{\text{alg}}$  of  $K$ .

**Proposition 1.3.8.** *If  $L|K$  is separable, letting  $\sigma : L \rightarrow \bar{K}$  vary over the  $K$ -embeddings of  $L$  into  $\bar{K}$ , we have*

1.  $\chi_z(t) = \prod_{\sigma} (t - \sigma z)$
2.  $\text{Tr}_{L|K}(z) = \sum_{\sigma} \sigma z$
3.  $\text{Nm}_{L|K}(z) = \prod_{\sigma} \sigma z$

*Proof.* Let  $d = [L : K(x)]$ . The characteristic polynomial is a power

$$\chi_z = \mu_z^d$$

where  $d = [L : K(z)]$ . Part 1 easily implies the others, by Vieta's formulas. ■

**Theorem 1.3.9.** *For a tower of finite extensions  $K \subseteq L \subseteq M$ , we have*

$$\begin{aligned} \text{Tr}_{L|K} \circ \text{Tr}_{M|L} &= \text{Tr}_{M|K} \\ \text{Nm}_{L|K} \circ \text{Nm}_{M|L} &= \text{Nm}_{M|K} \end{aligned}$$

### 1.4 Integral bases

**Definition 1.4.10.** *The **discriminant** of a basis  $\alpha_i$  of a separable extension  $L|K$  is defined by*

$$d(\alpha_1, \dots, \alpha_n) = \det((\sigma_i \alpha_j))^2$$

where the  $\sigma_i$  are the  $K$ -embeddings  $L \hookrightarrow \bar{K}$ .

**Proposition 1.4.11.** *For  $L|K$  a separable extension with basis  $\alpha_i$ , the function*

$$(x, y) = \text{Tr}_{L|K}(xy)$$

*yields a nondegenerate bilinear form on the  $K$ -vector space  $L$ .*

**Corollary 1.4.12.** *For  $L|K$  and  $\alpha_i$  as above,*

$$d(\alpha_1, \dots, \alpha_n) \neq 0.$$

*Proof.* The form has matrix

$$M = \text{Tr}_{L|K}((\alpha_i \alpha_j))$$

with respect to the given basis. The nondegeneracy of the form, which we have from Proposition 1.4.11, is equivalent to the statement that  $\det M \neq 0$ , whence the claim follows. ■

**Lemma 1.4.13.** *Let  $(\alpha_i)$  be a basis of  $L|K$  contained in  $B$ , with  $d = d(\alpha_1, \dots, \alpha_n)$ . Then*

$$dB \subseteq A\alpha_1 + \dots + A\alpha_n.$$

**Proposition 1.4.14.** *If  $L|K$  is separable and  $A$  is a PID, every finitely generated  $B$ -submodule  $M \neq 0$  of  $L$  is a free  $A$ -module of rank  $[L : K]$ .*

**Corollary 1.4.15.**  *$B$  admits an integral basis over  $A$ .*

**Proposition 1.4.16.** *Let  $M|K$  and  $N|K$  be two Galois extensions with  $M \cap N = K$ , with  $m = [M : K]$  and  $n = [N : K]$ . Fix integral bases  $(\alpha_i)_{1 \leq i \leq m}$  of  $M|K$  and  $(\beta_j)_{1 \leq j \leq n}$  of  $N|K$  respectively, with discriminants  $\mu$  and  $\nu$  respectively. If  $\mu$  and  $\nu$  are relatively prime, with  $x\mu + y\nu = 1$  for some  $x, y \in A$ , then  $(\alpha_i\beta_j)$  is an integral basis of  $MN$ , with discriminant  $m^\nu n^\mu$ .*

**Proposition 1.4.17.** *If  $i \subseteq j$  are two nonzero finite  $\mathcal{O}_K$ -submodules of  $K$ , then  $(j : i)$  is finite. Moreover,*

$$d(i) = (j : i)^2 d(j)$$

*holds.*

# Appendix A

## Complex analysis

### A.1 Holomorphy and complex differentiability

There are, broadly speaking, two criteria that we would like nice complex-valued functions to satisfy. The first is a notion of differentiability similar to the one from calculus, where every function can be linearly approximated by a *derivative*, while the second asks that every function be locally representable by a *power series* expansion.

This section will develop these notions, demonstrate relations between the two, and discuss some simple consequences of these conditions.

#### Initial definitions

Let  $\Omega \subseteq \mathbb{C}$  be an open set.

**Definition A.1.1.** A function  $f : \Omega \rightarrow \mathbb{C}$  is complex differentiable at  $z_0$  if the limit

$$f'(z_0) = \lim_{h \rightarrow 0} \frac{f(z_0 + h) - f(z_0)}{h}$$

exists.  $f$  is said to be complex differentiable on  $\Omega$  if it is complex differentiable at all  $z_0 \in \Omega$ .

**Definition A.1.2** (Narasimhan).  $f : \Omega \rightarrow \mathbb{C}$  is holomorphic on  $\Omega$  if, for all  $z_0 \in \Omega$ , there exists a neighborhood  $U \subseteq \Omega$  of  $z_0$  and a sequence  $\{c_n\}_{n \geq 0}$  of complex numbers such that, for all  $z \in U$ , the series

$$\sum_{n=0}^{\infty} c_n (z - z_0)^n$$

converges to  $f(z)$ .

These two definitions are in fact equivalent: holomorphy on  $\Omega$  is the same as  $\mathbb{C}$ -differentiability on  $\Omega$ . This is the content of the Cauchy-Goursat theorem, which we will prove later (TODO ref).



## Properties

Holomorphy and complex differentiability imply relations between the “ $x$ -behavior” and “ $y$ -behavior” of a function, so that there are certain rigidity properties we can be assured of. We now show a few properties which are all roughly similar in nature, culminating in the Definition A.1.6.

**Proposition A.1.3.** *Let  $f : \Omega \rightarrow \mathbf{C}$  be  $\mathbf{C}$ -differentiable at  $a \in \Omega$ . Then  $\partial_x f(a)$  and  $\partial_y f(a)$  exist, and*

$$\frac{\partial f}{\partial x}(a) = -i \frac{\partial f}{\partial y}(a) = f'(a)$$

*holds.*

*Proof.* In the Riemann tradition, write  $a = \sigma + it$ . We will calculate  $f'(a)$  in two ways, by approaching 0 along the real axis, then along the imaginary axis.

Taking  $0 \neq \xi \in \mathbf{R}$ ,

$$\begin{aligned} f'(a) &= \lim_{\xi \rightarrow 0} \frac{f(a + \xi) - f(a)}{\xi} \\ &= \lim_{\xi \rightarrow 0} \frac{f(\sigma + \xi, t) - f(\sigma, t)}{\xi} \\ &= \frac{\partial f}{\partial x}(a). \end{aligned}$$

Taking  $0 \neq \eta \in \mathbf{R}$ ,

$$\begin{aligned} f'(a) &= \lim_{\eta \rightarrow 0} \frac{f(a + i\eta) - f(a)}{i\eta} \\ &= \lim_{\eta \rightarrow 0} \frac{f(\sigma, t + \eta) - f(\sigma, t)}{i\eta} \\ &= \frac{1}{i} \frac{\partial f}{\partial y}(a). \end{aligned}$$

Equating these two expressions to  $f'(a)$  is then enough. ■

Note that  $x$  and  $y$  can be expressed in terms of  $z$  and  $\bar{z}$ :

$$\begin{aligned} x &= \frac{z + \bar{z}}{2} \\ y &= \frac{z - \bar{z}}{2i} \end{aligned}$$

This means one can (formally?) write, using the chain rule,

$$\frac{\partial f}{\partial z} = \frac{\partial f}{\partial x} \frac{\partial x}{\partial z} + \frac{\partial f}{\partial y} \frac{\partial y}{\partial z} = \frac{1}{2} \cdot \frac{\partial f}{\partial x} + \frac{1}{2i} \cdot \frac{\partial f}{\partial y} = \frac{1}{2}(f_z - if_y)$$

**Exercise A.1.A.** What is the analogous expression for  $\partial_{\bar{z}}$ ?

This motivates the following definition.

**Definition A.1.4.** The Wirtinger derivatives are differential operators defined as follows:

$$\begin{aligned}\partial_z &= \frac{\partial}{\partial z} = \frac{1}{2}(\partial_x - i\partial_y) \\ \partial_{\bar{z}} &= \frac{\partial}{\partial \bar{z}} = \frac{1}{2}(\partial_x + i\partial_y)\end{aligned}$$

**Proposition A.1.5.** If  $f : \Omega \rightarrow \mathbf{C}$  is  $\mathbf{C}$ -differentiable at  $a \in \Omega$ ,

$$\begin{aligned}\frac{\partial f}{\partial z}(a) &= f'(a) \\ \frac{\partial f}{\partial \bar{z}}(a) &= 0\end{aligned}$$

**Exercise A.1.B.** Prove this. (This is essentially a restatement of Proposition A.1.3 using the new notation.)

**Definition A.1.6.** Let  $f : \Omega \rightarrow \mathbf{C}$  be written as  $f = u + iv$ , where  $u, v : \Omega \rightarrow \mathbf{R}$ . Then the equations

$$\begin{aligned}\frac{\partial f}{\partial x} &= i \frac{\partial f}{\partial y} \\ \frac{\partial f}{\partial z} &= \frac{\partial f}{\partial \bar{z}} \\ \frac{\partial f}{\partial \bar{z}} &= 0\end{aligned}$$

are each equivalent to the following pair of equations:

$$u_x = v_y \tag{A.1}$$

$$-u_y = v_x \tag{A.2}$$

These differential equations are called the Cauchy-Riemann equations.

Define an  $\mathbf{R}$ -isomorphism of fields

$$\begin{aligned}\mu : \mathbf{C} &\rightarrow \mathbf{R}^2 \\ x + iy &\mapsto (x, y)\end{aligned}$$

Let  $f : \Omega \rightarrow \mathbf{C}$  have first partial derivatives at  $w$ . We have the *Jacobian map*, represented in the standard basis by

$$J_w(u, v) = \begin{bmatrix} u_x(w) & u_y(w) \\ v_x(w) & v_y(w) \end{bmatrix}$$

This is a local isomorphism of  $\mathbf{R}^2$  onto the tangent space  $T_w \mathbf{R}^2 \simeq \mathbf{R}^2$ . We “lift” this to  $\mathbf{C}$ :

**Definition A.1.7.** The tangent map of  $f = u + iv$  at  $w$  is

$$d_w f := \mu^{-1} \triangleleft J_w(u, v) \triangleleft \mu$$

**Proposition A.1.8.** We have  $\partial_{\bar{z}} f(w) = 0$  iff  $d_w f$  is  $\mathbf{C}$ -linear, that is, if

$$d_w f(\lambda \cdot z) = \lambda \cdot d_w f(z)$$

in which case

$$d_w f(z) = z \cdot \partial_z f(w) = z \cdot f'(w)$$

Notice that this says exactly that  $f$  is locally linear.

*Proof.* TODO. Pretty weird in Narasimhan. ■

$\mathbf{C}$ -differentiable functions satisfy the expected properties:

1. Given differentiable  $f, g : \Omega \rightarrow \mathbf{C}$  and  $\lambda \in \mathbf{C}$ ,

$$f + g : z \mapsto f(z) + g(z)$$

$$f \cdot g : z \mapsto f(z) \cdot g(z)$$

$$\lambda \cdot f : z \mapsto \lambda \cdot f(z)$$

are all  $\mathbf{C}$ -differentiable.

## A.2 Zeroes and poles

**Definition A.2.9** (Zeros).

**Definition A.2.10.** A function  $f : \mathbf{C} \rightarrow \mathbf{C}$  is elliptic if, for all  $\lambda$  in some lattice  $\Lambda$ ,  $f(z + \lambda) = f(z)$  for all  $z \in \mathbf{C}$ .

**Theorem A.2.11** (Liouville). Any bounded entire function is constant.

**Definition A.2.12.** If  $f$  is holomorphic on all of  $\mathbf{C}$ , it is said to be entire.

# Appendix B

## Fourier analysis

### B.1 Fourier expansions

Let  $g : \mathbb{C} \rightarrow \hat{\mathbb{C}}$  be a continuous function with period 1.

The  $n$ th *Fourier coefficient*  $a_n(y)$  is

$$a_n(y) = \hat{g}(n) = \int_0^1 g(z) \exp(-2\pi i n z) dx$$

Then we have the *Fourier expansion*

$$g(z) = \sum_{n=-\infty}^{\infty} a_n(y) \exp(2\pi i n z)$$

### B.2 Meromorphicity

The *nome* is a common building block for interesting functions.

$$q = q(z) := \exp(2\pi i z)$$

Let  $g$  be meromorphic in the notation of the previous section. Then there exists a unique meromorphic  $G : \mathbb{C}^\times \rightarrow \hat{\mathbb{C}}$  such that  $g(z) = G(q)$  (TODO why?): in other words, a period-1 meromorphic function of  $z$  is in fact a function of  $q(z)$ .

Note that  $G$  has a removable singularity at 0, so, by Theorem ???,  $G$  extends to a meromorphic function on  $\mathbb{C}$  iff

$$\lim_{q \rightarrow 0} G(q) |q|^m = 0$$

for some  $m$ . What does it mean for  $q$  to go to 0?

$$\begin{aligned} q \rightarrow 0 &\implies \exp(2\pi i(x + iy)) \rightarrow 0 \\ &\implies \exp(2\pi i x) e^{-2\pi y} \rightarrow 0 \\ &\implies y \rightarrow \infty \end{aligned}$$

so we have  $g(z)|q|^m \rightarrow 0$  as  $g(z)\exp(-2\pi my) \rightarrow 0$ , so we need

$$\text{as } \mathfrak{I}(z) \rightarrow \infty, \exists m \quad |g(z)| < \exp(2\pi my)$$

The meromorphy of  $G(q)$  at 0 thus requires  $\mathfrak{I}(z) \rightarrow \infty$ , in which case we say  $g$  is *meromorphic at  $i\infty$* . Then  $G$ , being meromorphic at 0, has a Laurent series expansion

$$g(z) = G(q) = \sum_{n=-m}^{\infty} c_n q^n = \sum_{n=-m}^{\infty} c_n e^{2\pi i n z}$$

Here  $m$  is the order of the pole of  $G$  at 0. However, we also have a Fourier expansion

$$g(z) = \sum_{n=-\infty}^{\infty} a_n(y) e^{2\pi i n z}$$

and, equating coefficients,

$$\begin{aligned} a_n(y) &= c_n & \text{for } n \geq -m \\ a_n(y) &= 0 & \text{for } n < -m \end{aligned}$$

# Appendix C

## Modular forms

### C.1 The hyperbolic plane

**Definition C.1.1.** *The upper half-plane in  $\mathbf{C}$  is*

$$\mathfrak{H} := \{h \in \mathbf{C} : \text{Im}(h) > 0\}$$

### C.2 Möbius transformations

$$\frac{az + b}{cz + d}$$

### C.3 The modular group

Define Möbius transformations

$$S = \begin{bmatrix} & 1 \\ -1 & \end{bmatrix}$$
$$T = \begin{bmatrix} & 1 \\ 1 & \end{bmatrix}$$

As before, the actions of these two matrices are as follows:

$$Sz = \frac{0z + 1}{-1z + 0} = -\frac{1}{z}$$
$$Tz = \frac{1z + 1}{0z + 1} = z + 1$$

$S$  is an inversion about the unit circle ( $z \mapsto 1/z$ ) followed by reflection across the imaginary axis ( $z \mapsto -z$ ), while  $T$  is a simple translation.

These form a “basis”, a generating set, for the modular group:

**Proposition C.3.2.**  $\text{PSL}_2(\mathbf{Z}) = \langle S, T \rangle$ .

## C.4 A fundamental domain for $\mathrm{PSL}_2(\mathbf{Z})$

**Definition C.4.3.** Let  $F \subset \mathfrak{H}$  be a closed set with connected interior, and let  $\Gamma$  be a subgroup of  $\mathrm{PSL}_2(\mathbf{Z})$ . We say  $F$  is a fundamental domain for  $\Gamma \backslash \mathfrak{H}$  or for  $\Gamma$  if

1. any  $h \in \mathfrak{H}$  is  $\Gamma$ -equivalent to some point in  $F$
2. no two interior points of  $F$  are equivalent under the  $\Gamma$  action
3. the boundary of  $F$  is piecewise smooth

Define  $\mathbf{M} = \mathrm{PSL}_2(\mathbf{Z})$ .

We now exhibit a fundamental domain for  $\mathrm{PSL}_2(\mathbf{Z})$ . Let

$$F = \{h \in \mathfrak{H} : |\Re(h)| \leq \frac{1}{2}, |h| \geq 1\}$$

**Proposition C.4.4.**  $F$  is a fundamental domain for  $\mathbf{M}$ .

## C.5 Congruence subgroups

**Definition C.5.5.** Let  $N \in \mathbf{Z}_{>0}$ . The modular group of level  $N$  is

$$\Gamma_0(N) = \left\{ \begin{bmatrix} a & b \\ c & d \end{bmatrix} \in \mathbf{M} : c \equiv 0 \pmod{N} \right\}$$

We also have

$$\Gamma_1(N) = \left\{ \begin{bmatrix} a & b \\ c & d \end{bmatrix} \in \mathbf{M} : \begin{bmatrix} a & b \\ c & d \end{bmatrix} \equiv \begin{bmatrix} 1 & b \\ 0 & 1 \end{bmatrix} \pmod{N} \right\}$$

and the principal congruence subgroups

$$\Gamma(N) = \left\{ \begin{bmatrix} a & b \\ c & d \end{bmatrix} \in \mathbf{M} : \begin{bmatrix} a & b \\ c & d \end{bmatrix} \equiv \begin{bmatrix} 1 & 0 \\ 0 & 1 \end{bmatrix} \pmod{N} \right\}$$