

M392C: Complex Geometry

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Abstract

Class taught by Berndt Siebert at the University of Texas at Austin, notes taken by Reese Lance. The notes are not live texed. They are post-mortem-texed, that is, taken by hand during class and typed later. Some of my own thoughts are interjected, but rarely. I initially thought to try to separate my thoughts from the professor's but it becomes too difficult. As such I will also try to expand on examples which are mentioned in passing in class and spell out proofs which are glossed over. This helps to justify the existence of this set of notes, as opposed to live-texed notes (which are often available for classes at this university), which are probably slightly better for a faithful representation of what is being taught in the classroom. Especially because some of my own content is interspersed throughout these notes, any corrections, questions, comments, suggestions, etc., can be sent via email (reese.lance@utexas.edu) or if you can find any other way to communicate with me, that is also fine. At the moment I'm trying to get the notes written, and worrying about the format later. I'm also not going to track theorem and lemma numbers, as I think that's mostly useless. If a proof somewhere says "Applying Theorem X", it can usually be determined from context what theorems need to be invoked, and if the reader doesn't find it readily apparent, then searching for the theorem in question will be a valuable experience. Also I always forget to write down the numbers. Also as I revisit and add in more stuff, the numbering becomes involved and I'd have to actually figure out how to number properly instead of just manually putting numbers, which is what would have been the plan. Bernd releases "Guiding Questions" periodically, which are meant to be concepts to refresh yourself with before the subsequent lectures. I will attempt to address some of these questions in the notes, but this is an intimidating proposition, as I'm already falling behind on notes, in classic fashion, so we will see if that actually ends up happening. Thanks to Arun Debray whose formatting choices inspired my own.

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Overview of Complex Analysis

Lecture 1, Aug 27. When working with several complex variables, the use of i as an index becomes problematic, since we want to assign $\sqrt{-1} = i$. So for instance, the standard coordinates

$$\mathbb{C}^n = \{z_1, \dots, z_n \mid z_i = x_i + iy_i, x_i, y_i \in \mathbb{R}\}$$

might appear confusing. However, it is always clear from context which i is meant: There is never an instance where you will be indexing over the complex numbers. One important concept we will discuss is the:

Theorem: (Cauchy Integral Formula, version 1) *If $f : U \rightarrow \mathbb{C}$ is a holomorphic function, and $\gamma : I \rightarrow U$ is a simple, closed curve in U , then for any $z \in \Omega_{int}$ defined by the Jordan Curve theorem applied to γ satisfies:*

$$f(z) = \frac{1}{2\pi i} \int_{\gamma} \frac{f(\zeta)}{(\zeta - z)} d\zeta$$

WLOG, we may treat γ as the unit disk, in which case CIF is written as

$$f(z) = \frac{1}{2\pi i} \int_{\|\zeta\|=1} \frac{f(\zeta)}{(\zeta - z)} d\zeta$$

Here, the curve we are interested in is γ , but we require f to be defined on a larger, open set U , but this is an unnecessary requirement: for CIF to hold, we need only require that f be continuous and holomorphic on Ω_{int} . We can see version 1 implies version 2 by taking increasingly large disks inside of Ω_{int} : Apply version 1 of CIF on a unit disk of radius $1 - \epsilon \subset \Omega_{int}$. Then we have the CIF formula for all points $z \in B(0, 1 - \epsilon)$. By perturbing the smaller disk slightly, a continuity argument says that the new integral,

$$\frac{1}{2\pi i} \int_{\|\zeta\|=1-\epsilon} \frac{f(\zeta)}{(\zeta - z)} d\zeta$$

converges. But of course, the value of $f(z)$ doesn't change, so neither does the integral. Then taking the limit $\epsilon \rightarrow 0$ shows the result. This is the more natural way to state the CIF, since it doesn't require the extra ad-hoc neighborhood U .

In a first encounter with complex analysis, we typically write down definitions formally and manipulate them calculus style:

$$\frac{\partial}{\partial z} := \frac{1}{2} \left(\frac{\partial}{\partial x} - i \frac{\partial}{\partial y} \right), \quad \text{and} \quad \frac{\partial}{\partial \bar{z}} := \frac{1}{2} \left(\frac{\partial}{\partial x} + i \frac{\partial}{\partial y} \right)$$

and similarly for the differential forms

$$dz = dx + i dy, \quad \text{and} \quad d\bar{z} = dx - i dy$$

If we consider a differentiable function $f = g + ih : \mathbb{C} \rightarrow \mathbb{C}$, we can write down its Jacobian

$$D_p f = \begin{pmatrix} g_x & g_y \\ h_x & h_y \end{pmatrix} (p)$$

If we wish to change coordinates to $z = x + iy$, we might want to write

$$D_p f = \begin{pmatrix} \partial_z f & \partial_{\bar{z}} f \\ \partial_z \bar{f} & \partial_{\bar{z}} \bar{f} \end{pmatrix}$$

But as we know, $\overline{(\partial_z f)} = \partial_{\bar{z}} \bar{f}$, so this Jacobian has inherent redundancy: the two entries which are diagonal from each other are related by the complex conjugate. In other words, this matrix has the form

$$D_p f = \begin{pmatrix} a & b \\ \bar{b} & \bar{a} \end{pmatrix}$$

So in order to compute the real jacobian, if you're working in complex coordinates, you must know $\partial_z f$ and $\partial_{\bar{z}} f$, it is not enough to just know one. If, in addition, f is holomorphic, then the Cauchy-Riemann equations imply $\partial_{\bar{z}} f = 0$, so all the information about the total derivative of f is contained in $\partial_z f$. So for $f \in \mathcal{O}(\mathbb{C})$ ¹,

$$D_p f = \begin{pmatrix} a & 0 \\ 0 & \bar{a} \end{pmatrix}$$

Now we move on to several variables, with the notion of holomorphicity:

Definition: If $U \subset \mathbb{C}^n$, then $f : U \rightarrow \mathbb{C}^n$ is holomorphic if the function

$$w \mapsto f(z_1, \dots, z_{i-1}, w, z_{i+1}, \dots, z_n) \text{ is holomorphic } \forall i \in \{1, \dots, n\}$$

i.e. if f is holomorphic in each of its variables. Then, if $f \in \mathcal{O}(U)$, we can write down a version of CIF by applying the 1-d version repeatedly to each entry:

$$f(z) = \frac{1}{(2\pi i)^n} \int_{\|z_i - w_i\| = \epsilon_i} \frac{f(\zeta_1, \dots, \zeta_n)}{(\zeta_1 - z_1) \dots (\zeta_n - z_n)} d\zeta_1 \dots d\zeta_n$$

here, the natural domain is

Definition: A poly disc, $D_\epsilon(w)$, is a product of disks,

$$\epsilon = (\epsilon_1, \dots, \epsilon_n), \quad D_\epsilon(w) := \{z \mid \|z_i - w_i\| < \epsilon_i\}$$

¹this denotes the set of holomorphic functions on \mathbb{C} , which is referred to by Siebert as a “structure sheaf”, which is a phrase that scares me

For $n > 1$, this is distinct from the unit disk, even when $\epsilon = (1, \dots, 1)$: in $n = 2$, $(i - \epsilon, i - \epsilon) \in D_{(1,1)}(0)$, but $(i - \epsilon, i - \epsilon)$ has distance ~ 2 from the origin, and so is not in the unit disk. As such, their boundaries are also not the same:

$$\partial\Delta^n = S^{2n-1} \not\cong \{z \in \mathbb{C}^n \mid \|z\| = 1\}$$

The RHS is sometimes referred to as the Shilov boundary.

As in one variable, the CIF gives us the power series expansions:

$$f(z) = \sum a_I z^I$$

where I is the multi-index $I = (i_1, \dots, i_n)$, and a_I are complex coefficients that can be obtained using the CIF.

Another interesting feature unique to $n > 1$ complex analysis is the concept of domains of holomorphy. A domain of holomorphy, Ω is characterized by the property that there is a holomorphic function on Ω which cannot be extended to a larger set.

In the case of Δ^2 , this means that

$$f \in \mathcal{O}(\Delta^2 \setminus \{0\}) \Rightarrow f \text{ extends over } 0$$

i.e. the map given by restriction

$$\mathcal{O}(\Delta^2) \rightarrow \mathcal{O}(\Delta^2 \setminus \{0\}) \text{ is a bijection}$$

This fact is a consequence of

Theorem (Hartogs): *Let f be a holomorphic function on a set $G \setminus K$ where G is an open subset of \mathbb{C}^n , for $(n \geq 2)$ and K is a compact subset of G . If the complement $G \setminus K$ is connected, then f can be extended to a unique holomorphic function on G .*

Switching gears again, suppose you have a holomorphic function f on a poly disc, and you want to study the zeros of this function.

Definition: A Weierstrass polynomial is a function on $\mathbb{C}^n = \{w, z\}$ for $w \in \mathbb{C}$ and $z \in \mathbb{C}^{n-1}$, of the form

$$w^d + a_1(z)w^{d-1} + \dots + a_d(z)$$

where a_i are holomorphic functions on $U \subset \mathbb{C}^{n-1}$. If you fix the z coordinate, f is a polynomial in w of degree d , with d roots. We will talk more about why this is a desirable form in the next lecture.

Hartogs Theorem and Domains of Holomorphy

Lecture 2, Sept 1. In this lecture, we will go through proofs for Hartogs' and Weierstrass Theorems.

Theorem (Hartog): Let $\epsilon, \epsilon' \in \mathbb{R}_{\geq 0}^n$, and $\epsilon'_i > \epsilon_i$, $n \geq 2$. Then the restriction map

$$\varphi : \mathcal{O}(D_\epsilon) \rightarrow \mathcal{O}(D_\epsilon \setminus \overline{D_{\epsilon'}})$$

is bijective.

Proof: To prove injectivity, we will use a result from 1 variable complex analysis:

Theorem (Identity Theorem): If $f, g : U \rightarrow \mathbb{C}$ are holomorphic functions on an open subset $U \subset \mathbb{C}$ such that there exists a non-empty open subset $V \subset U$ such that $f|_V = g|_V$, then $f = g$.

Knowing this, if we assume that $\varphi(f) = \varphi(g)$, then if we restrict to a single variable, say the first coordinate, and fix all the others, then $\varphi(f)(z) = \varphi(g)(z)$ for all $z \in D_\epsilon \setminus \overline{D_{\epsilon'}}$. An open set minus a closed set is open, and $\epsilon'_i > \epsilon_i$ implies $D_\epsilon \setminus \overline{D_{\epsilon'}}$ is non empty, so we satisfy the conditions to apply 1-dimensional Identity Theorem, setting $U = D_\epsilon$ and $V = D_\epsilon \setminus \overline{D_{\epsilon'}}$, so that $f = g$, and φ is injective.

To show surjectivity, define a projection

$$\pi : D_\epsilon \setminus \overline{D_{\epsilon'}} \rightarrow D_{\hat{\epsilon}} \setminus \overline{D_{\hat{\epsilon}'}}$$

where a hat denotes an omission of the first entry, so this is a projection from $\mathbb{C}^n \rightarrow \mathbb{C}^{n-1}$. Then we have

$$\pi^{-1}(\hat{z}) = \begin{cases} \Delta_{\epsilon_1} \setminus \overline{\Delta_{\epsilon'_1}} \times \{\hat{z}\} & z \in \overline{D_{\epsilon'}} \\ \Delta_{\epsilon_1} \times \{\hat{z}\} & \text{else} \end{cases}$$

Then take $f \in \mathcal{O}(D_\epsilon \setminus \overline{D_{\epsilon'}})$, and we want to find a function on D_ϵ which restricts to f on $D_\epsilon \setminus \overline{D_{\epsilon'}}$, and we will use CIF to define

$$g := \frac{1}{2\pi i} \int_{\|\zeta\| = \frac{\epsilon_1 + \epsilon'_1}{2}} \frac{f(\zeta, \hat{z})}{\zeta - z_1} d\zeta \in \mathcal{O}\left(D_{\left(\frac{\epsilon_1 + \epsilon'_1}{2}, \hat{\epsilon}\right)}\right)$$

So we are applying CIF in one dimension, specifically the first variable, so that g is holomorphic in the first variable, and f is holomorphic in the final $n - 1$ variables, so g is holomorphic in each of its variables, and thus is holomorphic on this restricted polydisc.

Then, if we restrict to the “annulus” polydisc $D\left(\frac{\epsilon_1 + \epsilon'_1}{2}, \hat{\epsilon}\right) \setminus \overline{D_{\epsilon'}}$, the Cauchy Integral Theorem tells us that $g = f$, as required. □

Question: What is the topology of the domain $D_\epsilon \setminus \overline{D_{\epsilon'}}$?
If we look in the 2D case, we can define a map

$$\begin{aligned} \mathbb{C}^2 &\rightarrow \mathbb{R}_{\geq 0}^2 \\ (z_1, z_2) &\mapsto (\|z_1\|, \|z_2\|) \end{aligned}$$

Then the fiber over each point in $\mathbb{R}_{\geq 0}^2$ is $S^1 \times S^1$.

To answer the question, we know that the poly discs have the topology of a ball. So if you remove a closed ball of smaller radius, you are left with a space that retracts onto the sphere, so $D_\epsilon \setminus \overline{D_{\epsilon'}}$ is topologically a sphere, S^{2n-1} . As such, it is connected and simply-connected.

Now we talk about domains of convergence of power series, $\sum_{k \in \mathbb{N}^n} a_k z^k$. Since the domain of convergence in $n = 1$ is a disk, we might guess that in n dimensions, the domain of convergence will be a polydisc. To prove this, we will use a lemma:

Lemma: If $w \in \mathbb{C}^{n*}$, and

$$\forall k, |a_k w^k| \leq C$$

for some C , i.e. there is a uniform bound on the norm of the terms, then $\sum_k a_k z^k$ converges absolutely in D_w .

Proof: Same as dimension 1 proof. □

So it turns out that the polydisc guess is not exactly right, but is close. Instead,

Definition: Given the map $\kappa : \mathbb{C}^n \rightarrow \mathbb{R}_{\geq 0}^n$ sending $(z_1, \dots, z_n) \mapsto (|z_1|, \dots, |z_n|)$, a Reinhardt domain is the pre-image of a domain $U \subset \mathbb{R}_{\geq 0}^n$ under κ . In this case, a domain is just a connected, open subset.

The fiber of a point $x \in \mathbb{R}_{\geq 0}^n$ with all non-zero entries is an n -torus. If not all the entries are non-zero, then the number of coordinate axes which x intersects determines the dimension of the resulting torus. These are just the fibers over single points, the Reinhardt domains are the pre-images over any connected, open set, so they can look very crazy.

Then we may define

Definition: The Logarithm map is the map

$$\begin{aligned} \text{Log} : \mathbb{C}^{n*} &\rightarrow \mathbb{R}^n \\ (z_1, \dots, z_n) &\mapsto (\log(|z_1|), \dots, \log(|z_n|)) \end{aligned}$$

Definition: If $U \subset \mathbb{C}^n$ is a Reinhardt domain, then U is logarithmically convex if $\text{Log}(U \cap \mathbb{C}^{n*})$ is a convex domain.

Proposition: If U is the domain of convergence for some power series $\sum_k a_k z^k$, then U is a logarithmically convex Reinhardt domain. In fact, this is an equivalence, but we'll only prove the

stated direction.

Proof: Assume that $\sum a_k z^k$ converges for $|z| = \delta$ and $|z| = \sigma$. Then there exists C such that

$$|a_k| \delta^k, |a_k| \sigma^k < C$$

Then let $\chi \in \mathbb{R}_{\geq 0}^n$ such that

$$\log \chi = \alpha \log \rho + \beta \log \sigma, \quad \alpha + \beta = 1$$

then, if $k = (k_1, \dots, k_n)$,

$$\begin{aligned} \log \chi^k &= \log \chi_1^{k_1} \dots \chi_n^{k_n} \\ &= \sum_{i=1}^n k_i \log \chi_i \\ &= \alpha \sum k_i \log \rho_i + \beta \sum \log \sigma_i \\ &= \alpha \log \rho^k + \beta \log \sigma^k \\ &\leq \alpha \log \frac{C}{|a_k|} \\ |a_k \chi^k| &< C \end{aligned}$$

The final line implies, by our earlier lemma, that $\sum_k a_k z^k$ converges on D_χ . □

Now we will discuss the Identity theorem in $n \geq 2$.

Theorem: If f, g are holomorphic functions on $U \subset \mathbb{C}^n$, and $\exists V \subset U$ open such that $f|_V = g|_V$, then $f = g$.

Recall the proof in $n = 1$: We wish to show Ω , the set of points on which f and g disagree, is closed, open, and nonempty in \mathbb{C} . To see it is closed, we note $f - g$ is a continuous function, and singletons are closed in T_1 spaces, the pre-image of a closed set is closed, so $(f - g)^{-1}(0)$ is a closed set, and $(f - g)^{-1}(0) = \Omega$. To show it is open, if we take some $c \in \Omega$, so that $f(c) \neq g(c)$, then because f and g are holomorphic, we can expand them into a Taylor series, each of which have non-zero radii of convergence (in this case, poly radius, since they will converge on a polydisc), and f and g agree on the polydisc, so Ω is open. Ω is non-empty by definition. □

Proposition: For $f : U \subset \mathbb{C}^n \rightarrow \mathbb{C}$ holomorphic and non-constant, f is open.

Proof: If f is non-constant, for any $z \in U$, there is some line through z on which f is non-constant, i.e. there is another point on the line, w , such that $f(z) \neq f(w)$. Applying the 1-dimensional open mapping theorem, we have the result. □

We will do some preparatory work to talk about Weierstrass preparation theorem next lecture:

Lemma: Let $\epsilon > \epsilon' > 0$, and $f \in \mathcal{O}(\Delta_\epsilon)$, and $\lambda_1, \dots, \lambda_d \in \Delta_{\epsilon'}$ are the zeros of f counted with multiplicities (so if λ_i has multiplicity r , then it appears r times in the list $\{\lambda_1, \dots, \lambda_d\}$). Then for $k \geq 0$,

$$\sum_i^n \lambda_i^k = \frac{1}{2\pi i} \int_{|z|=\epsilon'} z^k \frac{f'(z)}{f(z)} dz$$

The LHS is what is known as a Newton polynomial. Note that the case of $k = 0$, we recover the typical method of counting zeros.

Proof: By the residue theorem, the RHS is

$$\sum_{w \in \Delta_{\epsilon'}(0)} \text{Res} \left(z^k \frac{f'(z)}{f(z)}, w \right)$$

We will compute this for fixed λ . In this case, we can rewrite

$$f(z) = (z - \lambda)^r h(z)$$

where $h(z) \neq 0$. Then the residue is

$$\begin{aligned} \text{Res} \left(z^k \frac{f'(z)}{f(z)}, \lambda \right) &= \text{Res} \left(z^k \frac{r(z - \lambda)^{r-1} h(z) + (z - \lambda)^r h'(z)}{(z - \lambda)^r h(z)}, \lambda \right) \\ &= \text{Res} \left(z^k \left(\frac{1}{z - \lambda} + \frac{h'(z)}{h(z)} \right), \lambda \right) \\ &= \lambda^k r + 0 \end{aligned}$$

Then summing over all λ gives us the result. □

Weierstrass Preparation Theorem

Lecture 3, Sept 3. Consider the ring of symmetric polynomial functions. This has obvious basis of the elementary symmetric polynomials: if there are d variables, there is one elementary symmetric polynomial of degree n for each $n \leq d$, made by adding products of each variable. For example, the $n = 3$ elementary symmetric polynomial is

$$\sum_{1 \leq j < k < l \leq d} \lambda_j \lambda_k \lambda_l$$

It is a theorem that you can also pick a basis as these Newton polynomials, or “power sums”, i.e. there are transformations called “Newton identities” which express these Newton polynomials in terms of the elementary symmetric polynomials. One can look them up, it becomes complicated with higher degrees. For now, all we need to know is that it can be done.

Theorem (Weierstrass Preparation Theorem): *Adopt the coordinates (z, w) for $z \in \mathbb{C}$ and $w \in \mathbb{C}^{n-1}$. Then if $U \subset \mathbb{C}^{n-1}$, $f \in \mathcal{O}(U \times \Delta_\epsilon)$ such that $f|_{U \times \partial \Delta_\epsilon}$ has no zeros of f , then $\forall \epsilon' < \epsilon$ sufficiently close, \exists a Weierstrass polynomial g , such that $f|_{U \times \Delta_{\epsilon'}} = gh$, for some $h \in \mathcal{O}(U' \times \Delta_{\epsilon'})$, for $U' \subset U$, and $h(0) \neq 0$.*

Proof: We may assume that $f|_{U' \times \Delta_{\epsilon'}}$ has no zeros. Then for $w \in U$, let $\lambda_1(w), \dots, \lambda_{d(w)}(w)$ be the set of zeros of $f|_{\{w\} \times \Delta_{\epsilon'}}$. By the previous lemma, the function

$$d(w) = \sum_{i=1}^{d(w)} \lambda_i^0 = \frac{1}{2\pi i} \int_{|z|=\epsilon} \frac{\partial_z f}{f} dz$$

is locally constant wrt w , and if we take U connected, it is independent of w . So we define g pointwise:

$$g(z, w) = \prod_{i=1}^d (z - \lambda_i(w)) = z^d + \alpha_1(w)z^{d-1} + \dots + \alpha_d(w)$$

where $\alpha_i = (-1)^i \delta_i(\lambda_1, \dots, \lambda_d)$

where δ_i is the elementary symmetric polynomial. Then the fact from last lecture: The elementary symmetric functions can be expressed in terms of the Newton polynomials. But as the lemma showed, the Newton polynomials, for every k , are just integrals, and as

such are holomorphic. So the α_i are holomorphic, and so is g . Then defining $h = \frac{f}{g}$ gives the result. □

The reason this theorem is important in complex analysis is because it gives us tools to study the 0 loci of holomorphic functions, so-called analytic sets. It is saying that to study the zeros of f in this region, it is sufficient to study g , a Weierstrass polynomial, which is a factor of f .

Example: Apply WPT to the polynomial $f(z_1, z_2) = z_1^3 z_2 + z_1 z_2 + z_1^2 z_2^2 + z_2^2 + z_1 z_2^3$

We recall that given a real vector space, V , we can always get a complex vector space by complexifying V to $V_{\mathbb{C}} := V \otimes_{\mathbb{R}} \mathbb{C}$. Note that there is a canonical isomorphism $V_{\mathbb{C}} := V \otimes_{\mathbb{R}} \mathbb{C} = V \oplus V$ decomposing $V_{\mathbb{C}}$ into “real” and “complex” parts. Then the complex structure on $V \oplus V$ is the map $(v, w) \rightarrow (-w, v)$. If we think of (v, w) as $v + iw$, this choice of complex structure is natural. This is an important point because, given an abstract complex vector space, there is no notion of real or imaginary parts. The complex vector spaces which are obtained as tensor products of real vector spaces have this natural decomposition.

Definition: Given a complex vector space, V , a real structure on V is a conjugate linear involution $\kappa : V \rightarrow V$.

If an abstract complex vector has a real structure, then the fixed locus of κ , V^{κ} , defines an isomorphism

$$V^{\kappa} \otimes_{\mathbb{R}} \mathbb{C} \rightarrow V$$

which gives a decomposition of V into real and imaginary parts. Any complex vector space can be given a real structure, so that one may decompose vectors into real and imaginary parts, but only those which arise as a tensor product have this structure a priori. If handed an abstract complex vector space, one must impose this structure by picking a basis of V . But what would happen if you wanted to complexify a complex vector space? This happens already in complex manifold theory, where we study the complexified tangent bundle. But the tangent bundle already has a complex structure, that’s the definition of a complex manifold! If you complexify a complex vector space, V , you get an isomorphism $V \otimes_{\mathbb{R}} \mathbb{C} \cong V \oplus \bar{V}$, where \bar{V} is the complex vector space with the same underlying real vector space as V , but with the opposite complex structure, $-I$.

Finally, we will talk about analytic sets

Definition: An analytic set $Z \subset U$ is a set such that $\forall z \in U, \exists V \subset U$ open such that $\exists f_1, \dots, f_{\gamma} \in \mathcal{O}(\bar{V}), Z \cap V = \{z \in V \mid f_1(z) = \dots = f_{\gamma}(z) = 0\}$.

So an analytic set locally looks like the zero locus of finitely many holomorphic functions.

Definition: The germ of an analytic set at $0 \in \mathbb{C}^n$ is an equivalence class of pairs (U, Z) , where U is an open neighborhood of $0 \in \mathbb{C}^n$ and Z is an analytic subset of U . The equivalence relation imposed is $(U, Z) \sim (U', Z')$ iff $\exists V \subset (U \cap U')$ a neighborhood of 0 such that $Z \cap V = Z' \cap V$.

So the usual notion of a germ, that is setting things to be equivalent if they agree locally, but now locally means “on an analytic subset”. If α_i are elements of a ring, we denote $I(\alpha_i)$ to be the ideal generated by those elements. The germ of an analytic set Z is denoted $(Z, 0)$, so for example, the germ of a zero locus of an ideal I is written $(Z(I), 0)$.

Lemma: *The set of germs is in bijection with the set of finitely generated ideals in the ring of convergent power series of \mathbb{C}^n .*

Proof: Given a germ of an analytic set, let f_1, \dots, f_n be those functions which realize the definition of an analytic set, say at a point z . Then (f_1, \dots, f_n) generated a finitely generated ideal in the ring of convergent power series, since f_i are holomorphic. Conversely, given a finitely generated ideal let f_1, \dots, f_n be a finite generating set. Then $Z(f_1, \dots, f_n) := \{z \in U \mid f_i(z) = 0\}$ is analytic, and so we can take its germ. To do this, we made a choice of generators f_i , but the germ is well defined under this choice. This comes from the general fact that if I is finitely generated, then $Z(I)$ is well defined: If g_i is another generating set, then

$$\begin{aligned} \forall \alpha, f_\alpha &= \sum_i \lambda_i g_i, \quad \text{for } \lambda_i \in \mathcal{O}_{\mathbb{C}^n, 0} \\ &\Rightarrow Z(f_i) \subset Z(g_i) \end{aligned}$$

and vice versa¹.

□

This lemma is very reminiscent of Hilbert’s Nullstellensatz.

Definition: If I is an ideal, the radical of I is the set $\sqrt{I} := \{z \in R \mid \exists d > 0 \text{ s.t. } z^d \in I\}$, and if I is an ideal such that $I = \sqrt{I}$, then we call I a radical ideal.

Lemma: $Z(I) = Z(\sqrt{I})$.

Proof:

$$\begin{aligned} z_0 \in Z(I) &\Rightarrow f(z_0) = 0 \forall f \in I \\ \forall g \in \sqrt{I}, g^k &= f \text{ for some } f \in I \\ \Rightarrow g^k(z_0) &\equiv \underbrace{g(z_0) \dots g(z_0)}_{k \text{ times}} = f(z_0) = 0 \Rightarrow g(z_0) = 0 \end{aligned}$$

Conversely, $I \subset \sqrt{I} \Rightarrow Z(I) \subset Z(\sqrt{I})$.

□

Properties of $\mathcal{O}_{\mathbb{C}^n, 0}$: i) It is a UFD, ii) It is Noetherian, i.e. all ideals are finitely generated.

¹Siebert did not prove the statement that $(Z(I), 0)$ was well defined, and I doubt my own proof here. I used neither the fact that we took germs of $Z(I)$ instead of just $Z(I)$, nor the fact that it was finitely generated. If my proof is correct, then the bijection statement could have been made without any mention of germs, which is probably wrong.

In the future, we will study the basic correspondence we have just outlined:

$$\begin{aligned}
 \left\{ I \subset \mathcal{O}_{\mathbb{C}^n,0} \mid I = \sqrt{I} \right\} &\longleftrightarrow \{ (Z,0) \subset (\mathbb{C}^n,0) \text{ analytic} \} \\
 I &\longmapsto Z(I) \\
 \{ f \in \mathcal{O}_{\mathbb{C}^n,0} \mid Z \subset Z(f) \} &=: I(Z) \longleftarrow (Z,0)
 \end{aligned}$$

Differential Topology and Local Rings

Lecture 4, Sept 8. If we have a holomorphic function $f : U \subset \mathbb{C}^m \rightarrow \mathbb{C}^n$, and adopt the coordinates (z_i) on the domain and (w_i) on the codomain, with $z_i = x_i + iy_i$ and $w_i = u_i + iv_i$, and $f_i = g_i + ih_i$, the Jacobian of f is written as

$$J_{\mathbb{R}}(f)(z) = \begin{pmatrix} \partial_{x_j} g_i & \dots & \partial_{y_j} g_i \\ \vdots & & \vdots \\ \partial_{x_j} h_i & \dots & \partial_{y_j} h_i \end{pmatrix} \in M(2n \times 2m, \mathbb{R})$$

so this is a real matrix, and suppose we want to study its complexification. We recall that $\cdot \otimes_{\mathbb{R}} \mathbb{C}$ is a functor, so that we may apply it to maps as well, and it acts on maps by tensoring with the identity on the second component, so

$$J_{\mathbb{C}}(f)(z) = J_{\mathbb{R}}(f)(z) \otimes_{\mathbb{R}} Id_{\mathbb{C}} : V \otimes_{\mathbb{R}} \mathbb{C} \rightarrow W \otimes_{\mathbb{R}} \mathbb{C}$$

Note¹: $\det J_{\mathbb{R}}(f)(z) = \det J_{\mathbb{C}}(f)(z)$. We can write down this Jacobian for any real differentiable map. But what if f also happens to be holomorphic? As we discussed in the first lecture², we have

$$J_{\mathbb{C}}(f)(z) = \begin{pmatrix} \partial_{z_j} f_i & \dots & \partial_{\bar{z}_j} f_i \\ \vdots & & \vdots \\ \partial_{z_j} \bar{f}_i & \dots & \partial_{\bar{z}_j} \bar{f}_i \end{pmatrix} = \begin{pmatrix} \partial_{z_j} f_i & 0 \\ 0 & \partial_{z_j} \bar{f}_i \end{pmatrix}$$

So the Jacobian is really determined by one, smaller matrix in the upper left, which we will denote $J(f)(z)$. Then of course, $\det J_{\mathbb{C}}(f)(z) = |\det J(f)(z)|^2 \geq 0$. As a result, f is orientation preserving iff it is non-singular, and it cannot reverse orientation. We will see later that this is exactly the reason that complex manifolds are always orientable.

Definition: A map $f : U \rightarrow V$, for $U, V \subset \mathbb{C}^n$ open is biholomorphic if it is holomorphic and bijective, and its inverse is holomorphic.

This is the complex analogue of a diffeomorphism. We recall the

¹I actually don't know why this is. I know there's a formula to tensor two matrices together, but there has to be a better explanation than that. Maybe thinking about \det as the map on the top ext power

²Maybe it was second I don't remember

Theorem (Inverse Function Theorem): If f is a map $f : U \rightarrow V$ for $U, V \subset \mathbb{C}^n$, and z is a regular value (or equivalently³, $\det J(f)(z) \neq 0$), then there exist open sets $U', V' \subset \mathbb{C}^n$ such that $f|_{U'}$ is a biholomorphism onto V' .

Proof: We note that z is a regular value of f iff $\det J_{\mathbb{R}}(f)(z) \neq 0$, so the real inverse function theorem tells us that there exists a smooth inverse to f , $g : V' \rightarrow U'$. To see that it is also holomorphic, we show it satisfies the CR equations⁴:

$$\begin{aligned} 0 &= \partial_{\bar{z}_j}(f^{-1} \circ f) \\ &= \sum_{k=1}^n (\partial_{w_k} f^{-1} \cdot \partial_{\bar{z}_j} f_k + \partial_{\bar{w}_k} f^{-1} \cdot \partial_{\bar{z}_j} \bar{f}_k) \\ &= \sum_{k=1}^n \bar{\partial}_{z_j} \bar{f}_k \cdot \partial_{\bar{w}_k} f^{-1} \\ \Rightarrow J(f) &= (\partial_{z_j} f_k)_{ij} \text{ invertible} \Rightarrow \forall k \partial_{w_k} f^{-1} = 0 \end{aligned}$$

so that f^{-1} is holomorphic. □

Lemma: If $f : U \rightarrow V$ is holomorphic and bijective, then f is a biholomorphism.

We note that in the real case (replacing holomorphic and biholomorphic with smooth and diffeomorphic, respectively), this is not true. For example, take $y = x^3$.

Theorem (Implicit Function Theorem)⁵: Let $U \subset \mathbb{C}^m \times \mathbb{C}^n$ have coordinates (z, w) , and $f : U \rightarrow \mathbb{C}^n$ be holomorphic. If $J(f) = (\partial_{w_j} f_i)_{ij}$ is invertible at $(z_0, w_0) \in U$, and $f(z_0, w_0) = 0$, then $\exists U_1 \subset \mathbb{C}^m, U_2 \subset \mathbb{C}^n$ open, $U_1 \times U_2 \subset U$, $(z_0, w_0) \in U_1 \times U_2$, $g : U_1 \times U_2$ such that

$$\{(z, w) \in U_1 \times U_2 \mid f(z, w) = 0\} = \Gamma_g = \{(z, g(z)) \mid z \in U_1\}$$

Proof: Again we have the existence of such a function g from the real version of the implicit function theorem, and apply the chain rule to show that g is in fact holomorphic. □

We will now work our way up to the analytic Nullstellensatz. We recall that, essentially by definition, $\mathcal{O}_{\mathbb{C}^n, 0} = \mathbb{C}\{z_1, \dots, z_n\}$.

Definition: A ring is called local if it has a unique maximal ideal.

Lemma: The ring described above is local.

³It took me about 30 minutes to figure out, but they're equivalent because f is a map from $\mathbb{C}^n \rightarrow \mathbb{C}^n$, so the differential is a map between vector spaces of the same dimension, thus surjectivity implies isomorphism. Felt quite silly upon realizing this.

⁴The first equality holds because the identity map is holomorphic

⁵Bernd notes that in Germany, during oral exams for undergraduates, students would reliably fail at correctly formally stating the implicit function theorem. Indeed, this exact scenario happened to me in an undergraduate class last year. One part of one of our problems on an exam was to simply state the implicit function theorem, and I only received about half credit on that part.

Proof: Claim the unique maximal ideal is

$$\mathcal{M} := \left\{ f = \sum a_I z^I \mid a_0 = 0 \right\} = \{ f \in \mathcal{O}_{\mathbb{C}^n,0} \mid f(0) = 0 \}$$

This clearly defines an ideal. To show it is maximal, there is a well defined evaluation homomorphism:

$$\begin{aligned} ev_0 : \mathcal{O}_{\mathbb{C}^n,0} &\rightarrow \mathbb{C} \\ f &\mapsto f(0) = a_0 \end{aligned}$$

We see that the quotient $\mathcal{O}_{\mathbb{C}^n,0}/\mathcal{M} \cong \mathbb{C}$, and we recall that an ideal is maximal iff the quotient ring is a field, which \mathbb{C} is. □

We note that this property of locality is not unique to holomorphic functions: Smooth and continuous functions also share this property, for example.

Lemma: $\mathcal{O}_{\mathbb{C}^n,0}$ is an integral domain

Proof: There is a tedious, algebraic method to prove this lemma, but the geometric way is cleaner, via the identity theorem. WLOG, we may assume $f, g \in \mathcal{O}(D_\epsilon)$. Then if $f \not\equiv 0$, $\exists z \in D_\epsilon$ such that $f(z) \neq 0 \Rightarrow \exists D_{\epsilon'}(z)$ such that $f|_{D_{\epsilon'}(z)} \neq 0$. But then $g|_{D_{\epsilon'}(z)} \equiv 0$, so the identity theorem shows $g = 0$. □

Lemma: $\mathcal{O}_{\mathbb{C}^n,0}$ is a UFD.

Idea of Proof: The Wierstrass Preparation Theorem reduces this proof to a known statement about polynomials. □

Definition: A hypersurface is the set of zeros $Z(f) \subset U \subset \mathbb{C}^n$ for $f \in \mathcal{O}(U \setminus \{0\})$. The germ of a hypersurface is then the germ of the zero locus.

Corollary: Any germ of hypersurfaces decomposes uniquely into germs of irreducible hypersurfaces.

Proof: If we decompose f into irreducibles, $f = g_1 \dots g_r$, then $Z(f) = \cup Z(g_i)$. Then we must show the $Z(g_i)$ are irreducible. To do so, we must use a lemma about hypersurface containment:

Lemma: If $f, g \in \mathcal{O}_{\mathbb{C}^n,0}$ and g irreducible, then $Z(g) \subset Z(f) \Rightarrow g \mid f$.

With this lemma, the result can be found. □

Lec 5

Lecture 5, Sept 10. I missed lecture 5

Noether Normalization

Lecture 6, Sept 15. Theorem (Local Finite Mapping/Noether Normalizing): Let $X \subset (\mathbb{C}^n, 0)$ be a germ of analytic sets. Then there exists $d \leq n$, and a linear coordinate system $(z_1, \dots, z_d, z_{d+1}, \dots, z_n)$ such that the projection

$$\pi : \mathbb{C}^n \rightarrow \mathbb{C}^d, (z_1, \dots, z_n) \mapsto (z_1, \dots, z_d)$$

induces a proper, open surjection $X \rightarrow (\mathbb{C}^d, 0)$ with finite fibers, i.e. is a branched cover. We refer to d as the dimension of X .

Non-example: Take $\pi : \mathbb{C}^3 \rightarrow \mathbb{C}^2$ such that $(x, y, z) \mapsto (x, y)$, and take $X = Z(y - zx)$. Then X is irreducible, but π does not have finite fibers: $\pi^{-1}(0, 0) = \{0, 0\} \times \mathbb{C}$. This is an example of a blow-up. The theorem does not apply because $\mathcal{O}_{\mathbb{C}^2, 0} \rightarrow \mathcal{O}_{X, 0}$ is not an integral ring extension, which is a necessary step for the algebraic proof of the above theorem.

A priori, we may believe d depends on the choice of projection. There are 3 ways to show that this notion of dimension is well defined: a) We can show that outside of a nowhere dense, analytic subset, this is an unbranched covering. Then working locally, if we assume there is another covering π' of dimension d' , then we have a diffeomorphism from a subset of \mathbb{C}^d to $\mathbb{C}^{d'}$, and we are done. For the second method, we define

Definition: The Krull dimension of a commutative ring R is the supremum of the lengths of all chains of prime ideals:

$$\dim R = \sup\{k \mid f_0 \subset f_1 \subset \dots \subset f_k \in \mathcal{O}_{X, 0}\}$$

where f_i are prime ideals, and the containments are proper

So 0-dimensional rings correspond to maximal ideals. If you take a hypersurface inside X , that has codimension 1, then you continue until you reach an irreducible component. So if X is irreducible, $\mathcal{O}_{X, 0}$ is an integral domain. The final way of showing this is related to stratifications.

Definition: If X is an analytic subset, then the singular locus

$$X_{\text{sing}} = \{x \in X \mid X \text{ is not a submfld at } x\} \subset X$$

is the locus where X is not a submanifold.

X_{sing} can have any codimension. For example, take two copies of \mathbb{C}^2 and have them meet at the origin of \mathbb{C}^4 transversely, this has codim 2. Of course, the containment $X_{\text{sing}} \subset X$ must be a proper containment. Then we may take the series of successive singularities: the next step would be $(X_{\text{sing}})_{\text{sing}}$. Each set is also analytic, though that must be proved, and the dimension decreases with each step, so this series will always terminate. Then we define the dimension of X to be complex dimension $\dim X \setminus X_{\text{sing}}$.

So now we have an idea of what the zero loci of holomorphic functions in several complex variables looks like. They are closed, they admit these stratifications, locally they are coverings of domains in \mathbb{C}^d , and they have dimension d .

Onto complex manifolds.

Complex Projective Space

Lecture 7, Sept 17. We will start with projective space. Bernd remarks that $\mathbb{C}P^n$ is arguably the most important complex manifold, even moreso for compact complex manifolds. As a set,

$$\mathbb{C}P^n = (\mathbb{C}^{n+1} \setminus \{0\}) / \mathbb{C}^*$$

where \mathbb{C}^* acts diagonally. This induces the quotient topology. We note there is a surjection $S^{2n+1} \rightarrow \mathbb{C}P^n$, the antipodal map, which implies the latter is compact. The standard charts on $\mathbb{C}P^n$ is

$$U_i := \{[z_0, \dots, z_n] \in \mathbb{C}P^n \mid z_i \neq 0\}$$

Clearly these U_i cover $\mathbb{C}P^n$, and are biholomorphic to \mathbb{C}^n via

$$\begin{aligned} \varphi_i : U_i &\rightarrow \mathbb{C}^n \text{ parameterized by } w_0, \dots, \hat{w}_i, \dots, w_n \\ [z_0, \dots, z_n] &\mapsto \frac{1}{z_i}(z_0, \dots, \hat{z}_i, \dots, z_n) \end{aligned}$$

It is easily seen that the transition functions are biholomorphic.

Similarly, if V is a finite dimensional complex vector space, then can define $\mathbb{C}P(V)$ via charts, but we then need to worry about change of basis. If we have a change of basis $T : \mathbb{C}^{n+1} \rightarrow \mathbb{C}^{n+1}$, i.e. a linear isomorphism, we get a functor from complex vector spaces to complex manifolds, \mathbb{P} , which sends

$$T \mapsto \mathbb{P}(T) : \mathbb{C}P^n \rightarrow \mathbb{C}P^n$$

is a biholomorphism. Note we induce a map on the quotient simply by linearity. Another way to view this, we have a group homomorphism

$$F : GL(n+1, \mathbb{C}) \rightarrow \{\text{biholomorphisms of } \mathbb{C}P^n\}$$

What is $\text{Ker } F$? Clearly, all the scalar diagonal matrices, i.e. those of the form cA for $c \in \mathbb{C}^*$ and $A \in GL(n, \mathbb{C}^{n+1})$, because multiplication by a scalar preserves the $\mathbb{C}P^n$ class, by definition. It turns out this is the entire kernel, and the quotient by this normal subgroup $\text{Ker } F$ is called the projective linear group, $PGL(n+1, \mathbb{C}^{n+1})$.

There is an easy way to write down submanifolds of projective space via the homogeneous coordinates.

Definition: If $f \in \mathbb{C}[z_0, \dots, z_n]$, then it is called homogeneous of degree d if $\exists d$ s.t. $\forall \lambda \in \mathbb{C}^*$,

$$f(\lambda z_0, \dots, \lambda z_n) = \lambda^d f(z_0, \dots, z_n)$$

If f is as above, we may now define the zero locus $Z(f) := \{[z_0, \dots, z_n] \in \mathbb{CP}^n \mid f(z_0, \dots, z_n) = 0\}$. Note this is well defined because f is homogeneous: if

$$\begin{aligned} f(z_0, \dots, z_n) &= 0 \\ \Rightarrow f(\lambda z_0, \dots, \lambda z_n) &= \lambda^d f(z_0, \dots, z_n) = 0 \end{aligned}$$

Also note that f cannot be a function on \mathbb{CP}^n , by the maximum principal: a non-constant holomorphic function does not have local maxima. But because \mathbb{CP}^n is compact, every f would achieve a maximum, and thus must be constant. f would be holomorphic, since it is a polynomial in z_i , thus it must not be a function¹. Locally, we can easily compute its zero locus via charts: on U_i , $\varphi_i(Z(f) \cap U_i) = Z(f_i)$, where $f_i(w_0, \dots, \hat{w}_i, \dots, w_n) = f(w_0, \dots, 1, \dots, w_n)$ with the 1 in the i th index. This f_i is the “dehomogeneized” form of f .

Definition: A complex hypersurface in \mathbb{CP}^n is a set of the form $Z(f)$ for $f \in \mathbb{C}[z_0, \dots, z_n] \setminus \mathbb{C}$. As usual, given an ideal of homogeneous elements, $S \subset \mathbb{C}[z_0, \dots, z_n]$, we have a corresponding zero locus,

$$Z(S) := \bigcap_{f \in S} Z(f)$$

As before, these are analytic subsets, and the implicit function theorem shows that these are also submanifolds. We note that because $\mathbb{C}[z_0, \dots, z_n]$ is Noetherian, all ideals are finitely generated. Then because the zero locus of S is the same as the zero locus of the ideal it generates then pick a finite generating set, WLOG, we can treat S as being finite.

Theorem (Chow): $X \subset \mathbb{CP}^n$ is analytic $\Rightarrow X = Z(S)$ for some S .

i.e. every \mathbb{C} -analytic set in \mathbb{CP}^n is algebraic, i.e. defined globally by the zero set of homogeneous polynomials.

Question: Is every complex manifold algebraic?

In algebraic geometry, we look at the field $K(X)$, the field of rational functions on X , with transcendence degree² $\dim X$.

In complex analysis, there is a similar idea:

Definition: A meromorphic function on a complex manifold X is an equivalence class of pairs (U, f) with $U \subset X$ open and f a holomorphic function on U , $X \setminus U$ analytic and nowhere dense, such that $\forall x \in X, \exists V \ni x \subset X$ open, $g, h \in \mathcal{O}(V)$ such that h has no zeros on $V \setminus U$ and

$$f|_{V \cap U} = \frac{g}{h}|_{V \cap U}$$

¹I think this means that it would not be a well defined function, despite the fact that its zero locus is well defined. I don't know how else it could fail to be a function.

²don't really know what this is.

where $(U, f) \sim (V, g) \iff f|_{U \cap V} = g|_{U \cap V}$. Here we denote $\mathcal{K}(X)$ as the field of meromorphic functions on X .

Rather long winded definition, but the idea is that locally these should just look like quotients of holomorphic functions. Recall some of the behavior of single variable complex functions near singularities: Near an essential singularity, the function maps a neighborhood to the entire plane minus one or two points. Near a pole, you experience a blow up, and the Riemann extension theorem says that if you don't blow up³, then the singularity is removable.

Example: Take $f = \frac{z}{w}$ on $\mathbb{C}_{z,w}^2$. Clearly there is a pole along the $w = 0$ line. If you approach the horizontal axis ($w = 0$) from a vertical direction, you will get infinity. However, if you approach from an angle sufficiently close to vertical, say approaching the origin following the square root function, you get 0. So singularities are more complicated in multiple dimensions, because you cannot resolve this issue even by compactifying.

We note that multiplication on $K(X)$ is given by

$$(U, f) \cdot (V, g) = (U \cap V, f \cdot g)$$

and inversion given by

$$(U, f)^{-1} = (U \setminus Z(f), f^{-1})$$

and one can check that this is compatible with the equivalence classes.

Definition: The algebraic definition of X a compact complex manifold is the transcendence degree of $\mathcal{K}(X)$.

So given X , one can look at the field of rational functions, which is a subfield of $\mathcal{K}(X)$. It is not clear that these are the same, but another version of Chow's theorem says that any rational function is algebraic, so indeed they are the same field, i.e. that any meromorphic function can be written as the quotient of two homogeneous polynomials of the same degree. So

Theorem: (Remmert-Siegel) *If X is a compact, complex manifold, then the algebraic dimension is the transcendence degree of $\mathcal{K}(X)$, which is less than or equal to $\dim X = n$.*

Proof sketch: Cover the manifold with polydiscs. Then we take $f_1, \dots, f_{n+1} \in \mathcal{K}(X)$. Then we want to build an $F \in \mathbb{C}[x_1, \dots, x_{n+1}]$ such that $F(f_1, \dots, f_{n+1}) = 0$ and has zeros of high order at the centers of the polydiscs. So say the degree of F is m , and the order of vanishing is, say, m' . Then we recall the Schwarz' lemma: If you have a holomorphic function on the disk which vanishes at the origin to order k , then you can bound the growth of the function polynomially, with coefficients the norm of f . So for each polydisc, there is a smaller polydisc on which we can estimate the growth of F . You then argue that you can choose m and m' such that the Schwarz lemma must give a contradiction unless $F(f_1, \dots, f_{n+1}) = 0$, as desired.

In the next lecture, we will go through this proof in more detail.

³ie the function is locally bounded near all of its zeros

Proof of Remmert-Siegel

Lecture 8, Sept . We will now go through a proof of the previous theorem. So first, we pick f_i as $n + 1$ meromorphic functions on X . As we discussed last time, these functions locally look like quotients of holomorphic functions, i.e. $\forall x \in X, \exists U_x \subset X$ open such that $x \in U$ and $f_i|_{U_x} = \frac{g_{i,x}}{h_{i,x}}$, for some $g_{i,x} \in \mathcal{O}(U_x)$ and $h_{i,x} \in \mathcal{O}(U_x) \setminus \{0\}$. We also know if $g_{i,x}, h_{i,x} \in \mathcal{O}_{X,x}$ we may assume g and h are relatively prime, since $\mathcal{O}_{X,x}$ is a UFD. So at a point, the germs of these functions can be taken to be relatively prime.

Lemma: *If $g_{i,x}$ and $h_{i,x}$ are relatively prime at $\mathcal{O}(X, x)$, by possibly shrinking U_x , it holds that they are relatively prime on the whole neighborhood U_x , i.e. $\forall y \in U_x$, they are relatively prime in $\mathcal{O}_{X,y}$.*

Now the U_x give an open covering of X . By shrinking U_x to a chart, we choose refinements of $\{U_x\}$ such that

$$\begin{array}{ccccccc} \{x\} & \hookrightarrow & W_x & \hookrightarrow & V_x & \hookrightarrow & \overline{V_x} \hookrightarrow U_x \\ & & \downarrow \wr & & \downarrow \wr & & \\ & & D_{\frac{1}{2}}(0) & \hookrightarrow & D_1(0) & \hookrightarrow & \mathbb{C}^n \end{array}$$

Then because X is compact, the open cover $\{W_x\}$ has a finite subcover, $\{W_{x_i}\}$, with centers x_1, \dots, x_N , and we forget the double indexing, so define $\{W_k\} = \{W_{x_k}\}$, and the same for the V_x , and g and h . Now we will construct the F . We want it to satisfy

$$\partial^I F \left(\frac{g_{1,k}}{h_{1,k}}, \dots, \frac{g_{n+1,k}}{h_{n+1,k}} \right) (x_k) = 0 \quad \forall |I| < m'$$

so all the holomorphic derivatives, that is derivatives wrt z_i , not \bar{z}_i , as those vanish anyway, vanish. On the intersection $U_k \cap U_l$,

$$\begin{aligned} \frac{g_{i,k}}{h_{i,k}} &= \frac{g_{i,l}}{h_{i,l}} \\ \Rightarrow h_{i,k} &= h_{i,l} \cdot \varphi_{i,kl} \quad \text{where } \varphi_{i,kl} := \frac{g_{i,k}}{g_{i,l}} \end{aligned}$$

We then claim that $\varphi_{i,kl}$ has no zeros or poles¹. To see this, take some $y \in U_k \cap U_l$. Then by unique factorization, $\varphi_{i,kl} = \frac{\alpha}{\beta}$ for α, β coprime in $\mathcal{O}_{X,y}$. So $\varphi_{i,kl}$ is a unit in $\mathcal{O}(U_k \cap U_l)$. We then define

$$\varphi_{kl} = \prod_{i=1}^{n+1} \varphi_{i,kl}$$

which is holomorphic and bounded² on $\overline{V}_k \cap \overline{V}_l$. Also define

$$C := \max\{\|\varphi_{kl}\|\}_{k,l=1,\dots,N}$$

where $\|\cdot\|$ is the sup norm on $\overline{V}_k \cap \overline{V}_l$. Becuase this is closed, we have a max instead of sup. We note

$$\begin{aligned} \varphi_{kl} \cdot \varphi_{lk} &= \left(\prod_{i=1}^{n+1} \frac{g_{i,k}}{g_{i,l}} \right) \cdot \left(\prod_{i=1}^{n+1} \frac{g_{i,l}}{g_{i,k}} \right) = 1 \\ &\Rightarrow C \geq 1 \end{aligned}$$

Generally, we are going to use this φ_{kl} to clear denominators. Specifically, if we have $F\left(\frac{g_{1,k}}{h_{1,k}}, \dots, \frac{g_{n+1,k}}{h_{n+1,k}}\right)$ of degree m , then define $G_k \in \mathcal{O}(U_k)$ by

$$F = \frac{G_k}{\left(\prod_{i=1}^{n+1} h_{i,k}\right)^m}$$

We have that

$$G_k = \varphi_{lk}^m G_l$$

Now, assume we have

$$\binom{m+n+1}{m} > N \binom{m'-1+n}{m'-1}$$

and call this the $(*)$ condition. Claim $(*)$ implies there exists F of degree m satisfying

$$\partial^I F\left(\frac{g_{ik}}{h_{ik}}\right)(x_k) = 0$$

Indeed this is just a dimension count. There will be a nonzero kernel, so we can choose such an F . Now we may apply Schwarz' lemma, which says

$$\begin{aligned} f &\in \mathcal{O}(\overline{D_\epsilon(0)}), \partial_z^k f(0) = 0 \quad k = 0, \dots, m' - 1 \\ &\Rightarrow \forall z \in \overline{D_\epsilon(0)}, f(z) \leq \|f\| \cdot \left(\frac{|z|}{\epsilon}\right)^{m'} \end{aligned}$$

¹Can you not just see this from the equality $h_{i,k} = h_{i,l} \cdot \varphi_{i,kl}$? If φ had any zeros/poles, either $h_{i,k}$ has to have them as well, or $h_{i,l}$ has to cancel them somehow, which also shouldn't be possible.

²I think this is because we showed there are no poles/zeros.

This holds for higher dimensions. So we have that

$$\begin{aligned} \partial^I F(\dots)(x_k) = 0 &\Rightarrow \partial^I G_k(x_k) = 0 \forall k, \forall |I| < m' \\ \Rightarrow \forall x \in \overline{W}_k, |G_k(x_k)| &\leq \left(\frac{1}{1}\right)^{m'} \cdot C', \quad C' := \|G_k\|_{W_k} \end{aligned}$$

Now fix k , and $x \in \overline{W}_k$, which may not be in W_k , and write $C' = |G_k(x)|$. Now, X is covered by W_l , so $x \in W_l$ for some l . But on W_l , we have the estimate

$$\begin{aligned} C' = |G_k(x)| &= |G_l(x)| \cdot |\varphi_{lk}^m(x)| \leq \frac{C'}{2^{m'}} \cdot C^m \\ &\Rightarrow \left|1 - \frac{C^m}{2^{m'}}\right| C' \leq 0 \end{aligned}$$

this inequality will lead to a contradiction unless $G_k \equiv 0$ if we can show the first term is strictly larger than 0, forcing $C' = 0$. To do that, need to find m, m' such that

$$\frac{C^m}{2^{m'}} < 1$$

So because $C \geq 1$, we may write it as $C = 2^\lambda$ for some $\lambda \in \mathbb{R}_{\geq 0}$. Then

$$\frac{C^m}{2^{m'}} = 2^{\lambda m - m'} < 1 \iff m' > \lambda m$$

Recall we also wanted m and m' to satisfy

$$\binom{m+n+1}{m} > N \binom{m'-1+n}{m'-1}$$

Writing out the binomial expansion, one would see

$$\begin{aligned} \exists a, b \mid \binom{m+n+1}{m} &> a \cdot m^{n+1} \text{ for } m \gg 0 \\ N \cdot \binom{m'-1+n}{m'-1} &< b \cdot (m')^n \text{ for } m' \gg 0 \end{aligned}$$

Now take $m' = 2\lambda m$, so we satisfy $m' > \lambda m$. Then we can show

$$(*) \Leftarrow b \cdot (2\lambda)^n m^n = \frac{b \cdot (2\lambda)^n}{m} \cdot m^{n+1} < a \cdot m^{n+1}$$

The inequality follows because we can choose m so large that the coefficient in front of m^{n+1} is arbitrarily small. Working backwards from this conclusion, we can construct a proof of the theorem. □

We may consider complex tori $\pi : \mathbb{C}^n \setminus \mathbb{Z}^{2n}$, endowed with the quotient topology. In fact, π is a local homeomorphism, making the complex torus a complex manifold. More generally, if V is a complex vector space, and Γ is a free abelian discrete group of maximal rank $2n$. Then under a basis $b : V \rightarrow \mathbb{C}^n$,

$$\Gamma \mapsto \mathbb{Z}^n + \mathbb{Z}\tau_1 + \dots + \mathbb{Z}\tau_n$$

Then as a \mathbb{C}^∞ manifold, $X \cong (S^1)^n$, but $\mathbb{C}^n \setminus \Gamma \cong \mathbb{C}^n \setminus \Gamma' \iff \exists A \in GL(n, \mathbb{C}) \ A(\Gamma) = \Gamma'$.

Complex Tori and Holomorphic Group Actions

Lecture 9, . Weierstrass p -functions: We recall the Liouville theorem, which states that any bounded holomorphic function on \mathbb{C} is a constant. We want to consider periodic functions, but by Liouville, that would give us a function on the torus, which is bounded, and thus constant. However, meromorphic functions don't have this restriction, so we examine some such functions. One such function is the Weierstrass p -function, given by, for $\Gamma = \mathbb{Z} + \mathbb{Z}\tau \subset \mathbb{C}$

$$p(z) := \frac{1}{z^2} + \sum_{\omega \in \Gamma \setminus \{0\}} \left(\frac{1}{(z - \omega)^2} - \frac{1}{\omega^2} \right), \quad z \in \mathbb{C}, \omega \notin \Gamma$$

So we create a function with second order poles at each ω , with an essential singularity at the origin. It is not obvious that this converges, but it does. It is also Γ -periodic: $\forall \omega \in \Gamma, p(z + \omega) = p(z)$. It is easy to see its derivative is Γ periodic:

$$p'(z) = -2 \sum_{\omega \in \Gamma} \frac{1}{(z - \omega)^3}$$

But since the original function is just the integral of a periodic function, it is also periodic. So these are two meromorphic functions on a 1-dimensional manifold. By Remmert-Siegel, then, these are algebraically dependent. In fact,

$$p'^2 = 4p^3 - g_2p - g_3$$

where

$$g_2 = g_2(\tau) = 60 \cdot \sum_{(m,n) \neq (0,0)} \frac{1}{(m + n\tau)^4} \in \mathcal{O}(\mathbb{H})$$

$$g_3 = g_3(\tau) = 140 \cdot \sum_{(m,n) \neq (0,0)} \frac{1}{(m + n\tau)^6} \in \mathcal{O}(\mathbb{H})$$

These are called the elliptic modular forms of weight 4 and 6. Now, we may define the map $\mathbb{C} \setminus \Gamma \rightarrow \mathbb{C}^2 \rightarrow \mathbb{CP}^2$ via

$$z \mapsto (p(z), p'(z)) \mapsto [p(z), p'(z), 1]$$

which extends to a map from $\mathbb{C} \rightarrow \mathbb{CP}^2$. We can also quotient \mathbb{C} by Γ , call the result E , then we induce a map $E \rightarrow \mathbb{CP}^2$. Then we have

Theorem: $E = \mathbb{C} \setminus \mathbb{Z} + \mathbb{Z}\tau, \tau \in \mathbb{H}$ is biholomorphic to the plane cubic curve $C \subset \mathbb{CP}^2$ with

$$C \cap \mathbb{C}_{x,y}^2 = z(y^2 - 4x^3 - g_2(\tau)x - g_3(\tau))$$

This is all summarized in the diagram

$$\begin{array}{ccc}
 \mathbb{C} \setminus \Gamma & \xrightarrow{\quad} & \mathbb{C}^2 \\
 \downarrow & \nearrow \psi & \downarrow \\
 \mathbb{C} & \xrightarrow{\quad} & \mathbb{CP}^2 \\
 \downarrow & \nearrow \varphi & \downarrow \\
 E & \xrightarrow{\quad} & \{[x, y, z] \mid yz^2 = 4x^2 - g_2xz^2 - g_3z^3\}
 \end{array}$$

Then we may define

Definition: An elliptic curve is a smooth plane cubic curve, along with the data of a point on the curve, e , which corresponds to the unit in the torus.

“Definition”: A complex lie group is a lie group where the word differentiable is replaced with the word complex/holomorphic.

Examples: a) Linear groups: $GL(n, \mathbb{C})$, $SL(n, \mathbb{C})$, $Sp(n, \mathbb{C})$. Note $U(n, \mathbb{C})$ is not a complex Lie group, because the definition requires complex conjugation, which is not a holomorphic operation.

These are non-compact and non-abelian if $n > 1$.

b) $X = \mathbb{C}^n \setminus \Gamma$ is an abelian, compact, complex Lie group.

c) Finite and discrete groups is a 0-dimensional complex lie group.

Proposition: If G is a compact, complex Lie group, then G is abelian.

Proof: For $g \in G$, consider

$$\Phi_g : G \rightarrow G, \quad h \mapsto g^{-1} \cdot h \cdot g$$

Then G is abelian iff $\forall g \in G, \Phi_g = Id_G$. Denote \mathfrak{g} as the Lie algebra of G . Then apply the maximum principle to the components of

$$\begin{aligned}
 Ad : G &\rightarrow \text{Aut}(\mathfrak{g}) \cong GL(n, \mathbb{C}) \\
 g &\mapsto (\Phi_g)_A|_e
 \end{aligned}$$

This is a holomorphic map, so it must be constant. $Ad_e = Id$, so $Ad = Id_G$. Because G is compact, the exponential map is surjective. It also commutes with Φ_g , so that

$$\Phi_g(\exp(\xi)) = \exp((\Phi_g)_*\xi) = \exp\xi$$

so that Φ is the identity.

□

If X is a topological space with a continuous group action by G , then $X \rightarrow X/G$ makes X/G into a topological space. Two important notions we need are

Definition: An action is free if every stabilizer group is trivial.

Definition: If G is a topological group, the action is proper if the map

$$\begin{aligned} G \times X &\rightarrow X \times X \\ (g, x) &\mapsto (x, g \cdot x) \end{aligned}$$

is proper.

If G is compact, of course all actions are proper, but if G is not, then the notion of being proper is interesting.

Theorem: *If G is a complex Lie group which acts freely and properly on a complex manifold, X , then X/G is naturally a complex manifold, and the quotient map is holomorphic. In fact, $\pi : X \rightarrow X/G$ is a principal G -bundle.*

Examples of the previous theorem: a) The action of \mathbb{C}^* on $\mathbb{C}^{n+1} \setminus \{0\}$ which results in $\mathbb{C}P^n$.

b) The action of a full rank lattice Γ on \mathbb{C}^n that results in an elliptic curve.

Holomorphic group actions II

Lecture 10, Sept 29. Theorem: Let $G \times X \rightarrow X$ be a proper and free topological group action. Then

- i) X/G is Hausdorff.
- ii) $q : X \rightarrow X/G$ is a principal G -bundle.

Proofs of this fact are more available in the smooth and holomorphic categories. We have $\forall x \in X, q^{-1}(x) = \{G \cdot x\} \cong G$. The idea is then to find a “slice” of the group action, Z , which is a submanifold of X which hits each fiber in only one point, and has complementary dimension to the fiber, i.e. $\dim q^{-1}(x) + \dim Z = \dim X$. So locally we are just looking for a transverse submanifold to $\{G \cdot x\}$, which is done by splitting your chart into \mathbb{R}^{n-k} and \mathbb{R}^k , but showing this fully is non-trivial.

Example:

- i) $\mathbb{C}P^n = (\mathbb{C}^{n+1} / \{0\}) / \mathbb{C}^*$
- ii) the complex tori $\mathbb{C}^n \rightarrow \mathbb{C}^n / \Gamma$

Non-Example: For $\alpha \in \mathbb{R} \setminus \mathbb{Q}$, let $\mathbb{Z} \curvearrowright S^1 = \mathbb{R}/\mathbb{Z}$ by $x \mapsto x + n\alpha$. In this action, all orbits are dense.

Proposition: If G is discrete, X locally compact, then the action $G \curvearrowright X$ is proper and free iff G acts properly discontinuously, i.e. $\forall x \in X, \exists U \ni x$ such that $\forall g \in G \setminus \{e\}, g(U) \cap U = \emptyset$

Example: Linear \mathbb{C}^* actions: $\mathbb{C}^* \curvearrowright \mathbb{C}^n$ by

$$\lambda(z_1, \dots, z_n) = (\lambda^{a_1} z_1, \dots, \lambda^{a_n} z_n)$$

for integers a_i . When is $\mathbb{C}^n / \mathbb{C}^*$ a complex manifold?

Clearly if we leave the origin in, the result is not Hausdorff, because for any U containing the origin, and for any $\lambda \in \mathbb{C}^*, 0 \in \lambda \cdot U \cap U$, so the action cannot be proper and free, and the result is non Hausdorff. If we examine the $n = 2$ case, the case $a_2 < 0$, then $\mathbb{C}^2 / \mathbb{C}^* \cong \mathbb{C}P^1 \sqcup \{*\}$, which is non-Hausdorff, and $(\mathbb{C}^2 \setminus \{0\}) / \mathbb{C}^* \cong \mathbb{C}P^1$. If we have $a_1 < 0, a_2 = 0$, then we collapse one of the axes, namely every horizontal line in \mathbb{C}^2 collapses down to a point in the vertical line at the origin, yielding $\mathbb{C}^2 / \mathbb{C}^* \cong \mathbb{C} \sqcup \mathbb{C}$, so you would have to remove the whole axis to get something Hausdorff, namely \mathbb{C} . If one is positive and one negative, then you get $\mathbb{C}P^1 \sqcup \{*\}$, as before. These are under

the umbrella of geometric invariant theory (GIT), where we look at a polynomial ring $\mathbb{C}[x_1, \dots, x_n]$, and let G act on it by acting on \mathbb{C}^n . Typically in GIT we take G to be a complex reductive group, which is a group which is the complexification of a compact, real Lie group, for example GL or SL . These act on $V = \mathbb{C}^n$, and the quotient V/G is typically not a nice space. To fix this, you restrict to the stable locus, $\forall \lambda : \mathbb{C}^* \rightarrow G$ a homomorphism, $V^S := \{v \in V \mid \lambda(v) \equiv O_v \text{ is closed} \Rightarrow G_v \text{ is finite}\}$, and then V^S/G is typically a nice space.

We recall the uniformization theorem in 1D:

Theorem (Uniformization): *If Σ is a simply connected Riemann surface, it is conformally equivalent to either the open unit disk, the complex plane, or \mathbb{CP}^1 .*

If we take Σ to be a Riemann surface with universal cover D , the open unit disk, then there is a deck transformation group, Γ which acts on D , such that

$$\Sigma = \mathbb{H}/\Gamma \cong D/\Gamma$$

But as we know,

$$\text{Aut}(D) = SU(1, n) \cong SL(2, \mathbb{R})/\{\pm 1\}$$

with Γ as a discrete subgroup, and the group itself acts as mobius transformations. Such an example is called a ball quotient.

If we have a closed Riemann surface, C , and $E = \mathbb{C}/\Gamma$ is an elliptic curve, and $G \in (E, +)$ is a finite group acting on \mathbb{C} , then $G \times (C \times E) \rightarrow C \times E$ by

$$(g, (x, y)) \mapsto (g \cdot x, y + g)$$

is a free and proper action.

Definition: A hyper-elliptic curve is a Riemann surface which arises as a 2-to-1 cover of \mathbb{CP}^1 .

Of course, for any double branched cover, there is a \mathbb{Z}_2 action swapping the branches, and every elliptic curve has a \mathbb{Z}_2 action sending $z \mapsto -z$.

If we define $X := (C \times E)/G$, then $\pi : X \rightarrow C/G$ is holomorphic, with $\pi^{-1}([z]) \cong E/G_z$, and π is an example of what we call an elliptic fibration.

Calculus and Holomorphic Vector Bundles

Lecture 11, Oct 1. Some more examples of quotient manifolds:

Example (Hopf manifolds) The elliptic curve arose as a quotient of \mathbb{C}^*/\mathbb{Z} . We can generalize this to several variables in the following way: Given $\lambda \in \mathbb{C}^*$, $|\lambda| > 1$, there is an action $\mathbb{Z} \times \mathbb{C}^n \setminus \{0\} \rightarrow \mathbb{C}^n \setminus \{0\}$ by

$$k \cdot (z_1, \dots, z_n) \mapsto (\lambda^k z_1, \dots, \lambda^k z_n)$$

This action is proper and free. It looks like this is just a product of n elliptic curves, but this is an action of \mathbb{Z} , not \mathbb{Z}^n . So we define $X := (\mathbb{C}^n \setminus \{0\})/\mathbb{Z}$. Rather than a product of elliptic curves, there is a diffeomorphism $X \cong S^1 \times S^{2n-1}$. This construction is useful because, as a real manifold, it is not clear why $S^1 \times S^{2n-1}$ should have a complex structure, yet this diffeomorphism makes it clear by our previous theorem.

Example (Iwasawa manifolds) The real Heisenberg group is given by matrices of the form

$$H := \left\{ \begin{pmatrix} 1 & z_1 & z_2 \\ 0 & 1 & z_3 \\ 0 & 0 & 1 \end{pmatrix} \mid z_i \in \mathbb{R} \right\} \subset \mathrm{GL}(3, \mathbb{R})$$

and the complex Heisenberg group is the same, but with complex coefficients, $\cong \mathbb{C}^3$. By restricting our coefficients to the Gaussian integers, $\mathbb{Z} + \mathbb{Z}i$, we get a discrete subgroup $\Gamma = H \cap \mathrm{GL}(3, \mathbb{Z} + \mathbb{Z}i)$, with an action on H by multiplication. Of course, any such action is free, but it is also proper, so that $X := H/\Gamma$ is a complex manifold, called an Iwasawa manifold. We may define a projection $\pi : X \rightarrow E \times E$, where $E = \mathbb{C}/(\mathbb{Z} + \mathbb{Z}i)$. In that case, π is a holomorphic fiber bundle with fibers E , but $X \not\cong E \times E \times E$.

Example (Grassmannians): We define

$$G(k, n) := \{W \subset \mathbb{C}^n \text{ linear subspace with } \dim W = k\}$$

Note $G(1, n) \cong G(n-1, n) \cong \mathbb{C}P^{n-1}$. So the first interesting¹ Grassmannian is $G(2, 4)$, and it is the only one with $n = 4$. In general, how do we get charts? $\mathrm{GL}(n, \mathbb{C})$ acts transitively on $G(k, n)$, so it suffices to check on $W := \mathbb{C}^k \times \{0\}$. An element of $G(k, n)$ near

¹As in, not projective space

W is viewed as the graph of a map from a complimentary (not necessarily orthogonal) subspace of dimension $n - k$, so that we have charts given by

$$\begin{aligned} M((n - k) \times k, \mathbb{C}) &\rightarrow G(k, n) \\ A &\mapsto \text{im} \left(\begin{pmatrix} E \\ A \end{pmatrix} : \mathbb{C}^k \rightarrow \mathbb{C}^n \right) \end{aligned}$$

it can be checked that the transition functions are holomorphic.

We have the Plucker embedding: If $L \in G(k, n)$ is spanned by linearly independent vectors $v_1, \dots, v_k \in \mathbb{C}^n$, we can define a map

$$\begin{aligned} G(k, n) &\rightarrow \bigwedge^k \mathbb{C}^n \\ L(v_1, \dots, v_k) &\mapsto v_1 \wedge \dots \wedge v_k \end{aligned}$$

It is clear that changing the basis v_1, \dots, v_k only have the effect of multiplying $v_1 \wedge \dots \wedge v_k$ by a constant, so we get a map on the projective space:

$$\begin{aligned} G(k, n) &\rightarrow \mathbb{P} \left(\bigwedge^k \mathbb{C}^n \right) \cong \mathbb{C}P^{\binom{n}{k}-1} \\ L(v_1, \dots, v_k) &\mapsto [v_1 \wedge \dots \wedge v_k] \end{aligned}$$

There are still details to check, but this Plucker map is a holomorphic embedding, thus by Chow theorem, this manifold must be algebraic, given by the so-called Plucker relations: The easiest example is

$$G(2, 4) \cong V(z_{12}z_{34} - z_{13}z_{24} + z_{14}z_{23}) \subset \mathbb{C}P^5$$

Finally, we discuss calculus on manifolds: If we take M to be a complex manifold, $T_{\mathbb{R}}M$ is the real tangent bundle, which is locally spanned by, taking local coordinates $z_i = x_i + iy_i$,

$$\{\partial_{x_1}, \dots, \partial_{x_n}, \partial_{y_1}, \dots, \partial_{y_n}\}$$

If we complexify this, we get the holomorphic vector bundle

$$T_{\mathbb{C}}M = T_{\mathbb{R}}M \otimes \mathbb{C} \cong TM \oplus \overline{TM}$$

with local coordinates

$$\begin{aligned} \partial_{z_\mu} &= \frac{1}{2}(\partial_{x_\mu} - i\partial_{y_\mu}) \quad \text{spans } TM \\ \partial_{\bar{z}_\mu} &= \frac{1}{2}(\partial_{x_\mu} + i\partial_{y_\mu}) \quad \text{spans } \overline{TM} \end{aligned}$$

It turns out that $T_{\mathbb{C}}M$ is a holomorphic vector bundle of rank n , the complex dimension of M , and the above is a local frame.

Dually, we have $T_{\mathbb{R}}^*M$, spanned by $dx_1, \dots, dx_n, dy_1, \dots, dy_n$, and we complexify

$$T_{\mathbb{C}}^*M = T^*M \otimes \mathbb{C} \cong T^*M \oplus \overline{T^*M} \equiv T^{1,0}M \oplus T^{0,1}M$$

Then the (p, q) -forms are defined as

$$T^{(p,q)}M = \bigwedge^p T^*M \otimes_{\mathbb{C}} \bigwedge^q \overline{T^*M}$$

This is different from real differential forms because there, all components are the same, i.e. we can define k -forms, whereas here p denotes the number of holomorphic components, and q denotes the number of anti holomorphic components. We can of course define integration of these forms by integrating the real and imaginary part separately.

'Definition': A holomorphic vector bundle is a vector bundle in the holomorphic category.

So π is a map of complex manifolds, and local trivializations are biholomorphic.

Example: The tautological line bundle is defined by taking a finite dimensional complex vector space, V , and noting $\mathbb{P}(V) = \{\ell \subset V \mid \dim \ell = 1\}$, we may define the tautological line bundle over $\mathbb{P}(V)$, $L \rightarrow \mathbb{P}(V)$, as

$$\begin{array}{ccc} L = \{(\ell, v) \in \mathbb{P}(V) \times V \mid v \in \ell\} & & (\ell, v) \\ \downarrow & & \downarrow \\ \mathbb{P}(V) & & \ell \end{array}$$

Line bundles, of which the tautological line bundle is of course one, have the property that tensoring them together yields another line bundle, leading us to define the Picard group as the set of line bundles modulo isomorphism under the operation of tensor product. This is an abelian group, although $L \otimes L'$ need not be equal to $L' \otimes L$, they may not even type check, but they are always isomorphic. So the Picard group of a manifold, $\text{Pic}(X)$, is the group of line bundles.

Proposition:

- i) $\text{Pic}(\mathbb{C}^n) = 0$
- ii) $\text{Pic}(\mathbb{C}P^n) \cong \mathbb{Z} \cong \langle L \rangle$
- iii) If E is an elliptic curve, $\text{Pic}(E) \cong E \times \mathbb{Z}$

We will leave these unproved for now.

Holomorphic vector bundles II

Lecture 12, Oct 6. Lots of administrative content discussed today so shorter math part.

If M is a complex manifold and $L_i \in \text{Pic}(M)$, then $E = \bigoplus_i L_i$ gives us an easy way to generate higher rank vector bundles.

Theorem (Birkhoff-Grothendieck): *If $E \rightarrow \mathbb{CP}^1$ is a holomorphic vector bundle of rank r , then $E \cong \bigoplus L^{\otimes a_i}$ for $a_i \in \mathbb{Z}$ and $L \rightarrow \mathbb{CP}^1$ the tautological line bundle.*

Theorem (Atiyah): *If $M = \mathbb{C}/\Gamma$ is an elliptic curve, there exist indecomposable bundles E_1, E_2, \dots with rank $E_r = r$, such that if $E \rightarrow M$ is an indecomposable bundle, then for φ an automorphism of M ,*

$$E \cong \varphi^* E_i$$

for some i .

So up to translation, there is only one indecomposable bundle.

We are learning that there aren't many holomorphic vector bundles. However, we do always have the tangent and cotangent bundles to a complex manifold. If Y is a complex submanifold of X , a complex manifold, there is an inclusion $TY \rightarrow TX$, and as in the real case, we can define the holomorphic normal bundle as the quotient $N := TX|_Y / TY$. Locally, if we say $Y = Z(z_{m+1}, \dots, z_n)$, with local coordinates z_1, \dots, z_n , then

$$TY = \langle \partial_{z_1}, \dots, \partial_{z_m} \rangle \subset TX|_Y = \langle \partial_{z_1}, \dots, \partial_{z_n} \rangle$$

If q is the quotient map $TX|_Y \rightarrow N$, then $q(\partial_{z_{m+1}}), \dots, q(\partial_{z_n})$ form a frame for N .

There is an exact sequence

$$0 \longrightarrow TY \longrightarrow TX|_Y \longrightarrow N \longrightarrow 0$$

Lemma: *If $0 \rightarrow E \rightarrow F \rightarrow G \rightarrow 0$ is an exact sequence of holomorphic vector bundles, then*

$$\det F = \det E \otimes \det G \quad \text{via a canonical isomorphism}$$

□

Then the exact sequence for the normal bundle implies¹

$$\kappa_Y = \kappa_{X|Y} \otimes \det N$$

where if X is a complex manifold, κ_X denotes the determinant line bundle of its cotangent bundle. The above equation is known as the adjunction formula²

¹is the order not wrong here?

²There are a lot of adjunction formulas. We even saw one in the other notes on symplectic topology

Sheaf Theory

Lecture 13, Oct. 8. Our guiding example will be the situation: Let $E \rightarrow M$ be a holomorphic vector bundle. Then for every U open in M , we may assign an abelian group, $\xi(U) := \{s : U \rightarrow E|_U \mid s \text{ is a holomorphic section}\}$. If there is an open set $V \subset U$, the restriction map induces a restriction on the abelian group. This leads us to define

Definition: An (abelian group valued) presheaf on a topological space X is a functor $\mathcal{F} : \text{Op}(X) \rightarrow \text{AbGrp}$, where $\text{Op}(X)$ is the category of open sets of X , with morphisms given by inclusions, such that $\mathcal{F}(\emptyset) = 0$ and the restriction $U \subset U$ induces the trivial restriction on abelian groups.

Definition: An (abelian group valued) sheaf on a topological space X is an (abelian group valued) pre-sheaf on X , such that $\forall U = \cup V_i, V_i$ open in X ,

- i) for $s \in \mathcal{F}(U), \forall i, s|_{V_i} = 0 \Rightarrow s = 0$.
- ii) for $s_i \in \mathcal{F}(V_i), \forall i, j, s_i|_{V_i \cap V_j} = s_j|_{V_i \cap V_j} \Rightarrow \exists s \in \mathcal{F}(U) \mid s|_{V_i} = s_i$.

The first condition says that checking that a global section is 0 is equivalent to checking it locally, and the second condition says that given local sections that agree on intersections, we can glue together a global section. These conditions together imply that the resulting glued section is unique.

Examples:

- i) Sheaves of functions/sections: $\mathcal{O}_M, C_X^0, C_M^\infty, \xi$.
- ii) Ideal sheaves: Take $A \subset M$ an analytic subset. Then $\mathcal{I}_A = \{f \in \mathcal{O}_M(U) \mid A \subset Z(f)\}$.
- iii) Locally constant sheaves: If A is an abelian group, define $A_X(U) = \{U \rightarrow A \text{ locally constant}\}$. The more naive approach to iii) is just assigning $A_X(U) = A$, and all the restriction maps to be the identity. This is a presheaf, but not a sheaf. The first axiom of checking zeros is satisfied: in fact, if a section vanishes on any subspace then it vanishes on the whole open set. The second axiom is not satisfied, however. If we have some U which is disconnected, and $A \neq \{1\}$, then simply assign one component to 1 and one component to some $a \neq 1$, this defines a section on each component, but they do not glue together.

Of course, since a sheaf is defined as a certain type of functor, a sheaf morphism should be defined as a natural transformation:

Definition: A sheaf homomorphism $\varphi : \mathcal{F} \rightarrow \mathcal{G}$ is an assignment $\forall U, \exists \varphi_U : \mathcal{F}(U) \rightarrow$

$\mathcal{G}(U)$ such that $\forall V \supset U$, there is a commutative diagram

$$\begin{array}{ccc} \mathcal{F}(V) & \longrightarrow & \mathcal{F}(U) \\ \downarrow & & \downarrow \\ \mathcal{G}(V) & \longrightarrow & \mathcal{G}(U) \end{array}$$

where the horizontal maps are the restriction maps. So we have a set of objects (abelian sheafs) and a set of morphisms (natural transformations), which define a category, $Ab(X)$, the abelian sheafs over X .

To every $p \in X$, with a sheaf \mathcal{F} over X , we can associate the abelian group $\mathcal{F}_p := \varinjlim_{U \ni p} \mathcal{F}(U)$, called the stalk at p .

Proposition: *The result of this colimit is an abelian group consisting of germs of sections near p .*

Proof: To see this, we show

$$\mathcal{F}_x := \{x_U \in \mathcal{F}(U) \mid x \in U\} / (x_U \sim x_V \iff \exists W \subset U \cap V \mid x_U|_W = x_V|_W)$$

along with the maps, for every U , $\psi_U : \mathcal{F}(U) \rightarrow \mathcal{F}_x$ sending a section x_U to its germ, $[x_U]$, satisfies the universal property of colimits. First we note that taking germs of sections makes this diagram commute with respect to the restriction morphisms. To do this, assume there is another object, Z , with maps $\varphi_U : \mathcal{F}(U) \rightarrow Z$, then we need to show there is a unique morphism (group homomorphism) $\mathcal{F}_x \rightarrow Z$. We may define

$$x_U \mapsto \varphi_U(x_U)$$

which type checks because $x_U \in \mathcal{F}(U)$ and $\varphi_U : \mathcal{F}(U) \rightarrow Z$. □

For any sheaf morphism $\varphi : \mathcal{F} \rightarrow \mathcal{G}$, we have a map¹ $\varphi_p : \mathcal{F}_p \rightarrow \mathcal{G}_p$ sending

$$[x_U] \mapsto \varphi_U(x_U)$$

Proposition: If $\phi : \mathcal{F} \rightarrow \mathcal{G}$ is a sheaf morphism, then ϕ is an isomorphism iff $\forall p$, ϕ_p is an isomorphism. □

If $s \in \mathcal{F}(U)$, we can define $\text{supp}(s) = \{p \in U \mid s_p \neq 0\}$, s_p being the image of s in the stalk \mathcal{F}_p . Since $\mathcal{F}(U)$ is an abelian group, we can speak of a 0 element, hence this definition of support. In some settings we define the support to be the closure of the non-zero points, but here the support is automatically closed. To see this, take a point in

¹This is part of the statement that “take the stalk at x ” is a functor from the category of sheaves over X to the category of abelian groups.

the complement, i.e. for $U \ni p$, take $s' \in \mathcal{F}(U)$ such that $\psi_U(s') \equiv s'_p \equiv [s'] = [0]$. By definition of a germ, this means s' vanishes on some neighborhood of p , so being in the complement of the support is an open condition, so the support is closed.

In the same spirit, we can define the support of a sheaf, $\text{supp}(\mathcal{F}) := \{p \in X \mid \mathcal{F}_p \neq 0\}$, so the set of all non-zero stalks. In general, this is neither closed nor open. We will later see a class of sheaves called coherent sheaves, whose support is closed, and in fact analytic.

Now we will talk about sheaves of sets. The only thing that changes in the definition is that the target category is of course, Set , but also we need to replace our use of 0, since we are no longer in an abelian group. To do so, we replace the condition $F(\emptyset) = 0$ with the condition $F(\emptyset) = \emptyset$, and we also need to replace condition i). The spirit of this condition was that it suffices to check that two sections are equivalent (s and 0) locally, so we replace this with a more general locality statement: If $\{U_i\}$ is a cover of $U \subset X$ and $s, s' \in \mathcal{F}(U)$, if $s|_{U_i} = s'|_{U_i}$, then $s = s'$. Everything else is the same.

If \mathcal{F} is a presheaf of sets, we can form a topological space, the etale space,

$$\acute{\text{Et}}(\mathcal{F}) := \coprod_{p \in X} \mathcal{F}_p$$

and we can make it into a bundle with the projection

$$\begin{aligned} \pi : \acute{\text{Et}}(\mathcal{F}) &\rightarrow X \\ \mathcal{F}_x &\mapsto x \end{aligned}$$

How do we topologize the etale space? It's the largest topology making each section continuous. Spelling this out, we define $U \subset \acute{\text{Et}}(\mathcal{F})$ to be open iff $\forall V \subset X$ open, $\forall s \in \mathcal{F}(V)$, the set $\{p \in X \mid s_p \in U\}$ is open (in V and X).