String Theory

University of Cambridge Part III Mathematical Tripos

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Recommended Books and Resources

• J. Polchinski, *String Theory*

This two volume work is the standard introduction to the subject. Our lectures will more or less follow the path laid down in volume one covering the bosonic string. The book contains explanations and descriptions of many details that have been deliberately (and, I suspect, at times inadvertently) swept under a very large rug in these lectures. Volume two covers the superstring.

• M. Green, J. Schwarz and E. Witten, *Superstring Theory*

Another two volume set. It is now over 20 years old and takes a slightly old-fashioned route through the subject, with no explicit mention of conformal field theory. However, it does contain much good material and the explanations are uniformly excellent. Volume one is most relevant for these lectures.

• B. Zwiebach, *A First Course in String Theory*

This book grew out of a course given to undergraduates who had no previous exposure to general relativity or quantum field theory. It has wonderful pedagogical discussions of the basics of lightcone quantization. More surprisingly, it also has some very clear descriptions of several advanced topics, even though it misses out all the bits in between.

• P. Di Francesco, P. Mathieu and D. Sénéchal, *Conformal Field Theory*

This big yellow book is affectionately known as the yellow pages. It’s a great way to learn conformal field theory. At first glance, it comes across as slightly daunting because it’s big. (And yellow). But you soon realise that it’s big because it starts at the beginning and provides detailed explanations at every step. The material necessary for this course can be found in chapters 5 and 6.

Further References: “String Theory and M-Theory” by Becker, Becker and Schwarz and “String Theory in a Nutshell” (it’s a big nutshell) by Kiritsis both deal with the bosonic string fairly quickly, but include more advanced topics that may be of interest. The book “D-Branes” by Johnson has lively and clear discussions about the many joys of D-branes. Links to several excellent online resources, including video lectures by Shiraz Minwalla, are listed on the course webpage.
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Acknowledgements

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0. Introduction

String theory is an ambitious project. It purports to be an all-encompassing theory of the universe, unifying the forces of nature, including gravity, in a single quantum mechanical framework.

The premise of string theory is that, at the fundamental level, matter does not consist of point-particles but rather of tiny loops of string. From this slightly absurd beginning, the laws of physics emerge. General relativity, electromagnetism and Yang-Mills gauge theories all appear in a surprising fashion. However, they come with baggage. String theory gives rise to a host of other ingredients, most strikingly extra spatial dimensions of the universe beyond the three that we have observed. The purpose of this course is to understand these statements in detail.

These lectures differ from most other courses that you will take in a physics degree. String theory is speculative science. There is no experimental evidence that string theory is the correct description of our world and scant hope that hard evidence will arise in the near future. Moreover, string theory is very much a work in progress and certain aspects of the theory are far from understood. Unresolved issues abound and it seems likely that the final formulation has yet to be written. For these reasons, I’ll begin this introduction by suggesting some answers to the question: Why study string theory?

Reason 1. String theory is a theory of quantum gravity

String theory unifies Einstein’s theory of general relativity with quantum mechanics. Moreover, it does so in a manner that retains the explicit connection with both quantum theory and the low-energy description of spacetime.

But quantum gravity contains many puzzles, both technical and conceptual. What does spacetime look like at the shortest distance scales? How can we understand physics if the causal structure fluctuates quantum mechanically? Is the big bang truly the beginning of time? Do singularities that arise in black holes really signify the end of time? What is the microscopic origin of black hole entropy and what is it telling us? What is the resolution to the information paradox? Some of these issues will be reviewed later in this introduction.

Whether or not string theory is the true description of reality, it offers a framework in which one can begin to explore these issues. For some questions, string theory has given very impressive and compelling answers. For others, string theory has been almost silent.
Reason 2. String theory may be the theory of quantum gravity

With broad brush, string theory looks like an extremely good candidate to describe the real world. At low-energies it naturally gives rise to general relativity, gauge theories, scalar fields and chiral fermions. In other words, it contains all the ingredients that make up our universe. It also gives the only presently credible explanation for the value of the cosmological constant although, in fairness, I should add that the explanation is so distasteful to some that the community is rather amusingly split between whether this is a good thing or a bad thing. Moreover, string theory incorporates several ideas which do not yet have experimental evidence but which are considered to be likely candidates for physics beyond the standard model. Prime examples are supersymmetry and axions.

However, while the broad brush picture looks good, the finer details have yet to be painted. String theory does not provide unique predictions for low-energy physics but instead offers a bewildering array of possibilities, mostly dependent on what is hidden in those extra dimensions. Partly, this problem is inherent to any theory of quantum gravity: as we’ll review shortly, it’s a long way down from the Planck scale to the domestic energy scales explored at the LHC. Using quantum gravity to extract predictions for particle physics is akin to using QCD to extract predictions for how coffee makers work. But the mere fact that it’s hard is little comfort if we’re looking for convincing evidence that string theory describes the world in which we live.

While string theory cannot at present offer falsifiable predictions, it has nonetheless inspired new and imaginative proposals for solving outstanding problems in particle physics and cosmology. There are scenarios in which string theory might reveal itself in forthcoming experiments. Perhaps we’ll find extra dimensions at the LHC, perhaps we’ll see a network of fundamental strings stretched across the sky, or perhaps we’ll detect some feature of non-Gaussianity in the CMB that is characteristic of D-branes at work during inflation. My personal feeling however is that each of these is a long shot and we may not know whether string theory is right or wrong within our lifetimes. Of course, the history of physics is littered with naysayers, wrongly suggesting that various theories will never be testable. With luck, I’ll be one of them.

Reason 3. String theory provides new perspectives on gauge theories

String theory was born from attempts to understand the strong force. Almost forty years later, this remains one of the prime motivations for the subject. String theory provides tools with which to analyze down-to-earth aspects of quantum field theory that are far removed from high-falutin’ ideas about gravity and black holes.
Of immediate relevance to this course are the pedagogical reasons to invest time in string theory. At heart, it is the study of conformal field theory and gauge symmetry. The techniques that we’ll learn are not isolated to string theory, but apply to countless systems which have direct application to real world physics.

On a deeper level, string theory provides new and very surprising methods to understand aspects of quantum gauge theories. Of these, the most startling is the AdS/CFT correspondence, first conjectured by Juan Maldacena, which gives a relationship between strongly coupled quantum field theories and gravity in higher dimensions. These ideas have been applied in areas ranging from nuclear physics to condensed matter physics and have provided qualitative (and arguably quantitative) insights into strongly coupled phenomena.

Reason 4. String theory provides new results in mathematics

For the past 250 years, the close relationship between mathematics and physics has been almost a one-way street: physicists borrowed many things from mathematicians but, with a few noticeable exceptions, gave little back. In recent times, that has changed. Ideas and techniques from string theory and quantum field theory have been employed to give new “proofs” and, perhaps more importantly, suggest new directions and insights in mathematics. The most well known of these is mirror symmetry, a relationship between topologically different Calabi-Yau manifolds.

The four reasons described above also crudely characterize the string theory community: there are “relativists” and “phenomenologists” and “field theorists” and “mathematicians”. Of course, the lines between these different sub-disciplines are not fixed and one of the great attractions of string theory is its ability to bring together people working in different areas — from cosmology to condensed matter to pure mathematics — and provide a framework in which they can profitably communicate. In my opinion, it is this cross-fertilization between fields which is the greatest strength of string theory.

0.1 Quantum Gravity

This is a starter course in string theory. Our focus will be on the perturbative approach to the bosonic string and, in particular, why this gives a consistent theory of quantum gravity. Before we leap into this, it is probably best to say a few words about quantum gravity itself. Like why it’s hard. And why it’s important. (And why it’s not).

The Einstein Hilbert action is given by

\[ S_{EH} = \frac{1}{16\pi G_N} \int d^4x \sqrt{-g} R \]
Newton’s constant $G_N$ can be written as

$$8\pi G_N = \frac{hc}{M_{pl}^2}$$

Throughout these lectures we work in units with $\hbar = c = 1$. The Planck mass $M_{pl}$ defines an energy scale

$$M_{pl} \approx 2 \times 10^{18} \text{ GeV}.$$  

(This is sometimes referred to as the reduced Planck mass, to distinguish it from the scale without the factor of $8\pi$, namely $\sqrt{1/G_N} \approx 1 \times 10^{19} \text{ GeV}$).

There are a couple of simple lessons that we can already take from this. The first is that the relevant coupling in the quantum theory is $1/M_{pl}$. To see that this is indeed the case from the perspective of the action, we consider small perturbations around flat Minkowski space,

$$g_{\mu\nu} = \eta_{\mu\nu} + \frac{1}{M_{pl}} h_{\mu\nu}$$

The factor of $1/M_{pl}$ is there to ensure that when we expand out the Einstein-Hilbert action, the kinetic term for $h$ is canonically normalized, meaning that it comes with no powers of $M_{pl}$. This then gives the kind of theory that you met in your first course on quantum field theory, albeit with an infinite series of interaction terms,

$$S_{EH} = \int d^4x \left( \partial h \right)^2 + \frac{1}{M_{pl}} h \left( \partial h \right)^2 + \frac{1}{M_{pl}^2} h^2 \left( \partial h \right)^2 + \ldots$$

Each of these terms is schematic: if you were to do this explicitly, you would find a mess of indices contracted in different ways. We see that the interactions are suppressed by powers of $M_{pl}$. This means that quantum perturbation theory is an expansion in the dimensionless ratio $E^2/M_{pl}^2$, where $E$ is the energy associated to the process of interest. We learn that gravity is weak, and therefore under control, at low-energies. But gravitational interactions become strong as the energy involved approaches the Planck scale. In the language of the renormalization group, couplings of this type are known as irrelevant.

The second lesson to take away is that the Planck scale $M_{pl}$ is very very large. The LHC will probe the electroweak scale, $M_{EW} \sim 10^3 \text{ GeV}$. The ratio is $M_{EW}/M_{pl} \sim 10^{-15}$. For this reason, quantum gravity will not affect your daily life, even if your daily life involves the study of the most extreme observable conditions in the universe.
Gravity is Non-Renormalizable

Quantum field theories with irrelevant couplings are typically ill-behaved at high-energies, rendering the theory ill-defined. Gravity is no exception. Theories of this type are called non-renormalizable, which means that the divergences that appear in the Feynman diagram expansion cannot be absorbed by a finite number of counterterms. In pure Einstein gravity, the symmetries of the theory are enough to ensure that the one-loop S-matrix is finite. The first divergence occurs at two-loops and requires the introduction of a counterterm of the form,

$$\Gamma \sim \frac{1}{\epsilon M_{pl}^4} \int d^4 x \sqrt{-g} R_{\mu \nu} R^{\rho \sigma} R_{\lambda \kappa} R_{\mu \nu}^\lambda$$

with $\epsilon = 4 - D$. All indications point towards the fact that this is the first in an infinite number of necessary counterterms.

Coupling gravity to matter requires an interaction term of the form,

$$S_{int} = \int d^4 x \frac{1}{M_{pl}} h_{\mu \nu} T^{\mu \nu} + \mathcal{O}(h^2)$$

This makes the situation marginally worse, with the first divergence now appearing at one-loop. The Feynman diagram in the figure shows particle scattering through the exchange of two gravitons. When the momentum $k$ running in the loop is large, the diagram is badly divergent: it scales as

$$\frac{1}{M_{pl}^4} \int_{0}^{\infty} d^4 k$$

Non-renormalizable theories are commonplace in the history of physics, the most commonly cited example being Fermi’s theory of the weak interaction. The first thing to say about them is that they are far from useless! Non-renormalizable theories are typically viewed as effective field theories, valid only up to some energy scale $\Lambda$. One deals with the divergences by simply admitting ignorance beyond this scale and treating $\Lambda$ as a UV cut-off on any momentum integral. In this way, we get results which are valid to an accuracy of $E/\Lambda$ (perhaps raised to some power). In the case of the weak interaction, Fermi’s theory accurately predicts physics up to an energy scale of $\sqrt{1/G_F} \sim 100$ GeV. In the case of quantum gravity, Einstein’s theory works to an accuracy of $(E/M_{pl})^2$.
However, non-renormalizable theories are typically unable to describe physics at their cut-off scale $\Lambda$ or beyond. This is because they are missing the true ultra-violet degrees of freedom which tame the high-energy behaviour. In the case of the weak force, these new degrees of freedom are the W and Z bosons. We would like to know what missing degrees of freedom are needed to complete gravity.

**Singularities**

Only a particle physicist would phrase all questions about the universe in terms of scattering amplitudes. In general relativity we typically think about the geometry as a whole, rather than bastardizing the Einstein-Hilbert action and discussing perturbations around flat space. In this language, the question of high-energy physics turns into one of short distance physics. Classical general relativity is not to be trusted in regions where the curvature of spacetime approaches the Planck scale and ultimately becomes singular. A quantum theory of gravity should resolve these singularities.

The question of spacetime singularities is morally equivalent to that of high-energy scattering. Both probe the ultra-violet nature of gravity. A spacetime geometry is made of a coherent collection of gravitons, just as the electric and magnetic fields in a laser are made from a collection of photons. The short distance structure of spacetime is governed – after Fourier transform – by high momentum gravitons. Understanding spacetime singularities and high-energy scattering are different sides of the same coin.

There are two situations in general relativity where singularity theorems tell us that the curvature of spacetime gets large: at the big bang and in the center of a black hole. These provide two of the biggest challenges to any putative theory of quantum gravity.

**Gravity is Subtle**

It is often said that general relativity contains the seeds of its own destruction. The theory is unable to predict physics at the Planck scale and freely admits to it. Problems such as non-renormalizability and singularities are, in a Rumsfeldian sense, known unknowns. However, the full story is more complicated and subtle. On the one hand, the issue of non-renormalizability may not quite be the crisis that it first appears. On the other hand, some aspects of quantum gravity suggest that general relativity isn’t as honest about its own failings as is usually advertised. The theory hosts a number of unknown unknowns, things that we didn’t even know that we didn’t know. We won’t have a whole lot to say about these issues in this course, but you should be aware of them. Here I mention only a few salient points.
Firstly, there is a key difference between Fermi’s theory of the weak interaction and gravity. Fermi’s theory was unable to provide predictions for any scattering process at energies above $\sqrt{1/G_F}$. In contrast, if we scatter two objects at extremely high-energies in gravity — say, at energies $E \gg M_{pl}$ — then we know exactly what will happen: we form a big black hole. We don’t need quantum gravity to tell us this. Classical general relativity is sufficient. If we restrict attention to scattering, the crisis of non-renormalizability is not problematic at ultra-high energies. It’s troublesome only within a window of energies around the Planck scale.

Similar caveats hold for singularities. If you are foolish enough to jump into a black hole, then you’re on your own: without a theory of quantum gravity, no one can tell you what fate lies in store at the singularity. Yet, if you are smart and stay outside of the black hole, you’ll be hard pushed to see any effects of quantum gravity. This is because Nature has conspired to hide Planck scale curvatures from our inquisitive eyes. In the case of black holes this is achieved through cosmic censorship which is a conjecture in classical general relativity that says singularities are hidden behind horizons. In the case of the big bang, it is achieved through inflation, washing away any traces from the very early universe. Nature appears to shield us from the effects of quantum gravity, whether in high-energy scattering or in singularities. I think it’s fair to say that no one knows if this conspiracy is pointing at something deep, or is merely inconvenient for scientists trying to probe the Planck scale.

While horizons may protect us from the worst excesses of singularities, they come with problems of their own. These are the unknown unknowns: difficulties that arise when curvatures are small and general relativity says “trust me”. The entropy of black holes and the associated paradox of information loss strongly suggest that local quantum field theory breaks down at macroscopic distance scales. Attempts to formulate quantum gravity in de Sitter space, or in the presence of eternal inflation, hint at similar difficulties. Ideas of holography, black hole complimentarity and the AdS/CFT correspondence all point towards non-local effects and the emergence of spacetime. These are the deep puzzles of quantum gravity and their relationship to the ultra-violet properties of gravity is unclear.

As a final thought, let me mention the one observation that has an outside chance of being related to quantum gravity: the cosmological constant. With an energy scale of $\Lambda \sim 10^{-3}$ eV it appears to have little to do with ultra-violet physics. If it does have its origins in a theory of quantum gravity, it must either be due to some subtle “unknown unknown”, or because it is explained away as an environmental quantity as in string theory.
Is the Time Ripe?

Our current understanding of physics, embodied in the standard model, is valid up to energy scales of $10^3$ GeV. This is 15 orders of magnitude away from the Planck scale. Why do we think the time is now ripe to tackle quantum gravity? Surely we are like the ancient Greeks arguing about atomism. Why on earth do we believe that we’ve developed the right tools to even address the question?

The honest answer, I think, is hubris.

However, there is mild circumstantial evidence that the framework of quantum field theory might hold all the way to the Planck scale without anything very dramatic happening in between. The main argument is unification. The three coupling constants of Nature run logarithmically, meeting miraculously at the GUT energy scale of $10^{15}$ GeV. Just slightly later, the fourth force of Nature, gravity, joins them. While not overwhelming, this does provide a hint that perhaps quantum field theory can be taken seriously at these ridiculous scales. Historically I suspect this was what convinced large parts of the community that it was ok to speak about processes at $10^{18}$ GeV.

Finally, perhaps the most compelling argument for studying physics at the Planck scale is that string theory does provide a consistent unified quantum theory of gravity and the other forces. Given that we have this theory sitting in our laps, it would be foolish not to explore its consequences. The purpose of these lecture notes is to begin this journey.
1. The Relativistic String

All lecture courses on string theory start with a discussion of the point particle. Ours is no exception. We’ll take a flying tour through the physics of the relativistic point particle and extract a couple of important lessons that we’ll take with us as we move onto string theory.

1.1 The Relativistic Point Particle

We want to write down the Lagrangian describing a relativistic particle of mass $m$. In anticipation of string theory, we’ll consider $D$-dimensional Minkowski space $\mathbb{R}^{1,D-1}$. Throughout these notes, we work with signature

$$\eta_{\mu\nu} = \text{diag}(-1,+1,+1,\ldots,+1)$$

Note that this is the opposite signature to my quantum field theory notes.

If we fix a frame with coordinates $X^\mu = (t, \vec{x})$ the action is simple:

$$S = -m \int dt \sqrt{1 - \dot{\vec{x}} \cdot \dot{\vec{x}}}. \tag{1.1}$$

To see that this is correct we can compute the momentum $\vec{p}$, conjugate to $\vec{x}$, and the energy $E$ which is equal to the Hamiltonian,

$$\vec{p} = \frac{m\dot{\vec{x}}}{\sqrt{1 - \dot{\vec{x}} \cdot \dot{\vec{x}}}}, \quad E = \sqrt{m^2 + \vec{p}^2},$$

both of which should be familiar from courses on special relativity.

Although the Lagrangian (1.1) is correct, it’s not fully satisfactory. The reason is that time $t$ and space $\vec{x}$ play very different roles in this Lagrangian. The position $\vec{x}$ is a dynamical degree of freedom. In contrast, time $t$ is merely a parameter providing a label for the position. Yet Lorentz transformations are supposed to mix up $t$ and $\vec{x}$ and such symmetries are not completely obvious in (1.1). Can we find a new Lagrangian in which time and space are on equal footing?

One possibility is to treat both time and space as labels. This leads us to the concept of field theory. However, in this course we will be more interested in the other possibility: we will promote time to a dynamical degree of freedom. At first glance, this may appear odd: the number of degrees of freedom is one of the crudest ways we have to characterize a system. We shouldn’t be able to add more degrees of freedom
at will without fundamentally changing the system that we’re talking about. Another way of saying this is that the particle has the option to move in space, but it doesn’t have the option to move in time. It has to move in time. So we somehow need a way to promote time to a degree of freedom without it really being a true dynamical degree of freedom! How do we do this? The answer, as we will now show, is gauge symmetry.

Consider the action,

$$ S = -m \int d\tau \sqrt{\dot{X}^\mu \dot{X}^\nu \eta_{\mu\nu}} , \quad (1.2) $$

where $\mu = 0, \ldots, D - 1$ and $\dot{X}^\mu = dX^\mu/d\tau$. We’ve introduced a new parameter $\tau$ which labels the position along the worldline of the particle as shown by the dashed lines in the figure. This action has a simple interpretation: it is just the proper time $\int ds$ along the worldline.

Naively it looks as if we now have $D$ physical degrees of freedom rather than $D - 1$ because, as promised, the time direction $X^0 \equiv t$ is among our dynamical variables: $X^0 = X^0(\tau)$. However, this is an illusion. To see why, we need to note that the action (1.2) has a very important property: reparameterization invariance. This means that we can pick a different parameter $\tilde{\tau}$ on the worldline, related to $\tau$ by any monotonic function

$$ \tilde{\tau} = \tilde{\tau}(\tau) . $$

Let’s check that the action is invariant under transformations of this type. The integration measure in the action changes as $d\tau = d\tilde{\tau} |d\tau/d\tilde{\tau}|$. Meanwhile, the velocities change as $dX^\mu/d\tau = (dX^\mu/d\tilde{\tau})(d\tilde{\tau}/d\tau)$. Putting this together, we see that the action can just as well be written in the $\tilde{\tau}$ reparameterization,

$$ S = -m \int d\tilde{\tau} \sqrt{-\frac{dX^\mu}{d\tilde{\tau}} \frac{dX^\nu}{d\tilde{\tau}} \eta_{\mu\nu}} . $$

The upshot of this is that not all $D$ degrees of freedom $X^\mu$ are physical. For example, suppose you find a solution to this system, so that you know how $X^0$ changes with $\tau$ and how $X^1$ changes with $\tau$ and so on. Not all of that information is meaningful because $\tau$ itself is not meaningful. In particular, we could use our reparameterization invariance to simply set

$$ \tau = X^0(\tau) \equiv t \quad (1.3) $$
If we plug this choice into the action (1.2) then we recover our initial action (1.1). The reparameterization invariance is a *gauge symmetry* of the system. Like all gauge symmetries, it’s not really a symmetry at all. Rather, it is a redundancy in our description. In the present case, it means that although we seem to have $D$ degrees of freedom $X^\mu$, one of them is fake.

The fact that one of the degrees of freedom is a fake also shows up if we look at the momenta,

$$ p_\mu = \frac{\partial L}{\partial \dot{X}^\mu} = \frac{m\dot{X}^\nu \eta_{\mu\nu}}{\sqrt{-\dot{X}^\lambda \dot{X}^\rho \eta_{\lambda\rho}}} $$

These momenta aren’t all independent. They satisfy

$$ p_\mu p^\mu + m^2 = 0 $$

This is a constraint on the system. It is, of course, the mass-shell constraint for a relativistic particle of mass $m$. From the worldline perspective, it tells us that the particle isn’t allowed to sit still in Minkowski space: at the very least, it had better keep moving in a timelike direction with $(p^0)^2 \geq m^2$.

One advantage of the action (1.2) is that the Poincaré symmetry of the particle is now manifest, appearing as a global symmetry on the worldline

$$ X^\mu \rightarrow \Lambda^\mu_\nu X^\nu + e^\mu $$

where $\Lambda$ is a Lorentz transformation satisfying $\Lambda^\nu_\rho \eta^{\rho\sigma} \Lambda^\sigma_\rho = \eta^{\mu\sigma}$, while $c^\mu$ corresponds to a constant translation. We have made all the symmetries manifest at the price of introducing a gauge symmetry into our system. A similar gauge symmetry will arise in the relativistic string and much of this course will be devoted to understanding its consequences.

### 1.1.1 Quantization

It’s a trivial matter to quantize this action. We introduce a wavefunction $\Psi(X)$. This satisfies the usual Schrödinger equation,

$$ i \frac{\partial \Psi}{\partial \tau} = H\Psi. $$

But, computing the Hamiltonian $H = \dot{X}^\mu p_\mu - L$, we find that it vanishes: $H = 0$. This shouldn’t be surprising. It is simply telling us that the wavefunction doesn’t depend on
τ. Since the wavefunction is something physical while, as we have seen, τ is not, this is to be expected. Note that this doesn’t mean that time has dropped out of the problem. On the contrary, in this relativistic context, time $X^0$ is an operator, just like the spatial coordinates $\vec{x}$. This means that the wavefunction $\Psi$ is immediately a function of space and time. It is not like a static state in quantum mechanics, but more akin to the fully integrated solution to the non-relativistic Schrödinger equation.

The classical system has a constraint given by (1.5). In the quantum theory, we impose this constraint as an operator equation on the wavefunction, namely $(p^\mu p_\mu + m^2)\Psi = 0$. Using the usual representation of the momentum operator $p_\mu = -i \partial / \partial X^\mu$, we recognize this constraint as the Klein-Gordon equation

$$\left(-\frac{\partial}{\partial X^\mu} \frac{\partial}{\partial X^\nu} \eta^{\mu\nu} + m^2\right) \Psi(X) = 0$$

(1.7)

Although this equation is familiar from field theory, it’s important to realize that the interpretation is somewhat different. In relativistic field theory, the Klein-Gordon equation is the equation of motion obeyed by a scalar field. In relativistic quantum mechanics, it is the equation obeyed by the wavefunction. In the early days of field theory, the fact that these two equations are the same led people to think one should view the wavefunction as a classical field and quantize it a second time. This isn’t correct, but nonetheless the language has stuck and it is common to talk about the point particle perspective as “first quantization” and the field theory perspective as “second quantization”.

So far we’ve considered only a free point particle. How can we introduce interactions into this framework? We would have to first decide which interactions are allowed: perhaps the particle can split into two; perhaps it can fuse with other particles? Obviously, there is a huge range of options for us to choose from. We would then assign amplitudes for these processes to happen. There would be certain restrictions coming from the requirement of unitarity which, among other things, would lead to the necessity of anti-particles. We could draw diagrams associated to the different interactions — an example is given in the figure — and in this manner we would slowly build up the Feynman diagram expansion that is familiar from field theory. In fact, this was pretty much the way Feynman himself approached the topic of QED. However, in practice we rarely construct particle interactions in this way because the field theory framework provides a much better way of looking at things. In contrast, this way of building up interactions is exactly what we will later do for strings.

---

Figure 4:
1.1.2 Ein Einbein

There is another action that describes the relativistic point particle. We introduce yet another field on the worldline, \( e(\tau) \), and write

\[
S = \frac{1}{2} \int d\tau \left( e^{-1} \dot{X}^2 - em^2 \right),
\]

(1.8)

where we’ve used the notation \( \dot{X}^2 = \dot{X}^\mu \dot{X}^\nu \eta_{\mu\nu} \). For the rest of these lectures, terms like \( X^2 \) will always mean an implicit contraction with the spacetime Minkowski metric.

This form of the action makes it look as if we have coupled the worldline theory to 1d gravity, with the field \( e(\tau) \) acting as an einbein (in the sense of vierbeins that are introduced in general relativity). To see this, note that we could change notation and write this action in the more suggestive form

\[
S = -\frac{1}{2} \int d\tau \sqrt{-g_{\tau\tau}} \left( g_{\tau\tau} \dot{X}^2 + m^2 \right).
\]

(1.9)

where \( g_{\tau\tau} = (g^{\tau\tau})^{-1} \) is the metric on the worldline and \( e = \sqrt{-g_{\tau\tau}} \).

Although our action appears to have one more degree of freedom, \( e \), it can be easily checked that it has the same equations of motion as (1.2). The reason for this is that \( e \) is completely fixed by its equation of motion, \( \ddot{X}^2 + e^2 m^2 = 0 \). Substituting this into the action (1.8) recovers (1.2)

The action (1.8) has a couple of advantages over (1.2). Firstly, it works for massless particles with \( m = 0 \). Secondly, the absence of the annoying square root means that it’s easier to quantize in a path integral framework.

The action (1.8) retains invariance under reparameterizations which are now written in a form that looks more like general relativity. For transformations parameterized by an infinitesimal \( \eta \), we have

\[
\tau \rightarrow \tilde{\tau} = \tau - \eta(\tau), \quad \delta e = \frac{d}{d\tau}(\eta(\tau)e), \quad \delta X^\mu = \frac{dX^\mu}{d\tau} \eta(\tau)
\]

(1.10)

The einbein \( e \) transforms as a density on the worldline, while each of the coordinates \( X^\mu \) transforms as a worldline scalar.
1.2 The Nambu-Goto Action

A particle sweeps out a worldline in Minkowski space. A string sweeps out a worldsheet. We’ll parameterize this worldsheet by one timelike coordinate $\tau$, and one spacelike coordinate $\sigma$. In this section we’ll focus on closed strings and take $\sigma$ to be periodic, with range

$$\sigma \in [0, 2\pi) .$$  

We will sometimes package the two worldsheet coordinates together as $\sigma^\alpha = (\tau, \sigma)$, $\alpha = 0, 1$. Then the string sweeps out a surface in spacetime which defines a map from the worldsheet to Minkowski space, $X^\mu(\sigma, \tau)$ with $\mu = 0, \ldots, D-1$. For closed strings, we require

$$X^\mu(\sigma, \tau) = X^\mu(\sigma + 2\pi, \tau) .$$

In this context, spacetime is sometimes referred to as the target space to distinguish it from the worldsheet.

We need an action that describes the dynamics of this string. The key property that we will ask for is that nothing depends on the coordinates $\sigma^\alpha$ that we choose on the worldsheet. In other words, the string action should be reparameterization invariant. What kind of action does the trick? Well, for the point particle the action was proportional to the length of the worldline. The obvious generalization is that the action for the string should be proportional to the area, $A$, of the worldsheet. This is certainly a property that is characteristic of the worldsheet itself, rather than any choice of parameterization.

How do we find the area $A$ in terms of the coordinates $X^\mu(\sigma, \tau)$? The worldsheet is a curved surface embedded in spacetime. The induced metric, $\gamma_{\alpha\beta}$, on this surface is the pull-back of the flat metric on Minkowski space,

$$\gamma_{\alpha\beta} = \frac{\partial X^\mu}{\partial \sigma^\alpha} \frac{\partial X^\nu}{\partial \sigma^\beta} \eta_{\mu\nu} .$$  

Then the action which is proportional to the area of the worldsheet is given by,

$$S = -T \int d^2\sigma \sqrt{-\det \gamma} .$$  

Here $T$ is a constant of proportionality. We will see shortly that it is the tension of the string, meaning the mass per unit length.
We can write this action a little more explicitly. The pull-back of the metric is given by,

$$\gamma_{\alpha\beta} = \begin{pmatrix} \dot{X}^2 & \dot{X} \cdot X' \\ \dot{X} \cdot X' & (X')^2 \end{pmatrix}.$$

where $\dot{X}^\mu = \partial X^\mu / \partial \tau$ and $X'^\mu = \partial X^\mu / \partial \sigma$. The action then takes the form,

$$S = -T \int d^2\sigma \sqrt{-\dot{X}^2 (X')^2 + (\dot{X} \cdot X')^2}.$$  \hspace{1cm} (1.14)

This is the Nambu-Goto action for a relativistic string.

**Action = Area: A Check**

If you’re unfamiliar with differential geometry, the argument about the pull-back of the metric may be a bit slick. Thankfully, there’s a more pedestrian way to see that the action (1.14) is equal to the area swept out by the worldsheet. It’s slightly simpler to make this argument for a surface embedded in Euclidean space rather than Minkowski space. We choose some parameterization of the sheet in terms of $\tau$ and $\sigma$, as drawn in the figure, and we write the coordinates of Euclidean space as $\vec{X} (\sigma, \tau)$. We’ll compute the area of the infinitesimal shaded region. The vectors tangent to the boundary are,

$$\vec{dl}_1 = \frac{\partial \vec{X}}{\partial \sigma}, \quad \vec{dl}_2 = \frac{\partial \vec{X}}{\partial \tau}.$$

If the angle between these two vectors is $\theta$, then the area is then given by

$$ds^2 = |\vec{dl}_1| |\vec{dl}_2| \sin \theta = \sqrt{dl_1^2 dl_2^2 (1 - \cos^2 \theta)} = \sqrt{dl_1^2 dl_2^2 - (\vec{dl}_1 \cdot \vec{dl}_2)^2}$$  \hspace{1cm} (1.15)

which indeed takes the form of the integrand of (1.14).

**Tension and Dimension**

Let’s now see that $T$ has the physical interpretation of tension. We write Minkowski coordinates as $X^\mu = (t, \vec{x})$. We work in a gauge with $X^0 \equiv t = R \tau$, where $R$ is a constant that is needed to balance up dimensions (see below) and will drop out at the end of the argument. Consider a snapshot of a string configuration at a time when
\( d\vec{x}/d\tau = 0 \) so that the instantaneous kinetic energy vanishes. Evaluating the action for a time \( dt \) gives

\[
S = -T \int d\tau d\sigma R \sqrt{(d\vec{x}/d\sigma)^2} = -T \int dt \text{ (spatial length of string)} . \tag{1.16}
\]

But, when the kinetic energy vanishes, the action is proportional to the time integral of the potential energy,

\[
\text{potential energy} = T \times \text{ (spatial length of string)} .
\]

So \( T \) is indeed the energy per unit length as claimed. We learn that the string acts rather like an elastic band and its energy increases linearly with length. (This is different from the elastic bands you’re used to which obey Hooke’s law where energy increased quadratically with length). To minimize its potential energy, the string will want to shrink to zero size. We’ll see that when we include quantum effects this can’t happen because of the usual zero point energies.

There is a slightly annoying way of writing the tension that has its origin in ancient history, but is commonly used today

\[
T = \frac{1}{2\pi \alpha'} \tag{1.17}
\]

where \( \alpha' \) is pronounced “alpha-prime”. In the language of our ancestors, \( \alpha' \) is referred to as the “universal Regge slope”. We’ll explain why later in this course.

At this point, it’s worth pointing out some conventions that we have, until now, left implicit. The spacetime coordinates have dimension \([X] = -1\). In contrast, the worldsheet coordinates are taken to be dimensionless, \([\sigma] = 0\). (This can be seen in our identification \( \sigma \equiv \sigma + 2\pi \)). The tension is equal to the mass per unit length and has dimension \([T] = 2\). Obviously this means that \([\alpha'] = -2\). We can therefore associate a length scale, \( l_s \), by

\[
\alpha' = l_s^2 \tag{1.18}
\]

The string scale \( l_s \) is the natural length that appears in string theory. In fact, in a certain sense (that we will make more precise later in the course) this length scale is the only parameter of the theory.
Actual Strings vs. Fundamental Strings

There are several situations in Nature where string-like objects arise. Prime examples include magnetic flux tubes in superconductors and chromo-electric flux tubes in QCD. Cosmic strings, a popular speculation in cosmology, are similar objects, stretched across the sky. In each of these situations, there are typically two length scales associated to the string: the tension, $T$ and the width of the string, $L$. For all these objects, the dynamics is governed by the Nambu-Goto action as long as the curvature of the string is much greater than $L$. (In the case of superconductors, one should work with a suitable non-relativistic version of the Nambu-Goto action).

However, in each of these other cases, the Nambu-Goto action is not the end of the story. There will typically be additional terms in the action that depend on the width of the string. The form of these terms is not universal, but often includes a rigidity piece of form $L \int K^2$, where $K$ is the extrinsic curvature of the worldsheet. Other terms could be added to describe fluctuations in the width of the string.

The string scale, $l_s$, or equivalently the tension, $T$, depends on the kind of string that we’re considering. For example, if we’re interested in QCD flux tubes then we would take

$$T \sim (1 \text{ GeV})^2$$  \hspace{1cm} (1.19)

In this course we will consider fundamental strings which have zero width. What this means in practice is that we take the Nambu-Goto action as the complete description for all configurations of the string. These strings will have relevance to quantum gravity and the tension of the string is taken to be much larger, typically an order of magnitude or so below the Planck scale.

$$T \lesssim M_{pl}^2 = (10^{18} \text{ GeV})^2$$  \hspace{1cm} (1.20)

However, I should point out that when we try to view string theory as a fundamental theory of quantum gravity, we don’t really know what value $T$ should take. As we will see later in this course, it depends on many other aspects, most notably the string coupling and the volume of the extra dimensions.

1.2.1 Symmetries of the Nambu-Goto Action

The Nambu-Goto action has two types of symmetry, each of a different nature.

- Poincaré invariance of the spacetime (1.6). This is a global symmetry from the perspective of the worldsheet, meaning that the parameters $\Lambda^{\mu}_{\nu}$ and $c^\mu$ which label
the symmetry transformation are constants and do not depend on worldsheet coordinates $\sigma^\alpha$.

- Reparameterization invariance, $\sigma^\alpha \to \tilde{\sigma}^\alpha(\sigma)$. As for the point particle, this is a gauge symmetry. It reflects the fact that we have a redundancy in our description because the worldsheet coordinates $\sigma^\alpha$ have no physical meaning.

1.2.2 Equations of Motion

To derive the equations of motion for the Nambu-Goto string, we first introduce the momenta which we call $\Pi$ because there will be countless other quantities that we want to call $p$ later,

$$
\Pi^\tau_\mu = \frac{\partial L}{\partial \dot{X}^\mu} = -T \frac{(\dot{X} \cdot X')X'_\mu - (X'^2)\dot{X}_\mu}{\sqrt{(\dot{X} \cdot X')^2 - \dot{X}^2 X'^2}}.
$$

$$
\Pi^\sigma_\mu = \frac{\partial L}{\partial X'^{\mu}} = -T \frac{(\dot{X} \cdot X')\dot{X}_\mu - (\dot{X}^2)X'_\mu}{\sqrt{(\dot{X} \cdot X')^2 - \dot{X}^2 X'^2}}.
$$

The equations of motion are then given by,

$$
\frac{\partial \Pi^\tau_\mu}{\partial \tau} + \frac{\partial \Pi^\sigma_\mu}{\partial \sigma} = 0
$$

These look like nasty, non-linear equations. In fact, there’s a slightly nicer way to write these equations, starting from the earlier action (1.13). Recall that the variation of a determinant is $\delta \sqrt{-\gamma} = \frac{1}{2} \sqrt{-\gamma} \gamma^{\alpha\beta} \delta \gamma_{\alpha\beta}$. Using the definition of the pull-back metric $\gamma_{\alpha\beta}$, this gives rise to the equations of motion

$$
\partial_\alpha (\sqrt{-\det \gamma} \gamma^{\alpha\beta} \partial_\beta X^\mu) = 0,
$$

Although this notation makes the equations look a little nicer, we’re kidding ourselves. Written in terms of $X^\mu$, they are still the same equations. Still nasty.

1.3 The Polyakov Action

The square-root in the Nambu-Goto action means that it’s rather difficult to quantize using path integral techniques. However, there is another form of the string action which is classically equivalent to the Nambu-Goto action. It eliminates the square root at the expense of introducing another field,

$$
S = -\frac{1}{4\pi\alpha'} \int d^2\sigma \sqrt{-g} g^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu \eta_{\mu\nu}
$$

where $g \equiv \det g$. This is the Polyakov action. (Polyakov didn’t discover the action, but he understood how to work with it in the path integral and for this reason it carries his name. The path integral treatment of this action will be the subject of Chapter 5).
The new field is $g_{\alpha\beta}$. It is a dynamical metric on the worldsheet. From the perspective of the worldsheet, the Polyakov action is a bunch of scalar fields $X$ coupled to 2d gravity.

The equation of motion for $X^\mu$ is

$$\partial_\alpha(\sqrt{-g}g^{\alpha\beta}\partial_\beta X^\mu) = 0 , \quad (1.23)$$

which coincides with the equation of motion (1.21) from the Nambu-Goto action, except that $g_{\alpha\beta}$ is now an independent variable which is fixed by its own equation of motion. To determine this, we vary the action (remembering again that $\delta\sqrt{-g} = -\frac{1}{2}\sqrt{-g}g_{\alpha\beta}\delta g^{\alpha\beta} = +\frac{1}{2}\sqrt{-g}g^{\alpha\beta}\delta g_{\alpha\beta}$),

$$\delta S = -\frac{T}{2} \int d^2\sigma \, \delta g^{\alpha\beta} \left(\sqrt{-g} \partial_\alpha X^\mu \partial_\beta X^\nu - \frac{1}{2}\sqrt{-g}g_{\alpha\beta}g^{\rho\sigma}\partial_\rho X^\mu \partial_\sigma X^\nu\right) \eta_{\mu\nu} = 0 . \quad (1.24)$$

The worldsheet metric is therefore given by,

$$g_{\alpha\beta} = 2f(\sigma) \partial_\alpha X \cdot \partial_\beta X , \quad (1.25)$$

where the function $f(\sigma)$ is given by,

$$f^{-1} = g^{\rho\sigma} \partial_\rho X \cdot \partial_\sigma X$$

A comment on the potentially ambiguous notation: here, and below, any function $f(\sigma)$ is always short-hand for $f(\sigma, \tau)$: it in no way implies that $f$ depends only on the spatial worldsheet coordinate.

We see that $g_{\alpha\beta}$ isn’t quite the same as the pull-back metric $\gamma_{\alpha\beta}$ defined in equation (1.12); the two differ by the conformal factor $f$. However, this doesn’t matter because, rather remarkably, $f$ drops out of the equation of motion (1.23). This is because the $\sqrt{-g}$ term scales as $f$, while the inverse metric $g^{\alpha\beta}$ scales as $f^{-1}$ and the two pieces cancel. We therefore see that Nambu-Goto and the Polyakov actions result in the same equation of motion for $X$.

In fact, we can see more directly that the Nambu-Goto and Polyakov actions coincide. We may replace $g_{\alpha\beta}$ in the Polyakov action (1.22) with its equation of motion $g_{\alpha\beta} = 2f \gamma_{\alpha\beta}$. The factor of $f$ also drops out of the action for the same reason that it dropped out of the equation of motion. In this manner, we recover the Nambu-Goto action (1.13).
### 1.3.1 Symmetries of the Polyakov Action

The fact that the presence of the factor $f(\sigma, \tau)$ in (1.25) didn’t actually affect the equations of motion for $X^\mu$ reflects the existence of an extra symmetry which the Polyakov action enjoys. Let’s look more closely at this. Firstly, the Polyakov action still has the two symmetries of the Nambu-Goto action,

- Poincaré invariance. This is a global symmetry on the worldsheet.

$$X^\mu \to \Lambda^\mu_\nu X^\nu + c^\mu .$$

- Reparameterization invariance, also known as diffeomorphisms. This is a gauge symmetry on the worldsheet. We may redefine the worldsheet coordinates as $\sigma^\alpha \to \tilde{\sigma}^\alpha(\sigma)$. The fields $X^\mu$ transform as worldsheet scalars, while $g_{\alpha\beta}$ transforms in the manner appropriate for a 2d metric.

$$X^\mu(\sigma) \to \tilde{X}^\mu(\tilde{\sigma}) = X^\mu(\sigma)$$

$$g_{\alpha\beta}(\sigma) \to \tilde{g}_{\alpha\beta}(\tilde{\sigma}) = \frac{\partial \sigma^\gamma}{\partial \tilde{\sigma}^\alpha} \frac{\partial \sigma^\delta}{\partial \tilde{\sigma}^\beta} g_{\gamma\delta}(\sigma)$$

It will sometimes be useful to work infinitesimally. If we make the coordinate change $\sigma^\alpha \to \tilde{\sigma}^\alpha = \sigma^\alpha - \eta^\alpha(\sigma)$, for some small $\eta$. The transformations of the fields then become,

$$\delta X^\mu(\sigma) = \eta^\alpha \partial_\alpha X^\mu$$

$$\delta g_{\alpha\beta}(\sigma) = \nabla_\alpha \eta_\beta + \nabla_\beta \eta_\alpha$$

where the covariant derivative is defined by $\nabla_\alpha \eta_\beta = \partial_\alpha \eta_\beta - \Gamma^\sigma_{\alpha\beta} \eta_\sigma$ with the Levi-Civita connection associated to the worldsheet metric given by the usual expression,

$$\Gamma^\sigma_{\alpha\beta} = \frac{1}{2} g^{\sigma\rho}(\partial_\alpha g_{\beta\rho} + \partial_\beta g_{\rho\alpha} - \partial_\rho g_{\alpha\beta})$$

Together with these familiar symmetries, there is also a new symmetry which is novel to the Polyakov action. It is called Weyl invariance.

- Weyl Invariance. Under this symmetry, $X^\mu(\sigma) \to X^\mu(\sigma)$, while the metric changes as

$$g_{\alpha\beta}(\sigma) \to \Omega^2(\sigma) g_{\alpha\beta}(\sigma) .$$

Or, infinitesimally, we can write $\Omega^2(\sigma) = e^{2\phi(\sigma)}$ for small $\phi$ so that

$$\delta g_{\alpha\beta}(\sigma) = 2\phi(\sigma) g_{\alpha\beta}(\sigma) .$$
It is simple to see that the Polyakov action is invariant under this transformation: the factor of $\Omega^2$ drops out just as the factor of $f$ did in equation (1.25), canceling between $\sqrt{-g}$ and the inverse metric $g^{\alpha\beta}$. This is a gauge symmetry of the string, as seen by the fact that the parameter $\Omega$ depends on the worldsheet coordinates $\sigma$. This means that two metrics which are related by a Weyl transformation (1.26) are to be considered as the same physical state.

![Figure 7: An example of a Weyl transformation](image)

How should we think of Weyl invariance? It is not a coordinate change. Instead it is the invariance of the theory under a local change of scale which preserves the angles between all lines. For example the two worldsheet metrics shown in the figure are viewed by the Polyakov string as equivalent. This is rather surprising! And, as you might imagine, theories with this property are extremely rare. It should be clear from the discussion above that the property of Weyl invariance is special to two dimensions, for only there does the scaling factor coming from the determinant $\sqrt{-g}$ cancel that coming from the inverse metric. But even in two dimensions, if we wish to keep Weyl invariance then we are strictly limited in the kind of interactions that can be added to the action. For example, we would not be allowed a potential term for the worldsheet scalars of the form,

$$\int d^2\sigma \sqrt{-g} V(X).$$

These break Weyl invariance. Nor can we add a worldsheet cosmological constant term,

$$\mu \int d^2\sigma \sqrt{-g}.$$

This too breaks Weyl invariance. We will see later in this course that the requirement of Weyl invariance becomes even more stringent in the quantum theory. We will also see what kind of interactions terms can be added to the worldsheet. Indeed, much of this course can be thought of as the study of theories with Weyl invariance.
1.3.2 Fixing a Gauge

As we have seen, the equation of motion (1.23) looks pretty nasty. However, we can use the redundancy inherent in the gauge symmetry to choose coordinates in which they simplify. Let’s think about what we can do with the gauge symmetry.

Firstly, we have two reparameterizations to play with. The worldsheet metric has three independent components. This means that we expect to be able to set any two of the metric components to a value of our choosing. We will choose to make the metric locally conformally flat, meaning

\[ g_{\alpha\beta} = e^{2\phi} \eta_{\alpha\beta} , \tag{1.27} \]

where \( \phi(\sigma, \tau) \) is some function on the worldsheet. You can check that this is possible by writing down the change of the metric under a coordinate transformation and seeing that the differential equations which result from the condition (1.27) have solutions, at least locally. Choosing a metric of the form (1.27) is known as conformal gauge.

We have only used reparameterization invariance to get to the metric (1.27). We still have Weyl transformations to play with. Clearly, we can use these to remove the last independent component of the metric and set \( \phi = 0 \) such that,

\[ g_{\alpha\beta} = \eta_{\alpha\beta} . \tag{1.28} \]

We end up with the flat metric on the worldsheet in Minkowski coordinates.

A Diversion: How to make a metric flat

The fact that we can use Weyl invariance to make any two-dimensional metric flat is an important result. Let’s take a quick diversion from our main discussion to see a different proof that isn’t tied to the choice of Minkowski coordinates on the worldsheet. We’ll work in 2d Euclidean space to avoid annoying minus signs. Consider two metrics related by a Weyl transformation, \( g'_{\alpha\beta} = e^{2\phi} g_{\alpha\beta} \). One can check that the Ricci scalars of the two metrics are related by,

\[ \sqrt{g} R' = \sqrt{g} (R - 2\nabla^2 \phi) . \tag{1.29} \]

We can therefore pick a \( \phi \) such that the new metric has vanishing Ricci scalar, \( R' = 0 \), simply by solving this differential equation for \( \phi \). However, in two dimensions (but not in higher dimensions) a vanishing Ricci scalar implies a flat metric. The reason is simply that there aren’t too many indices to play with. In particular, symmetry of the Riemann tensor in two dimensions means that it must take the form,

\[ R_{\alpha\beta\gamma\delta} = \frac{R}{2} (g_{\alpha\gamma} g_{\beta\delta} - g_{\alpha\delta} g_{\beta\gamma}) . \]
So $R' = 0$ is enough to ensure that $R'_{\alpha\beta\gamma\delta} = 0$, which means that the manifold is flat. In equation (1.28), we’ve further used reparameterization invariance to pick coordinates in which the flat metric is the Minkowski metric.

**The equations of motion and the stress-energy tensor**

With the choice of the flat metric (1.28), the Polyakov action simplifies tremendously and becomes the theory of $D$ free scalar fields. (In fact, this simplification happens in any conformal gauge).

$$S = -\frac{1}{4\alpha'} \int d^2\sigma \partial_\alpha X \cdot \partial^\alpha X ,$$

and the equations of motion for $X^\mu$ reduce to the free wave equation,

$$\partial_\alpha \partial^\alpha X^\mu = 0 .$$

Now that looks too good to be true! Are the horrible equations (1.23) really equivalent to a free wave equation? Well, not quite. There is something that we’ve forgotten: we picked a choice of gauge for the metric $g_{\alpha\beta}$. But we must still make sure that the equation of motion for $g_{\alpha\beta}$ is satisfied. In fact, the variation of the action with respect to the metric gives rise to a rather special quantity: it is the stress-energy tensor, $T_{\alpha\beta}$.

With a particular choice of normalization convention, we define the stress-energy tensor to be

$$T_{\alpha\beta} = -\frac{2}{T} \frac{1}{\sqrt{-g}} \frac{\partial S}{\partial g^{\alpha\beta}} .$$

We varied the Polyakov action with respect to $g_{\alpha\beta}$ in (1.24). When we set $g_{\alpha\beta} = \eta_{\alpha\beta}$ we get

$$T_{\alpha\beta} = \partial_\alpha X \cdot \partial_\beta X - \frac{1}{2} \eta_{\alpha\beta} \eta^{\rho\sigma} \partial_\rho X \cdot \partial_\sigma X .$$

The equation of motion associated to the metric $g_{\alpha\beta}$ is simply $T_{\alpha\beta} = 0$. Or, more explicitly,

$$T_{01} = \dot{X} \cdot X' = 0$$
$$T_{00} = T_{11} = \frac{1}{2}(\dot{X}^2 + X'^2) = 0 .$$

We therefore learn that the equations of motion of the string are the free wave equations (1.31) subject to the two constraints (1.33) arising from the equation of motion $T_{\alpha\beta} = 0$. 

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Getting a feel for the constraints

Let’s try to get some intuition for these constraints. There is a simple meaning of the first constraint in (1.33): we must choose our parameterization such that lines of constant $\sigma$ are perpendicular to the lines of constant $\tau$, as shown in the figure.

But we can do better. To gain more physical insight, we need to make use of the fact that we haven’t quite exhausted our gauge symmetry. We will discuss this more in Section 2.2, but for now one can check that there is enough remnant gauge symmetry to allow us to go to static gauge,

$$X^0 \equiv t = R\tau ,$$

so that $(X^0)' = 0$ and $\dot{X}^0 = R$, where $R$ is a constant that is needed on dimensional grounds. The interpretation of this constant will become clear shortly. Then, writing $X^\mu = (t, \vec{x})$, the equation of motion for spatial components is the free wave equation,

$$\ddot{\vec{x}} - \vec{x}'' = 0$$

while the constraints become

$$\dot{\vec{x}} \cdot \vec{x}' = 0$$
$$\dot{\vec{x}}^2 + \vec{x}'^2 = R^2$$

The first constraint tells us that the motion of the string must be perpendicular to the string itself. In other words, the physical modes of the string are transverse oscillations. There is no longitudinal mode. We’ll also see this again in Section 2.2.

From the second constraint, we can understand the meaning of the constant $R$: it is related to the length of the string when $\dot{x} = 0$,

$$\int d\sigma \sqrt{(d\vec{x}/d\sigma)^2} = 2\pi R .$$

Of course, if we have a stretched string with $\dot{x} = 0$ at one moment of time, then it won’t stay like that for long. It will contract under its own tension. As this happens, the second constraint equation relates the length of the string to the instantaneous velocity of the string.
1.4 Mode Expansions

Let’s look at the equations of motion and constraints more closely. The equations of motion (1.31) are easily solved. We introduce lightcone coordinates on the worldsheet,

$$\sigma^\pm = \tau \pm \sigma ,$$

in terms of which the equations of motion simply read

$$\partial_+ \partial_- X^\mu = 0$$

The most general solution is,

$$X^\mu(\sigma, \tau) = X_L^\mu(\sigma^+) + X_R^\mu(\sigma^-)$$

for arbitrary functions $X_L^\mu$ and $X_R^\mu$. These describe left-moving and right-moving waves respectively. Of course the solution must still obey both the constraints (1.33) as well as the periodicity condition,

$$X^\mu(\sigma, \tau) = X^\mu(\sigma + 2\pi, \tau) . \quad (1.35)$$

The most general, periodic solution can be expanded in Fourier modes,

$$X_L^\mu(\sigma^+) = \frac{1}{2} x^\mu + \frac{1}{2} \alpha' p^\mu \sigma^+ + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \bar{\alpha}_n e^{-in\sigma^+} ,$$

$$X_R^\mu(\sigma^-) = \frac{1}{2} x^\mu + \frac{1}{2} \alpha' p^\mu \sigma^- + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \alpha_n e^{-in\sigma^-} . \quad (1.36)$$

This mode expansion will be very important when we come to the quantum theory. Let’s make a few simple comments here.

- Various normalizations in this expression, such as the $\alpha'$ and factor of $1/n$ have been chosen for later convenience.

- $X_L$ and $X_R$ do not individually satisfy the periodicity condition (1.35) due to the terms linear in $\sigma^\pm$. However, the sum of them is invariant under $\sigma \to \sigma + 2\pi$ as required.

- The variables $x^\mu$ and $p^\mu$ are the position and momentum of the center of mass of the string. This can be checked, for example, by studying the Noether currents arising from the spacetime translation symmetry $X^\mu \to X^\mu + c^\mu$. One finds that the conserved charge is indeed $p^\mu$.

- Reality of $X^\mu$ requires that the coefficients of the Fourier modes, $\alpha_n^\mu$ and $\tilde{\alpha}_n^\mu$, obey

$$\alpha_n^\mu = (\alpha_{-n}^\mu)^* , \quad \tilde{\alpha}_n^\mu = (\tilde{\alpha}_{-n}^\mu)^* . \quad (1.37)$$
1.4.1 The Constraints Revisited

We still have to impose the two constraints (1.33). In the worldsheet lightcone coordinates $\sigma^\pm$, these become,

$$(\partial_+ X)^2 = (\partial_- X)^2 = 0 \ .$$

These equations give constraints on the momenta $p^\mu$ and the Fourier modes $\alpha_n^\mu$ and $\tilde{\alpha}_n^\mu$.

To see what these are, let’s look at

$$\partial_- X^\mu = \partial_- X_R^\mu = \frac{\alpha'}{2} p^\mu + \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \alpha_n^\mu e^{-i n \sigma^-}$$

where in the second line the sum is over all $n \in \mathbb{Z}$ and we have defined $\alpha_0^\mu$ to be

$$\alpha_0^\mu \equiv \sqrt{\frac{\alpha'}{2}} p^\mu \ .$$

The constraint (1.38) can then be written as

$$(\partial_- X)^2 = \frac{\alpha'}{2} \sum_{m,p} \alpha_m \cdot \alpha_p e^{-i(m+p)\sigma^-}$$

$$= \frac{\alpha'}{2} \sum_{m,n} \alpha_m \cdot \alpha_{n-m} e^{-i n \sigma^-}$$

$$\equiv \alpha' \sum_n L_n e^{-i n \sigma^-} = 0 \ .$$

where we have defined the sum of oscillator modes,

$$L_n = \frac{1}{2} \sum_m \alpha_{n-m} \cdot \alpha_m \ .$$

We can also do the same for the left-moving modes, where we again define an analogous sum of operator modes,

$$\tilde{L}_n = \frac{1}{2} \sum_m \tilde{\alpha}_{n-m} \cdot \tilde{\alpha}_m \ .$$

with the zero mode defined to be,

$$\tilde{\alpha}_0^\mu \equiv \sqrt{\frac{\alpha'}{2}} p^\mu \ .$$
The fact that $\tilde{\alpha}_0^\mu = \alpha_0^\mu$ looks innocuous but is a key point to remember when we come to quantize the string. The $L_n$ and $\tilde{L}_n$ are the Fourier modes of the constraints. Any classical solution of the string of the form (1.36) must further obey the infinite number of constraints,

$$L_n = \tilde{L}_n = 0 \quad n \in \mathbb{Z}.$$ 

We’ll meet these objects $L_n$ and $\tilde{L}_n$ again in a more general context when we come to discuss conformal field theory.

The constraints arising from $L_0$ and $\tilde{L}_0$ have a rather special interpretation. This is because they include the square of the spacetime momentum $p^\mu$. But, the square of the spacetime momentum is an important quantity in Minkowski space: it is the square of the rest mass of a particle,

$$p_\mu p^\mu = -M^2.$$ 

So the $L_0$ and $\tilde{L}_0$ constraints tell us the effective mass of a string in terms of the excited oscillator modes, namely

$$M^2 = \frac{4}{\alpha'} \sum_{n>0} \alpha_n \cdot \alpha_{-n} = \frac{4}{\alpha'} \sum_{n>0} \tilde{\alpha}_n \cdot \tilde{\alpha}_{-n} \quad (1.41)$$

Because both $\alpha_0^\mu$ and $\tilde{\alpha}_0^\mu$ are equal to $\sqrt{\alpha'/2} p^\mu$, we have two expressions for the invariant mass: one in terms of right-moving oscillators $\alpha_n^\mu$ and one in terms of left-moving oscillators $\tilde{\alpha}_n^\mu$. And these two terms must be equal to each other. This is known as level matching. It will play an important role in the next section where we turn to the quantum theory.
2. The Quantum String

Our goal in this section is to quantize the string. We have seen that the string action involves a gauge symmetry and whenever we wish to quantize a gauge theory we’re presented with a number of different ways in which we can proceed. If we’re working in the canonical formalism, this usually boils down to one of two choices:

- We could first quantize the system and then subsequently impose the constraints that arise from gauge fixing as operator equations on the physical states of the system. For example, in QED this is the Gupta-Bleuler method of quantization that we use in Lorentz gauge. In string theory it consists of treating all fields $X^\mu$, including time $X^0$, as operators and imposing the constraint equations (1.33) on the states. This is usually called covariant quantization.

- The alternative method is to first solve all of the constraints of the system to determine the space of physically distinct classical solutions. We then quantize these physical solutions. For example, in QED, this is the way we proceed in Coulomb gauge. Later in this chapter, we will see a simple way to solve the constraints of the free string.

Of course, if we do everything correctly, the two methods should agree. Usually, each presents a slightly different challenge and offers a different viewpoint.

In these lectures, we’ll take a brief look at the first method of covariant quantization. However, at the slightest sign of difficulties, we’ll bail! It will be useful enough to see where the problems lie. We’ll then push forward with the second method described above which is known as lightcone quantization in string theory. Although we’ll succeed in pushing quantization through to the end, our derivations will be a little cheap and unsatisfactory in places. In Section 5 we’ll return to all these issues, armed with more sophisticated techniques from conformal field theory.

2.1 A Lightning Look at Covariant Quantization

We wish to quantize $D$ free scalar fields $X^\mu$ whose dynamics is governed by the action (1.30). We subsequently wish to impose the constraints

\[ \dot{X} \cdot X' = \dot{X}^2 + X'^2 = 0. \]  

The first step is easy. We promote $X^\mu$ and their conjugate momenta $\Pi_\mu = (1/2\pi\alpha')\dot{X}_\mu$ to operator valued fields obeying the canonical equal-time commutation relations,

\[
[X^\mu(\sigma, \tau), \Pi_\nu(\sigma', \tau)] = i\delta(\sigma - \sigma') \delta^\mu_\nu, \\
[X^\mu(\sigma, \tau), X^\nu(\sigma', \tau)] = [\Pi_\mu(\sigma, \tau), \Pi_\nu(\sigma', \tau)] = 0.
\]
We translate these into commutation relations for the Fourier modes $x^\mu$, $p^\mu$, $\alpha_n^\mu$ and $\tilde{\alpha}_n^\mu$. Using the mode expansion (1.36) we find
\[
[x^\mu, p_\nu] = i\delta_\mu^\nu \quad \text{and} \quad [\alpha_n^\mu, \alpha_m^\nu] = [\tilde{\alpha}_n^\mu, \tilde{\alpha}_m^\nu] = n \eta^\mu_\nu \delta_{n+m,0} ,
\] (2.2)
with all others zero. The commutation relations for $x^\mu$ and $p^\mu$ are expected for operators governing the position and momentum of the center of mass of the string. The commutation relations of $\alpha_n^\mu$ and $\tilde{\alpha}_n^\mu$ are those of harmonic oscillator creation and annihilation operators in disguise. And the disguise isn’t that good. We just need to define (ignoring the $\mu$ index for now)
\[
a_n = \frac{\alpha_n}{\sqrt{n}} \quad , \quad a_n^\dagger = \frac{\alpha_{-n}}{\sqrt{n}} \quad n > 0
\] (2.3)
Then (2.2) gives the familiar $[a_n, a_m^\dagger] = \delta_{mn}$. So each scalar field gives rise to two infinite towers of creation and annihilation operators, with $\alpha_n$ acting as a rescaled annihilation operator for $n > 0$ and as a creation operator for $n < 0$. There are two towers because we have right-moving modes $\alpha_n$ and left-moving modes $\tilde{\alpha}_n$.

With these commutation relations in hand we can now start building the Fock space of our theory. We introduce a vacuum state of the string $|0\rangle$, defined to obey
\[
\alpha_n^\mu |0\rangle = 0 \quad \text{for} \quad n > 0
\] (2.4)
The vacuum state of string theory has a different interpretation from the analogous object in field theory. This is not the vacuum state of spacetime. It is instead the vacuum state of a single string. This is reflected in the fact that the operators $x^\mu$ and $p^\mu$ give extra structure to the vacuum. The true ground state of the string is $|0\rangle$, tensored with a spatial wavefunction $\Psi(x)$. Alternatively, if we work in momentum space, the vacuum carries another quantum number, $p^\mu$, which is the eigenvalue of the momentum operator. We should therefore write the vacuum as $|0; p\rangle$, which still obeys (2.4), but now also
\[
\hat{p}^\mu |0; p\rangle = p^\mu |0; p\rangle
\] (2.5)
where (for the only time in these lecture notes) we’ve put a hat on the momentum operator $\hat{p}^\mu$ on the left-hand side of this equation to distinguish it from the eigenvalue $p^\mu$ on the right-hand side.

We can now start to build up the Fock space by acting with creation operators $\alpha_n^\mu$ and $\tilde{\alpha}_n^\mu$ with $n < 0$. A generic state comes from acting with any number of these creation operators on the vacuum,
\[
(|0\rangle)^{\nu_1} (\alpha_{-1}^\mu)^{\nu_2} \cdots (\alpha_{-2}^\mu)^{\nu_2} \cdots |0; p\rangle
\]
Each state in the Fock space is a different excited state of the string. Each has the interpretation of a different species of particle in spacetime. We’ll see exactly what particles they are shortly. But for now, notice that because there’s an infinite number of ways to excite a string there are an infinite number of different species of particles in this theory.

2.1.1 Ghosts

There’s a problem with the Fock space that we’ve constructed: it doesn’t have positive norm. The reason for this is that one of the scalar fields, $X^0$, comes with the wrong sign kinetic term in the action (1.30). From the perspective of the commutation relations, this issue raises its head in presence of the spacetime Minkowski metric in the expression

$$[\alpha^\mu_n, \alpha^\nu_m] = n \eta^{\mu\nu} \delta_{n,m}.$$ 

This gives rise to the offending negative norm states, which come with an odd number of timelike oscillators excited, for example

$$\langle p'; 0 | \alpha^0_1 \alpha^0_{-1} | 0; p \rangle \sim -\delta^D(p - p')$$

This is the first problem that arises in the covariant approach to quantization. States with negative norm are referred to as ghosts. To make sense of the theory, we have to make sure that they can’t be produced in any physical processes. Of course, this problem is familiar from attempts to quantize QED in Lorentz gauge. In that case, gauge symmetry rides to the rescue since the ghosts are removed by imposing the gauge fixing constraint. We must hope that the same happens in string theory.

2.1.2 Constraints

Although we won’t push through with this programme at the present time, let us briefly look at what kind of constraints we have in string theory. In terms of Fourier modes, the classical constraints can be written as $L_n = \bar{L}_n = 0$, where

$$L_n = \frac{1}{2} \sum_m \alpha_{n-m} \cdot \alpha_m$$

and similar for $\bar{L}_n$. As in the Gupta-Bleuler quantization of QED, we don’t impose all of these as operator equations on the Hilbert space. Instead we only require that the operators $L_n$ and $\bar{L}_n$ have vanishing matrix elements when sandwiched between two physical states $|\text{phys}\rangle$ and $|\text{phys}'\rangle$,

$$\langle \text{phys}' | L_n | \text{phys} \rangle = \langle \text{phys}' | \bar{L}_n | \text{phys} \rangle = 0$$
Because \( L_n^1 = L_{-n} \), it is therefore sufficient to require
\[
L_n |\text{phys}\rangle = \tilde{L}_n |\text{phys}\rangle = 0 \quad \text{for } n > 0 \quad (2.6)
\]
However, we still haven’t explained how to impose the constraints \( L_0 \) and \( \tilde{L}_0 \). And these present a problem that doesn’t arise in the case of QED. The problem is that, unlike for \( L_n \) with \( n \neq 0 \), the operator \( L_0 \) is not uniquely defined when we pass to the quantum theory. There is an operator ordering ambiguity arising from the commutation relations (2.2). Commuting the \( \alpha_n^\mu \) operators past each other in \( L_0 \) gives rise to extra constant terms.

**Question:** How do we know what order to put the \( \alpha_n^\mu \) operators in the quantum operator \( L_0 \)? Or the \( \tilde{\alpha}_n^\mu \) operators in \( \tilde{L}_0 \)?

**Answer:** We don’t! Yet. Naively it looks as if each different choice will define a different theory when we impose the constraints. To make this ambiguity manifest, for now let’s just pick a choice of ordering. We define the quantum operators to be normal ordered, with the annihilation operators \( \alpha_n^\mu \), \( n > 0 \), moved to the right,
\[
L_0 = \sum_{m=1}^{\infty} \alpha_{-m} \cdot \alpha_m + \frac{1}{2} \alpha_0^2 \quad , \quad \tilde{L}_0 = \sum_{m=1}^{\infty} \tilde{\alpha}_{-m} \cdot \tilde{\alpha}_m + \frac{1}{2} \tilde{\alpha}_0^2
\]
Then the ambiguity rears its head in the different constraint equations that we could impose, namely
\[
(L_0 - a) |\text{phys}\rangle = (\tilde{L}_0 - a) |\text{phys}\rangle = 0 \quad (2.7)
\]
for some constant \( a \).

As we saw classically, the operators \( L_0 \) and \( \tilde{L}_0 \) play an important role in determining the spectrum of the string because they include a term quadratic in the momentum \( \alpha_n^\mu = \tilde{\alpha}_n^\mu = \sqrt{\alpha'/2} p^\mu \). Combining the expression (1.41) with our constraint equation for \( L_0 \) and \( \tilde{L}_0 \), we find the spectrum of the string is given by,
\[
M^2 = \frac{4}{\alpha'} \left( -a + \sum_{m=1}^{\infty} \alpha_{-m} \cdot \alpha_m \right) = \frac{4}{\alpha'} \left( -a + \sum_{m=1}^{\infty} \tilde{\alpha}_{-m} \cdot \tilde{\alpha}_m \right)
\]
We learn therefore that the undetermined constant \( a \) has a direct physical effect: it changes the mass spectrum of the string. In the quantum theory, the sums over \( \alpha_n^\mu \) modes are related to the number operators for the harmonic oscillator: they count the number of excited modes of the string. The level matching in the quantum theory tells us that the number of left-moving modes must equal the number of right-moving modes.
Ultimately, we will find that the need to decouple the ghosts forces us to make a unique choice for the constant $a$. (Spoiler alert: it turns out to be $a = 1$). In fact, the requirement that there are no ghosts is much stronger than this. It also restricts the number of scalar fields that we have in the theory. (Another spoiler: $D = 26$). If you’re interested in how this works in covariant formulation then you can read about it in the book by Green, Schwarz and Witten. Instead, we’ll show how to quantize the string and derive these values for $a$ and $D$ in lightcone gauge. However, after a trip through the world of conformal field theory, we’ll come back to these ideas in a context which is closer to the covariant approach.

### 2.2 Lightcone Quantization

We will now take the second path described at the beginning of this section. We will try to find a parameterization of all classical solutions of the string. This is equivalent to finding the classical phase space of the theory. We do this by solving the constraints (2.1) in the classical theory, leaving behind only the physical degrees of freedom.

Recall that we fixed the gauge to set the worldsheet metric to

$$g_{\alpha \beta} = \eta_{\alpha \beta}.$$  

However, this isn’t the end of our gauge freedom. There still remain gauge transformations which preserve this choice of metric. In particular, any coordinate transformation $\sigma \to \tilde{\sigma}(\sigma)$ which changes the metric by

$$\eta_{\alpha \beta} \to \Omega^2(\sigma) \eta_{\alpha \beta}, \quad (2.8)$$

can be undone by a Weyl transformation. What are these coordinate transformations? It’s simplest to answer this using lightcone coordinates on the worldsheet,

$$\sigma^\pm = \tau \pm \sigma, \quad (2.9)$$

where the flat metric on the worldsheet takes the form,

$$ds^2 = -d\sigma^+ d\sigma^-$$

In these coordinates, it’s clear that any transformation of the form

$$\sigma^+ \to \tilde{\sigma}^+(\sigma^+) \quad , \quad \sigma^- \to \tilde{\sigma}^-(\sigma^-), \quad (2.10)$$

simply multiplies the flat metric by an overall factor (2.8) and so can be undone by a compensating Weyl transformation. Some quick comments on this surviving gauge symmetry:
Recall that in Section 1.3.2 we used the argument that 3 gauge invariances (2 reparameterizations + 1 Weyl) could be used to fix 3 components of the worldsheet metric $g_{\alpha\beta}$. What happened to this argument? Why do we still have some gauge symmetry left? The reason is that $\tilde{\sigma}^{\pm}$ are functions of just a single variable, not two. So we did fix nearly all our gauge symmetries. What is left is a set of measure zero amongst the full gauge symmetry that we started with.

The remaining reparameterization invariance (2.10) has an important physical implication. Recall that the solutions to the equations of motion are of the form $X^\mu_L(\sigma^+) + X^\mu_R(\sigma^-)$ which looks like 2D functions worth of solutions. Of course, we still have the constraints which, in terms of $\sigma^{\pm}$, read

$$\left( \partial_+ X \right)^2 = \left( \partial_- X \right)^2 = 0 ,$$

which seems to bring the number down to $2(D - 1)$ functions. But the reparameterization invariance (2.10) tells us that even some of these are fake since we can always change what we mean by $\sigma^{\pm}$. The physical solutions of the string are therefore actually described by $2(D - 2)$ functions. But this counting has a nice interpretation: the degrees of freedom describe the transverse fluctuations of the string.

The above comment reaches the same conclusion as the discussion in Section 1.3.2. There, in an attempt to get some feel for the constraints, we claimed that we could go to static gauge $X^0 = R\tau$ for some dimensionful parameter $R$. It is easy to check that this is simple to do using reparameterizations of the form (2.10). However, to solve the string constraints in full, it turns out that static gauge is not that useful. Rather we will use something called “lightcone gauge”.

### 2.2.1 Lightcone Gauge

We would like to gauge fix the remaining reparameterization invariance (2.10). The best way to do this is called lightcone gauge. In counterpoint to the worldsheet lightcone coordinates (2.9), we introduce the spacetime lightcone coordinates,

$$X^\pm = \sqrt{\frac{1}{2}} (X^0 \pm X^{D-1}) .$$

Note that this choice picks out a particular time direction and a particular spatial direction. It means that any calculations that we do involving $X^\pm$ will not be manifestly Lorentz invariant. You might think that we needn’t really worry about this. We could try to make the following argument: “The equations may not look Lorentz invariant
but, since we started from a Lorentz invariant theory, at the end of the day any physical process is guaranteed to obey this symmetry”. Right?! Well, unfortunately not. One of the more interesting and subtle aspects of quantum field theory is the possibility of anomalies: these are symmetries of the classical theory that do not survive the journey of quantization. When we come to the quantum theory, if our equations don’t look Lorentz invariant then there’s a real possibility that it’s because the underlying physics actually isn’t Lorentz invariant. Later we will need to spend some time figuring out under what circumstances our quantum theory keeps the classical Lorentz symmetry.

In lightcone coordinates, the spacetime Minkowski metric reads

\[ ds^2 = -2dX^+dX^- + \sum_{i=1}^{D-2} dX^i dX^i \]

This means that indices are raised and lowered with \( A_+ = -A^- \) and \( A_- = -A^+ \) and \( A_i = A^i \). The product of spacetime vectors reads \( A \cdot B = -A^+ B^- - A^- B^+ + A^i B^i \).

Let’s look at the solution to the equation of motion for \( X^+ \). It reads,

\[ X^+ = X_L^+(\sigma^+) + X_R^+(\sigma^-) \]

We now gauge fix. We use our freedom of reparameterization invariance to choose coordinates such that

\[ X_L^+ = \frac{1}{2} x^+ + \frac{1}{2} \alpha' p^+ \sigma^+ , \quad X_R^+ = \frac{1}{2} x^+ + \frac{1}{2} \alpha' p^+ \sigma^- . \]

You might think that we could go further and eliminate \( p^+ \) and \( x^+ \) but this isn’t possible because we don’t quite have the full freedom of reparameterization invariance since all functions should remain periodic in \( \sigma \). The upshot of this choice of gauge is that

\[ X^+ = x^+ + \alpha' p^+ \tau . \quad (2.13) \]

This is lightcone gauge. Notice that, as long as \( p^+ \neq 0 \), we can always shift \( x^+ \) by a shift in \( \tau \).

There’s something a little disconcerting about the choice (2.13). We’ve identified a timelike worldsheet coordinate with a null spacetime coordinate. Nonetheless, as you can see from the figure, it seems to be a good parameterization of the worldsheet. One could imagine that the parameterization might break if the string is actually massless and travels in the \( X^- \) direction, with \( p^+ = 0 \). But otherwise, all should be fine.
Solving for $X^-$

The choice (2.13) does the job of fixing the reparameterization invariance (2.10). As we will now see, it also renders the constraint equations trivial. The first thing that we have to worry about is the possibility of extra constraints arising from this new choice of gauge fixing. This can be checked by looking at the equation of motion for $X^+$,

$$\partial_+ \partial_- X^- = 0$$

But we can solve this by the usual ansatz,

$$X^- = X^+_L (\sigma^+) + X^+_R (\sigma^-).$$

We’re still left with all the other constraints (2.11). Here we see the real benefit of working in lightcone gauge (which is actually what makes quantization possible at all): $X^-$ is completely determined by these constraints. For example, the first of these reads

$$2 \partial_+ X^- \partial_+ X^+ = \sum_{i=1}^{D-2} \partial_+ X^i \partial_+ X^i$$

(2.14)

which, using (2.13), simply becomes

$$\partial_+ X^-_L = \frac{1}{\alpha' p^+} \sum_{i=1}^{D-2} \partial_+ X^i \partial_+ X^i.$$  

(2.15)

Similarly,

$$\partial_- X^-_R = \frac{1}{\alpha' p^+} \sum_{i=1}^{D-2} \partial_- X^i \partial_- X^i.$$  

(2.16)

So, up to an integration constant, the function $X^- (\sigma^+, \sigma^-)$ is completely determined in terms of the other fields. If we write the usual mode expansion for $X^+_L/R$

$$X^+_L (\sigma^+) = \frac{1}{2} x^- + \frac{1}{2} \alpha' p^- \sigma^+ + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \tilde{\alpha}_n e^{-i n \sigma^+},$$

$$X^-_R (\sigma^-) = \frac{1}{2} x^- + \frac{1}{2} \alpha' p^- \sigma^- + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \alpha_n e^{-i n \sigma^-}.$$  

then $x^-$ is the undetermined integration constant, while $p^-$, $\alpha_n$ and $\tilde{\alpha}_n$ are all fixed by the constraints (2.15) and (2.16). For example, the oscillator modes $\alpha_n^-$ are given by,

$$\alpha_n^- = \sqrt{\frac{1}{2 \alpha' p^+}} \sum_{m=-\infty}^{+\infty} \sum_{i=1}^{D-2} \alpha_{n-m}^i \alpha^i_m,$$  

(2.17)
A special case of this is the \( \alpha_0^- = \sqrt{\alpha'/2} p^- \) equation, which reads

\[
\frac{\alpha' p^-}{2} = \frac{1}{2p^\pm} \sum_{i=1}^{D-2} \left( \frac{1}{2} \alpha'^i p'^i + \sum_{n \neq 0} \alpha_n^i \alpha_{-n}^i \right) .
\]  

We also get another equation for \( p^- \) from the \( \tilde{\alpha}_0^- \) equation arising from (2.15)

\[
\frac{\alpha' p^-}{2} = \frac{1}{2p^\pm} \sum_{i=1}^{D-2} \left( \frac{1}{2} \alpha'^i p'^i + \sum_{n \neq 0} \tilde{\alpha}_n^i \tilde{\alpha}_{-n}^i \right) .
\]  

From these two equations, we can reconstruct the old, classical, level matching conditions (1.41). But now with a difference:

\[
M^2 = 2p^+ p^- - \sum_{i=1}^{D-2} p^i p^i = \frac{4}{\alpha'} \sum_{i=1}^{D-2} \sum_{n > 0} \alpha_{-n}^i \alpha_n^i = \frac{4}{\alpha'} \sum_{i=1}^{D-2} \sum_{n > 0} \tilde{\alpha}_{-n}^i \tilde{\alpha}_n^i .
\]  

The difference is that now the sum is over oscillators \( \alpha^i \) and \( \tilde{\alpha}^i \) only, with \( i = 1, \ldots, D-2 \). We’ll refer to these as transverse oscillators. Note that the string isn’t necessarily living in the \( X^0-X^{D-1} \) plane, so these aren’t literally the transverse excitations of the string. Nonetheless, if we specify the \( \alpha^i \) then all other oscillator modes are determined. In this sense, they are the physical excitation of the string.

Let’s summarize the state of play so far. The most general classical solution is described in terms of \( 2(D-2) \) transverse oscillator modes \( \alpha_n^i \) and \( \tilde{\alpha}_n^i \), together with a number of zero modes describing the center of mass and momentum of the string: \( x^i, p^i, p^+ \) and \( x^- \). But \( x^+ \) can be absorbed by a shift of \( \tau \) in (2.13) and \( p^- \) is constrained to obey (2.18) and (2.19). In fact, \( p^- \) can be thought of as (proportional to) the lightcone Hamiltonian. Indeed, we know that \( p^- \) generates translations in \( x^+ \), but this is equivalent to shifts in \( \tau \).

### 2.2.2 Quantization

Having identified the physical degrees of freedom, let’s now quantize. We want to impose commutation relations. Some of these are easy:

\[
[x^i, p^j] = i \delta^{ij} , \quad [x^-, p^+] = -i \\
[\alpha_n^i, \alpha_m^j] = [\tilde{\alpha}_n^i, \tilde{\alpha}_m^j] = n \delta^{ij} \delta_{n+m,0} .
\]  

\[
(2.21)
\]
all of which follow from the commutation relations (2.2) that we saw in covariant quantization1.

What to do with $x^+$ and $p^-$? We could implement $p^-$ as the Hamiltonian acting on states. In fact, it will prove slightly more elegant (but equivalent) if we promote both $x^+$ and $p^-$ to operators with the expected commutation relation,

$$[x^+, p^-] = -i .$$

This is morally equivalent to writing $[t, H] = -i$ in non-relativistic quantum mechanics, which is true on a formal level. In the present context, it means that we can once again choose states to be eigenstates of $p^\mu$, with $\mu = 0, \ldots, D$, but the constraints (2.18) and (2.19) must still be imposed as operator equations on the physical states. We’ll come to this shortly.

The Hilbert space of states is very similar to that described in covariant quantization: we define a vacuum state, $|0; p\rangle$ such that

$$\hat{p}^\mu |0; p\rangle = p^\mu |0; p\rangle , \quad \alpha^i_n |0; p\rangle = \tilde{\alpha}^i_n |0; p\rangle = 0 \quad \text{for } n > 0$$

and we build a Fock space by acting with the creation operators $\alpha^i_n$ and $\tilde{\alpha}^i_n$ with $n > 0$. The difference with the covariant quantization is that we only act with transverse oscillators which carry a spatial index $i = 1, \ldots, D - 2$. For this reason, the Hilbert space is, by construction, positive definite. We don’t have to worry about ghosts.

---

1Mea Culpa: We’re not really supposed to do this. The whole point of the approach that we’re taking is to quantize just the physical degrees of freedom. The resulting commutation relations are not, in general, inherited from the larger theory that we started with simply by closing our eyes and forgetting about all the other fields that we’ve gauge fixed. We can see the problem by looking at (2.17), where $\alpha^i_-$ is determined in terms of $\alpha^i_+$. This means that the commutation relations for $\alpha^i_+$ might be infected by those of $\alpha^i_-$ which could potentially give rise to extra terms. The correct procedure to deal with this is to figure out the Poisson bracket structure of the physical degrees of freedom in the classical theory. Or, in fancier language, the symplectic form on the phase space which schematically looks like

$$\omega \sim \int d\sigma - d\dot{X}^+ \wedge dX^- - d\dot{X}^- \wedge dX^+ + 2d\dot{X}^i \wedge dX^i ,$$

The reason that the commutation relations (2.21) do not get infected is because the $\alpha^-$ terms in the symplectic form come multiplying $X^+$. Yet $X^+$ is given in (2.13). It has no oscillator modes. That means that the symplectic form doesn’t pick up the Fourier modes of $X^-$ and so doesn’t receive any corrections from $\alpha^-_n$. The upshot of this is that the naive commutation relations (2.21) are actually right.
The Constraints

Because $p^{-}$ is not an independent variable in our theory, we must impose the constraints (2.18) and (2.19) by hand as operator equations which define the physical states. In the classical theory, we saw that these constraints are equivalent to mass-shell conditions (2.20).

But there’s a problem when we go to the quantum theory. It’s the same problem that we saw in covariant quantization: there’s an ordering ambiguity in the sum over oscillator modes on the right-hand side of (2.20). If we choose all operators to be normal ordered then this ambiguity reveals itself in an overall constant, $a$, which we have not yet determined. The final result for the mass of states in lightcone gauge is:

$$M^2 = \frac{4}{\alpha'} \left( \sum_{i=1}^{D-2} \sum_{n>0} \alpha_i^i \alpha_n^i - a \right) = \frac{4}{\alpha'} \left( \sum_{i=1}^{D-2} \sum_{n>0} \tilde{\alpha}_i^i \tilde{\alpha}_n^i - a \right)$$

Since we’ll use this formula quite a lot in what follows, it’s useful to introduce quantities related to the number operators of the harmonic oscillator,

$$N = \sum_{i=1}^{D-2} \sum_{n>0} \alpha_i^i \alpha_n^i \quad , \quad \tilde{N} = \sum_{i=1}^{D-2} \sum_{n>0} \tilde{\alpha}_i^i \tilde{\alpha}_n^i . \quad (2.24)$$

These are not quite number operators because of the factor of $1/\sqrt{n}$ in (2.3). The value of $N$ and $\tilde{N}$ is often called the level. Which, if nothing else, means that the name “level matching” makes sense. We now have

$$M^2 = \frac{4}{\alpha'} (N - a) = \frac{4}{\alpha'} (\tilde{N} - a) . \quad (2.25)$$

How are we going to fix $a$? Later in the course we’ll see the correct way to do it. For now, I’m just going to give you a quick and dirty derivation.

The Casimir Energy

“I told him that the sum of an infinite no. of terms of the series: 1 + 2 + 3 + 4 + \ldots = -\frac{1}{12} under my theory. If I tell you this you will at once point out to me the lunatic asylum as my goal. ”

* Ramanujan, in a letter to G.H.Hardy.

What follows is a heuristic derivation of the normal ordering constant $a$. Suppose that we didn’t notice that there was any ordering ambiguity and instead took the naive classical result directly over to the quantum theory, that is

$$\frac{1}{2} \sum_{n \neq 0} \alpha_{-n}^i \alpha_n^i = \frac{1}{2} \sum_{n<0} \alpha_{-n}^i \alpha_n^i + \frac{1}{2} \sum_{n>0} \alpha_{-n}^i \alpha_n^i .$$
where we’ve left the sum over \( i = 1, \ldots, D - 2 \) implicit. We’ll now try to put this in normal ordered form, with the annihilation operators \( \alpha_n^i \) with \( n > 0 \) on the right-hand side. It’s the first term that needs changing. We get

\[
\frac{1}{2} \sum_{n<0} [\alpha_n^i \alpha_{-n}^i - n(D - 2)] + \frac{1}{2} \sum_{n>0} \alpha_{-n}^i \alpha_n^i = \sum_{n>0} \alpha_{-n}^i \alpha_n^i + \frac{D-2}{2} \sum_{n>0} n .
\]

The final term clearly diverges. But it at least seems to have a physical interpretation: it is the sum of zero point energies of an infinite number of harmonic oscillators. In fact, we came across exactly the same type of term in the course on quantum field theory where we learnt that, despite the divergence, one can still extract interesting physics from this. This is the physics of the Casimir force.

Let’s recall the steps that we took to derive the Casimir force. Firstly, we introduced an ultra-violet cut-off \( \epsilon \ll 1 \), probably muttering some words about no physical plates being able to withstand very high energy quanta. Unfortunately, those words are no longer available to us in string theory, but let’s proceed regardless. We replace the divergent sum over integers by the expression,

\[
\sum_{n=1}^{\infty} n \rightarrow \sum_{n=1}^{\infty} ne^{-\epsilon n} = -\frac{\partial}{\partial \epsilon} \sum_{n=1}^{\infty} e^{-\epsilon n} = -\frac{\partial}{\partial \epsilon} (1 - e^{-\epsilon})^{-1} = \frac{1}{\epsilon^2} - \frac{1}{12} + O(\epsilon).
\]

Obviously the \( 1/\epsilon^2 \) piece diverges as \( \epsilon \rightarrow 0 \). This term should be renormalized away. In fact, this is necessary to preserve the Weyl invariance of the Polyakov action since it contributes to a cosmological constant on the worldsheet. After this renormalization, we’re left with the wonderful answer, first intuited by Ramanujan

\[
\sum_{n=1}^{\infty} n = -\frac{1}{12} .
\]

While heuristic, this argument does predict the correct physical Casimir energy measured in one-dimensional systems. For example, this effect is seen in simulations of quantum spin chains.

What does this mean for our string? It means that we should take the unknown constant \( a \) in the mass formula (2.25) to be,

\[
M^2 = \frac{4}{\alpha'} \left( N - \frac{D-2}{24} \right) = \frac{4}{\alpha'} \left( \tilde{N} - \frac{D-2}{24} \right) .
\]

This is the formula that we will use to determine the spectrum of the string.
Zeta Function Regularization

I appreciate that the preceding argument is not totally convincing. We could spend some time making it more robust at this stage, but it’s best if we wait until later in the course when we will have the tools of conformal field theory at our disposal. We will eventually revisit this issue and provide a respectable derivation of the Casimir energy in Section 4.4.1. For now I merely offer an even less convincing argument, known as zeta-function regularization.

The zeta-function is defined, for \( \text{Re}(s) > 1 \), by the sum

\[
\zeta(s) = \sum_{n=1}^{\infty} n^{-s}.
\]

But \( \zeta(s) \) has a unique analytic continuation to all values of \( s \). In particular,

\[
\zeta(-1) = -\frac{1}{12}.
\]

Good? Good. This argument is famously unconvincing the first time you meet it! But it’s actually a very useful trick for getting the right answer.

2.3 The String Spectrum

Finally, we’re in a position to analyze the spectrum of a single, free string.

2.3.1 The Tachyon

Let’s start with the ground state \( |0; p \rangle \) defined in (2.23). With no oscillators excited, the mass formula (2.26) gives

\[
M^2 = -\frac{1}{\alpha'} \frac{D - 2}{6}.
\]

(2.27)

But that’s a little odd. It’s a negative mass-squared. Such particles are called tachyons.

In fact, tachyons aren’t quite as pathological as you might think. If you’ve heard of these objects before, it’s probably in the context of special relativity where they’re strange beasts which always travel faster than the speed of light. But that’s not the right interpretation. Rather we should think more in the language of quantum field theory. Suppose that we have a field in spacetime — let’s call it \( T(X) \) — whose quanta will give rise to this particle. The mass-squared of the particle is simply the quadratic term in the action, or

\[
M^2 = \left. \frac{\partial^2 V(T)}{\partial T^2} \right|_{T=0}
\]
So the negative mass-squared in (2.27) is telling us that we’re expanding around a maximum of the potential for the tachyon field as shown in the figure. Note that from this perspective, the Higgs field in the standard model at $H = 0$ is also a tachyon.

The fact that string theory turns out to sit at an unstable point in the tachyon field is unfortunate. The natural question is whether the potential has a good minimum elsewhere, as shown in the figure to the right. No one knows the answer to this! Naive attempts to understand this don’t work. We know that around $T = 0$, the leading order contribution to the potential is negative and quadratic. But there are further terms that we can compute using techniques that we’ll describe in Section 6. An expansion of the tachyon potential around $T = 0$ looks like

$$V(T) = \frac{1}{2}M^2T^2 + c_3T^3 + c_4T^4 + \ldots$$

It turns out that the $T^3$ term in the potential does give rise to a minimum. But the $T^4$ term destabilizes it again. Moreover, the $T$ field starts to mix with other scalar fields in the theory that we will come across soon. The ultimate fate of the tachyon in the bosonic string is not yet understood.

The tachyon is a problem for the bosonic string. It may well be that this theory makes no sense — or, at the very least, has no time-independent stable solutions. Or perhaps we just haven’t worked out how to correctly deal with the tachyon. Either way, the problem does not arise when we introduce fermions on the worldsheet and study the superstring. This will involve several further technicalities which we won’t get into in this course. Instead, our time will be put to better use if we continue to study the bosonic string since all the lessons that we learn will carry over directly to the superstring. However, one should be aware that the problem of the unstable vacuum will continue to haunt us throughout this course.

Although we won’t describe it in detail, at several times along our journey we’ll make an aside about how calculations work out for the superstring.

2.3.2 The First Excited States

We now look at the first excited states. If we act with a creation operator $\alpha_{-1}^j$, then the level matching condition (2.25) tells us that we also need to act with a $\tilde{\alpha}_{-1}^i$ operator. This gives us $(D - 2)^2$ particle states,

$$\tilde{\alpha}_{-1}^i \alpha_{-1}^j |0; p\rangle , \quad (2.28)$$
each of which has mass

\[ M^2 = \frac{4}{\alpha'} \left( 1 - \frac{D - 2}{24} \right). \]

But now we seem to have a problem. Our states have space indices \( i, j = 1, \ldots, D - 2 \).
The operators \( \alpha^i \) and \( \tilde{\alpha}^i \) each transform in the vector representation of \( SO(D - 2) \subset SO(1, D - 1) \) which is manifest in lightcone gauge. But ultimately we want these states to fit into some representation of the full Lorentz \( SO(1, D - 1) \) group. That looks as if it’s going to be hard to arrange. This is the first manifestation of the comment that we made after equation (2.12): it’s tricky to see Lorentz invariance in lightcone gauge.

To proceed, let’s recall Wigner’s classification of representations of the Poincaré group. We start by looking at massive particles in \( R^{1,D-1} \). After going to the rest frame of the particle by setting \( p^\mu = (p, 0, \ldots, 0) \), we can watch how any internal indices transform under the little group \( SO(D - 1) \) of spatial rotations. The upshot of this is that any massive particle must form a representation of \( SO(D - 1) \). But the particles described by (2.28) have \((D - 2)^2\) states. There’s no way to package these states into a representation of \( SO(D - 1) \) and this means that there’s no way that the first excited states of the string can form a massive representation of the \( D \)-dimensional Poincaré group.

It looks like we’re in trouble. Thankfully, there’s a way out. If the states are massless, then we can’t go to the rest frame. The best that we can do is choose a spacetime momentum for the particle of the form \( p^\mu = (p, 0, \ldots, 0, p) \). In this case, the particles fill out a representation of the little group \( SO(D - 2) \). This means that massless particles get away with having fewer internal states than massive particles. For example, in four dimensions the photon has two polarization states, but a massive spin-1 particle must have three.

The first excited states (2.28) happily sit in a representation of \( SO(D - 2) \). We learn that if we want the quantum theory to preserve the \( SO(1, D - 1) \) Lorentz symmetry that we started with, then these states will have to be massless. And this is only the case if the dimension of spacetime is

\[ D = 26. \]

This is our first derivation of the critical dimension of the bosonic string.

Moreover, we’ve found that our theory contains a bunch of massless particles. And massless particles are interesting because they give rise to long range forces. Let’s look
more closely at what massless particles the string has given us. The states (2.28) transform in the $24 \otimes 24$ representation of SO(24). These decompose into three irreducible representations:

traceless symmetric $\oplus$ anti-symmetric $\oplus$ singlet (=trace)

To each of these modes, we associate a massless field in spacetime such that the string oscillation can be identified with a quantum of these fields. The fields are:

$$G_{\mu\nu}(X) , \quad B_{\mu\nu}(X) , \quad \Phi(X) \quad (2.29)$$

Of these, the first is the most interesting and we shall have more to say momentarily. The second is an anti-symmetric tensor field which is usually called the anti-symmetric tensor field. It also goes by the names of the “Kalb-Ramond field” or, in the language of differential geometry, the “2-form”. The scalar field is called the dilaton. These three massless fields are common to all string theories. We’ll learn more about the role these fields play later in the course.

The particle in the symmetric traceless representation of SO(24) is particularly interesting. This is a massless spin 2 particle. However, there are general arguments, due originally to Feynman and Weinberg, that any theory of interacting massless spin two particles must be equivalent to general relativity. We should therefore identify the field $G_{\mu\nu}(X)$ with the metric of spacetime. Let’s pause briefly to review the thrust of these arguments.

**Why Massless Spin 2 = General Relativity**

Let’s call the spacetime metric $G_{\mu\nu}(X)$. We can expand around flat space by writing

$$G_{\mu\nu} = \eta_{\mu\nu} + h_{\mu\nu}(X) .$$

Then the Einstein-Hilbert action has an expansion in powers of $h$. If we truncate to quadratic order, we simply have a free theory which we may merrily quantize in the usual canonical fashion: we promote $h_{\mu\nu}$ to an operator and introduce the associated creation and annihilation operators $a_{\mu\nu}$ and $a^\dagger_{\mu\nu}$. This way of looking at gravity is anathema to those raised in the geometrical world of general relativity. But from a particle physics language it is very standard: it is simply the quantization of a massless spin 2 field, $h_{\mu\nu}$.

---

2A very readable description of this can be found in the first few chapters of the Feynman Lectures on Gravitation.
However, even on this simple level, there is a problem due to the indefinite signature of the spacetime Minkowski metric. The canonical quantization relations of the creation and annihilation operators are schematically of the form,

\[ [a_{\mu\nu}, a_{\rho\sigma}^\dagger] \sim \eta_{\mu\rho}\eta_{\nu\sigma} + \eta_{\mu\sigma}\eta_{\nu\rho} \]

But this will lead to a Hilbert space with negative norm states coming from acting with time-like creation operators. For example, the one-graviton state of the form,

\[ a_{0i}^\dagger |0\rangle \quad (2.30) \]

suffers from a negative norm. This should be becoming familiar by now: it is the usual problem that we run into if we try to covariantly quantize a gauge theory. And, indeed, general relativity is a gauge theory. The gauge transformations are diffeomorphisms. We would hope that this saves the theory of quantum gravity from these negative norm states.

Let’s look a little more closely at what the gauge symmetry looks like for small fluctuations \( h_{\mu\nu} \). We’ve butchered the Einstein-Hilbert action and left only terms quadratic in \( h \). Including all the index contractions, we find

\[
S_{EH} = \frac{M_{pl}^2}{2} \int d^4x \left[ \partial_\mu h_{\rho}^\rho \partial_\nu h^{\mu\nu} - \partial_\nu h^{\rho\mu} \partial_\mu h_{\rho\nu} + \frac{1}{2} \partial_\mu h_{\rho\nu} \partial^\rho h^{\mu\nu} - \frac{1}{2} \partial_\mu h_{\nu\rho} \partial^\rho h^{\mu\nu} \right] + \ldots
\]

One can check that this truncated action is invariant under the gauge symmetry,

\[ h_{\mu\nu} \longrightarrow h_{\mu\nu} + \partial_\mu \xi_\nu + \partial_\nu \xi_\mu \quad (2.31) \]

for any function \( \xi_\mu(X) \). The gauge symmetry is the remnant of diffeomorphism invariance, restricted to small deviations away from flat space. With this gauge invariance in hand one can show that, just like QED, the negative norm states decouple from all physical processes.

To summarize, theories of massless spin 2 fields only make sense if there is a gauge symmetry to remove the negative norm states. In general relativity, this gauge symmetry descends from diffeomorphism invariance. The argument of Feynman and Weinberg now runs this logic in reverse. It goes as follows: suppose that we have a massless, spin 2 particle. Then, at the linearized level, it must be invariant under the gauge symmetry (2.31) in order to eliminate the negative norm states. Moreover, this symmetry must survive when interaction terms are introduced. But the only way to do this is to ensure that the resulting theory obeys diffeomorphism invariance. That means the theory of any interacting, massless spin 2 particle is Einstein gravity, perhaps supplemented by higher derivative terms.
We haven’t yet shown that string theory includes interactions for $h_{\mu\nu}$ but we will come to this later in the course. More importantly, we will also explicitly see how Einstein’s field equations arise directly in string theory.

**A Comment on Spacetime Gauge Invariance**

We’ve surreptitiously put $\mu, \nu = 0, \ldots, 25$ indices on the spacetime fields, rather than $i, j = 1, \ldots, 24$. The reason we’re allowed to do this is because both $G_{\mu\nu}$ and $B_{\mu\nu}$ enjoy a spacetime gauge symmetry which allows us to eliminate appropriate modes. Indeed, this is exactly the gauge symmetry (2.31) that entered the discussion above. It isn’t possible to see these spacetime gauge symmetries from the lightcone formalism of the string since, by construction, we find only the physical states (although, by consistency alone, the gauge symmetries must be there). One of the main advantages of pushing through with the covariant calculation is that it does allow us to see how the spacetime gauge symmetry emerges from the string worldsheet. Details can be found in Green, Schwarz and Witten. We’ll also briefly return to this issue in Section 5.

**2.3.3 Higher Excited States**

We rescued the Lorentz invariance of the first excited states by choosing $D = 26$ to ensure that they are massless. But now we’ve used this trick once, we still have to worry about all the other excited states. These also carry indices that take the range $i, j = 1, \ldots, D - 2 = 24$ and, from the mass formula (2.26), they will all be massive and so must form representations of $SO(D - 1)$. It looks like we’re in trouble again.

Let’s examine the string at level $N = \tilde{N} = 2$. In the right-moving sector, we now have two different states: $\alpha_{-1}^i \alpha_{-1}^j |0\rangle$ and $\alpha_{-2}^i |0\rangle$. The same is true for the left-moving sector, meaning that the total set of states at level 2 is (in notation that is hopefully obvious, but probably technically wrong)

$$(\alpha_{-1}^i \alpha_{-1}^j \oplus \alpha_{-2}^i) \otimes (\tilde{\alpha}_{-1}^i \tilde{\alpha}_{-1}^j \oplus \tilde{\alpha}_{-2}^i) |0; p\rangle.$$ 

These states have mass $M^2 = 4/\alpha'$. How many states do we have? In the left-moving sector, we have,

$$\frac{1}{2}(D - 2)(D - 1) + (D - 2) = \frac{1}{2}D(D - 1) - 1.$$ 

But, remarkably, that does fit nicely into a representation of $SO(D - 1)$, namely the traceless symmetric tensor representation.

In fact, one can show that all excited states of the string fit nicely into $SO(D - 1)$ representations. The only consistency requirement that we need for Lorentz invariance is to fix up the first excited states: $D = 26$. 

– 45 –
Note that if we are interested in a fundamental theory of quantum gravity, then all these excited states will have masses close to the Planck scale so are unlikely to be observable in particle physics experiments. Nonetheless, as we shall see when we come to discuss scattering amplitudes, it is the presence of this infinite tower of states that tames the ultra-violet behaviour of gravity.

2.4 Lorentz Invariance Revisited

The previous discussion allowed us to derive both the critical dimension and the spectrum of string theory in the quickest fashion. But the derivation creaks a little in places. The calculation of the Casimir energy is unsatisfactory the first time one sees it. Similarly, the explanation of the need for massless particles at the first excited level is correct, but seems rather cheap considering the huge importance that we’re placing on the result.

As I’ve mentioned a few times already, we’ll shortly do better and gain some physical insight into these issues, in particular the critical dimension. But here I would just like to briefly sketch how one can be a little more rigorous within the framework of lightcone quantization. The question, as we’ve seen, is whether one preserves spacetime Lorentz symmetry when we quantize in lightcone gauge. We can examine this more closely.

Firstly, let’s go back to the action for free scalar fields (1.30) before we imposed lightcone gauge fixing. Here the full Poincaré symmetry was manifest: it appears as a global symmetry on the worldsheet,

$$X^\mu \to \Lambda^\mu_\nu X^\nu + e^\mu$$

But recall that in field theory, global symmetries give rise to Noether currents and their associated conserved charges. What are the Noether currents associated to this Poincaré transformation? We can start with the translations $X^\mu \to X^\mu + e^\mu$. A quick computation shows that the current is,

$$P^\alpha_\mu = T \partial^\alpha X_\mu$$

which is indeed a conserved current since $\partial_\alpha P^\alpha_\mu = 0$ is simply the equation of motion. Similarly, we can compute the $\frac{1}{2} D(D - 1)$ currents associated to Lorentz transformations. They are,

$$J^{\alpha}_{\mu\nu} = P^{\alpha}_\mu X_\nu - P^{\alpha}_\nu X_\mu$$

It’s not hard to check that $\partial_\alpha J^{\alpha}_{\mu\nu} = 0$ when the equations of motion are obeyed.
The conserved charges arising from this current are given by $M_{\mu\nu} = \int d\sigma J^\tau_{\mu\nu}$. Using the mode expansion (1.36) for $X^\mu$, these can be written as

$$M_{\mu\nu} = \int d\sigma J^\tau_{\mu\nu} = \sum_{n=1}^{\infty} \frac{1}{n} (\alpha_n^\mu \alpha_n^\nu - \alpha_n^\nu \alpha_n^\mu) - i \sum_{n=1}^{\infty} \frac{1}{n} (\bar{\alpha}_n^\mu \bar{\alpha}_n^\nu - \bar{\alpha}_n^\nu \bar{\alpha}_n^\mu)$$

$$\equiv l_{\mu\nu} + S_{\mu\nu} + \tilde{S}_{\mu\nu}$$

The first piece, $l_{\mu\nu}$, is the orbital angular momentum of the string while the remaining pieces $S_{\mu\nu}$ and $\tilde{S}_{\mu\nu}$ tell us the angular momentum due to excited oscillator modes. Classically, these obey the Poisson brackets of the Lorentz algebra. Moreover, if we quantize in the covariant approach, the corresponding operators obey the commutation relations of the Lorentz Lie algebra, namely

$$[M^{\rho\sigma}, M^{\tau\nu}] = \eta^{\rho\tau} M^{\sigma\nu} - \eta^{\rho\nu} M^{\sigma\tau} - \eta^{\sigma\nu} M^{\rho\tau}$$

However, things aren’t so easy in lightcone gauge. Lorentz invariance is not guaranteed and, in general, is not there. The right way to go about looking for it is to make sure that the Lorentz algebra above is reproduced by the generators $M_{\mu\nu}$. It turns out that the smoking gun lies in the commutation relation,

$$[M^{i-}, M^{j-}] = 0$$

Does this equation hold in lightcone gauge? The problem is that it involves the operators $p^-$ and $\alpha_n^-$, both of which are fixed by (2.17) and (2.18) in terms of the other operators. So the task is to compute this commutation relation $[M^{i-}, M^{j-}]$, given the commutation relations (2.21) for the physical degrees of freedom, and check that it vanishes. To do this, we re-instate the ordering ambiguity $a$ and the number of spacetime dimension $D$ as arbitrary variables and proceed.

The part involving orbital angular momenta $l_{\mu\nu}$ is fairly straightforward. (Actually, there’s a small subtlety because we must first make sure that the operator $l_{\mu\nu}$ is Hermitian by replacing $x^\mu p^\nu$ with $\frac{1}{2} (x^\mu p^\nu + p^\nu x^\mu)$). The real difficulty comes from computing the commutation relations $[S^{i-}, S^{j-}]$. This is messy$^3$. After a tedious computation, one finds,

$$[M^{i-}, M^{j-}] = \frac{2}{(p^+)^2} \sum_{n>0} \left[ \frac{D-2}{24} - 1 \right] n + \frac{1}{n} \left[ a - \frac{D-2}{24} \right] (\alpha_n^i \alpha_n^j - \alpha_n^j \alpha_n^i) + (\alpha \leftrightarrow \bar{\alpha})$$

$^3$The original, classic, paper where lightcone quantization was first implemented is Goddard, Goldstone, Rebbi and Thorn “Quantum Dynamics of a Massless Relativistic String”, Nucl. Phys. B56 (1973). A pedestrian walkthrough of this calculation can be found in the lecture notes by Gleb Arutyunov. A link is given on the course webpage.
The right-hand side does not, in general, vanish. We learn that the relativistic string can only be quantized in flat Minkowski space if we pick,

\[ D = 26 \quad \text{and} \quad a = 1. \]

2.5 A Nod to the Superstring

We won’t provide details of the superstring in this course, but will pause occasionally to make some pertinent comments. Although what follows is nothing more than a list of facts, it will hopefully be helpful in orienting you when you do come to study this material.

The key difference between the bosonic string and the superstring is the addition of fermionic modes on its worldsheet. The resulting worldsheet theory is supersymmetric. (At least in the so-called Neveu-Schwarz-Ramond formalism). Hence the name “superstring”. Applying the kind of quantization procedure we’ve discussed in this section, one finds the following results:

- The critical dimension of the superstring is \( D = 10 \).
- There is no tachyon in the spectrum.
- The massless bosonic fields \( G_{\mu\nu}, B_{\mu\nu} \) and \( \Phi \) are all part of the spectrum of the superstring. In this context, \( B_{\mu\nu} \) is sometimes referred to as the Neveu-Schwarz 2-form. There are also massless spacetime fermions, as well as further massless bosonic fields. As we now discuss, the exact form of these extra bosonic fields depends on exactly what superstring theory we consider.

While the bosonic string is unique, there are a number of discrete choices that one can make when adding fermions to the worldsheet. This gives rise to a handful of different perturbative superstring theories. (Although later developments reveal that they are actually all part of the same framework which sometimes goes by the name of \( M\)-theory). The most important of these discrete options is whether we add fermions in both the left-moving and right-moving sectors of the string, or whether we choose the fermions to move only in one direction, usually taken to be right-moving. This gives rise to two different classes of string theory.

- Type II strings have both left and right-moving worldsheet fermions. The resulting spacetime theory in \( D = 10 \) dimensions has \( \mathcal{N} = 2 \) supersymmetry, which means 32 supercharges.

- Heterotic strings have just right-moving fermions. The resulting spacetime theory has \( \mathcal{N} = 1 \) supersymmetry, or 16 supercharges.
In each of these cases, there is then one further discrete choice that we can make. This leaves us with four superstring theories. In each case, the massless bosonic fields include $G_{\mu\nu}$, $B_{\mu\nu}$ and $\Phi$ together with a number of extra fields. These are:

- **Type IIA**: In the type II theories, the extra massless bosonic excitations of the string are referred to as *Ramond-Ramond* fields. For Type IIA, they are a 1-form $C_{\mu}$ and a 3-form $C_{\mu\nu\rho}$. Each of these is to be thought of as a gauge field. The gauge invariant information lies in the field strengths which take the form $F = dC$.

- **Type IIB**: The Ramond-Ramond gauge fields consist of a scalar $C$, a 2-form $C_{\mu\nu}$ and a 4-form $C_{\mu\nu\rho\sigma}$. The 4-form is restricted to have a self-dual field strength: $F_5 = *F_5$. (Actually, this statement is almost true...we'll look a little closer at this in Section 7.3.3).

- **Heterotic $SO(32)$**: The heterotic strings do not have Ramond-Ramond fields. Instead, each comes with a non-Abelian gauge field in spacetime. The heterotic strings are named after the gauge group. For example, the Heterotic $SO(32)$ string gives rise to an $SO(32)$ Yang-Mills theory in ten dimensions.

- **Heterotic $E_8 \times E_8$**: The clue is in the name. This string gives rise to an $E_8 \times E_8$ Yang-Mills field in ten-dimensions.

It is sometimes said that there are five perturbative superstring theories in ten dimensions. Here we’ve only mentioned four. The remaining theory is called Type I and includes open strings moving in flat ten dimensional space as well as closed strings. We’ll mention it in passing in the following section.
3. Open Strings and D-Branes

In this section we discuss the dynamics of open strings. Clearly their distinguishing feature is the existence of two end points. Our goal is to understand the effect of these end points. The spatial coordinate of the string is parameterized by

\[ \sigma \in [0, \pi] . \]

The dynamics of a generic point on a string is governed by local physics. This means that a generic point has no idea if it is part of a closed string or an open string. The dynamics of an open string must therefore still be described by the Polyakov action. But this must now be supplemented by something else: boundary conditions to tell us how the end points move. To see this, let’s look at the Polyakov action in conformal gauge

\[
S = -\frac{1}{4\pi \alpha'} \int d^2 \sigma \, \partial^\alpha X \cdot \partial^\alpha X .
\]

As usual, we derive the equations of motion by finding the extrema of the action. This involves an integration by parts. Let’s consider the string evolving from some initial configuration at \( \tau = \tau_i \) to some final configuration at \( \tau = \tau_f \):

\[
\delta S = -\frac{1}{2\pi \alpha'} \int_{\tau_i}^{\tau_f} d\tau \int_0^\pi d\sigma \, \partial^\alpha X \cdot \partial^\alpha \delta X
\]

\[= \frac{1}{2\pi \alpha'} \int d^2 \sigma \left( \partial^\alpha \partial^\alpha X \right) \cdot \delta X + \text{total derivative}
\]

For an open string the total derivative picks up the boundary contributions

\[
\frac{1}{2\pi \alpha'} \left[ \int_0^\pi d\sigma \cdot \delta X \right]_{\tau = \tau_f}^{\tau = \tau_i} - \frac{1}{2\pi \alpha'} \left[ \int_{\tau_i}^{\tau_f} d\tau \cdot X' \cdot \delta X \right]_{\sigma = 0}^{\sigma = \pi}
\]

The first term is the kind that we always get when using the principle of least action. The equations of motion are derived by requiring that \( \delta X^\mu = 0 \) at \( \tau = \tau_i \) and \( \tau_f \) and so it vanishes. However, the second term is novel. In order for it too to vanish, we require

\[ \partial_\sigma X^\mu \delta X^\mu = 0 \quad \text{at } \sigma = 0, \pi \]

There are two different types of boundary conditions that we can impose to satisfy this:
Neumann boundary conditions.

\[ \partial_\sigma X^\mu = 0 \quad \text{at } \sigma = 0, \pi \]  

(3.1)

Because there is no restriction on \( \delta X^\mu \), this condition allows the end of the string to move freely. To see the consequences of this, it’s useful to repeat what we did for the closed string and work in static gauge with \( X^0 \equiv t = R \tau \), for some dimensionful constant \( R \). Then, as in equations (1.34), the constraints read

\[ \ddot{\vec{x}} \cdot \vec{x}' = 0 \quad \text{and} \quad \dot{\vec{x}}^2 + \vec{x}'^2 = R^2 \]

But at the end points of the string, \( \vec{x}' = 0 \). So the second equation tells us that \( |d\vec{x}/dt| = 1 \). Or, in other words, the end point of the string moves at the speed of light.

Dirichlet boundary conditions

\[ \delta X^\mu = 0 \quad \text{at } \sigma = 0, \pi \]  

(3.2)

This means that the end points of the string lie at some constant position, \( X^\mu = c^\mu \), in space.

At first sight, Dirichlet boundary conditions may seem a little odd. Why on earth would the strings be fixed at some point \( c^\mu \)? What is special about that point? Historically people were pretty hung up about this and Dirichlet boundary conditions were rarely considered until the mid-1990s. Then everything changed due to an insight of Polchinski...

Let’s consider Dirichlet boundary conditions for some coordinates and Neumann for the others. This means that at both end points of the string, we have

\[ \partial_\sigma X^a = 0 \quad \text{for } a = 0, \ldots, p \]

\[ X^I = c^I \quad \text{for } I = p + 1, \ldots, D - 1 \]  

(3.3)

This fixes the end-points of the string to lie in a \((p + 1)\)-dimensional hypersurface in spacetime such that the \( SO(1,D - 1) \) Lorentz group is broken to,

\[ SO(1,D - 1) \rightarrow SO(1,p) \times SO(D - p - 1) \].

This hypersurface is called a \( D\)-brane or, when we want to specify its dimension, a \( Dp \)-brane. Here \( D \) stands for Dirichlet, while \( p \) is the number of spatial dimensions of the brane. So, in this language, a \( D0 \)-brane is a particle; a \( D1 \)-brane is itself a string; a \( D2 \)-brane a membrane and so on. The brane sits at specific positions \( c^I \) in the transverse space. But what is the interpretation of this hypersurface?
It turns out that the D-brane hypersurface should be thought of as a new, dynamical object in its own right. This is a conceptual leap that is far from obvious. Indeed, it took decades for people to fully appreciate this fact. String theory is not just a theory of strings: it also contains higher dimensional branes. In Section 7.5 we will see how these D-branes develop a life of their own. Some comments:

- We’ve defined D-branes that are infinite in space. However, we could just as well define finite D-branes by specifying closed surfaces on which the string can end.

- There are many situations where we want to describe strings that have Neumann boundary conditions in all directions, meaning that the string is free to move throughout spacetime. It’s best to understand this in terms of a space-filling D-brane. No Dirichlet conditions means D-branes are everywhere!

- The D$p$-brane described above always has Neumann boundary conditions in the $X^0$ direction. What would it mean to have Dirichlet conditions for $X^0$? Obviously this is a little weird since the object is now localized at a fixed point in time. But there is an interpretation of such an object: it is an instanton. This “D-instanton” is usually referred to as a D($-1$)-brane. It is related to tunneling effects in the quantum theory.

Mode Expansion

We take the usual mode expansion for the string, with $X^\mu = X^\mu_L(\sigma^+) + X^\mu_R(\sigma^-)$ and

$$X^\mu_L(\sigma^+) = \frac{1}{2} x^\mu + \alpha^\mu \cdot p^\mu \sigma^+ + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \tilde{\alpha}_n e^{-i n \sigma^+},$$

$$X^\mu_R(\sigma^-) = \frac{1}{2} x^\mu + \alpha^\mu \cdot p^\mu \sigma^- + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \alpha_n e^{-i n \sigma^-}. \quad (3.4)$$

The boundary conditions impose relations on the modes of the string. They are easily checked to be:

- Neumann boundary conditions, $\partial_\sigma X^a = 0$, at the end points require that
  $$\alpha_n^a = \tilde{\alpha}_n^a \quad (3.5)$$

- Dirichlet boundary conditions, $X^I = c^I$, at the end points require that
  $$x^I = c^I, \quad p^I = 0, \quad \alpha_n^I = -\tilde{\alpha}_n^I$$

So for both boundary conditions, we only have one set of oscillators, say $\alpha_n$. The $\tilde{\alpha}_n$ are then determined by the boundary conditions.
It’s worth pointing out that there is a factor of 2 difference in the $p^\mu$ term between the open string (3.4) and the closed string (1.36). This is to ensure that $p^\mu$ for the open string retains the interpretation of the spacetime momentum of the string when $\sigma \in [0, \pi]$. To see this, one needs to check the Noether current associated to translations of $X^\mu$ on the worldsheet: it was given in (2.33). The conserved charge is then

$$P^\mu = \int_0^\pi d\sigma (P\tau)^\mu = \frac{1}{2\pi\alpha'} \int_0^\pi d\sigma \dot{X}'^\mu = p^\mu$$

as advertised. Note that we’ve needed to use the Neumann conditions (3.5) to ensure that the Fourier modes don’t contribute to this integral.

### 3.1 Quantization

To quantize, we promote the fields $x^a$ and $p^a$ and $\alpha^\mu_n$ to operators. The other elements in the mode expansion are fixed by the boundary conditions. An obvious, but important, point is that the position and momentum degrees of freedom, $x^a$ and $p^a$, have a spacetime index that takes values $a = 0, \ldots, p$. This means that the spatial wavefunctions only depend on the coordinates of the brane not the whole spacetime. Said another, quantizing an open string gives rise to states which are restricted to lie on the brane.

To determine the spectrum, it is again simplest to work in lightcone gauge. The spacetime lightcone coordinate is chosen to lie within the brane,

$$X^\pm = \sqrt{\frac{1}{2}} (X^0 \pm X^p)$$

Quantization now proceeds in the same manner as for the closed string until we arrive at the mass formula for states which is a sum over the transverse modes of the string.

$$M^2 = \frac{1}{\alpha'} \left( \sum_{i=1}^{p-1} \sum_{n>0} \alpha^{i}_{-n} \alpha^{i}_{n} + \sum_{i=p+1}^{D-1} \sum_{n>0} \alpha^{i}_{-n} \alpha^{i}_{n} - a \right)$$

The first sum is over modes parallel to the brane, the second over modes perpendicular to the brane. It’s worth commenting on the differences with the closed string formula. Firstly, there is an overall factor of 4 difference. This can be traced to the lack of the factor of 1/2 in front of $p^\mu$ in the mode expansion that we discussed above. Secondly, there is a sum only over $\alpha$ modes. The $\tilde{\alpha}$ modes are not independent because of the boundary conditions.
Open and Closed

In the mass formula, we have once again left the normal ordering constant $a$ ambiguous. As in the closed string case, requiring the Lorentz symmetry of the quantum theory — this time the reduced symmetry $SO(1, p) \times SO(D - p - 1)$ — forces us to choose

$$D = 26 \quad \text{and} \quad a = 1 .$$

These are the same values that we found for the closed string. This reflects an important fact: the open string and closed string are not different theories. They are both different states inside the same theory.

More precisely, theories of open strings necessarily contain closed strings. This is because, once we consider interactions, an open string can join to form a closed string as shown in the figure. We’ll look at interactions in Section 6. The question of whether this works the other way — meaning whether closed string theories require open strings — is a little more involved and is cleanest to state in the context of the superstring. For type II superstrings, the open strings and D-branes are necessary ingredients. For heterotic superstrings, there appear to be no open strings and no D-branes. For the bosonic theory, it seems likely that the open strings are a necessary ingredient although I don’t know of a killer argument. But since we’re not sure whether the theory exists due to the presence of the tachyon, the point is probably moot. In the remainder of these lectures, we’ll view the bosonic string in the same manner as the type II string and assume that the theory includes both closed strings and open strings with their associated D-branes.

3.1.1 The Ground State

The ground state is defined by

$$\alpha_n^i |0; p\rangle = 0 \quad n > 0$$

The spatial index now runs over $i = 1, \ldots, p − 1, p + 1, \ldots, D − 1$. The ground state has mass

$$M^2 = -\frac{1}{\alpha'}$$

It is again tachyonic. Its mass is half that of the closed string tachyon. As we commented above, this time the tachyon is confined to the brane. In contrast to the closed string tachyon, the open string tachyon is now fairly well understood and its potential
is of the form shown in the figure. The interpretation is that the brane is unstable. It will decay, much like a resonance state in field theory. It does this by dissolving into closed string modes. The end point of this process — corresponding to the minimum at $T > 0$ in the figure — is simply a state with no D-brane. The difference between the value of the potential at the minimum and at $T = 0$ is the tension of the D-brane.

Notice that although there is a minimum of the potential at $T > 0$, it is not a global minimum. The potential seems to drop off without bound to the left. This is still not well understood. There are suggestions that it is related in some way to the closed string tachyon.

3.1.2 First Excited States: A World of Light

The first excited states are massless. They fall into two classes:

- Oscillators longitudinal to the brane,
  \[ \alpha^a_{-1}|0; p\rangle \quad a = 1, \ldots, p - 1 \]
  The spacetime indices $a$ lie within the brane so this state transforms under the $SO(1, p)$ Lorentz group. It is a spin 1 particle on the brane or, in other words, it is a photon. We introduce a gauge field $A_a$ with $a = 0, \ldots, p$ lying on the brane whose quanta are identified with this photon.

- Oscillators transverse to the brane,
  \[ \alpha^I_{-1}|0; p\rangle \quad I = p + 1, \ldots, D - 1 \]
  These states are scalars under the $SO(1, p)$ Lorentz group of the brane. They can be thought of as arising from scalar fields $\phi^I$ living on the brane. These scalars have a nice interpretation: they are fluctuations of the brane in the transverse directions. This is our first hint that the D-brane is a dynamical object. Note that although the $\phi^I$ are scalar fields under the $SO(1, p)$ Lorentz group of the brane, they do transform as a vector under the $SO(D - p - 1)$ rotation group transverse to the brane. This appears as a global symmetry on the brane worldvolume.
3.1.3 Higher Excited States and Regge Trajectories

At level \( N \), the mass of the string state is

\[
M^2 = \frac{1}{\alpha'}(N - 1)
\]

The maximal spin of these states arises from the symmetric tensor. It is

\[
J_{\text{max}} = N = \alpha' M^2 + 1
\]

Plotting the spin vs. the mass-squared, we find straight lines. These are usually called Regge trajectories. (Or sometimes Chew-Frautschi trajectories). They are seen in Nature in both the spectrum of mesons and baryons. Some examples involving \( \rho \)-mesons are shown in the figure. These stringy Regge trajectories suggest a naive cartoon picture of mesons as two rotating quarks connected by a confining flux tube.

The value of the string tension required to match the hadron spectrum of QCD is \( T \sim 1 \text{ GeV} \). This relationship between the strong interaction and the open string was one of the original motivations for the development of string theory and it is from here that the parameter \( \alpha' \) gets its (admittedly rarely used) name “Regge slope”. In these enlightened modern times, the connection between the open string and quarks lives on in the AdS/CFT correspondence.

3.1.4 Another Nod to the Superstring

Just as supersymmetry eliminates the closed string tachyon, so it removes the open string tachyon. Open strings are an ingredient of the type II string theories. The possible D-branes are

- Type IIA string theory has stable \( D_p \)-branes with \( p \) even.
- Type IIB string theory has stable \( D_p \)-branes with \( p \) odd.

The most important reason that D-branes are stable in the type II string theories is that they are charged under the Ramond-Ramond fields. (This was actually Polchinski’s insight that made people take D-branes seriously). However, type II string theories also contain unstable branes, with \( p \) odd in type IIA and \( p \) even in type IIB.
The fifth string theory (which was actually the first to be discovered) is called Type I. Unlike the other string theories, it contains both open and closed strings moving in flat ten-dimensional Lorentz-invariant spacetime. It can be thought of as the Type IIB theory with a bunch of space-filling D9-branes, together with something called an orientifold plane. You can read about this in Polchinski.

As we mentioned above, the heterotic string doesn’t have (finite energy) D-branes. This is due to an inconsistency in any attempt to reflect left-moving modes into right-moving modes.

### 3.2 Brane Dynamics: The Dirac Action

We have introduced D-branes as fixed boundary conditions for the open string. However, we’ve already seen a hint that these objects are dynamical in their own right, since the massless scalar excitations \( \phi^I \) have a natural interpretation as transverse fluctuations of the brane. Indeed, if a theory includes both open strings and closed strings, then the D-branes have to be dynamical because there can be no rigid objects in a theory of gravity. The dynamical nature of D-branes will become clearer as the course progresses.

But any dynamical object should have an action which describes how it moves. Moreover, after our discussion in Section 1, we already know what this is! On grounds of Lorentz invariance and reparameterization invariance alone, the action must be a higher dimensional extension of the Nambu-Goto action. This is

\[
S_{Dp} = -T_p \int d^{p+1} \xi \sqrt{-\det \gamma} \tag{3.6}
\]

where \( T_p \) is the tension of the Dp-brane which we will determine later, while \( \xi^a, a = 0, \ldots, p \), are the worldvolume coordinates of the brane. \( \gamma_{ab} \) is the pull back of the spacetime metric onto the worldvolume,

\[
\gamma_{ab} = \frac{\partial X^\mu}{\partial \xi^a} \frac{\partial X^\nu}{\partial \xi^b} \eta_{\mu\nu}.
\]

This is called the *Dirac action*. It was first written down by Dirac for a membrane some time before Nambu and Goto rediscovered it in the context of the string.

To make contact with the fields \( \phi^I \), we can use the reparameterization invariance of the Dirac action to go to static gauge. For an infinite, flat Dp-brane we can choose

\[
X^a = \xi^a \quad a = 0, \ldots, p.
\]
The dynamical transverse coordinates are then identified with the fluctuations $\phi^I$ through

$$X^I(\xi) = 2\pi\alpha' \phi^I(\xi) \quad I = p + 1, \ldots, D - 1$$

However, the Dirac action can’t be the whole story. It describes the transverse fluctuations of the D-brane, but has nothing to say about the $U(1)$ gauge field $A_\mu$ which lives on the D-brane. There must be some action which describes how this gauge field moves as well. We will return to this in Section 7.

What’s Special About Strings?

We could try to quantize the Dirac action (3.6) for a D-brane in the same manner that we quantized the action for the string. Is this possible? The answer, at present, is no. There appear to be both technical and conceptual obstacles. The technical issue is just that it’s hard. Weyl invariance was one of our chief weapons in attacking the string, but it doesn’t hold for higher dimensional objects.

The conceptual issue is that quantizing a membrane, or higher dimensional object, would not give rise to a discrete spectrum of states which have the interpretation of particles. In this way, they appear to be fundamentally different from the string.

Let’s get some intuition for why this is the case. The energy of a string is proportional to its length. This ensures that strings behave more or less like familiar elastic bands. What about D2-branes? Now the energy is proportional to the area. In the back of your mind, you might be thinking of a rubber-like sheet. But membranes, and higher dimensional objects, governed by the Dirac action don’t behave as household rubber sheets. They are more flexible. This is because a membrane can form many different shapes with the same area. For example, a tubular membrane of length $L$ and radius $1/L$ has the same area for all values of $L$; short and stubby, or long and thin. This means that long thin spikes can develop on a membrane at no extra cost of energy. In particular, objects connected by long thin tubes have the same energy, regardless of their separation. After quantization, this property gives rise to a continuous spectrum of states. A quantum membrane, or higher dimensional object, does not have the single particle interpretation that we saw for the string. The expectation is that the quantum membrane should describe multi-particle states.
3.3 Multiple Branes: A World of Glue

Consider two parallel D*p-branes. An open string now has options. It could either end on the same brane, or stretch between the two branes. Let’s consider the string that stretches between the two. It obeys

\[ X^I(0, \tau) = c^I \quad \text{and} \quad X^I(\pi, \tau) = d^I \]

where \( c^I \) and \( d^I \) are the positions of the two branes. In terms of the mode expansion, this requires

\[ X^I = c^I + \frac{(d^I - c^I)\sigma}{\pi} + \text{oscillator modes} \]

The classical constraints then read

\[ \partial_+ X \cdot \partial_+ X = \alpha'^2 p^2 + \frac{|\vec{d} - \vec{c}|^2}{4\pi^2} + \text{oscillator modes} = 0 \]

which means the classical mass-shell condition is

\[ M^2 = \frac{|\vec{d} - \vec{c}|^2}{(2\pi\alpha')^2} + \text{oscillator modes} \]

The extra term has an obvious interpretation: it is the mass of a classical string stretched between the two branes. The quantization of this string proceeds as before. After we include the normal ordering constant, the ground state of this string is only tachyonic if \(|\vec{d} - \vec{c}|^2 < 4\pi^2\alpha'\). Or in other words, the ground state is tachyonic if the branes approach to a sub-stringy distance.

There is an obvious generalization of this to the case of \( N \) parallel branes. Each end point of the string has \( N \) possible places on which to end. We can label each end point with a number \( m, n = 1, \ldots, N \) which tell us which brane it ends on. This label is sometimes referred to as a Chan-Paton factor.

Consider now the situation where all branes lie at the same position in spacetime. Each end point of the string has \( N \) possible places on which to end. We can label each end point with a number \( m, n = 1, \ldots, N \) which tell us which brane it ends on. This label is sometimes referred to as a Chan-Paton factor.

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Consider now the situation where all branes lie at the same position in spacetime. Each end point can lie on one of \( N \) different branes, giving \( N^2 \) possibilities in total. Each of these strings has the mass spectrum of an open string, meaning that there are now \( N^2 \) different particles of each type. It’s natural to arrange the associated fields to sit inside \( N \times N \) Hermitian matrices. We then have the open string tachyon \( T^m_n \) and the massless fields

\[ (\phi^I)^m_n, \quad (A_a)^m_n \]  

(3.7)

Here the components of the matrix tell us which string the field came from. Diagonal components arise from strings which have both ends on the same brane.
The gauge field $A_a$ is particularly interesting. Written in this way, it looks like a $U(N)$ gauge connection. We will later see that this is indeed the case. One can show that as $N$ branes coincide, the $U(1)^N$ gauge symmetry of the branes is enhanced to $U(N)$. The scalar fields $\phi^I$ transform in the adjoint of this symmetry.
4. Introducing Conformal Field Theory

The purpose of this section is to get comfortable with the basic language of two-dimensional conformal field theory\(^4\). This is a topic which has many applications outside of string theory, most notably in statistical physics where it offers a description of critical phenomena. Moreover, it turns out that conformal field theories in two dimensions provide rare examples of interacting, yet exactly solvable, quantum field theories. In recent years, attention has focussed on conformal field theories in higher dimensions due to their role in the AdS/CFT correspondence.

A **conformal transformation** is a change of coordinates \(\sigma^\alpha \rightarrow \tilde{\sigma}^\alpha(\sigma)\) such that the metric changes by

\[
g_{\alpha\beta}(\sigma) \rightarrow \Omega^2(\sigma)g_{\alpha\beta}(\sigma)\]  

(4.1)

A **conformal field theory** (CFT) is a field theory which is invariant under these transformations. This means that the physics of the theory looks the same at all length scales. Conformal field theories care about angles, but not about distances.

A transformation of the form (4.1) has a different interpretation depending on whether we are considering a fixed background metric \(g_{\alpha\beta}\), or a dynamical background metric. When the metric is dynamical, the transformation is a diffeomorphism; this is a gauge symmetry. When the background is fixed, the transformation should be thought of as an honest, physical symmetry, taking the point \(\sigma^\alpha\) to point \(\tilde{\sigma}^\alpha\). This is now a global symmetry with the corresponding conserved currents.

In the context of string theory in the Polyakov formalism, the metric is dynamical and the transformations (4.1) are residual gauge transformations: diffeomorphisms which can be undone by a Weyl transformation.

In contrast, in this section we will be primarily interested in theories defined on fixed backgrounds. Apart from a few noticeable exceptions, we will usually take this background to be flat. This is the situation that we are used to when studying quantum field theory.

\(^4\)Much of the material covered in this section was first described in the ground breaking paper by Belavin, Polyakov and Zamalodchikov, “Infinite Conformal Symmetry in Two-Dimensional Quantum Field Theory”, Nucl. Phys. B241 (1984). The application to string theory was explained by Friedan, Martinec and Shenker in “Conformal Invariance, Supersymmetry and String Theory”, Nucl. Phys. B271 (1986). The canonical reference for learning conformal field theory is the excellent review by Ginsparg. A link can be found on the course webpage.
Of course, we can alternate between thinking of theories as defined on fixed or fluctuating backgrounds. Any theory of 2d gravity which enjoys both diffeomorphism and Weyl invariance will reduce to a conformally invariant theory when the background metric is fixed. Similarly, any conformally invariant theory can be coupled to 2d gravity where it will give rise to a classical theory which enjoys both diffeomorphism and Weyl invariance. Notice the caveat “classical”! In some sense, the whole point of this course is to understand when this last statement also holds at the quantum level.

Even though conformal field theories are a subset of quantum field theories, the language used to describe them is a little different. This is partly out of necessity. Invariance under the transformation (4.1) can only hold if the theory has no preferred length scale. But this means that there can be nothing in the theory like a mass or a Compton wavelength. In other words, conformal field theories only support massless excitations. The questions that we ask are not those of particles and S-matrices. Instead we will be concerned with correlation functions and the behaviour of different operators under conformal transformations.

4.0.1 Euclidean Space

Although we’re ultimately interested in Minkowski signature worldsheets, it will be much simpler and elegant if we work instead with Euclidean worldsheets. There’s no funny business here — everything we do could also be formulated in Minkowski space.

The Euclidean worldsheet coordinates are \((\sigma^1, \sigma^2) = (\sigma^1, i\sigma^0)\) and it will prove useful to form the complex coordinates,

\[ z = \sigma^1 + i\sigma^2 \quad \text{and} \quad \bar{z} = \sigma^1 - i\sigma^2 \]

which are the Euclidean analogue of the lightcone coordinates. Motivated by this analogy, it is common to refer to holomorphic functions as “left-moving” and anti-holomorphic functions as “right-moving”.

The holomorphic derivatives are

\[ \partial_z \equiv \partial = \frac{1}{2}(\partial_1 - i\partial_2) \quad \text{and} \quad \partial_{\bar{z}} \equiv \bar{\partial} = \frac{1}{2}(\partial_1 + i\partial_2) \]

These obey \(\partial z = \bar{\partial} \bar{z} = 1\) and \(\partial \bar{z} = \bar{\partial} z = 0\). We will usually work in flat Euclidean space, with metric

\[ ds^2 = (d\sigma^1)^2 + (d\sigma^2)^2 = dz d\bar{z} \quad (4.2) \]
In components, this flat metric reads

\[ g_{zz} = g_{\bar{z}\bar{z}} = 0 \quad \text{and} \quad g_{z\bar{z}} = \frac{1}{2} \]

With this convention, the measure factor is \( dzd\bar{z} = 2d\sigma^1 d\sigma^2 \). We define the delta-function such that \( \int d^2 \sigma \delta(\sigma) = 1 \). Notice that because we also have \( \int d^2 \sigma \delta(\sigma) = 1 \), this means that there is a factor of 2 difference between the two delta functions. Vectors naturally have their indices up: \( v^z = (v^1 + iv^2) \) and \( v^{\bar{z}} = (v^1 - iv^2) \). When indices are down, the vectors are \( v_z = \frac{1}{2}(v^1 - iv^2) \) and \( v_{\bar{z}} = \frac{1}{2}(v^1 + iv^2) \).

### 4.0.2 The Holomorphy of Conformal Transformations

In the complex Euclidean coordinates \( z \) and \( \bar{z} \), conformal transformations of flat space are simple: they are any holomorphic change of coordinates,

\[ z \to z' = f(z) \quad \text{and} \quad \bar{z} \to \bar{z}' = \bar{f}(\bar{z}) \]

Under this transformation, \( ds^2 = dzd\bar{z} \to |df/dz|^2 dzd\bar{z} \), which indeed takes the form (4.1). Note that we have an infinite number of conformal transformations — in fact, a whole functions worth \( f(z) \). This is special to conformal field theories in two dimensions. In higher dimensions, the space of conformal transformations is a finite dimensional group. For theories defined on \( \mathbb{R}^{p,q} \), the conformal group is \( SO(p+1,q+1) \) when \( p + q > 2 \).

A couple of particularly simple and important examples of 2d conformal transformations are

- \( z \to z + a \): This is a translation.
- \( z \to \zeta z \): This is a rotation for \( |\zeta| = 1 \) and a scale transformation (also known as a *dilatation*) for real \( \zeta \neq 1 \).

For many purposes, it’s simplest to treat \( z \) and \( \bar{z} \) as independent variables. In doing this, we’re really extending the worldsheet from \( \mathbb{R}^2 \) to \( \mathbb{C}^2 \). This will allow us to make use of various theorems from complex methods. However, at the end of the day we should remember that we’re really sitting on the real slice \( \mathbb{R}^2 \subset \mathbb{C}^2 \) defined by \( \bar{z} = z^* \).

### 4.1 Classical Aspects

We start by deriving some properties of classical theories which are invariant under conformal transformations (4.1).
4.1.1 The Stress-Energy Tensor

One of the most important objects in any field theory is the stress-energy tensor (also known as the energy-momentum tensor). This is defined in the usual way as the matrix of conserved currents which arise from translational invariance,

$$\delta \sigma^\alpha = \epsilon^\alpha.$$ 

In flat spacetime, a translation is a special case of a conformal transformation.

There’s a cute way to derive the stress-energy tensor in any theory. Suppose for the moment that we are in flat space $g_{\alpha\beta} = \eta_{\alpha\beta}$. Recall that we can usually derive conserved currents by promoting the constant parameter $\epsilon$ that appears in the symmetry to a function of the spacetime coordinates. The change in the action must then be of the form,

$$\delta S = \int d^2 \sigma \ J^\alpha \partial_\alpha \epsilon \quad (4.3)$$

for some function of the fields, $J^\alpha$. This ensures that the variation of the action vanishes when $\epsilon$ is constant, which is of course the definition of a symmetry. But when the equations of motion are satisfied, we must have $\delta S = 0$ for all variations $\epsilon(\sigma)$, not just constant $\epsilon$. This means that when the equations of motion are obeyed, $J^\alpha$ must satisfy

$$\partial_\alpha J^\alpha = 0$$

The function $J^\alpha$ is our conserved current.

Let’s see how this works for translational invariance. If we promote $\epsilon$ to a function of the worldsheet variables, the change of the action must be of the form (4.3). But what is $J^\alpha$? At this point we do the cute thing. Consider the same theory, but now coupled to a dynamical background metric $g_{\alpha\beta}(\sigma)$. In other words, coupled to gravity. Then we could view the transformation

$$\delta \sigma^\alpha = \epsilon^\alpha(\sigma)$$

as a diffeomorphism and we know that the theory is invariant as long as we make the corresponding change to the metric

$$\delta g_{\alpha\beta} = \partial_\alpha \epsilon_\beta + \partial_\beta \epsilon_\alpha .$$

This means that if we just make the transformation of the coordinates in our original theory, then the change in the action must be the opposite of what we get if we just
transform the metric. (Because doing both together leaves the action invariant). So we have

\[ \delta S = - \int d^2 \sigma \frac{\partial S}{\partial g_{\alpha \beta}} \delta g_{\alpha \beta} = -2 \int d^2 \sigma \frac{\partial S}{\partial g_{\alpha \beta}} \partial \epsilon_{\beta} \]

Note that \( \partial S/\partial g_{\alpha \beta} \) in this expression is really a functional derivatives but we won’t be careful about using notation to indicate this. We now have the conserved current arising from translational invariance. We will add a normalization constant which is standard in string theory (although not necessarily in other areas) and define the stress-energy tensor to be

\[ T_{\alpha \beta} = -\frac{4\pi}{\sqrt{g}} \frac{\partial S}{\partial g_{\alpha \beta}} \]

(4.4)

If we have a flat worldsheet, we evaluate \( T_{\alpha \beta} \) on \( g_{\alpha \beta} = \delta_{\alpha \beta} \) and the resulting expression obeys \( \partial^\alpha T_{\alpha \beta} = 0 \). If we’re working on a curved worldsheet, then the energy-momentum tensor is covariantly conserved, \( \nabla^\alpha T_{\alpha \beta} = 0 \).

**The Stress-Energy Tensor is Traceless**

In conformal theories, \( T_{\alpha \beta} \) has a very important property: its trace vanishes. To see this, let’s vary the action with respect to a scale transformation which is a special case of a conformal transformation,

\[ \delta g_{\alpha \beta} = \epsilon g_{\alpha \beta} \]

(4.5)

Then we have

\[ \delta S = \int d^2 \sigma \frac{\partial S}{\partial g_{\alpha \beta}} \delta g_{\alpha \beta} = -\frac{1}{4\pi} \int d^2 \sigma \sqrt{g} \epsilon T^\alpha_{\alpha} \]

But this must vanish in a conformal theory because scaling transformations are a symmetry. So

\[ T^\alpha_{\alpha} = 0 \]

This is the key feature of a conformal field theory in any dimension. Many theories have this feature at the classical level, including Maxwell theory and Yang-Mills theory in four-dimensions. However, it is much harder to preserve at the quantum level. (The weight of the world rests on the fact that Yang-Mills theory fails to be conformal at the quantum level). Technically the difficulty arises due to the need to introduce a scale when regulating the theories. Here we will be interested in two-dimensional theories
which succeed in preserving the conformal symmetry at the quantum level.

Looking Ahead: Even when the conformal invariance survives in a 2d quantum theory, the vanishing trace $T^\alpha_\alpha = 0$ will only turn out to hold in flat space. We will derive this result in section 4.4.2.

The Stress-Tensor in Complex Coordinates

In complex coordinates, $z = \sigma^1 + i\sigma^2$, the vanishing of the trace $T^\alpha_\alpha = 0$ becomes

$$T_{zz} = 0$$

Meanwhile, the conservation equation $\partial_\alpha T^\alpha_\beta = 0$ becomes $\partial T^{zz} = \bar{\partial} T^{\bar{z}\bar{z}} = 0$. Or, lowering the indices on $T$,

$$\bar{\partial} T_{zz} = 0 \quad \text{and} \quad \partial T_{\bar{z}\bar{z}} = 0$$

In other words, $T_{zz} = T_{zz}(z)$ is a holomorphic function while $T_{\bar{z}\bar{z}} = T_{\bar{z}\bar{z}}(\bar{z})$ is an anti-holomorphic function. We will often use the simplified notation

$$T_{zz}(z) \equiv T(z) \quad \text{and} \quad T_{\bar{z}\bar{z}}(\bar{z}) \equiv \bar{T}(\bar{z})$$

4.1.2 Noether Currents

The stress-energy tensor $T_{\alpha\beta}$ provides the Noether currents for translations. What are the currents associated to the other conformal transformations? Consider the infinitesimal change,

$$z' = z + \epsilon(z) \quad , \quad \bar{z}' = \bar{z} + \bar{\epsilon}(\bar{z})$$

where, making contact with the two examples above, constant $\epsilon$ corresponds to a translation while $\epsilon(z) \sim z$ corresponds to a rotation and dilatation. To compute the current, we’ll use the same trick that we saw before: we promote the parameter $\epsilon$ to depend on the worldsheet coordinates. But it’s already a function of half of the worldsheet coordinates, so this now means $\epsilon(z) \rightarrow \epsilon(z, \bar{z})$. Then we can compute the change in the action, again using the fact that we can make a compensating change in the metric,

$$\delta S = - \int d^2\sigma \frac{\partial S}{\partial g^{\alpha\beta}} \delta g^{\alpha\beta}$$

$$= \frac{1}{2\pi} \int d^2\sigma \ T_{\alpha\beta} (\partial^\alpha \delta \sigma^\beta)$$

$$= \frac{1}{2\pi} \int d^2z \ 1 \ 2 [T_{zz} (\partial^z \delta \bar{z}) + T_{\bar{z}\bar{z}} (\partial^{\bar{z}} \delta \bar{z})]$$

$$= \frac{1}{2\pi} \int d^2z \ [T_{zz} \partial_z \epsilon + T_{\bar{z}\bar{z}} \partial_{\bar{z}} \bar{\epsilon}]$$ (4.6)
Firstly note that if $\epsilon$ is holomorphic and $\bar{\epsilon}$ is anti-holomorphic, then we immediately have $\delta S = 0$. This, of course, is the statement that we have a symmetry on our hands. (You may wonder where in the above derivation we used the fact that the theory was conformal. It lies in the transition to the third line where we needed $T_{z\bar{z}} = 0$).

At this stage, let’s use the trick of treating $z$ and $\bar{z}$ as independent variables. We look at separate currents that come from shifts in $z$ and shifts $\bar{z}$. Let’s first look at the symmetry

$$\delta z = \epsilon(z) \quad , \quad \delta \bar{z} = 0$$

We can read off the conserved current from (4.6) by using the standard trick of letting the small parameter depend on position. Since $\epsilon(z)$ already depends on position, this means promoting $\epsilon \to \epsilon(z)f(\bar{z})$ for some function $f$ and then looking at the $\bar{\partial}f$ terms in (4.6). This gives us the current

$$J^z = 0 \quad \text{and} \quad J^{\bar{z}} = T_{z\bar{z}}(z)\epsilon(z) \equiv T(z)\epsilon(z) \quad (4.7)$$

Importantly, we find that the current itself is also holomorphic. We can check that this is indeed a conserved current: it should satisfy $\partial_{\alpha}J^\alpha = \partial_z J^z + \partial_{\bar{z}}J^{\bar{z}} = 0$. But in fact it does so with room to spare: it satisfies the much stronger condition $\partial_{\bar{z}}J^{\bar{z}} = 0$.

Similarly, we can look at transformations $\delta \bar{z} = \bar{\epsilon}(\bar{z})$ with $\delta z = 0$. We get the anti-holomorphic current $\bar{J}$,

$$\bar{J}^z = \bar{T}(\bar{z})\bar{\epsilon}(\bar{z}) \quad \text{and} \quad \bar{J}^{\bar{z}} = 0 \quad (4.8)$$

### 4.1.3 An Example: The Free Scalar Field

Let’s illustrate some of these ideas about classical conformal theories with the free scalar field,

$$S = \frac{1}{4\pi\alpha'} \int d^2\sigma \, \partial_\alpha X \, \partial^\alpha X$$

Notice that there’s no overall minus sign, in contrast to our earlier action (1.30). That’s because we’re now working with a Euclidean worldsheet metric. The theory of a free scalar field is, of course, dead easy. We can compute anything we like in this theory. Nonetheless, it will still exhibit enough structure to provide an example of all the abstract concepts that we will come across in CFT. For this reason, the free scalar field will prove a good companion throughout this part of the lectures.
Firstly, let’s just check that this free scalar field is actually conformal. In particular, we can look at rescaling \( \sigma^\alpha \rightarrow \lambda \sigma^\alpha \). If we view this in the sense of an active transformation, the coordinates remain fixed but the value of the field at point \( \sigma \) gets moved to point \( \lambda \sigma \). This means,

\[
X(\sigma) \rightarrow X(\lambda^{-1}\sigma) \quad \text{and} \quad \frac{\partial X(\sigma)}{\partial \sigma^\alpha} \rightarrow \frac{\partial X(\lambda^{-1}\sigma)}{\partial \sigma^\alpha} = \frac{1}{\lambda} \frac{\partial X(\tilde{\sigma})}{\partial \tilde{\sigma}}
\]

where we’ve defined \( \tilde{\sigma} = \lambda^{-1}\sigma \). The factor of \( \lambda^{-2} \) coming from the two derivatives in the Lagrangian then cancels the Jacobian factor from the measure \( d^2\sigma = \lambda^2 d^2\tilde{\sigma} \), leaving the action invariant. Note that any polynomial interaction term for \( X \) would break conformal invariance.

The stress-energy tensor for this theory is defined using (4.4),

\[
T_{\alpha\beta} = -\frac{1}{\alpha'} \left( \partial_\alpha X \partial_\beta X - \frac{1}{2} \delta_{\alpha\beta} (\partial X)^2 \right),
\]

which indeed satisfies \( T^\alpha_{\alpha} = 0 \) as it should. The stress-energy tensor looks much simpler in complex coordinates. It is simple to check that \( T_z\bar{z} = 0 \) while

\[
T = -\frac{1}{\alpha'} \partial X \partial X \quad \text{and} \quad \bar{T} = -\frac{1}{\alpha'} \bar{\partial} X \bar{\partial} X
\]

The equation of motion for \( X \) is \( \partial \bar{\partial} X = 0 \). The general classical solution decomposes as,

\[
X(z, \bar{z}) = X(z) + \bar{X}(\bar{z})
\]

When evaluated on this solution, \( T \) and \( \bar{T} \) become holomorphic and anti-holomorphic functions respectively.

### 4.2 Quantum Aspects

So far our discussion has been entirely classical. We now turn to the quantum theory. The first concept that we want to discuss is actually a feature of any quantum field theory. But it really comes into its own in the context of CFT: it is the operator product expansion.

#### 4.2.1 Operator Product Expansion

Let’s first describe what we mean by a local operator in a CFT. We will also refer to these objects as fields. There is a slight difference in terminology between CFTs and more general quantum field theories. Usually in quantum field theory, one reserves the
term “field” for the objects $\phi$ which sit in the action and are integrated over in the path integral. In contrast, in CFT the term “field” refers to any local expression that we can write down. This includes $\phi$, but also includes derivatives $\partial^n \phi$ or composite operators such as $e^{i\phi}$. All of these are thought of as different fields in a CFT. It should be clear from this that the set of all “fields” in a CFT is always infinite even though, if you were used to working with quantum field theory, you would talk about only a finite number of fundamental objects $\phi$. Obviously, this is nothing to be scared about. It’s just a change of language: it doesn’t mean that our theory got harder.

We now define the operator product expansion (OPE). It is a statement about what happens as local operators approach each other. The idea is that two local operators inserted at nearby points can be closely approximated by a string of operators at one of these points. Let’s denote all the local operators of the CFT by $O_i$, where $i$ runs over the set of all operators. Then the OPE is

$$O_i(z, \bar{z}) O_j(w, \bar{w}) = \sum_k C^k_{ij}(z - w, \bar{z} - \bar{w}) O_k(w, \bar{w}) \tag{4.10}$$

Here $C^k_{ij}(z - w, \bar{z} - \bar{w})$ are a set of functions which, on grounds of translational invariance, depend only on the separation between the two operators. We will write a lot of operator equations of the form (4.10) and it’s important to clarify exactly what they mean: they are always to be understood as statements which hold as operator insertions inside time-ordered correlation functions,

$$\langle O_i(z, \bar{z}) O_j(w, \bar{w}) \ldots \rangle = \sum_k C^k_{ij}(z - w, \bar{z} - \bar{w}) \langle O_k(w, \bar{w}) \ldots \rangle$$

where the $\ldots$ can be any other operator insertions that we choose. Obviously it would be tedious to continually write $\langle \ldots \rangle$. So we don’t. But it’s always implicitly there. There are further caveats about the OPE that are worth stressing

- The correlation functions are always assumed to be time-ordered. (Or something similar that we will discuss in Section 4.5.1). This means that as far as the OPE is concerned, everything commutes since the ordering of operators is determined inside the correlation function anyway. So we must have $O_i(z, \bar{z}) O_j(w, \bar{w}) = O_j(w, \bar{w}) O_i(z, \bar{z})$. (There is a caveat here: if the operators are Grassmann objects, then they pick up an extra minus sign when commuted, even inside time-ordered products).
• The other operator insertions in the correlation function (denoted \ldots above) are arbitrary. Except they should be at a distance large compared to $|z - w|$. It turns out — rather remarkably — that in a CFT the OPEs are exact statements and have a radius of convergence equal to the distance to the nearest other insertion. We will return to this in Section 4.6. The radius of convergence is denoted in the figure by the dotted line.

• The OPEs have singular behaviour as $z \rightarrow w$. In fact, this singular behaviour will really be the only thing we care about! It will turn out to contain the same information as commutation relations, as well as telling us how operators transform under symmetries. Indeed, in many equations we will simply write the singular terms in the OPE and denote the non-singular terms as $+ \ldots$.

4.2.2 Ward Identities

The spirit of Noether’s theorem in quantum field theories is captured by operator equations known as Ward Identities. Here we derive the Ward identities associated to conformal invariance. We start by considering a general theory with a symmetry. Later we will restrict to conformal symmetries.

Games with Path Integrals

We’ll take this opportunity to get comfortable with some basic techniques using path integrals. Schematically, the path integral takes the form

$$Z = \int \mathcal{D}\phi \ e^{-S[\phi]}$$

where $\phi$ collectively denote all the fields (in the path integral sense...not the CFT sense!). A symmetry of the quantum theory is such that an infinitesimal transformation $\phi' = \phi + \epsilon \delta \phi$

leaves both the action and the measure invariant,

$$S[\phi'] = S[\phi] \quad \text{and} \quad \mathcal{D}\phi' = \mathcal{D}\phi$$

(In fact, we only really need the combination $\mathcal{D}\phi e^{-S[\phi]}$ to be invariant but this subtlety won’t matter in this course). We use the same trick that we employed earlier in the classical theory and promote $\epsilon \rightarrow \epsilon(\sigma)$. Then, typically, neither the action nor the measure are invariant but, to leading order in $\epsilon$, the change has to be proportional to
∂ε. We have

\[ Z \longrightarrow \int D\phi' \exp(-S[\phi']) \]
\[ = \int D\phi \exp \left( -S[\phi] - \frac{1}{2\pi} \int J^\alpha \partial_\alpha \epsilon \right) \]
\[ = \int D\phi e^{-S[\phi]} \left( 1 - \frac{1}{2\pi} \int J^\alpha \partial_\alpha \epsilon \right) \]

where the factor of \(1/2\pi\) is merely a convention and \(\int\) is shorthand for \(\int d^2\sigma \sqrt{g}\). Notice that the current \(J^\alpha\) may now also have contributions from the measure transformation as well as the action.

Now comes the clever step. Although the integrand has changed, the actual value of the partition function can’t have changed at all. After all, we just redefined a dummy integration variable \(\phi\). So the expression above must be equal to the original \(Z\). Or, in other words,

\[ \int D\phi e^{-S[\phi]} \left( \int J^\alpha \partial_\alpha \epsilon \right) = 0 \]

Moreover, this must hold for all \(\epsilon\). This gives us the quantum version of Noether’s theorem: the vacuum expectation value of the divergence of the current vanishes:

\[ \langle \partial_\alpha J^\alpha \rangle = 0 \]

We can repeat these tricks of this sort to derive some stronger statements. Let’s see what happens when we have other insertions in the path integral. The time-ordered correlation function is given by

\[ \langle O_1(\sigma_1) \ldots O_n(\sigma_n) \rangle = \frac{1}{Z} \int D\phi e^{-S[\phi]} O_1(\sigma_1) \ldots O_n(\sigma_n) \]

We can think of these as operators inserted at particular points on the plane as shown in the figure. As we described above, the operators \(O_i\) are any general expressions that we can form from the \(\phi\) fields. Under the symmetry of interest, the operator will change in some way, say

\[ O_i \rightarrow O_i + \epsilon \delta O_i \]

We once again promote \(\epsilon \rightarrow \epsilon(\sigma)\). As our first pass, let’s pick a choice of \(\epsilon(\sigma)\) which only has support away from the operator insertions as shown in the Figure 20. Then,

\[ \delta O_i(\sigma_i) = 0 \]
and the above derivation goes through in exactly the same way to give

$$\langle \partial_\alpha J^\alpha(\sigma) O_1(\sigma_1) \ldots O_n(\sigma_n) \rangle = 0 \quad \text{for } \sigma \neq \sigma_i$$

Because this holds for any operator insertions away from \(\sigma\), from the discussion in Section 4.2.1 we are entitled to write the operator equation

$$\partial_\alpha J^\alpha = 0$$

But what if there are operator insertions that lie at the same point as \(J^\alpha\)? In other words, what happens as \(\sigma\) approaches one of the insertion points? The resulting formulae are called Ward identities. To derive these, let’s take \(\epsilon(\sigma)\) to have support in some region that includes the point \(\sigma_1\), but not the other points as shown in Figure 21. The simplest choice is just to take \(\epsilon(\sigma)\) to be constant inside the shaded region and zero outside. Now using the same procedure as before, we find that the original correlation function is equal to,

$$\frac{1}{Z} \int D\phi e^{-S[\phi]} \left( 1 - \frac{1}{2\pi} \int J^\alpha \partial_\alpha \epsilon \right) (O_1 + \epsilon \delta O_1) O_2 \ldots O_n$$

Working to leading order in \(\epsilon\), this gives

$$-\frac{1}{2\pi} \int \epsilon \partial_\alpha \langle J^\alpha(\sigma) O_1(\sigma_1) \ldots \rangle = \langle \delta O_1(\sigma_1) \ldots \rangle$$

where the integral on the left-hand-side is only over the region of non-zero \(\epsilon\). This is the Ward Identity.

**Ward Identities for Conformal Transformations**

Ward identities (4.11) hold for any symmetries. Let’s now see what they give when applied to conformal transformations. There are two further steps needed in the derivation. The first simply comes from the fact that we’re working in two dimensions and we can use Stokes’ theorem to convert the integral on the left-hand-side of (4.11) to a line integral around the boundary. Let \(\hat{n}^\alpha\) be the unit vector normal to the boundary. For any vector \(J^\alpha\), we have

$$\int \epsilon \partial_\alpha J^\alpha = \oint \partial_\epsilon (J_1 d\sigma^2 - J_2 d\sigma^1) = -i \oint \partial_\epsilon (J_z dz - J_{\bar{z}} d\bar{z})$$

Figure 20:

Figure 21:
where we have written the expression both in Cartesian coordinates $\sigma^\alpha$ and complex coordinates on the plane. As described in Section 4.0.1, the complex components of the vector with indices down are defined as $J_z = \frac{1}{2}(J_1 - iJ_2)$ and $J_{\bar{z}} = \frac{1}{2}(J_1 + iJ_2)$. So, applying this to the Ward identity (4.11), we find for two dimensional theories

$$\frac{i}{2\pi} \oint_{\partial \epsilon} dz \langle J_z(z, \bar{z}) O_1(\sigma_1) \ldots \rangle = - \frac{i}{2\pi} \oint_{\partial \epsilon} d\bar{z} \langle J_{\bar{z}}(z, \bar{z}) O_1(\sigma_1) \ldots \rangle = \langle \delta O_1(\sigma_1) \ldots \rangle$$

So far our derivation holds for any conserved current $J$ in two dimensions. At this stage we specialize to the currents that arise from conformal transformations (4.7) and (4.8). Here something nice happens because $J_z$ is holomorphic while $J_{\bar{z}}$ is anti-holomorphic. This means that the contour integral simply picks up the residue,

$$\frac{i}{2\pi} \oint_{\partial \epsilon} dz \langle J_z(z) O_1(\sigma_1) \rangle = - \text{Res} [J_z O_1]$$

where this means the residue in the OPE between the two operators,

$$J_z(z) O_1(w, \bar{w}) = \ldots + \frac{\text{Res}[J_z O_1(w, \bar{w})]}{z - w} + \ldots$$

So we find a rather nice way of writing the Ward identities for conformal transformations. If we again view $z$ and $\bar{z}$ as independent variables, the Ward identities split into two pieces. From the change $\delta z = \epsilon(z)$, we get

$$\delta O_1(\sigma_1) = - \text{Res} [J_z(z) O_1(\sigma_1)] = - \text{Res} [\epsilon(z) T(z) O_1(\sigma_1)]$$

where, in the second equality, we have used the expression for the conformal current (4.7). Meanwhile, from the change $\delta \bar{z} = \bar{\epsilon}(\bar{z})$, we have

$$\delta O_1(\sigma_1) = - \text{Res} [\bar{J}_{\bar{z}}(z) O_1(\sigma_1)] = - \text{Res} [\bar{\epsilon}(\bar{z}) \bar{T}(\bar{z}) O_1(\sigma_1)]$$

where the minus sign comes from the fact that the $\oint d\bar{z}$ boundary integral is taken in the opposite direction.

This result means that if we know the OPE between an operator and the stress-tensors $T(z)$ and $\bar{T}(\bar{z})$, then we immediately know how the operator transforms under conformal symmetry. Or, standing this on its head, if we know how an operator transforms then we know at least some part of its OPE with $T$ and $\bar{T}$.

### 4.2.3 Primary Operators

The Ward identity allows us to start piecing together some OPEs by looking at how operators transform under conformal symmetries. Although we don’t yet know the
action of general conformal symmetries, we can start to make progress by looking at the two simplest examples.

**Translations:** If $\delta z = \epsilon$, a constant, then all operators transform as

$$\mathcal{O}(z - \epsilon) = \mathcal{O}(z) - \epsilon \partial \mathcal{O}(z) + \ldots$$

The Noether current for translations is the stress-energy tensor $T$. The Ward identity in the form (4.12) tells us that the OPE of $T$ with any operator $\mathcal{O}$ must be of the form,

$$T(z) \mathcal{O}(w, \bar{w}) = \ldots + \frac{\partial \mathcal{O}(w, \bar{w})}{z-w} + \ldots$$  (4.13)

Similarly, the OPE with $\bar{T}$ is

$$\bar{T}(\bar{z}) \mathcal{O}(w, \bar{w}) = \ldots + \frac{\bar{\partial} \mathcal{O}(w, \bar{w})}{\bar{z}-\bar{w}} + \ldots$$  (4.14)

**Rotations and Scaling:** The transformation

$$z \rightarrow z + \epsilon z \quad \text{and} \quad \bar{z} \rightarrow \bar{z} + \bar{\epsilon} \bar{z}$$  (4.15)

describes rotation for $\epsilon$ purely imaginary and scaling (dilatation) for $\epsilon$ real. Not all operators have good transformation properties under these actions. This is entirely analogous to the statement in quantum mechanics that not all states transform nicely under the Hamiltonian $H$ and angular momentum operator $L$. However, in quantum mechanics we know that the eigenstates of $H$ and $L$ can be chosen as a basis of the Hilbert space provided, of course, that $[H, L] = 0$.

The same statement holds for operators in a CFT: we can choose a basis of local operators that have good transformation properties under rotations and dilatations. In fact, we will see in Section 4.6 that the statement about local operators actually follows from the statement about states.

**Definition:** An operator $\mathcal{O}$ is said to have weight $(h, \tilde{h})$ if, under $\delta z = \epsilon z$ and $\delta \bar{z} = \bar{\epsilon} \bar{z}$, $\mathcal{O}$ transforms as

$$\delta \mathcal{O} = -\epsilon (h \mathcal{O} + z \partial \mathcal{O}) - \bar{\epsilon} (\tilde{h} \mathcal{O} + \bar{z} \bar{\partial} \mathcal{O})$$  (4.16)

The terms $\partial \mathcal{O}$ in this expression would be there for any operator. They simply come from expanding $\mathcal{O}(z - \epsilon z, \bar{z} - \bar{\epsilon} \bar{z})$. The terms $h \mathcal{O}$ and $\tilde{h} \mathcal{O}$ are special to operators which are eigenstates of dilatations and rotations. Some comments:
• Both \( h \) and \( \tilde{h} \) are real numbers. In a unitary CFT, all operators have \( h, \tilde{h} \geq 0 \). We will prove this is Section 4.5.4.

• The weights are not as unfamiliar as they appear. They simply tell us how operators transform under rotations and scalings. But we already have names for these concepts from undergraduate days. The eigenvalue under rotation is usually called the \textit{spin}, \( s \), and is given in terms of the weights as

\[
s = h - \tilde{h}
\]

Meanwhile, the \textit{scaling dimension} \( \Delta \) of an operator is

\[
\Delta = h + \tilde{h}
\]

• To motivate these definitions, it’s worth recalling how rotations and scale transformations act on the underlying coordinates. Rotations are implemented by the operator

\[
L = -i(\sigma^1 \partial_2 - \sigma^2 \partial_1) = z\partial - \bar{z}\bar{\partial}
\]

while the dilation operator \( D \) which gives rise to scalings is

\[
D = \sigma^\alpha \partial_\alpha = z\partial + \bar{z}\bar{\partial}
\]

• The scaling dimension is nothing more than the familiar “dimension” that we usually associate to fields and operators by dimensional analysis. For example, worldsheet derivatives always increase the dimension of an operator by one: \( \Delta[\partial] = +1 \). The tricky part is that the naive dimension that fields have in the classical theory is not necessarily the same as the dimension in the quantum theory.

Let’s compare the transformation law (4.16) with the Ward identity (4.12). The Noether current arising from rotations and scaling \( \delta z = \epsilon z \) was given in (4.7): it is \( J(z) = zT(z) \). This means that the residue of the \( J\mathcal{O} \) OPE will determine the \( 1/z^2 \) term in the \( T\mathcal{O} \) OPE. Similar arguments hold, of course, for \( \delta\bar{z} = \bar{\epsilon}\bar{z} \) and \( \bar{T} \). So, the upshot of this is that, for an operator \( \mathcal{O} \) with weight \((h, \tilde{h})\), the OPE with \( T \) and \( \bar{T} \) takes the form

\[
T(z) \mathcal{O}(w, \bar{w}) = \ldots + h \frac{\mathcal{O}(w, \bar{w})}{(z - w)^2} + \frac{\partial \mathcal{O}(w, \bar{w})}{z - w} + \ldots
\]

\[
\bar{T}(\bar{z}) \mathcal{O}(w, \bar{w}) = \ldots + \tilde{h} \frac{\mathcal{O}(w, \bar{w})}{(\bar{z} - \bar{w})^2} + \frac{\bar{\partial} \mathcal{O}(w, \bar{w})}{\bar{z} - \bar{w}} + \ldots
\]
Primary Operators

A primary operator is one whose OPE with $T$ and $\bar{T}$ truncates at order $(z-w)^{-2}$ or order $(\bar{z}-\bar{w})^{-2}$ respectively. There are no higher singularities:

$$T(z) \mathcal{O}(w, \bar{w}) = h \frac{\mathcal{O}(w, \bar{w})}{(z-w)^2} + \frac{\partial \mathcal{O}(w, \bar{w})}{z-w} + \text{non-singular}$$

$$\bar{T}(\bar{z}) \mathcal{O}(w, \bar{w}) = \bar{h} \frac{\mathcal{O}(w, \bar{w})}{(\bar{z}-\bar{w})^2} + \frac{\bar{\partial} \mathcal{O}(w, \bar{w})}{\bar{z}-\bar{w}} + \text{non-singular}$$

Since we now know all singularities in the $T\mathcal{O}$ OPE, we can reconstruct the transformation under all conformal transformations. The importance of primary operators is that they have particularly simple transformation properties. Focussing on $\delta z = \epsilon(z)$, we have

$$\delta \mathcal{O}(w, \bar{w}) = -\text{Res} \left[ \epsilon(z) T(z) \mathcal{O}(w, \bar{w}) \right]$$

$$= -\text{Res} \left[ \epsilon(z) \left( h \frac{\mathcal{O}(w, \bar{w})}{(z-w)^2} + \frac{\partial \mathcal{O}(w, \bar{w})}{z-w} + \ldots \right) \right]$$

We want to look at smooth conformal transformations and so require that $\epsilon(z)$ itself has no singularities at $z = w$. We can then Taylor expand

$$\epsilon(z) = \epsilon(w) + \epsilon'(w) (z-w) + \ldots$$

We learn that the infinitesimal change of a primary operator under a general conformal transformation $\delta z = \epsilon(z)$ is

$$\delta \mathcal{O}(w, \bar{w}) = -h \epsilon'(w) \mathcal{O}(w, \bar{w}) - \epsilon(w) \partial \mathcal{O}(w, \bar{w})$$

(4.17)

There is a similar expression for the anti-holomorphic transformations $\delta \bar{z} = \bar{\epsilon}(\bar{z})$.

Equation (4.17) holds for infinitesimal conformal transformations. It is a simple matter to integrate up to find how primary operators change under a finite conformal transformation,

$$z \rightarrow \tilde{z}(z) \quad \text{and} \quad \bar{z} \rightarrow \tilde{\bar{z}}(\bar{z})$$

The general transformation of a primary operator is given by

$$\mathcal{O}(z, \bar{z}) \rightarrow \tilde{\mathcal{O}}(\tilde{z}, \tilde{\bar{z}}) = \left( \frac{\partial \tilde{z}}{\partial z} \right)^{-h} \left( \frac{\partial \tilde{\bar{z}}}{\partial \bar{z}} \right)^{-\tilde{h}} \mathcal{O}(z, \bar{z})$$

(4.18)

It will turn out that one of the main objects of interest in a CFT is the spectrum of weights $(h, \tilde{h})$ of primary fields. This will be equivalent to computing the particle mass spectrum in a quantum field theory. In the context of statistical mechanics, the weights of primary operators are the critical exponents.
4.3 An Example: The Free Scalar Field

Let’s look at how all of this works for the free scalar field. We’ll start by familiarizing ourselves with some techniques using the path integral. The action is,

\[ S = \frac{1}{4\pi\alpha'} \int d^2\sigma \, \partial_\alpha X \partial^\alpha X \]  

(4.19)

The classical equation of motion is \( \partial^2 X = 0 \). Let’s start by seeing how to derive the analogous statement in the quantum theory using the path integral. The key fact that we’ll need is that the integral of a total derivative vanishes in the path integral just as it does in an ordinary integral. From this we have,

\[
0 = \int D X \, \frac{\delta}{\delta X(\sigma)} e^{-S} = \int D X \, e^{-S} \left[ \frac{1}{2\pi\alpha'} \partial^2 X(\sigma) \right]
\]

But this is nothing more than the Ehrenfest theorem which states that expectation values of operators obey the classical equations of motion,

\[
\langle \partial^2 X(\sigma) \rangle = 0
\]

4.3.1 The Propagator

The next thing that we want to do is compute the propagator for \( X \). We could do this using canonical quantization, but it will be useful to again see how it works using the path integral. This time we look at,

\[
0 = \int D X \, \frac{\delta}{\delta X(\sigma)} \left[ e^{-S} X(\sigma') \right] = \int D X \, e^{-S} \left[ \frac{1}{2\pi\alpha'} \partial^2 X(\sigma) X(\sigma') + \delta(\sigma - \sigma') \right]
\]

So this time we learn that

\[
\langle \partial^2 X(\sigma) X(\sigma') \rangle = -2\pi\alpha' \delta(\sigma - \sigma')
\]

(4.20)

Note that if we’d computed this in the canonical approach, we would have found the same answer: the \( \delta \)-function arises in this calculation because all correlation functions are time-ordered.

We can now treat (4.20) as a differential equation for the propagator \( \langle X(\sigma)X(\sigma') \rangle \). To solve this equation, we need the following standard result

\[
\partial^2 \ln(\sigma - \sigma')^2 = 4\pi\delta(\sigma - \sigma')
\]

(4.21)
Since this is important, let’s just quickly check that it’s true. It’s a simple application of Stokes’ theorem. Set $\sigma' = 0$ and integrate over $\int d^2\sigma$. We obviously get $4\pi$ from the right-hand-side. The left-hand-side gives

$$\int d^2\sigma \partial^2 \ln(\sigma_1^2 + \sigma_2^2) = \int d^2\sigma \partial^2 \left( \frac{2\sigma_1}{\sigma_1^2 + \sigma_2^2} \right) = 2 \oint \frac{(\sigma_1 d\sigma_2 - \sigma_2 d\sigma_1)}{\sigma_1^2 + \sigma_2^2}$$

Switching to polar coordinates $\sigma_1 + i\sigma_2 = re^{i\theta}$, we can rewrite this expression as

$$2 \int \frac{r^2 d\theta}{r^2} = 4\pi$$

confirming (4.21). Applying this result to our equation (4.20), we get the propagator of a free scalar in two-dimensions,

$$\langle X(\sigma)X(\sigma') \rangle = -\frac{\alpha'}{2} \ln(\sigma - \sigma')^2$$

The propagator has a singularity as $\sigma \to \sigma'$. This is an ultra-violet divergence and is common to all field theories. It also has a singularity as $|\sigma - \sigma'| \to \infty$. This is telling us something important that we will mention in Section 4.3.2.

Finally, we could repeat our trick of looking at total derivatives in the path integral, now with other operator insertions $O_1(\sigma_1), \ldots O_n(\sigma_n)$ in the path integral. As long as $\sigma, \sigma' \neq \sigma_i$, then the whole analysis goes through as before. But this is exactly our criterion to write the operator product equation,

$$X(\sigma)X(\sigma') = -\frac{\alpha'}{2} \ln(\sigma - \sigma')^2 + \ldots \quad (4.22)$$

We can also write this in complex coordinates. The classical equation of motion $\partial \bar{\partial}X = 0$ allows us to split the operator $X$ into left-moving and right-moving pieces,

$$X(z, \bar{z}) = X(z) + \bar{X}(\bar{z})$$

We’ll focus just on the left-moving piece. This has the operator product expansion,

$$X(z)X(w) = -\frac{\alpha'}{2} \ln(z - w) + \ldots$$

The logarithm means that $X(z)$ doesn’t have any nice properties under the conformal transformations. For this reason, the “fundamental field” $X$ is not really the object of interest in this theory! However, we can look at the derivative of $X$. This has a rather nice looking OPE,

$$\partial X(z) \partial X(w) = -\frac{\alpha'}{2} \frac{1}{(z - w)^2} + \text{non-singular} \quad (4.23)$$
4.3.2 An Aside: No Goldstone Bosons in Two Dimensions

The infra-red divergence in the propagator has an important physical implication. Let’s start by pointing out one of the big differences between quantum mechanics and quantum field theory in $d = 3 + 1$ dimensions. Since the language used to describe these two theories is rather different, you may not even be aware that this difference exists.

Consider the quantum mechanics of a particle on a line. This is a $d = 0 + 1$ dimensional theory of a free scalar field $X$. Let’s prepare the particle in some localized state—say a Gaussian wavefunction $\Psi(X) \sim \exp(-X^2/L^2)$. What then happens? The wavefunction starts to spread out. And the spreading doesn’t stop. In fact, the would-be ground state of the system is a uniform wavefunction of infinite width, which isn’t a state in the Hilbert space because it is non-normalizable.

Let’s now compare this to the situation of a free scalar field $X$ in a $d = 3 + 1$ dimensional field theory. Now we think of this as a scalar without potential. The physics is very different: the theory has an infinite number of ground states, determined by the expectation value $\langle X \rangle$. Small fluctuations around this vacuum are massless: they are Goldstone bosons for broken translational invariance $X \to X + c$.

We see that the physics is very different in field theories in $d = 0 + 1$ and $d = 3 + 1$ dimensions. The wavefunction spreads along flat directions in quantum mechanics, but not in higher dimensional field theories. But what happens in $d = 1 + 1$ and $d = 2 + 1$ dimensions? It turns out that field theories in $d = 1 + 1$ dimensions are more like quantum mechanics: the wavefunction spreads. Theories in $d = 2 + 1$ dimensions and higher exhibit the opposite behaviour: they have Goldstone bosons. The place to see this is the propagator. In $d$ spacetime dimensions, it takes the form

$$\langle X(r) X(0) \rangle \sim \begin{cases} 1/r^{d-2} & d \neq 2 \\ \ln r & d = 2 \end{cases}$$

which diverges at large $r$ only for $d = 1$ and $d = 2$. If we perturb the vacuum slightly by inserting the operator $X(0)$, this correlation function tells us how this perturbation falls off with distance. The infra-red divergence in low dimensions is telling us that the wavefunction wants to spread.

The spreading of the wavefunction in low dimensions means that there is no spontaneous symmetry breaking and no Goldstone bosons. It is usually referred to as the Coleman-Mermin-Wagner theorem. Note, however, that it certainly doesn’t prohibit massless excitations in two dimensions: it only prohibits Goldstone-like massless excitations.
4.3.3 The Stress-Energy Tensor and Primary Operators

We want to compute the OPE of $T$ with other operators. Firstly, what is $T$? We computed it in the classical theory in (4.9). It is,

$$T = -\frac{1}{\alpha'} \partial X \partial X$$

(4.24)

But we need to be careful about what this means in the quantum theory. It involves the product of two operators defined at the same point and this is bound to mean divergences if we just treat it naively. In canonical quantization, we would be tempted to normal order by putting all annihilation operators to the right. This guarantees that the vacuum has zero energy. Here we do something that is basically equivalent, but without reference to creation and annihilation operators. We write

$$T = -\frac{1}{\alpha'} : \partial X \partial X : \equiv -\frac{1}{\alpha'} \lim_{z \to w} (\partial X(z) \partial X(w) - \langle \partial X(z) \partial X(w) \rangle)$$

(4.25)

which, by construction, has $\langle T \rangle = 0$.

With this definition of $T$, let’s start to compute the OPEs to determine the primary fields in the theory.

Claim 1: $\partial X$ is a primary field with weight $h = 1$ and $\tilde{h} = 0$.

Proof: We need to figure out how to take products of normal ordered operators

$$T(z) \partial X(w) = -\frac{1}{\alpha'} : \partial X(z) \partial X(z) : \partial X(w)$$

The operators on the left-hand side are time-ordered (because all operator expressions of this type are taken to live inside time-ordered correlation functions). In contrast, the right-hand side is a product of normal-ordered operators. But we know how to change normal ordered products into time ordered products: this is the content of Wick’s theorem. Although we have defined normal ordering in (4.25) without reference to creation and annihilation operators, Wick’s theorem still holds. We must sum over all possible contractions of pairs of operators, where the term “contraction” means that we replace the pair by the propagator,

$$\langle \partial X(z) \partial X(w) \rangle = -\frac{\alpha'}{2} \frac{1}{(z - w)^2}$$

Using this, we have

$$T(z) \partial X(w) = -\frac{2}{\alpha'} \partial X(z) \left( -\frac{\alpha'}{2} \frac{1}{(z - w)^2} + \text{non-singular} \right)$$
Here the “non-singular” piece includes the totally normal ordered term: $T(z) \partial X(w)$. It is only the singular part that interests us. Continuing, we have

$$T(z) \partial X(w) = \frac{\partial X(z)}{(z-w)^2} + \ldots = \frac{\partial X(w)}{(z-w)^2} + \frac{\partial^2 X(w)}{z-w} + \ldots$$

This is indeed the OPE for a primary operator of weight $h = 1$. □

Note that higher derivatives $\partial^n X$ are not primary for $n > 1$. For example, $\partial^2 X$ has weight $(h, \tilde{h}) = (2, 0)$, but is not a primary operator, as we see from the OPE,

$$T(z) \partial^2 X(w) = \partial_w \left[ \frac{\partial X(w)}{(z-w)^2} + \ldots \right] = \frac{2\partial X(w)}{(z-w)^3} + \frac{2\partial^2 X(w)}{(z-w)^2} + \ldots$$

The fact that the field $\partial^n X$ has weight $(h, \tilde{h}) = (n, 0)$ fits our natural intuition: each derivative provides spin $s = 1$ and dimension $\Delta = 1$, while the field $X$ does not appear to be contributing, presumably reflecting the fact that it has naive, classical dimension zero. However, in the quantum theory, it is not correct to say that $X$ has vanishing dimension: it has an ill-defined dimension due to the logarithmic behaviour of its OPE (4.22). This is responsible for the following, more surprising, result

Claim 2: The field $e^{ikX}$ is primary with weight $h = \tilde{h} = \alpha'k^2/4$.

This result is not what we would guess from the classical theory\footnote{We could, however, guess it with a little knowledge of renormalisation. Indeed, we previously derived this result in the lectures on Statistical Field Theory where we computed RG flows in the Sine-Gordon model; see Section 4.4.3 of those lectures.}. Indeed, it’s obvious that it has a quantum origin because the weight is proportional to $\alpha'$, which sits outside the action in the same place that $\hbar$ would (if we hadn’t set it to one). Note also that this means that the spectrum of the free scalar field is continuous. This is related to the fact that the range of $X$ is non-compact. Generally, CFTs will have a discrete spectrum.

Proof: Let’s first compute the OPE with $\partial X$. We have

$$\partial X(z) : e^{ikX(w)} : = \sum_{n=0}^{\infty} \frac{(ik)^n}{n!} \partial X(z) : X(w)^n : = \sum_{n=1}^{\infty} \frac{(ik)^n}{(n-1)!} : X(w)^{n-1} : \left( -\frac{\alpha'}{2} \frac{1}{z-w} \right) + \ldots = -\frac{i\alpha'k}{2} \frac{e^{ikX(w)}}{z-w} + \ldots (4.26)$$
From this, we can compute the OPE with $T$.

\[
T(z) : e^{ikX(w)} : = \frac{1}{\alpha'} : \partial X(z) \partial X(z) : : e^{ikX(w)} :
\]
\[
= \frac{\alpha' k^2}{4} : e^{ikX(w)} : + \frac{ik}{z - w} : \partial X(z) e^{ikX(w)} : + \ldots
\]

where the first term comes from two contractions, while the second term comes from a single contraction. Replacing $\partial_z$ by $\partial_w$ in the final term we get

\[
T(z) : e^{ikX(w)} : = \frac{\alpha' k^2}{4} : e^{ikX(w)} : + \frac{\partial_w : e^{ikX(w)} :}{z - w} + \ldots
\]

(4.27)

showing that $e^{ikX(w)} :$ is indeed primary. We will encounter this operator frequently later, but will choose to simplify notation and drop the normal ordering colons. Normal ordering will just be assumed from now on. \square.

Finally, let's check to see the OPE of $T$ with itself. This is again just an exercise in Wick contractions.

\[
T(z) T(w) = \frac{1}{\alpha'^2} : \partial X(z) \partial X(z) : : \partial X(w) \partial X(w) :
\]
\[
= \frac{2}{\alpha'^2} \left( -\frac{\alpha'}{2} \frac{1}{(z - w)^2} \right)^2 - \frac{4}{\alpha'^2} \frac{\alpha' : \partial X(z) \partial X(w) :}{(z - w)^2} + \ldots
\]

The factor of 2 in front of the first term comes from the two ways of performing two contractions; the factor of 4 in the second term comes from the number of ways of performing a single contraction. Continuing,

\[
T(z) T(w) = \frac{1/2}{(z - w)^4} + \frac{2T(w)}{(z - w)^2} - \frac{2}{\alpha'} \frac{\partial^2 X(w) \partial X(w)}{z - w} + \ldots
\]
\[
= \frac{1/2}{(z - w)^4} + \frac{2T(w)}{(z - w)^2} + \frac{\partial T(w)}{z - w} + \ldots
\]

(4.28)

We learn that $T$ is not a primary operator in the theory of a single free scalar field. It is an operator of weight $(h, \tilde{h}) = (2, 0)$, but it fails the primary test on account of the $(z - w)^{-4}$ term. In fact, this property of the stress energy tensor a general feature of all CFTs which we now explore in more detail.

### 4.4 The Central Charge

In any CFT, the most prominent example of an operator which is not primary is the stress-energy tensor itself.
For the free scalar field, we have already seen that $T$ is an operator of weight $(h, \tilde{h}) = (2, 0)$. This remains true in any CFT. The reason for this is simple: $T_{\alpha\beta}$ has dimension $\Delta = 2$ because we obtain the energy by integrating over space. It has spin $s = 2$ because it is a symmetric 2-tensor. But these two pieces of information are equivalent to the statement that $T$ is has weight $(2, 0)$. Similarly, $\bar{T}$ has weight $(0, 2)$. This means that the $TT$ OPE takes the form,

$$T(z)T(w) = \ldots + \frac{2T(w)}{(z-w)^2} + \frac{\partial T(w)}{z-w} + \ldots$$

and similar for $\bar{T}\bar{T}$. What other terms could we have in this expansion? Since each term has dimension $\Delta = 4$, any operators that appear on the right-hand-side must be of the form

$$\frac{\mathcal{O}_n}{(z-w)^n}$$

where $\Delta[\mathcal{O}_n] = 4 - n$. But, in a unitary CFT there are no operators with $h, \tilde{h} < 0$. (We will prove this shortly). So the most singular term that we can have is of order $(z-w)^{-4}$. Such a term must be multiplied by a constant. We write,

$$T(z)T(w) = \frac{c/2}{(z-w)^4} + \frac{2T(w)}{(z-w)^2} + \frac{\partial T(w)}{z-w} + \ldots$$

and, similarly,

$$\bar{T}(\bar{z})\bar{T}(\bar{w}) = \frac{\tilde{c}/2}{(\bar{z}-\bar{w})^4} + \frac{2\bar{T}(\bar{w})}{(\bar{z}-\bar{w})^2} + \frac{\bar{\partial}\bar{T}(\bar{w})}{\bar{z}-\bar{w}} + \ldots$$

The constants $c$ and $\tilde{c}$ are called the central charges. (Sometimes they are referred to as left-moving and right-moving central charges). They are perhaps the most important numbers characterizing the CFT. We can already get some intuition for the information contained in these two numbers. Looking back at the free scalar field (4.28) we see that it has $c = \tilde{c} = 1$. If we instead considered $D$ non-interacting free scalar fields, we would get $c = \tilde{c} = D$. This gives us a hint: $c$ and $\tilde{c}$ are somehow measuring the number of degrees of freedom in the CFT. This is true in a deep sense! However, be warned: $c$ is not necessarily an integer.

Before moving on, it’s worth pausing to explain why we didn’t include a $(z-w)^{-3}$ term in the $TT$ OPE. The reason is that the OPE must obey $T(z)T(w) = T(w)T(z)$ because, as explained previously, these operator equations are all taken to hold inside time-ordered correlation functions. So the quick answer is that a $(z-w)^{-3}$ term would
not be invariant under \( z \leftrightarrow w \). However, you may wonder how the \((z - w)^{-1}\) term manages to satisfy this property. Let’s see how this works:

\[
T(w) T(z) = \frac{c/2}{(z-w)^4} + \frac{2T(z)}{(z-w)^2} + \frac{\partial T(z)}{w-z} + \ldots
\]

Now we can Taylor expand \( T(z) = T(w) + (z-w)\partial T(w) + \ldots \) and \( \partial T(z) = \partial T(w) + \ldots \). Using this in the above expression, we find

\[
T(w) T(z) = \frac{c/2}{(z-w)^4} + \frac{2T(w)}{(z-w)^2} - \frac{\partial T(w)}{z-w} + \ldots = T(z) T(w)
\]

This trick of Taylor expanding saves the \((z - w)^{-1}\) term. It wouldn’t work for the \((z - w)^{-3}\) term.

**The Transformation of Energy**

So \( T \) is not primary unless \( c = 0 \). And we will see shortly that all theories have \( c > 0 \). What does this mean for the transformation of \( T \)?

\[
\delta T(w) = -\text{Res} [\epsilon(z) T(z) T(w)] = -\text{Res} \left[ \epsilon(z) \left( \frac{c/2}{(z-w)^4} + \frac{2T(w)}{(z-w)^2} + \frac{\partial T(w)}{z-w} + \ldots \right) \right]
\]

If \( \epsilon(z) \) contains no singular terms, we can expand

\[
\epsilon(z) = \epsilon(w) + \epsilon'(w)(z-w) + \frac{1}{2} \epsilon''(z-w)^2 + \frac{1}{6} \epsilon'''(w)(z-w)^3 + \ldots
\]

from which we find

\[
\delta T(w) = -\epsilon(w) \partial T(w) - 2\epsilon'(w) T(w) - \frac{c}{12} \epsilon'''(w) \tag{4.30}
\]

This is the infinitesimal version. We would like to know what becomes of \( T \) under the finite conformal transformation \( z \rightarrow \tilde{z}(z) \). The answer turns out to be

\[
\tilde{T}(\tilde{z}) = \left( \frac{\partial \tilde{z}}{\partial z} \right)^{-2} \left[ T(z) - \frac{c}{12} S(\tilde{z}, z) \right] \tag{4.31}
\]

where \( S(\tilde{z}, z) \) is known as the **Schwarzian** and is defined by

\[
S(\tilde{z}, z) = \left( \frac{\partial^3 \tilde{z}}{\partial z^3} \right) \left( \frac{\partial \tilde{z}}{\partial z} \right)^{-1} - \frac{3}{2} \left( \frac{\partial^2 \tilde{z}}{\partial z^2} \right)^2 \left( \frac{\partial \tilde{z}}{\partial z} \right)^{-2} \tag{4.32}
\]

It is simple to check that the Schwarzian has the right infinitesimal form to give \((4.30)\). Its key property is that it preserves the group structure of successive conformal transformations.
4.4.1 \( c \) is for Casimir

Note that the extra term in the transformation (4.31) of \( T \) does not depend on \( T \) itself. In particular, it will be the same evaluated on all states. It only affects the constant term — or zero mode — in the energy. In other words, it is the Casimir energy of the system.

Let’s look at an example that will prove to be useful later for the string. Consider the Euclidean cylinder, parameterized by

\[
    w = \sigma + i\tau, \quad \sigma \in [0, 2\pi)
\]

We can make a conformal transformation from the cylinder to the complex plane by

\[
    z = e^{-iw}
\]

The fact that the cylinder and the plane are related by a conformal map means that if we understand a given CFT on the cylinder, then we immediately understand it on the plane. And vice-versa. Notice that constant time slices on the cylinder are mapped to circles of constant radius. The origin, \( z = 0 \), is the distant past, \( \tau \to -\infty \).

What becomes of \( T \) under this transformation? The Schwarzian can be easily calculated to be \( S(z, w) = 1/2 \). So we find,

\[
    T_{\text{cylinder}}(w) = -z^2 T_{\text{plane}}(z) + \frac{c}{24} \tag{4.33}
\]

Suppose that the ground state energy vanishes when the theory is defined on the plane: \( \langle T_{\text{plane}} \rangle = 0 \). What happens on the cylinder? We want to look at the Hamiltonian, which is defined by

\[
    H \equiv \int d\sigma \, T_{\tau\tau} = -\int d\sigma \, (T_{ww} + \bar{T}_{\bar{w}\bar{w}})
\]

The conformal transformation then tells us that the ground state energy on the cylinder is

\[
    E = -\frac{2\pi(c + \tilde{c})}{24} \tag{4.34}
\]

This is indeed the (negative) Casimir energy on a cylinder. For a free scalar field, we have \( c = \tilde{c} = 1 \) and the energy density \( E/2\pi = -1/12 \). This is the same result that we got in Section 2.2.2, but this time with no funny business where we throw out infinities.
An Application: The Lüscher Term

If we’re looking at a physical system, the cylinder will have a radius \( L \). In this case, the Casimir energy is given by \( E = \frac{-2\pi (c + \tilde{c})}{24L} \). There is an application of this to QCD-like theories. Consider two quarks in a confining theory, separated by a distance \( L \). If the tension of the confining flux tube is \( T \), then the string will be stable as long as \( TL \lesssim m \), the mass of the lightest quark. The energy of the stretched string as a function of \( L \) is given by

\[
E(L) = TL + a - \frac{\pi c}{24L} + \ldots
\]

Here \( a \) is an undetermined constant, while \( c \) counts the number of degrees of freedom of the QCD flux tube. (There is no analog of \( \tilde{c} \) here because of the reflecting boundary conditions at the end of the string). If the string has no internal degrees of freedom, then \( c = 2 \) for the two transverse fluctuations. This contribution to the string energy is known as the Lüscher term.

4.4.2 The Weyl Anomaly

There is another way in which the central charge affects the stress-energy tensor. Recall that in the classical theory, one of the defining features of a CFT was the vanishing of the trace of the stress tensor,

\[
T^\alpha_\alpha = 0
\]

However, things are more subtle in the quantum theory. While \( \langle T^\alpha_\alpha \rangle \) indeed vanishes in flat space, it will not longer be true if we place the theory on a curved background. The purpose of this section is to show that

\[
\langle T^\alpha_\alpha \rangle = -\frac{c}{12} R
\]  

(4.35)

where \( R \) is the Ricci scalar of the 2d worldsheet. Before we derive this formula, some quick comments:

- Equation (4.35) holds for any state in the theory — not just the vacuum. This reflects the fact that it comes from regulating short distant divergences in the theory. But, at short distances all finite energy states look basically the same.

- Because \( \langle T^\alpha_\alpha \rangle \) is the same for any state it must be equal to something that depends only on the background metric. This something should be local and must be dimension 2. The only candidate is the Ricci scalar \( R \). For this reason, the formula \( \langle T^\alpha_\alpha \rangle \sim R \) is the most general possibility. The only question is: what is the coefficient. And, in particular, is it non-zero?
By a suitable choice of coordinates, we can always put any 2d metric in the form 
\[ g_{\alpha\beta} = e^{2\omega} \delta_{\alpha\beta}. \] 
In these coordinates, the Ricci scalar is given by 
\[ R = -2e^{-2\omega} \partial^2 \omega \] 
which depends explicitly on the function \( \omega \). Equation (4.35) is then telling us 
that any conformal theory with \( c \neq 0 \) has at least one physical observable, \( \langle T_\alpha^\alpha \rangle \), 
which takes different values on backgrounds related by a Weyl transformation \( \omega \). 
This result is referred to as the Weyl anomaly, or sometimes as the trace anomaly.

There is also a Weyl anomaly for conformal field theories in higher dimensions. For example, 4d CFTs are characterized by two numbers, \( a \) and \( c \), which appear 
as coefficients in the Weyl anomaly, 
\[ \langle T^\mu_\mu \rangle_{4d} = \frac{c}{16\pi^2} C_{\rho\sigma\kappa\lambda} C^{\rho\sigma\kappa\lambda} - \frac{a}{16\pi^2} \tilde{R}_{\rho\sigma\kappa\lambda} \tilde{R}^{\rho\sigma\kappa\lambda} \] 
where \( C \) is the Weyl tensor and \( \tilde{R} \) is the dual of the Riemann tensor.

Equation (4.35) involves only the left-moving central charge \( c \). You might wonder 
what’s special about the left-moving sector. The answer, of course, is nothing. We also have 
\[ \langle T_\alpha^\alpha \rangle = \frac{-\bar{c}}{12} R \] 
In flat space, conformal field theories with different \( c \) and \( \bar{c} \) are perfectly accept-
able. However, if we wish these theories to be consistent in fixed, curved back-
grounds, then we require \( c = \bar{c} \). This is an example of a gravitational anomaly.

The fact that Weyl invariance requires \( c = 0 \) will prove crucial in string theory. 
We shall return to this in Chapter 5.

We will now prove the Weyl anomaly formula (4.35). Firstly, we need to derive 
an intermediate formula: the \( T_{z\bar{z}} T_{w\bar{w}} \) OPE. Of course, in the classical theory we found 
that conformal invariance requires \( T_{z\bar{z}} = 0 \). We will now show that it’s a little more 
subtle in the quantum theory.

Our starting point is the equation for energy conservation, 
\[ \partial T_{z\bar{z}} = -\bar{\partial} T_{z\bar{z}} \] 
Using this, we can express our desired OPE in terms of the familiar \( TT \) OPE, 
\[ \partial \bar{z} T_{z\bar{z}}(z, \bar{z}) \partial \bar{w} T_{w\bar{w}}(w, \bar{w}) = \bar{\partial} \bar{z} T_{z\bar{z}}(z, \bar{z}) \bar{\partial} \bar{w} T_{w\bar{w}}(w, \bar{w}) = \bar{\partial} \bar{z} \bar{\partial} \bar{w} \left[ \frac{c/2}{(z-w)^4} + \ldots \right] \] (4.37)
Now you might think that the right-hand-side just vanishes: after all, it is an anti-holomorphic derivative \(\bar{\partial}\) of a holomorphic quantity. But we shouldn’t be so cavalier because there is a singularity at \(z = w\). For example, consider the following equation,

\[
\bar{\partial}_z \partial_{\bar{w}} \ln |z - w|^2 = \bar{\partial}_z \frac{1}{z - w} = 2\pi \delta(z - w, \bar{z} - \bar{w}) \tag{4.38}
\]

We proved this statement after equation (4.21). (The factor of 2 difference from (4.21) can be traced to the conventions we defined for complex coordinates in Section 4.0.1). Looking at the intermediate step in (4.38), we again have an anti-holomorphic derivative of a holomorphic function and you might be tempted to say that this also vanishes. But you’d be wrong: subtle things happen because of the singularity and equation (4.38) tells us that the function \(1/z\) secretly depends on \(\bar{z}\). (This should really be understood as a statement about distributions, with the delta function integrated against arbitrary test functions). Using this result, we can write

\[
\bar{\partial}_z \partial_{\bar{w}} \frac{1}{(z - w)^4} = \frac{1}{6} \bar{\partial}_z \partial_{\bar{w}} \left( \partial^2_z \partial_w \frac{1}{z - w} \right) = \frac{\pi}{3} \partial^2_z \partial_w \partial_{\bar{w}} \delta(z - w, \bar{z} - \bar{w})
\]

Inserting this into the correlation function (4.37) and stripping off the \(\partial_z \partial_w\) derivatives on both sides, we end up with what we want,

\[
T_{zz}(z, \bar{z}) T_{ww}(w, \bar{w}) = \frac{c\pi}{6} \partial_z \partial_{\bar{w}} \delta(z - w, \bar{z} - \bar{w}) \tag{4.39}
\]

So the OPE of \(T_{zz}\) and \(T_{ww}\) almost vanishes, but there’s some strange singular behaviour going on as \(z \to w\). This is usually referred to as a contact term between operators and, as we have shown, it is needed to ensure the conservation of energy-momentum. We will now see that this contact term is responsible for the Weyl anomaly.

We assume that \(\langle T^\alpha_\alpha \rangle = 0\) in flat space. Our goal is to derive an expression for \(\langle T^\alpha_\alpha \rangle\) close to flat space. Firstly, consider the change of \(\langle T^\alpha_\alpha \rangle\) under a general shift of the metric \(\delta g_{\alpha\beta}\). Using the definition of the energy-momentum tensor (4.4), we have

\[
\delta \langle T^\alpha_\alpha(\sigma) \rangle = \delta \int D\phi \ e^{-S} T^\alpha_\alpha(\sigma) = \frac{1}{4\pi} \int D\phi \ e^{-S} \left( T^\alpha_\alpha(\sigma) \int d^2\sigma' \sqrt{g} \ \delta g^{\beta\gamma} T^\beta_\gamma(\sigma') \right)
\]

If we now restrict to a Weyl transformation, the change to a flat metric is \(\delta g_{\alpha\beta} = 2\omega \delta_{\alpha\beta}\), so the change in the inverse metric is \(\delta g^{\alpha\beta} = -2\omega \delta^{\alpha\beta}\). This gives

\[
\delta \langle T^\alpha_\alpha(\sigma) \rangle = -\frac{1}{2\pi} \int D\phi \ e^{-S} \left( T^\alpha_\alpha(\sigma) \int d^2\sigma' \ \omega(\sigma') T^\beta_\beta(\sigma') \right) \tag{4.40}
\]
Now we see why the OPE (4.39) determines the Weyl anomaly. We need to change between complex coordinates and Cartesian coordinates, keeping track of factors of 2. We have

\[ T^\alpha_\alpha(\sigma) T^\beta_\beta(\sigma') = 16 T_{z\bar{z}}(z, \bar{z}) T_{w\bar{w}}(w, \bar{w}) \]

Meanwhile, using the conventions laid down in 4.0.1, we have

\[ 8 \partial_z \bar{\partial}_w \delta(z - w, \bar{z} - \bar{w}) = -\partial^2 \delta(\sigma - \sigma') \]

This gives us the OPE in Cartesian coordinates

\[ T^\alpha_\alpha(\sigma) T^\beta_\beta(\sigma') = -\frac{c\pi}{3} \partial^2 \delta(\sigma - \sigma') \]

We now plug this into (4.40) and integrate by parts to move the two derivatives onto the conformal factor \( \omega \). We’re left with,

\[ \delta \langle T^\alpha_\alpha \rangle = \frac{c}{6} \partial^2 \omega \quad \Rightarrow \quad \langle T^\alpha_\alpha \rangle = -\frac{c}{12} R \]

where, to get to the final step, we’ve used (4.36) and, since we’re working infinitesimally, we can replace \( e^{-2\omega} \approx 1 \). This completes the proof of the Weyl anomaly, at least for spaces infinitesimally close to flat space. The fact that \( R \) remains on the right-hand-side for general 2d surfaces follows simply from the comments after equation (4.35), most pertinently the need for the expression to be reparameterization invariant.

4.4.3 c is for Cardy

The Casimir effect and the Weyl anomaly have a similar smell. In both, the central charge provides an extra contribution to the energy. We now demonstrate a different avatar of the central charge: it tells us the density of high energy states.

We will study conformal field theory on a Euclidean torus. We’ll keep our normalization \( \sigma \in [0, 2\pi) \), but now we also take \( \tau \) to be periodic, lying in the range

\[ \tau \in [0, \beta) \]

The partition function of a theory with periodic Euclidean time has a very natural interpretation: it is related to the free energy of the theory at temperature \( T = 1/\beta \).

\[ Z[\beta] = \text{Tr} e^{-\beta H} = e^{-\beta F} \quad (4.41) \]

At very low temperatures, \( \beta \rightarrow \infty \), the free energy is dominated by the lowest energy state. All other states are exponentially suppressed. But we saw in 4.4.1 that the vacuum state on the cylinder has Casimir energy \( H = -c/12 \). In the limit of low temperature, the partition function is therefore approximated by

\[ Z \rightarrow e^{c\beta/12} \quad \text{as} \quad \beta \rightarrow \infty \quad (4.42) \]
Now comes the trick. In Euclidean space, both directions of the torus are on equal footing. We’re perfectly at liberty to decide that \( \sigma \) is “time” and \( \tau \) is “space”. This can’t change the value of the partition function. So let’s make the swap.

To compare to our original partition function, we want the spatial direction to have range \([0, 2\pi]\). Happily, due to the conformal nature of our theory, we arrange this through the scaling

\[
\tau \to \frac{2\pi}{\beta} \tau, \quad \sigma \to \frac{2\pi}{\beta} \sigma
\]

Now we’re back where we started, but with the temporal direction taking values in \( \sigma \in [0, 4\pi^2/\beta] \). This tells us that the high-temperature and low-temperature partition functions are related,

\[
Z[4\pi^2/\beta] = Z[\beta]
\]

This is called modular invariance. We’ll come across it again in Section 6.4. Writing \( \beta' = 4\pi^2/\beta \), this tells us the very high temperature behaviour of the partition function

\[
Z[\beta'] \to e^{c\pi^2/3\beta'} \quad \text{as} \quad \beta' \to 0
\]

But the very high temperature limit of the partition function is sampling all states in the theory. On entropic grounds, this sampling is dominated by the high energy states. So this computation is telling us how many high energy states there are.

To see this more explicitly, let’s do some elementary manipulations in statistical mechanics. Any system has a density of states \( \rho(E) = e^{S(E)} \), where \( S(E) \) is the entropy. The free energy is given by

\[
e^{-\beta F} = \int dE \rho(E) e^{-\beta E} = \int dE e^{S(E) - \beta E}
\]

In two dimensions, all systems have an entropy which scales at large energy as

\[
S(E) \to N\sqrt{E} \quad \text{(4.43)}
\]

The coefficient \( N \) counts the number of degrees of freedom. The fact that \( S \sim \sqrt{E} \) is equivalent to the fact that \( F \sim T^2 \), as befits an energy density in a theory with one
spatial dimension. To see this, we need only approximate the integral by the saddle point $S'(E_\star) = \beta$. From (4.43), this gives us the free energy

$$F \sim N^2 T^2$$

We can now make the statement about the central charge more explicit. In a conformal field theory, the entropy of high energy states is given by

$$S(E) \sim \sqrt{cE}$$

This is Cardy’s formula. A more careful analysis of the coefficients shows that the high energy density of states scales as

$$S(E) \rightarrow 2\pi \sqrt{\frac{c}{6}} \left(ER - \frac{c}{24}\right)$$

(4.44)

where the offset is the Casimir energy (4.34) that we derived previously. This is the contribution from left-movers. There is a similar contribution from right-movers, depending on $\tilde{c}$.

4.4.4 $c$ has a Theorem

The connection between the central charge and the degrees of freedom in a theory is given further weight by a result of Zamalodchikov, known as the $c$-theorem. The idea of the $c$-theorem is to stand back and look at the space of all theories and the renormalization group (RG) flows between them.

Conformal field theories are special. They are the fixed points of the renormalization group, looking the same at all length scales. One can consider perturbing a conformal field theory by adding an extra term to the action,

$$S \rightarrow S + \alpha \int d^2 \sigma \mathcal{O}(\sigma)$$

Here $\mathcal{O}$ is a local operator of the theory, while $\alpha$ is some coefficient. These perturbations fall into three classes, depending on the dimension $\Delta$ of $\mathcal{O}$.

- $\Delta < 2$: In this case, $\alpha$ has positive dimension: $[\alpha] = 2 - \delta$. Such deformations are called relevant because they are important in the infra-red. RG flow takes us away from our original CFT. We only stop flowing when we hit a new CFT (which could be trivial with $c = 0$).

- $\Delta = 2$: The constant $\alpha$ is dimensionless. Such deformations are called marginal. The deformed theory defines a new CFT.
• $\Delta > 2$: The constant $\alpha$ has negative dimension. These deformations are irrelevant. The infra-red physics is still described by the original CFT. But the ultra-violet physics is altered.

We expect information is lost as we flow from an ultra-violet theory to the infra-red. The c-theorem makes this intuition precise. The theorem exhibits a function $c$ on the space of all theories which monotonically decreases along RG flows. At the fixed points, $c$ coincides with the central charge of the CFT.

**A Thermodynamic Proof of the c-Theorem**

There are a number of different proofs of the c-theorem. Here we give one that is particularly physical. The basic idea is to heat up the system to a finite temperature $T$ and compute the speed of sound. The c-theorem follows from the requirement that the speed of sound does not exceed the speed of light (which, in our conventions, is simply 1). I should warn you that the style of argument in this section is somewhat different from the rest of these lectures. But, if nothing else, it reminds you that just because you’re learning string theory, you shouldn’t neglect basic physics!

Let’s first start with a CFT. For simplicity, we assume that $c = \tilde{c}$. Then, from (4.44), we have the asymptotic behaviour

$$S(E) \rightarrow 4\pi \sqrt{\frac{cER}{6}}$$

where we have dropped the $c/24$ offset, and the overall coefficient is $4\pi$ rather than $2\pi$ because we are including both left- and right-moving sectors. To compare with familiar, thermodynamic formulae we write this in terms of the spatial volume $V = 2\pi R$, so

$$S(E) \rightarrow 4\pi \sqrt{\frac{\pi cEV}{3}}$$

Now, the temperature is defined to be

$$\frac{1}{T} = \frac{\partial S}{\partial E} = 2\pi \sqrt{\frac{\pi cV}{3E}} \quad \Rightarrow \quad \sqrt{E} = 2\pi T \sqrt{\frac{\pi cV}{3}}$$

From this, we can compute the entropy of a CFT as a function of temperature, rather than as a function of energy

$$S(T) = \frac{8\pi^3 cVT}{3} \quad \Rightarrow \quad s(T) = \frac{8\pi^3 c}{3} T$$

where $s = S/V$ is the entropy density.
Now we’ll consider a more general situation. We’ll flow from some CFT in the UV with central charge \( c_{UV} \) to another CFT in the IR with central charge \( c_{IR} \). It may be that the final theory is gapped – meaning that everything is massless – in which case \( c_{IR} = 0 \). Our goal is to prove that, regardless of the flow, we always have \( c_{UV} \geq c_{IR} \) (with equality if there is no flow at all). To achieve this, we need to play around with some thermodynamic identities. In particular, we need to following result

**Claim:**

\[
  s = \left. \frac{\partial P}{\partial T} \right|_V
\]

(4.46)

with \( P \) the pressure.

**Proof:** Given the energy \( E = E(S,V) \), the first law of thermodynamics tells us

\[
dE = TdS - PdV
\]

The free energy is then defined as \( F(T,V) = E - TS \) and obeys

\[
dF = -SdT - PdV
\]

(4.47)

But the free energy is extensive and this means that it must, in fact, be proportional to \( V \) since this is the only extensive quantity that it can depend on. So

\[
F(T,V) = -P(T)V
\]

From this we learn that

\[
dF = -\frac{\partial P}{\partial T}VdT - PdV
\]

Comparing to (4.47) gives us the claimed result (4.46). \( \square \)

Finally, we recall that the speed of sound in a system is given by (see, for example, the lectures on Fluid Mechanics)

\[
c_s^2 = \frac{dP}{d\epsilon}
\]

where \( \epsilon = E/V \) is the energy density. At fixed volume, we have

\[
dE = TdS \quad \Rightarrow \quad d\epsilon = Tds
\]
All of which means that we can express the speed of sound as

\[ c_s^2 = \frac{1}{T} \frac{dP}{ds} = \frac{1}{T} \frac{dP}{dT} \frac{dT}{ds} = \frac{s}{T} \frac{dT}{ds} = \frac{d \log T}{d \log s} \]

This is the key result that we need. Now we define a thermal \( c \)-function

\[ \chi = \frac{s}{T} \]

As we’ve seen in (4.45), when we have a CFT the function \( \chi \) is proportional to the central charge: \( \chi = 8\pi^3 c/3 \). If we flow from a CFT in the UV, with central charge \( c_{UV} \), to a different CFT in the IR with central charge \( c_{IR} \), then \( \chi \) will interpolate between these two values (multiplied by \( 8\pi^3/3 \)) as we vary the temperature. To prove the \( c \)-theorem, we need to show that as we decrease the temperature, and so excite lower energy degrees of freedom, the function \( \chi \) necessarily decreases. We do this by relating \( \chi \) to the speed of sound,

\[ \frac{1}{c_s^2} = \frac{d \log s}{d \log T} = \frac{d \log (\chi T)}{d \log T} = 1 + \frac{d \log \chi}{d \log T} \]

By causality, we must have \( c_s^2 \leq 1 \) (with equality when we have a CFT) and so

\[ \frac{d \log \chi}{d \log T} \geq 0 \quad \Rightarrow \quad \frac{d \chi}{dT} \geq 0 \]

But this is what we wanted. We learn that we necessarily have \( c_{UV} \geq c_{IR} \). This is the \( c \)-theorem.

4.5 The Virasoro Algebra

So far our discussion has been limited to the operators of the CFT. We haven’t said anything about states. We now remedy this. We start by taking a closer look at the map between the cylinder and the plane.

4.5.1 Radial Quantization

To discuss states in a quantum field theory we need to think about where they live and how they evolve. For example, consider a two dimensional quantum field theory defined on the plane. Traditionally, when quantizing this theory, we parameterize the plane by Cartesian coordinates \((t, x)\) which we’ll call “time” and “space”. The states live on spatial slices. The Hamiltonian generates time translations and hence governs the evolution of states.
However, the map between the cylinder and the plane suggests a different way to quantize a CFT on the plane. The complex coordinate on the cylinder is taken to be $\omega$, while the coordinate on the plane is $z$. They are related by,

$$\omega = \sigma + i\tau, \quad z = e^{-i\omega}$$

On the cylinder, states live on spatial slices of constant $\sigma$ and evolve by the Hamiltonian,

$$H = \partial_{\tau}$$

After the map to the plane, the Hamiltonian becomes the dilatation operator

$$D = z\partial + \bar{z}\bar{\partial}$$

If we want the states on the plane to remember their cylindrical roots, they should live on circles of constant radius. Their evolution is governed by the dilatation operator $D$. This approach to a theory is known as radial quantization.

Usually in a quantum field theory, we’re interested in time-ordered correlation functions. Time ordering on the cylinder becomes radial ordering on the plane. Operators in correlation functions are ordered so that those inserted at larger radial distance are moved to the left.

**Virasoro Generators**

Let’s look at what becomes of the stress tensor $T(z)$ evaluated on the plane. On the cylinder, we would decompose $T$ in a Fourier expansion.

$$T_{\text{cylinder}}(w) = -\sum_{m=-\infty}^{\infty} L_m e^{imw} + \frac{c}{24}$$

---

**Figure 24:** The map from the cylinder to the plane.

**Figure 25:**
After the transformation (4.33) to the plane, this becomes the Laurent expansion

\[ T(z) = \sum_{m=-\infty}^{\infty} \frac{L_m}{z^{m+2}} \]

As always, a similar statement holds for the right-moving sector

\[ \bar{T}(\bar{z}) = \sum_{m=-\infty}^{\infty} \frac{\bar{L}_m}{\bar{z}^{m+2}} \]

We can invert these expressions to get \( L_m \) in terms of \( T(z) \). We need to take a suitable contour integral

\[ L_n = \frac{1}{2\pi i} \oint dz \, z^{n+1} T(z) \quad , \quad \bar{L}_n = \frac{1}{2\pi i} \oint d\bar{z} \, \bar{z}^{n+1} \bar{T}(\bar{z}) \] (4.48)

where, if we just want \( L_n \) or \( \bar{L}_n \), we must make sure that there are no other insertions inside the contour.

In radial quantization, \( L_n \) is the conserved charge associated to the conformal transformation \( \delta z = z^{n+1} \). To see this, recall that the corresponding Noether current, given in (4.7), is \( J(z) = z^{n+1}T(z) \). Moreover, the contour integral \( \oint dz \) maps to the integral around spatial slices on the cylinder. This tells us that \( L_n \) is the conserved charge where “conserved” means that it is constant under time evolution on the cylinder, or under radial evolution on the plane. Similarly, \( \bar{L}_n \) is the conserved charge associated to the conformal transformation \( \delta \bar{z} = \bar{z}^{n+1} \).

When we go to the quantum theory, conserved charges become generators for the transformation. Thus the operators \( L_n \) and \( \bar{L}_n \) generate the conformal transformations \( \delta z = z^{n+1} \) and \( \delta \bar{z} = \bar{z}^{n+1} \). They are known as the Virasoro generators. In particular, our two favorite conformal transformations are

- \( L_{-1} \) and \( \bar{L}_{-1} \) generate translations in the plane.
- \( L_0 \) and \( \bar{L}_0 \) generate scaling and rotations.

The Hamiltonian of the system — which measures the energy of states on the cylinder — is mapped into the dilatation operator on the plane. When acting on states of the theory, this operator is represented as

\[ D = L_0 + \bar{L}_0 \]
4.5.2 The Virasoro Algebra

If we have some number of conserved charges, the first thing that we should do is compute their algebra. Representations of this algebra then classify the states of the theory. (For example, think angular momentum in the hydrogen atom). For conformal symmetry, we want to determine the algebra obeyed by the $L_n$ generators. It’s a nice fact that the commutation relations are actually encoded $TT$ OPE. Let’s see how this works.

We want to compute $[L_m, L_n]$. Let’s write $L_m$ as a contour integral over $\oint dz$ and $L_n$ as a contour integral over $\oint dw$. (Note: both $z$ and $w$ denote coordinates on the complex plane now). The commutator is

$$[L_m, L_n] = \left( \oint \frac{dz}{2\pi i} \oint \frac{dw}{2\pi i} - \oint \frac{dw}{2\pi i} \oint \frac{dz}{2\pi i} \right) z^{m+1} w^{n+1} T(z) T(w)$$

What does this actually mean?! We need to remember that all operator equations are to be viewed as living inside time-ordered correlation functions. Except, now we’re working on the $z$-plane, this statement has transmuted into radially ordered correlation functions: outies to the left, innies to the right.

So $L_m L_n$ means $\oint \int z$ while $L_n L_m$ means $\oint \int z$.

The trick to computing the commutator is to first fix $w$ and do the $\oint dz$ integrations. The resulting contour is,

$$[L_m, L_n] = \oint \frac{dw}{2\pi i} \oint \frac{dz}{2\pi i} z^{m+1} w^{n+1} T(z) T(w)$$

In other words, we do the $z$-integration around a fixed point $w$, to get

$$[L_m, L_n] = \oint \frac{dw}{2\pi i} \oint \frac{dz}{2\pi i} z^{m+1} w^{n+1} T(z) T(w)$$
\[
\oint \frac{dw}{2\pi i} \text{Res} \left[ z^{m+1}w^{n+1}\left( \frac{c/2}{(z-w)^4} + \frac{2T(w)}{(z-w)^2} + \frac{\partial T(w)}{z-w} + \ldots \right) \right]
\]

To compute the residue at \( z = w \), we first need to Taylor expand \( z^{m+1} \) about the point \( w \),

\[
z^{m+1} = w^{m+1} + (m+1)w^m(z-w) + \frac{1}{2}m(m+1)w^{m-1}(z-w)^2
\]

\[
+ \frac{1}{6}m(m^2-1)w^{m-2}(z-w)^3 + \ldots
\]

The residue then picks up a contribution from each of the three terms,

\[
[L_m, L_n] = \oint \frac{dw}{2\pi i} w^{n+1} \left[ w^{m+1}\partial T(w) + 2(m+1)w^m T(w) + \frac{c}{12}m(m^2-1)w^{m-2} \right]
\]

To proceed, it is simplest to integrate the first term by parts. Then we do the \( w \)-integral. But for both the first two terms, the resulting integral is of the form (4.48) and gives us \( L_{m+n} \). For the third term, we pick up the pole. The end result is

\[
[L_m, L_n] = (m-n)L_{m+n} + \frac{c}{12}m(m^2-1)\delta_{m+n,0}
\]

This is the Virasoro algebra. It’s quite famous. The \( \tilde{L}_n \)'s satisfy exactly the same algebra, but with \( c \) replaced by \( \tilde{c} \). Of course, \([L_n, \tilde{L}_m] = 0\). The appearance of \( c \) as an extra term in the Virasoro algebra is the reason it is called the “central charge”. In general, a central charge is an extra term in an algebra that commutes with everything else.

**Conformal = Diffeo + Weyl**

We can build some intuition for the Virasoro algebra. We know that the \( L_n \)'s generate conformal transformations \( \delta z = z^{n+1} \). Let’s consider something closely related: a coordinate transformation \( \delta z = z^{n+1} \). These are generated by the vector fields

\[
l_n = z^{n+1}\partial_z
\]

But it’s a simple matter to compute their commutation relations:

\[
[l_n, l_m] = (m-n)l_{m+n}
\]

So this is giving us the first part of the Virasoro algebra. But what about the central term? The key point to remember is that, as we stressed at the beginning of this chapter, a conformal transformation is not just a reparameterization of the coordinates: it is a reparameterization, followed by a compensating Weyl rescaling. The central term in the Virasoro algebra is due to the Weyl rescaling.
4.5.3 Representations of the Virasoro Algebra

With the algebra of conserved charges at hand, we can now start to see how the conformal symmetry classifies the states into representations.

Suppose that we have some state \(|\psi\rangle\) that is an eigenstate of \(L_0\) and \(\tilde{L}_0\).

\[
L_0 |\psi\rangle = h |\psi\rangle, \quad \tilde{L}_0 |\psi\rangle = \tilde{h} |\psi\rangle
\]

Back on the cylinder, this corresponds to some state with energy

\[
\frac{E}{2\pi} = h + \tilde{h} - \frac{c + \tilde{c}}{24}
\]

For this reason, we’ll refer to the eigenvalues \(h\) and \(\tilde{h}\) as the energy of the state. By acting with the \(L_n\) operators, we can get further states with eigenvalues

\[
L_0L_n |\psi\rangle = (L_nL_0 - nL_n) |\psi\rangle = (h - n)L_n |\psi\rangle
\]

This tells us that \(L_n\) are raising and lowering operators depending on the sign of \(n\). When \(n > 0\), \(L_n\) lowers the energy of the state and \(L_{-n}\) raises the energy of the state. If the spectrum is to be bounded below, there must be some states which are annihilated by all \(L_n\) and \(\tilde{L}_n\) for \(n > 0\). Such states are called primary. They obey

\[
L_n |\psi\rangle = \tilde{L}_n |\psi\rangle = 0 \quad \text{for all } n > 0
\]

In the language of representation theory, they are also called highest weight states. They are the states of lowest energy.

Representations of the Virasoro algebra can now be built by acting on the primary states with raising operators \(L_{-n}\) with \(n > 0\). Obviously this results in an infinite tower of states. All states obtained in this way are called descendants. From an initial primary state \(|\psi\rangle\), the tower fans out...

\[
|\psi\rangle, \quad L_{-1} |\psi\rangle, \quad L^2_{-1} |\psi\rangle, \quad L_{-2} |\psi\rangle, \quad L^3_{-1} |\psi\rangle, \quad L_{-1}L_{-2} |\psi\rangle, \quad L_{-3} |\psi\rangle
\]

The whole set of states is called a Verma module. They are the irreducible representations of the Virasoro algebra. This means that if we know the spectrum of primary states, then we know the spectrum of the whole theory.
Some comments:

- The vacuum state $|0\rangle$ has $h = 0$. This state obeys
  
  $$L_n |0\rangle = 0 \quad \text{for all } n \geq -1$$  
  \hspace{1cm} (4.50)

  Note that this state preserves the maximum number of symmetries: like all primary states, it is annihilated by $L_n$ with $n > 0$, but it is also annihilated by $L_0$ and $L_{-1}$. This fits with our intuition that the vacuum state should be invariant under as many symmetries as possible. You might think that we could go further and require that the vacuum state obeys $L_n |0\rangle = 0$ for all $n$. But that isn’t consistent with the central charge term in Virasoro algebra. The requirements (4.50) are the best we can do.

- This discussion should be ringing bells. We saw something very similar in the covariant quantization of the string, where we imposed conditions (2.6) as constraints. We will see the connection between the primary states and the spectrum of the string in Section 5.

- There’s a subtlety that you should be aware of: the states in the Verma module are not necessarily all independent. It could be that some linear combination of the states vanishes. This linear combination is known as a null state. The existence of null states depends on the values of $h$ and $c$. For example, suppose that we are in a theory in which the central charge is $c = 2h(5 - 8h)/(2h + 1)$, where $h$ is the energy of a primary state $|\psi\rangle$. Then it is simple to check that the following combination has vanishing norm:

  $$L_{-2} |\psi\rangle - \frac{3}{2(2h + 1)} L_{-1}^2 |\psi\rangle$$  
  \hspace{1cm} (4.51)

- There is a close relationship between the primary states and the primary operators defined in Section 4.2.3. In fact, the energies $h$ and $\tilde{h}$ of primary states will turn out to be exactly the weights of primary operators in the theory. This connection will be described in Section 4.6.

4.5.4 Consequences of Unitarity

There is one physical requirement that a theory must obey which we have so far neglected to mention: unitarity. This is the statement that probabilities are conserved when we are in Minkowski signature spacetime. Unitarity follows immediately if we have a Hermitian Hamiltonian which governs time evolution. But so far our discussion has been somewhat algebraic and we’ve not enforced this condition. Let’s do so now.
We retrace our footsteps back to the Euclidean cylinder and then back again to the Minkowski cylinder where we can ask questions about time evolution. Here the Hamiltonian density takes the form

\[ \mathcal{H} = T_{ww} + T_{\bar{w}\bar{w}} = \sum_n L_n e^{-i \sigma^+} + \bar{L}_n e^{-i \sigma^-} \]

So for the Hamiltonian to be Hermitian, we require

\[ L_n = L_n^\dagger \]

This requirement imposes some strong constraints on the structure of CFTs. Here we look at a couple of trivial, but important, constraints that arise due to unitarity and the requirement that the physical Hilbert space does not contain negative norm states.

- \( h \geq 0 \): This fact follows from looking at the norm,

\[ |L_{-1} | \psi \rangle|^2 = \langle \psi | L_{-1} L_{+1} | \psi \rangle = \langle \psi | [L_{+1}, L_{-1}] | \psi \rangle = 2h \langle \psi | \psi \rangle \geq 0 \]

The only state with \( h = 0 \) is the vacuum state \( |0 \rangle \).

- \( c > 0 \): To see this, we can look at

\[ |L_{-n} |0 \rangle|^2 = \langle 0 | [L_n, L_{-n}] | 0 \rangle = \frac{c}{12} n(n^2 - 1) \geq 0 \quad (4.52) \]

So \( c \geq 0 \). If \( c = 0 \), the only state in the vacuum module is the vacuum itself. It turns out that, in fact, the only state in the whole theory is the vacuum itself. Any non-trivial CFT has \( c > 0 \).

There are many more requirements of this kind that constrain the theory. In fact, it turns out that for CFTs with \( c < 1 \) these requirements are enough to classify and solve all theories.

4.6 The State-Operator Map

In this section we describe one particularly important aspect of conformal field theories: a map between states and local operators.

Firstly, let’s get some perspective. In a typical quantum field theory, the states and local operators are very different objects. While local operators live at a point in spacetime, the states live over an entire spatial slice. This is most clear if we write down a Schrödinger-style wavefunction. In field theory, this object is actually a wavefunctional, \( \Psi[\phi(\sigma)] \), describing the probability for every field configuration \( \phi(\sigma) \) at each point \( \sigma \) in space (but at a fixed time).
Given that states and local operators are such very different beasts, it’s a little surprising that in a CFT there is an isomorphism between them: it’s called the state-operator map. The key point is that the distant past in the cylinder gets mapped to a single point $z = 0$ in the complex plane. So specifying a state on the cylinder in the far past is equivalent to specifying a local disturbance at the origin.

To make this precise, we need to recall how to write down wavefunctions using path integrals. Different states are computed by putting different boundary conditions on the functional integral. Let’s start by returning to quantum mechanics and reviewing a few simple facts. The propagator for a particle to move from position $x_i$ at time $\tau_i$ to position $x_f$ at time $\tau_f$ is given by

$$G(x_f, x_i) = \int_{x(\tau_i) = x_i}^{x(\tau_f) = x_f} Dx \, e^{iS}$$

This means that if our system starts off in some state described by the wavefunction $\psi_i(x_i)$ at time $\tau_i$ then (ignoring the overall normalization) it evolves to the state

$$\psi_f(x_f, \tau_f) = \int dx_i G(x_f, x_i) \psi_i(x_i, \tau_i)$$

There are two lessons to take from this. Firstly, to determine the value of the wavefunction at a given point $x_f$, we evaluate the path integral restricting to paths which satisfy $x(\tau_f) = x_f$. Secondly, the initial state $\psi(x_i)$ acts as a weighting factor for the integral over initial boundary conditions.

Let’s now write down the same formula in a field theory, where we’re dealing with wavefunctionals. We’ll work with the Euclidean path integral on the cylinder. If we start with some state $\Psi_i[\phi_i(\sigma)]$ at time $\tau_i$, then it will evolve to the state

$$\Psi_f[\phi_f(\sigma), \tau_f] = \int D\phi_i \int_{\phi(\tau_i) = \phi_i}^{\phi(\tau_f) = \phi_f} D\phi \, e^{-S[\phi]} \, \Psi_i[\phi_i(\sigma), \tau_i]$$

How do we write a similar expression for states after the map to the complex plane? Now the states are defined on circles of constant radius, say $|z| = r$, and evolution is governed by the dilatation operator. Suppose the initial state is defined at $|z| = r_i$. In the path integral, we integrate over all fields with fixed boundary conditions $\phi(r_i) = \phi_i$ and $\phi(r_f) = \phi_f$ on the two edges of the annulus shown in the figure,

$$\Psi_f[\phi_f(\sigma), r_f] = \int D\phi_i \int_{\phi(r_i) = \phi_i}^{\phi(r_f) = \phi_f} D\phi \, e^{-S[\phi]} \, \Psi_i[\phi_i(\sigma), r_i]$$

This is the traditional way to define a state in field theory, albeit with a slight twist because we’re working in radial quantization. We see that the effect of the initial state is to change the weighting of the path integral over the inner ring at $|z| = r_i$. 


Let’s now see what happens as we take the initial state back to the far past and, ultimately, to $z = 0$? We must now integrate over the whole disc $|z| \leq r_f$, rather than the annulus. The only effect of the initial state is now to change the weighting of the path integral at the point $z = 0$. But that’s exactly what we mean by a local operator inserted at that point. This means that each local operator $\mathcal{O}(z=0)$ defines a different state in the theory,

$$\Psi[\phi_f; r] = \int_{\phi(r)=\phi_f}^{\phi(r)=\phi_f} \mathcal{D}\phi \ e^{-S[\phi]} \mathcal{O}(z=0)$$

We’re now integrating over all field configurations within the disc, including all possible values of the field at $z = 0$, which is analogous to integrating over the boundary conditions $\int \mathcal{D}\phi_i$ on the inner circle.

- The state-operator map is only true in conformal field theories where we can map the cylinder to the plane. It also holds in conformal field theories in higher dimensions (where $\mathbb{R} \times S^{D-1}$ can be mapped to the plane $\mathbb{R}^D$). In non-conformal field theories, a typical local operator creates many different states.

- The state-operator map does not say that the number of states in the theory is equal to the number of operators: this is never true. It does say that the states are in one-to-one correspondence with the local operators.

- You might think that you’ve seen something like this before. In the canonical quantization of free fields, we create states in a Fock space by acting with creation operators. That’s not what’s going on here! The creation operators are just about as far from local operators as you can get. They are the Fourier transforms of local operators.

- There’s a special state that we can create this way: the vacuum. This arises by inserting the identity operator $1$ into the path integral. Back in the cylinder.
picture, this just means that we propagate the state back to time \( \tau = -\infty \) which is a standard trick used in the Euclidean path integral to project out all but the ground state. For this reason the vacuum is sometimes referred to, in operator notation, as \(|1\rangle\).

### 4.6.1 Some Simple Consequences

Let’s use the state-operator map to wrap up a few loose ends that have arisen in our study of conformal field theory.

Firstly, we’ve defined two objects that we’ve called “primary”: states and operators. The state-operator map relates the two. Consider the state \(|\mathcal{O}\rangle\), built from inserting a primary operator \(\mathcal{O}\) into the path integral at \(z = 0\). We can look at,

\[
L_n |\mathcal{O}\rangle = \oint \frac{dz}{2\pi i} z^{n+1} T(z) \mathcal{O}(z = 0) \\
= \oint \frac{dz}{2\pi i} z^{n+1} \left( \frac{h \mathcal{O}}{z^2} + \frac{\partial \mathcal{O}}{z} + \ldots \right)
\]

You may wonder what became of the path integral \(\int D\phi e^{-S[\phi]}\) in this expression. The answer is that it’s still implicitly there. Remember that operator expressions such as (4.48) are always taken to hold inside correlation functions. But putting an operator in the correlation function is the same thing as putting it in the path integral, weighted with \(e^{-S[\phi]}\).

From (4.53) we can see the effect of various generators on states

- \(L_{-1} |\mathcal{O}\rangle = |\partial \mathcal{O}\rangle\): In fact, this is true for all operators, not just primary ones. It is expected since \(L_{-1}\) is the translation generator.
- \(L_0 |\mathcal{O}\rangle = h |\mathcal{O}\rangle\): This is true of any operator with well defined transformation under scaling.
- \(L_n |\mathcal{O}\rangle = 0\) for all \(n > 0\). This is true only of primary operators \(\mathcal{O}\). Moreover, it is our requirement for \(|\mathcal{O}\rangle\) to be a primary state.

This has an important consequence. We stated earlier that one of the most important things to compute in a CFT is the spectrum of weights of primary operators. This seems like a slightly obscure thing to do. But now we see that it has a much more direct, physical meaning. It is the spectrum of energy and angular momentum of states of the theory defined on the cylinder.
Another loose end: when defining operators which carry specific weight, we made the statement that we could always work in a basis of operators which have specified eigenvalues under $D$ and $L$. This follows immediately from the statement that we can always find a basis of eigenstates of $H$ and $L$ on the cylinder.

Finally, we can use this idea of the state-operator map to understand why the OPE works so well in conformal field theories. Suppose that we’re interested in some correlation function, with operator insertions as shown in the figure. The statement of the OPE is that we can replace the two inner operators by a sum of operators at $z = 0$, independent of what’s going on outside of the dotted line. As an operator statement, that sounds rather surprising. But this follows by computing the path integral up to the dotted line, by which point the only effect of the two operators is to determine what state we have. This provides us a way of understanding why the OPE is exact in CFTs, with a radius of convergence equal to the next-nearest insertion.

**4.6.2 Our Favourite Example: The Free Scalar Field**

Let’s illustrate the state-operator map by returning yet again to the free scalar field. On a Euclidean cylinder, we have the mode expansion

$$X(w, \bar{w}) = x + \alpha' p \tau + i \sqrt{\alpha'} \sum_{n \neq 0} \frac{1}{n} \left( \alpha_n e^{i \alpha w} + \tilde{\alpha}_n e^{i \alpha \bar{w}} \right)$$

where we retain the requirement of reality in Minkowski space, which gave us $\alpha_n^* = \alpha_{-n}$ and $\tilde{\alpha}_n^* = \tilde{\alpha}_{-n}$. We saw in Section 4.3 that $X$ does not have good conformal properties. Before transforming to the $z = e^{-i w}$ plane, we should work with the primary field on the cylinder,

$$\partial_w X(w, \bar{w}) = -\sqrt{\alpha'} \sum_n \alpha_n e^{i \alpha w} \quad \text{with} \quad \alpha_0 \equiv i \sqrt{\alpha' / 2} p$$
Since $\partial X$ is a primary field of weight $h = 1$, its transformation to the plane is given by (4.18) and reads

$$\partial_z X(z) = \left( \frac{\partial}{\partial w} \right)^{-1} \partial_w X(w) = -i \sqrt{\frac{\alpha'}{2}} \sum_n \frac{\alpha_n}{z^{n+1}}$$

and similar for $\bar{\partial} X$. Inverting this gives an equation for $\alpha_n$ as a contour integral,

$$\alpha_n = i \sqrt{\frac{2}{\alpha'}} \oint \frac{dz}{2\pi i} z^n \partial X(z)$$

(4.54)

Just as the TT OPE allowed us to determine the $[L_m, L_n]$ commutation relations in the previous section, so the $\partial X \partial X$ OPE contains the information about the $[\alpha_m, \alpha_n]$ commutation relations. The calculation is straightforward,

$$[\alpha_m, \alpha_n] = -\frac{2}{\alpha'} \left( \oint \frac{dz}{2\pi i} \oint \frac{dw}{2\pi i} - \oint \frac{dw}{2\pi i} \oint \frac{dz}{2\pi i} \right) z^m w^n \partial X(z) \partial X(w)$$

$$= -\frac{2}{\alpha'} \oint \frac{dw}{2\pi i} \operatorname{Res}_{z=w} \left[ z^m w^n \left( -\frac{\alpha'/2}{(z-w)^2} + \ldots \right) \right]$$

$$= m \oint \frac{dw}{2\pi i} w^{m+n+1} = m \delta_{m+n,0}$$

where, in going from the second to third line, we have Taylor expanded $z$ around $w$. Hearteningly, the final result agrees with the commutation relation (2.2) that we derived in string theory using canonical quantization.

The State-Operator Map for the Free Scalar Field

Let’s now look at the map between states and local operators. We know from canonical quantization that the Fock space is defined by acting with creation operators $\alpha_{-m}$ with $m > 0$ on the vacuum $|0\rangle$. The vacuum state itself obeys $\alpha_m |0\rangle = 0$ for $m > 0$. Finally, there is also the zero mode $\alpha_0 \sim p$ which provides all states with another quantum number. A general state is given by

$$\prod_{m=1}^{\infty} \alpha_{-m}^k |0; p\rangle$$

Let’s try and recover these states by inserting operators into the path integral. Our first task is to check whether the vacuum state is indeed equivalent to the insertion of the identity operator. In other words, is the ground state wavefunctional of the theory on the circle $|z| = r$ really given by

$$\Psi_0[X_f] = \int_{X_f(r)} \mathcal{D}X \ e^{-S[X]}$$

(4.55)
We want to check that this satisfies the definition of the vacuum state, namely \( \alpha_m|0\rangle = 0 \) for \( m > 0 \). How do we act on the wavefunctional with an operator? We should still integrate over all field configurations \( X(z, \bar{z}) \), subject to the boundary conditions at \( X(|z| = r) = X_f \). But now we should insert the contour integral (4.54) at some \( |w| < r \) (because, after all, the state is only going to vanish after we’ve hit it with \( \alpha_m \), not before!). So we look at

\[
\alpha_m \Psi_0[X_f] = \int^{X_f} D X e^{-S[X]} \oint dw \frac{2\pi i}{w^m} \partial X(w)
\]

The path integral is weighted by the action (4.19) for a free scalar field. If a given configuration diverges somewhere inside the disc \( |z| < r \), then the action also diverges. This ensures that only smooth functions \( \partial X(z) \), which have no singularity inside the disc, contribute. But for such functions we have

\[
\oint dw \frac{2\pi i}{w^m} \partial X(z) = 0 \quad \text{for all } m \geq 0
\]

So the state (4.55) is indeed the vacuum state. In fact, since \( \alpha_0 \) also annihilates this state, it is identified as the vacuum state with vanishing momentum.

What about the excited states of the theory?

**Claim:** \( \alpha_{-m}|0\rangle = |\partial^m X\rangle \). By which we mean that the state \( \alpha_{-m}|0\rangle \) can be built from the path integral,

\[
\alpha_{-m}|0\rangle = \int D X e^{-S[X]} \partial^m X(z = 0)
\] (4.56)

**Proof:** We can check this by acting on \( |\partial^m X\rangle \) with the annihilation operators \( \alpha_n \).

\[
\alpha_n |\partial^m X\rangle \sim \int^{X_f(r)} D X e^{-S[X]} \oint dw \frac{2\pi i}{w^n} \partial X(w) \partial^m X(z = 0)
\]

We can focus on the operator insertions and use the OPE (4.23). We drop the path integral and just focus on the operator equation (because, after all, operator equations only make sense in correlation functions which is the same thing as in path integrals). We have

\[
\oint dw \frac{2\pi i}{w^n} w^m \partial^m - \frac{1}{(w - z)^2} \bigg|_{z = 0} = m! \oint dw \frac{2\pi i}{w^{n-1}} w^{n-1} = 0 \quad \text{unless } m = n
\]

This confirms that the state (4.56) has the right properties. \( \square \)
Finally, we should worry about the zero mode, or momentum $\alpha_0 \sim p$. It is simple to show using the techniques above (together with the OPE (4.26)) that the momentum of a state arises by the insertion of the primary operator $e^{ipX}$. For example,

$$|0;p\rangle \sim \int \mathcal{D}X \ e^{-S[X]} \ e^{ipX(z=0)} .$$

### 4.7 Brief Comments on Conformal Field Theories with Boundaries

The open string lives on the infinite strip with spatial coordinate $\sigma \in [0, \pi]$. Here we make just a few brief comments on the corresponding conformal field theories.

As before, we can define the complex coordinate $w = \sigma + i\tau$ and make the conformal map

$$z = e^{-iw}$$

This time the map takes us to the upper-half plane: $\text{Im} \ z \geq 0$. The end points of the string are mapped to the real axis, $\text{Im} \ z = 0$.

Much of our previous discussion goes through as before. But now we need to take care of boundary conditions at $\text{Im} \ z = 0$. Let’s first look at $T_{\alpha\beta}$. Recall that the stress-energy tensor exists because of translational invariance. We still have translational invariance in the direction parallel to the boundary — let’s call the associated tangent vector $t^\alpha$. But translational invariance is broken perpendicular to the boundary — we call the normal vector $n^\alpha$. The upshot of this is that $T_{\alpha\beta} t^\beta$ remains a conserved current.

To implement Neumann boundary conditions, we insist that none of the current flows out of the boundary. The condition is

$$T_{\alpha\beta} n^\alpha t^\beta = 0 \quad \text{at} \quad \text{Im} \ z = 0$$

In complex coordinates, this becomes

$$T_{zz} = T_{\bar{z}\bar{z}} \quad \text{at} \quad \text{Im} \ z = 0$$

There’s a simple way to implement this: we extend the definition of $T_{zz}$ from the upper-half plane to the whole complex plane by defining

$$T_{zz}(z) = T_{\bar{z}\bar{z}}(\bar{z})$$
For the closed string we had both functions $T$ and $\bar{T}$ in the whole plane. But for the open string, we have just one of these — say, $T$, — in the whole plane. This contains the same information as both $T$ and $\bar{T}$ in the upper-half plane. It’s simpler to work in the whole plane and focus just on $T$. Correspondingly, we now have just a single set of Virasoro generators,

$$L_n = \oint \frac{dz}{2\pi i} z^{n+1} T_{zz}(z)$$

There is no independent $\bar{L}_n$ for the open string.

A similar doubling trick works when computing the propagator for the free scalar field. The scalar field $X(z, \bar{z})$ is only defined in the upper-half plane. Suppose we want to implement Neumann boundary conditions. Then the propagator is defined by

$$\langle X(z, \bar{z}) X(w, \bar{w}) \rangle = G(z, \bar{z}; w, \bar{w})$$

which obeys $\partial^2 G = -2\pi \alpha' \delta(z - w, \bar{z} - \bar{w})$ subject to the boundary condition

$$\partial_\sigma G(z, \bar{z}; w, \bar{w})|_{\sigma=0} = 0$$

But we solve problems like this in our electrodynamics courses. A useful way of proceeding is to introduce an “image charge” in the lower-half plane. We now let $X(z, \bar{z})$ vary over the whole complex plane with its dynamics governed by the propagator

$$G(z, \bar{z}; w, \bar{w}) = -\frac{\alpha'}{2} \ln |z - w|^2 - \frac{\alpha'}{2} \ln |z - \bar{w}|^2$$

(4.57)

Much of the remaining discussion of CFTs carries forward with only minor differences. However, there is one point that is simple but worth stressing because it will be of importance later. This concerns the state-operator map. Recall the logic that leads us to this idea: we consider a state at fixed time on the strip and propagate it back to past infinity $\tau \to -\infty$. After the map to the half-plane, past infinity is again the origin. But now the origin lies on the boundary. We learn that the state-operator map relates states to local operators defined on the boundary.

This fact ensures that theories on a strip have fewer states than those on the cylinder. For example, for a free scalar field, Neumann boundary conditions require $\partial X = \bar{\partial} X$ at $\text{Im} z = 0$. (This follows from the requirement that $\partial_\sigma X = 0$ at $\sigma = 0, \pi$ on the strip). On the cylinder, the operators $\partial X$ and $\bar{\partial} X$ give rise to different states; on the strip they give rise to the same state. This, of course, mirrors what we’ve seen for the quantization of the open string where boundary conditions mean that we have only half the oscillator modes to play with.
5. The Polyakov Path Integral and Ghosts

At the beginning of the last chapter, we stressed that there are two very different interpretations of conformal symmetry depending on whether we’re thinking of a fixed 2d background or a dynamical 2d background. In applications to statistical physics, the background is fixed and conformal symmetry is a global symmetry. In contrast, in string theory the background is dynamical. Conformal symmetry is a gauge symmetry, a remnant of diffeomorphism invariance and Weyl invariance.

But gauge symmetries are not symmetries at all. They are redundancies in our description of the system. As such, we can’t afford to lose them and it is imperative that they don’t suffer an anomaly in the quantum theory. At worst, theories with gauge anomalies make no sense. (For example, Yang-Mills theory coupled to only left-handed fundamental fermions is a nonsensical theory for this reason). At best, it may be possible to recover the quantum theory, but it almost certainly has nothing to do with the theory that you started with.

Piecing together some results from the previous chapter, it looks like we’re in trouble. We saw that the Weyl symmetry is anomalous since the expectation value of the stress-energy tensor takes different values on backgrounds related by a Weyl symmetry:

\[ \langle T_\alpha \rangle = -\frac{c}{12} R \]

On fixed backgrounds, that’s merely interesting. On dynamical backgrounds, it’s fatal. What can we do? It seems that the only way out is to ensure that our theory has \( c = 0 \). But we’ve already seen that \( c > 0 \) for all non-trivial, unitary CFTs. We seem to have reached an impasse. In this section we will discover the loophole. It turns out that we do indeed require \( c = 0 \), but there’s a way to achieve this that makes sense.

5.1 The Path Integral

In Euclidean space the Polyakov action is given by,

\[ S_{\text{Poly}} = \frac{1}{4\pi\alpha'} \int d^2 \sigma \sqrt{g} \; g^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu \delta_{\mu\nu} \]

From now on, our analysis of the string will be in terms of the path integral\(^6\). We integrate over all embedding coordinates \( X^\mu \) and all worldsheet metrics \( g_{\alpha\beta} \). Schematically,

the path integral is given by,

\[ Z = \frac{1}{\text{Vol}} \int \mathcal{D}g \mathcal{D}X \ e^{-S_{\text{Poly}}[X,g]} \]

The “Vol” term is all-important. It refers to the fact that we shouldn’t be integrating over all field configurations, but only those physically distinct configurations not related by diffeomorphisms and Weyl symmetries. Since the path integral, as written, sums over all fields, the “Vol” term means that we need to divide out by the volume of the gauge action on field space.

To make the situation more explicit, we need to split the integration over all field configurations into two pieces: those corresponding to physically distinct configurations — schematically depicted as the dotted line in the figure — and those corresponding to gauge transformations — which are shown as solid lines. Dividing by “Vol” simply removes the piece of the partition function which comes from integrating along the solid-line gauge orbits.

In an ordinary integral, if we change coordinates then we pick up a Jacobian factor for our troubles. The path integral is no different. We want to decompose our integration variables into physical fields and gauge orbits. The tricky part is to figure out what Jacobian we get. Thankfully, there is a standard method to determine the Jacobian, first introduced by Faddeev and Popov. This method works for all gauge symmetries, including Yang-Mills and you will also learn about it in the “Advanced Quantum Field Theory” course.

5.1.1 The Faddeev-Popov Method

We have two gauge symmetries: diffeomorphisms and Weyl transformations. We will schematically denote both of these by \( \zeta \). The change of the metric under a general gauge transformation is \( g \rightarrow g^\zeta \). This is shorthand for,

\[ g_{\alpha\beta}(\sigma) \rightarrow g^\zeta_{\alpha\beta}(\sigma') = e^{2\omega(\sigma)} \frac{\partial \sigma^\gamma}{\partial \sigma'^\alpha} \frac{\partial \sigma'^\delta}{\partial \sigma^\gamma} g_{\gamma\delta}(\sigma) \]

In two dimensions these gauge symmetries allow us to put the metric into any form that we like — say, \( \hat{g} \). This is called the fiducial metric and will represent our choice of gauge fixing. Two caveats:

- Firstly, it’s not true that we can put any 2d metric into the form \( \hat{g} \) of our choosing. This is only true locally. Globally, it remains true if the worldsheet has the
topology of a cylinder or a sphere, but not for higher genus surfaces. We’ll revisit this issue in Section 6.

- Secondly, fixing the metric locally to \( \hat{g} \) does not fix all the gauge symmetries. We still have the conformal symmetries to deal with. We’ll revisit this in the Section 6 as well.

Our goal is to only integrate over physically inequivalent configurations. To achieve this, first consider the integral over the gauge orbit of \( \hat{g} \). For some value of the gauge transformation \( \zeta \), the configuration \( g^\zeta \) will coincide with our original metric \( g \). We can put a delta-function in the integral to get

\[
\int D\zeta \delta(g - \hat{g}^\zeta) = \Delta^{-1}_{FP}[g]
\] (5.1)

This integral isn’t equal to one because we need to take into account the Jacobian factor. This is analogous to the statement that \( \int dx \delta(f(x)) = 1/|f'| \), evaluated at points where \( f(x) = 0 \). In the above equation, we have written this Jacobian factor as \( \Delta^{-1}_{FP} \). The inverse of this, namely \( \Delta_{FP} \), is called the Faddeev-Popov determinant. We will evaluate it explicitly shortly. Some comments:

- This whole procedure is rather formal and runs into the usual difficulties with trying to define the path integral. Just as for Yang-Mills theory, we will find that it results in sensible answers.

- We will assume that our gauge fixing is good, meaning that the dotted line in the previous figure cuts through each physically distinct configuration exactly once. Equivalently, the integral over gauge transformations \( D\zeta \) clicks exactly once with the delta-function and we don’t have to worry about discrete ambiguities (known as Gribov copies in QCD).

- The measure is taken to be the analogue of the Haar measure for Lie groups, invariant under left and right actions

\[
D\zeta = D(\zeta' \zeta) = D(\zeta \zeta')
\]

When gauge fixing in Yang-Mills theory, the first thing we do is prove that the Faddeev-Popov determinant \( \Delta_{FP} \) is gauge invariant. However, our route here is a little more subtle. As we’ve stressed above, the Weyl anomaly means that our original theory actually fails to be gauge invariant. We will see that the Faddeev-Popov determinant also fails but can, in certain circumstances, cancel the original failure leaving behind a well-defined theory.
The Faddeev-Popov procedure starts by inserting a factor of unity into the path integral, in the guise of

\[ 1 = \Delta_{FP}[g] \int \mathcal{D}\zeta \, \delta(g - \hat{g} \zeta) \]

We’ll call the resulting path integral expression \( Z[\hat{g}] \) since it depends on the choice of fiducial metric \( \hat{g} \). The first thing we do is use the \( \delta(g - \hat{g} \zeta) \) delta-function to do the integral over metrics,

\[
Z[\hat{g}] = \frac{1}{\text{Vol}} \int \mathcal{D}\zeta \mathcal{D}X \mathcal{D}g \, \Delta_{FP}[g] \delta(g - \hat{g} \zeta) \, e^{-S_{\text{Poly}}[X,g]}
= \frac{1}{\text{Vol}} \int \mathcal{D}\zeta \mathcal{D}X \, \Delta_{FP}[\hat{g} \zeta] \, e^{-S_{\text{Poly}}[X,\hat{g} \zeta]} \tag{5.2}
\]

At this stage the integrand depends on \( \hat{g} \), where \( \zeta \) is shorthand for a diffeomorphism and Weyl transformation. Everything in the equation is invariant under diffeomorphisms, but Weyl transformations are another matter. We know that quantum theory \( \int \mathcal{D}X \, e^{-S_{\text{Poly}}} \) suffers a Weyl anomaly. The action \( S_{\text{Poly}} \) is invariant under Weyl rescalings, so the subtlety must come from the measure. Meanwhile, anticipating what’s to come, we will find a similar issue with the Faddeev-Popov determinant \( \Delta_{FP} \).

If, however, we find ourselves in the fortunate situation where the problems cancel then things would work out nicely. In that situation, everything on the right-hand side of (5.2) would conspire to be invariant under both diffeomorphisms and Weyl transformations and we could write

\[
Z[\hat{g}] = \frac{1}{\text{Vol}} \int \mathcal{D}\zeta \mathcal{D}X \, \Delta_{FP}[\hat{g}] \, e^{-S_{\text{Poly}}[X,\hat{g}]}
\]

But now, nothing depends on the gauge transformation \( \zeta \). Indeed, this is precisely the integration over the gauge orbits that we wanted to isolate and it cancels the “Vol” factor sitting outside. We’re left with

\[
Z[\hat{g}] = \int \mathcal{D}X \, \Delta_{FP}[\hat{g}] \, e^{-S_{\text{Poly}}[X,\hat{g}]} \tag{5.3}
\]

This is the integral over physically distinct configurations — the dotted line in the previous figure. We see that the Faddeev-Popov determinant is precisely the Jacobian factor that we need.

Clearly the above discussion only flies if we find ourselves in a situation in which the theory (5.2) is genuinely Weyl invariant. Our next task is to understand when this happens which means that we need to figure out what becomes of \( \Delta_{FP} \) when we do a Weyl transformation.
5.1.2 The Faddeev-Popov Determinant

We still need to compute $\Delta_{FP}[\hat{g}]$. It’s defined in (5.1). Let’s look at gauge transformations $\zeta$ which are close to the identity. In this case, the delta-function $\delta(g - \hat{g} \xi)$ is going to be non-zero when the metric $g$ is close to the fiducial metric $\hat{g}$. In fact, it will be sufficient to look at the delta-function $\delta(\hat{g} - \hat{g} \xi)$, which is only non-zero when $\zeta = 0$. We take an infinitesimal Weyl transformation parameterized by $\omega(\sigma)$ and an infinitesimal diffeomorphism $\delta \sigma^\alpha = v^\alpha(\sigma)$. The change in the metric is

$$\delta \hat{g}_{\alpha\beta} = 2\omega \hat{g}_{\alpha\beta} + \nabla_\alpha v_\beta + \nabla_\beta v_\alpha$$

Plugging this into the delta-function, the expression for the Faddeev-Popov determinant becomes

$$\Delta_{FP}^{-1}[\hat{g}] = \int D\omega Dv \delta(2\omega \hat{g}_{\alpha\beta} + \nabla_\alpha v_\beta + \nabla_\beta v_\alpha)$$

(5.4)

where we’ve replaced the integral $D\zeta$ over the gauge group with the integral $D\omega Dv$ over the Lie algebra of group since we’re near the identity. (We also suppress the subscript on $v_\alpha$ in the measure factor to keep things looking tidy).

At this stage it’s useful to represent the delta-function in its integral, Fourier form. For a single delta-function, this is $\delta(x) = \int dp \exp(2\pi i px)$. But the delta-function in (5.4) is actually a delta-functional: it restricts a whole function. Correspondingly, the integral representation is in terms of a functional integral,

$$\Delta_{FP}^{-1}[\hat{g}] = \int D\omega Dv D\beta \exp \left( 2\pi i \int d^2 \sigma \sqrt{\hat{g}} \beta^{\alpha\beta}[2\omega \hat{g}_{\alpha\beta} + \nabla_\alpha v_\beta + \nabla_\beta v_\alpha] \right)$$

where $\beta^{\alpha\beta}$ is a symmetric 2-tensor on the worldsheet.

We now simply do the $\int D\omega$ integral. It doesn’t come with any derivatives, so it merely acts as a Lagrange multiplier, setting

$$\beta^{\alpha\beta} \hat{g}_{\alpha\beta} = 0$$

In other words, after performing the $\omega$ integral, $\beta^{\alpha\beta}$ is symmetric and traceless. We’ll take this to be the definition of $\beta^{\alpha\beta}$ from now on. So, finally we have

$$\Delta_{FP}^{-1}[\hat{g}] = \int Dv D\beta \exp \left( 4\pi i \int d^2 \sigma \sqrt{\hat{g}} \beta^{\alpha\beta} \nabla_\alpha v_\beta \right)$$
5.1.3 Ghosts

The previous manipulations give us an expression for $\Delta_{FP}^{-1}$. But we want to invert it to get $\Delta_{FP}$. Thankfully, there’s a simple way to achieve this. Because the integrand is quadratic in $v$ and $\beta$, we know that the integral computes the inverse determinant of the operator $\nabla_\alpha$. (Strictly speaking, it computes the inverse determinant of the projection of $\nabla_\alpha$ onto symmetric, traceless tensors. This observation is important because it means the relevant operator is a square matrix which is necessary to talk about a determinant). But we also know how to write down an expression for the determinant $\Delta_{FP}$, instead of its inverse, in terms of path integrals: we simply need to replace the commuting integration variables with anti-commuting fields,

$$
\beta_{\alpha\beta} \rightarrow b_{\alpha\beta} \\
v^\alpha \rightarrow c^\alpha
$$

where $b$ and $c$ are both Grassmann-valued fields (i.e. anti-commuting). They are known as *ghost fields*. This gives us our final expression for the Faddeev-Popov determinant,

$$
\Delta_{FP}[g] = \int \mathcal{D}b \mathcal{D}c \exp[iS_{\text{ghost}}]
$$

where the ghost action is defined to be

$$
S_{\text{ghost}} = \frac{1}{2\pi} \int d^2\sigma \sqrt{g} \ b_{\alpha\beta} \nabla_\alpha c^\beta \tag{5.5}
$$

and we have chosen to rescale the $b$ and $c$ fields at this last step to get a factor of $1/2\pi$ sitting in front of the action. (This only changes the normalization of the partition function which doesn’t matter). Rotating back to Euclidean space, the factor of $i$ disappears. The expression for the full partition function (5.3) is

$$
Z[g] = \int \mathcal{D}X \mathcal{D}b \mathcal{D}c \exp (-S_{\text{Poly}}[X, \hat{g}] - S_{\text{ghost}}[b, c, \hat{g}])
$$

Something lovely has happened. Although the ghost fields were introduced as some auxiliary constructs, they now appear on the same footing as the dynamical fields $X$. We learn that gauge fixing comes with a price: our theory has extra ghost fields.

The role of these ghost fields is to cancel the unphysical gauge degrees of freedom, leaving only the $D - 2$ transverse modes of $X^\mu$. Unlike lightcone quantization, they achieve this in a way which preserves Lorentz invariance.
Simplifying the Ghost Action

The ghost action (5.5) looks fairly simple. But it looks even simpler if we work in conformal gauge,

\[ \hat{g}_{\alpha\beta} = e^{2\omega} \delta_{\alpha\beta} \]

The determinant is \( \sqrt{\hat{g}} = e^{2\omega} \). Recall that in complex coordinates, the measure is \( d^2\sigma = \frac{1}{2} d^2z \), while we can lower the index on the covariant derivative using \( \nabla^z = g^{\bar{z}\bar{z}} \nabla_{\bar{z}} = 2e^{-2\omega} \nabla_{\bar{z}} \). We have

\[ S_{\text{ghost}} = \frac{1}{2\pi} \int d^2z \left( b_{zz} \nabla_{\bar{z}} c^z + b_{\bar{z}\bar{z}} \nabla_z c^{\bar{z}} \right) \]

In deriving this, remember that there is no field \( b_{\bar{z}z} \) because \( b_{\alpha\beta} \) is traceless. Now comes the nice part: the covariant derivatives are actually just ordinary derivatives. To see why this is the case, look at

\[ \nabla_z c^z = \partial_z c^z + \Gamma^z_{\bar{z}\alpha} c^\alpha \]

But the Christoffel symbols are given by

\[ \Gamma^z_{\bar{z}\alpha} = \frac{1}{2} g^{\bar{z}\bar{z}} \left( \partial_{\bar{z}} g_{\alpha\bar{z}} + \partial_{\bar{z}} g_{\bar{z}\alpha} - \partial_{\bar{z}} g_{\bar{z}\bar{z}} \right) = 0 \quad \text{for } \alpha = z, \bar{z} \]

So in conformal gauge, the ghost action factorizes into two free theories,

\[ S_{\text{ghost}} = \frac{1}{2\pi} \int d^2z \left( b_{zz} \partial_{\bar{z}} c^z + b_{\bar{z}\bar{z}} \partial_z c^{\bar{z}} \right) \]

The action doesn’t depend on the conformal factor \( \omega \). In other words, it is Weyl invariant without any need to change \( b \) and \( c \): these are therefore both neutral under Weyl transformations.

(It’s worth pointing out that \( b_{\alpha\beta} \) and \( c^\alpha \) are neutral under Weyl transformations. But if we raise or lower these indices, then the fields pick up factors of the metric. So \( b^{\alpha\beta} \) and \( c_{\alpha} \) would not be neutral under Weyl transformations).

5.2 The Ghost CFT

Fixing the Weyl and diffeomorphism gauge symmetries has left us with two new dynamical ghost fields, \( b \) and \( c \). Both are Grassmann (i.e. anti-commuting) variables. Their dynamics is governed by a CFT. Define

\[ b = b_{zz}, \quad \bar{b} = b_{\bar{z}\bar{z}} \]
\[ c = c^z, \quad \bar{c} = c^{\bar{z}} \]
The ghost action is given by

\[ S_{\text{ghost}} = \frac{1}{2\pi} \int d^2z \ (b \bar{\partial}c + \bar{b} \partial\bar{c}) \]

Which gives the equations of motion

\[ \bar{\partial}b = \partial\bar{b} = \bar{\partial}c = \partial\bar{c} = 0 \]

So we see that \( b \) and \( c \) are holomorphic fields, while \( \bar{b} \) and \( \bar{c} \) are anti-holomorphic.

Before moving onto quantization, there’s one last bit of information we need from the classical theory: the stress tensor for the \( bc \) ghosts. The calculation is a little bit fiddly. We use the general definition of the stress tensor \((4.4)\), which requires us to return to the theory \((5.5)\) on a general background and vary the metric \( g^{\alpha\beta} \). The complications are twofold. Firstly, we pick up a contribution from the Christoffel symbol that is lurking inside the covariant derivative \( \nabla^\alpha \). Secondly, we must also remember that \( b_{\alpha\beta} \) is traceless. But this is a condition which itself depends on the metric: \( b_{\alpha\beta} g^{\alpha\beta} = 0 \). To account for this we should add a Lagrange multiplier to the action imposing tracelessness. After correctly varying the metric, we may safely retreat back to flat space where the end result is rather simple. We have \( T_{zz} = 0 \), as we must for any conformal theory. Meanwhile, the holomorphic and anti-holomorphic parts of the stress tensor are given by,

\[ T = 2(\partial c) b + c \partial b \quad , \quad \bar{T} = 2(\bar{\partial}\bar{c}) \bar{b} + \bar{c} \bar{\partial}\bar{b}. \]  

\[(5.6)\]

**Operator Product Expansions**

We can compute the OPEs of these fields using the standard path integral techniques that we employed in the last chapter. In what follows, we’ll just focus on the holomorphic piece of the CFT. We have, for example,

\[ 0 = \int \mathcal{D}b\mathcal{D}c \ \frac{\delta}{\delta b(\sigma)} \ [e^{-S_{\text{ghost}} b(\sigma')} \] = \[ \int \mathcal{D}b\mathcal{D}c \ e^{-S_{\text{ghost}}} \left[ -\frac{1}{2\pi} \bar{\partial}c(\sigma) b(\sigma') + \delta(\sigma - \sigma') \right] \]

which tells us that

\[ \bar{\partial}c(\sigma) b(\sigma') = 2\pi \delta(\sigma - \sigma') \]

Similarly, looking at \( \delta/\delta c(\sigma) \) gives

\[ \bar{\partial}b(\sigma) c(\sigma') = 2\pi \delta(\sigma - \sigma') \]
We can integrate both of these equations using our favorite formula \( \bar{\partial}(1/z) = 2\pi \delta(z, \bar{z}) \).

We learn that the OPEs between fields are given by

\[
\begin{align*}
    b(z) c(w) &= \frac{1}{z - w} + \ldots \\
    c(w) b(z) &= \frac{1}{w - z} + \ldots
\end{align*}
\]

In fact the second equation follows from the first equation and Fermi statistics. The OPEs of \( b(z) b(w) \) and \( c(z) c(w) \) have no singular parts. They vanish as \( z \to w \).

Finally, we need the stress tensor of the theory. After normal ordering, it is given by

\[
T(z) = 2 : \partial c(z) b(z) : + : c(z) \partial b(z) :
\]

We will shortly see that with this choice, \( b \) and \( c \) carry appropriate weights for tensor fields which are neutral under Weyl rescaling.

**Primary Fields**

We will now show that both \( b \) and \( c \) are primary fields, with weights \( h = 2 \) and \( h = -1 \) respectively. Let’s start by looking at \( c \). The OPE with the stress tensor is

\[
T(z) c(w) = 2 : \partial c(z) b(z) : c(w) + : c(z) \partial b(z) : c(w) = 2 \partial c(z) - \frac{c(z)}{(z - w)^2} + \ldots = -\frac{c(w)}{(z - w)^2} + \frac{\partial c(w)}{z - w} + \ldots
\]

confirming that \( c \) has weight \(-1\). When taking the OPE with \( b \), we need to be a little more careful with minus signs. We get

\[
T(z) b(w) = 2 : \partial c(z) b(z) : b(w) + : c(z) \partial b(z) : b(w) = -2b(z) \left( -\frac{1}{(z - w)^2} \right) - \frac{\partial b(z)}{z - w} = \frac{2b(w)}{(z - w)^2} + \frac{\partial b(w)}{z - w} + \ldots
\]

showing that \( b \) has weight \( 2 \). As we’ve pointed out a number of times, conformal = diffeo + Weyl. We mentioned earlier that the fields \( b \) and \( c \) are neutral under Weyl transformations. This is reflected in their weights, which are due solely to diffeomorphisms as dictated by their index structure: \( b_{zz} \) and \( c^z \).

**The Central Charge**

Finally, we can compute the \( TT \) OPE to determine the central charge of the \( bc \) ghost system.

\[
T(z) T(w) = 4 : \partial c(z)b(z) : \partial c(w)b(w) : + 2 : \partial c(z)b(z) : c(w)\partial b(w) : + 2 : c(z)\partial b(z) : \partial c(w)b(w) : + : c(z)\partial b(z) : c(w)\partial b(w) :
\]
For each of these terms, making two contractions gives a $(z - w)^{-4}$ contribution to the OPE. There are also two ways to make a single contraction. These give $(z - w)^{-1}$ or $(z - w)^{-2}$ or $(z - w)^{-3}$ contributions depending on what the derivatives hit. The end result is

$$T(z) T(w) = \frac{-4}{(z - w)^4} + \frac{4}{(z - w)^2} \partial c(z) b(w) : - \frac{4}{(z - w)^2} b(z) \partial c(w) :$$

$$- \frac{4}{(z - w)^4} + \frac{2}{z - w} \partial c(z) \partial b(w) : - \frac{4}{(z - w)^3} b(z) c(w) :$$

$$- \frac{4}{(z - w)^4} - \frac{4}{(z - w)^3} c(z) b(w) : + \frac{2}{z - w} \partial b(z) \partial c(w) :$$

$$- \frac{1}{(z - w)^4} - \frac{c(z) \partial b(w)}{(z - w)^2} + \frac{\partial b(z) c(w)}{(z - w)^2} + \ldots$$

After some Taylor expansions to turn $f(z)$ functions into $f(w)$ functions, together with a little collecting of terms, this can be written as,

$$T(z) T(w) = \frac{-13}{(z - w)^4} + \frac{2T(w)}{(z - w)^2} + \frac{\partial T(w)}{z - w} + \ldots$$

The first thing to notice is that it indeed has the form expected of $TT$ OPE. The second, and most important, thing to notice is the central charge of the $bc$ ghost system: it is

$$c = -26$$

5.3 The Critical “Dimension” of String Theory

Let’s put the pieces together. We’ve learnt that gauge fixing the diffeomorphisms and Weyl gauge symmetries results in the introduction of ghosts which contribute central charge $c = -26$. We’ve also learnt that the Weyl symmetry is anomalous unless $c = 0$. Since the Weyl symmetry is a gauge symmetry, it’s crucial that we keep it. We’re forced to add exactly the right degrees of freedom to the string to cancel the contribution from the ghosts.

The simplest possibility is to add $D$ free scalar fields. Each of these contributes $c = 1$ to the central charge, so the whole procedure is only consistent if we pick

$$D = 26$$

This agrees with the result we found in Chapter 2: it is the critical dimension of string theory.
However, there’s no reason that we have to work with free scalar fields. The consistency requirement is merely that the degrees of freedom of the string are described by a CFT with \( c = 26 \). Any CFT will do. Each such CFT describes a different background in which a string can propagate. If you like, the space of CFTs with \( c = 26 \) can be thought of as the space of classical solutions of string theory.

We learn that the “critical dimension” of string theory is something of a misnomer: it is really a “critical central charge”. Only for rather special CFTs can this central charge be thought of as a spacetime dimension.

For example, if we wish to describe strings moving in 4d Minkowski space, we can take \( D = 4 \) free scalars (one of which will be timelike) together with some other \( c = 22 \) CFT. This CFT may have a geometrical interpretation, or it may be something more abstract. The CFT with \( c = 22 \) is sometimes called the “internal sector” of the theory. It is what we really mean when we talk about the “extra hidden dimensions of string theory”. We’ll see some examples of CFTs describing curved spaces in Section 7.

There’s one final subtlety: we need to be careful with the transition back to Minkowski space. After all, we want one of the directions of the CFT, \( X^0 \), to have the wrong sign kinetic term. One safe way to do this is to keep \( X^0 \) as a free scalar field, with the remaining degrees of freedom described by some \( c = 25 \) CFT. This doesn’t seem quite satisfactory though since it doesn’t allow for spacetimes which evolve in time — and, of course, these are certainly necessary if we wish to understand early universe cosmology. There are still some technical obstacles to understanding the worldsheet of the string in time-dependent backgrounds. To make progress, and discuss string cosmology, we usually bi-pass this issue by working with the low-energy effective action which we will derive in Section 7.

5.3.1 The Usual Nod to the Superstring

The superstring has another gauge symmetry on the worldsheet: supersymmetry. This gives rise to more ghosts, the so-called \( \beta\gamma \) system, which turns out to have central charge +11. Consistency then requires that the degrees of freedom of the string have central charge \( c = 26 - 11 = 15 \).

However, now the CFTs must themselves be invariant under supersymmetry, which means that bosons come matched with fermions. If we add \( D \) bosons, then we also need to add \( D \) fermions. A free boson has \( c = 1 \), while a free fermion has \( c = 1/2 \). So, the total number of free bosons that we should add is \( D(1 + 1/2) = 15 \), giving us the
critical dimension of the superstring:

\[ D = 10 \]

5.3.2 An Aside: Non-Critical Strings

Although it’s a slight departure from the our main narrative, it’s worth pausing to mention what Polyakov actually did in his four page paper. His main focus was not critical strings, with \( D = 26 \), but rather non-critical strings with \( D \neq 26 \). From the discussion above, we know that these suffer from a Weyl anomaly. But it turns out that there is a way to make sense of the situation.

The starting point is to abandon Weyl invariance from the beginning. We start with \( D \) free scalar fields coupled to a dynamical worldsheet metric \( g_{\alpha\beta} \). (More generally, we could have any CFT). We still want to keep reparameterization invariance, but now we ignore the constraints of Weyl invariance. Of course, it seems likely that this isn’t going to have too much to do with the Nambu-Goto string, but let’s proceed anyway. Without Weyl invariance, there is one extra term that it is natural to add to the 2d theory: a worldsheet cosmological constant \( \mu \),

\[
S_{\text{non-critical}} = \frac{1}{4\pi\alpha'} \int d^2\sigma \sqrt{g} \left( g^{\alpha\beta} \partial_{\alpha} X^\mu \partial_{\beta} X_\mu + \mu \right)
\]

Our goal will be to understand how the partition function changes under a Weyl rescaling. There will be two contributions: one from the explicit \( \mu \) dependence and one from the Weyl anomaly. Consider two metrics related by a Weyl transformation

\[ \hat{g}_{\alpha\beta} = e^{2\omega} g_{\alpha\beta} \]

As we vary \( \omega \), the partition function \( Z[\hat{g}] \) changes as

\[
\frac{1}{Z} \frac{\partial Z}{\partial \omega} = \frac{1}{Z} \int D X \ e^{-S} \left( - \frac{\partial S}{\partial \hat{g}_{\alpha\beta}} \frac{\partial \hat{g}_{\alpha\beta}}{\partial \omega} \right) = \frac{1}{Z} \int D X \ e^{-S} \left( - \frac{1}{2\pi} \sqrt{\hat{g}} T^\alpha_{\alpha} \right) = \frac{c}{24\pi} \sqrt{\hat{g}} \hat{R} - \frac{1}{2\pi\alpha'} \mu e^{2\omega}
\]

where, in the last two lines, we used the Weyl anomaly (4.35) and the relationship between Ricci curvatures (1.29). The central charge appearing in these formulae includes the contribution from the ghosts,

\[ c = D - 26 \]
We can now just treat this as a differential equation for the partition function $Z$ and solve. This allows us to express the partition function $Z[\hat{g}]$, defined on one worldsheet metric, in terms of $Z[g]$, defined on another. The relationship is,

$$Z[\hat{g}] = Z[g] \exp \left[ -\frac{1}{4\pi\alpha'} \int d^2\sigma \sqrt{\hat{g}} \left( 2\mu e^{2\omega} - \frac{\epsilon\alpha'}{6} \left( g^{\alpha\beta} \partial_\alpha \omega \partial_\beta \omega + R \omega \right) \right) \right]$$

We see that the scaling mode $\omega$ inherits a kinetic term. It now appears as a new dynamical scalar field in the theory. It is often called the Liouville field on account of the exponential potential term multiplying $\mu$. Solving this theory is quite hard\textsuperscript{7}. Notice also that our new scalar field $\omega$ appears in the final term multiplying the Ricci scalar $R$. We will describe the significance of this in Section 7.2.1. We'll also see another derivation of this kind of Lagrangian in Section 7.4.4.

### 5.4 States and Vertex Operators

In Chapter 2 we determined the spectrum of the string in flat space. What is the spectrum for a general string background? The theory consists of the $b$ and $c$ ghosts, together with a $c = 26$ CFT. At first glance, it seems that we have a greatly enlarged Hilbert space since we can act with creation operators from all fields, including the ghosts. However, as you might expect, not all of these states will be physical. After correctly accounting for the gauge symmetry, only some subset survives.

The elegant method to determine the physical Hilbert space in a gauge fixed action with ghosts is known as **BRST quantization**. You will learn about it in the “Advanced Quantum Field Theory” course where you will apply it to Yang-Mills theory. Although a correct construction of the string spectrum employs the BRST method, we won’t describe it here for lack of time. A very clear description of the general method and its application to the string can be found in Section 4.2 of Polchinski’s book.

Instead, we will make do with a poor man’s attempt to determine the spectrum of the string. Our strategy is to simply pretend that the ghosts aren’t there and focus on the states created by the fields of the matter CFT (i.e. the $X^\mu$ fields if we’re talking about flat space). As we’ll explain in the next section, if we’re only interested in tree-level scattering amplitudes then this will suffice.

To illustrate how to compute the spectrum of the string, let’s go back to flat $D = 26$ dimensional Minkowski space and the discussion of covariant quantization in Section

\textsuperscript{7}A good review can be found Seiberg’s article “Notes on Quantum Liouville Theory and Quantum Gravity”, Prog. Theor. Phys. Suppl. 102 (1990) 319.
2.1. We found that physical states $|\Psi\rangle$ are subject to the Virasoro constraints (2.6) and (2.7) which read

$$L_n |\Psi\rangle = 0 \quad \text{for } n > 0$$
$$L_0 |\Psi\rangle = c |\Psi\rangle$$

and similar for $\tilde{L}_n$,

$$\tilde{L}_n |\Psi\rangle = 0 \quad \text{for } n > 0$$
$$\tilde{L}_0 |\Psi\rangle = \tilde{c} |\Psi\rangle$$

where we have, just briefly, allowed for the possibility of different normal ordering coefficients $c$ and $\tilde{c}$ for the left- and right-moving sectors. But there’s a name for states in a conformal field theory obeying these requirements: they are primary states of weight $(c, \tilde{c})$.

So how do we fix the normal ordering ambiguities $c$ and $\tilde{c}$? A simple way is to first replace the states with operator insertions on the worldsheet using the state-operator map: $|\Psi\rangle \to \mathcal{O}$. But we have a further requirement on the operators $\mathcal{O}$: gauge invariance. There are two gauge symmetries: reparameterization invariance and Weyl symmetry. Both restrict the possible states.

Let’s start by considering reparameterization invariance. In the last section, we happily placed operators at specific points on the worldsheet. But in a theory with a dynamical metric, this doesn’t give rise to a diffeomorphism invariant operator. To make an object that is invariant under reparameterizations of the worldsheet coordinates, we should integrate over the whole worldsheet. Our operator insertions (in conformal gauge) are therefore of the form,

$$V \sim \int d^2z \, \mathcal{O}$$

(5.7)

Here the $\sim$ sign reflects the fact that we’ve dropped an overall normalization constant which we’ll return to in the next section.

Integrating over the worldsheet takes care of diffeomorphisms. But what about Weyl symmetries? The measure $d^2z$ has weight $(-1, -1)$ under rescaling. To compensate, the operator $\mathcal{O}$ must have weight $(+1, +1)$. This is how we fix the normal ordering ambiguity: we require $c = \tilde{c} = 1$. Note that this agrees with the normal ordering coefficient $a = 1$ that we derived in lightcone quantization in Chapter 2.
This, then, is the rather rough derivation of the string spectrum. The physical states are the primary states of the CFT with weight $(+1, +1)$. The operators (5.7) associated to these states are called *vertex operators*.

### 5.4.1 An Example: Closed Strings in Flat Space

Let's use this new language to rederive the spectrum of the closed string in flat space. We start with the ground state of the string, which was previously identified as a tachyon. As we saw in Section 4, the vacuum of a CFT is associated to the identity operator. But we also have the zero modes. We can give the string momentum $p^\mu$ by acting with the operator $e^{ip \cdot X}$. The vertex operator associated to the ground state of the string is therefore

$$V_{\text{tachyon}} \sim \int d^2 z : e^{ip \cdot X} :$$  \hspace{1cm} (5.8)

In Section 4.3.3, we showed that the operator $e^{ip \cdot X}$ is primary with weight $h = \tilde{h} = \alpha' p^2 / 4$. But Weyl invariance requires that the operator has weight $(+1, +1)$. This is only true if the mass of the state is

$$M^2 \equiv -p^2 = -\frac{4}{\alpha'}$$

This is precisely the mass of the tachyon that we saw in Section 2.

Let's now look at the first excited states. In covariant quantization, these are of the form $\zeta_{\mu \nu} \alpha^\mu_{-1} \tilde{\alpha}^\nu_{-1} |0; p\rangle$, where $\zeta_{\mu \nu}$ is a constant tensor that determines the type of state, together with its polarization. (Recall: traceless symmetric $\zeta_{\mu \nu}$ corresponds to the graviton, anti-symmetric $\zeta_{\mu \nu}$ corresponds to the $B_{\mu \nu}$ field and the trace of $\zeta_{\mu \nu}$ corresponds to the scalar known as the dilaton). From (4.56), the vertex operator associated to this state is,

$$V_{\text{excited}} \sim \int d^2 z : e^{ip \cdot X} \partial X^\mu \bar{\partial} X^\nu : \zeta_{\mu \nu}$$ \hspace{1cm} (5.9)

where $\partial X^\mu$ gives us a $\alpha^\mu_{-1}$ excitation, while $\bar{\partial} X^\mu$ gives a $\tilde{\alpha}^\mu_{-1}$ excitation. It’s easy to check that the weight of this operator is $h = \tilde{h} = 1 + \alpha' p^2 / 4$. Weyl invariance therefore requires that

$$p^2 = 0$$

confirming that the first excited states of the string are indeed massless. However, we still need to check that the operator in (5.9) is actually primary. We know that $\partial X$ is
primary and we know that \( e^{ip \cdot X} \) is primary, but now we want to consider them both sitting together inside the normal ordering. This means that there are extra terms in the Wick contraction which give rise to \( 1/(z-w)^3 \) terms in the OPE, potentially ruining the primacy of our operator. One such term arises from a double contraction, one of which includes the \( e^{ip \cdot X} \) operator. This gives rise to an offending term proportional to \( p^\mu \zeta_{\mu\nu} \). The same kind of contraction with \( \bar{T} \) gives rise to a term proportional to \( p^\nu \zeta_{\mu\nu} \).

In order for these terms to vanish, the polarization tensor must satisfy

\[
p^\mu \zeta_{\mu\nu} = p^\nu \zeta_{\mu\nu} = 0
\]

which is precisely the transverse polarization condition expected for a massless particle.

### 5.4.2 An Example: Open Strings in Flat Space

As explained in Section 4.7, vertex operators for the open-string are inserted on the boundary \( \partial M \) of the worldsheet. We still need to ensure that these operators are diffeomorphism invariant which is achieved by integrating over \( \partial M \). The vertex operator for the open string tachyon is

\[
V_{\text{tachyon}} \sim \int_{\partial M} ds : e^{ip \cdot X} :
\]

We need to figure out the dimension of the boundary operator \( : e^{ip \cdot X} : \). It’s not the same as for the closed string. The reason is due to presence of the image charge in the propagator (4.57) for a free scalar field on a space with boundary. This propagator appears in the Wick contractions in the OPEs and affects the weights. Let’s see why this is the case. Firstly, we look at a single scalar field \( X \),

\[
\partial X(z) : e^{ipX(w,\bar{w})} : = \sum_{n=1}^{\infty} \frac{(ip)^n}{(n-1)!} : X(w,\bar{w})^{n-1} : \left( -\frac{\alpha'}{2} \frac{1}{z-w} - \frac{\alpha'}{2} \frac{1}{z-\bar{w}} \right) + \ldots
\]

\[
= -\frac{i\alpha'p}{2} : e^{ipX(w,\bar{w})} : \left( \frac{1}{z-w} + \frac{1}{z-\bar{w}} \right) + \ldots
\]

With this result, we can now compute the OPE with \( T \),

\[
T(z) : e^{ipX(w,\bar{w})} : = \frac{\alpha'p^2}{4} : e^{ipX} : \left( \frac{1}{z-w} + \frac{1}{z-\bar{w}} \right)^2 + \ldots
\]

When the operator \( : e^{ipX(w,\bar{w})} : \) is placed on the boundary \( w = \bar{w} \), this becomes

\[
T(z) : e^{ipX(w,\bar{w})} : = \frac{\alpha'p^2}{(z-w)^2} : e^{ipX(w,\bar{w})} : + \ldots
\]

This tells us that the boundary operator \( : e^{ip \cdot X} : \) is indeed primary, with weight \( \alpha'p^2 \).
For the open string, Weyl invariance requires that operators have weight +1 in order to cancel the scaling dimension of −1 coming from the boundary integral $\int ds$. So the mass of the open string ground state is

$$M^2 \equiv -p^2 = -\frac{1}{\alpha'}$$

in agreement with the mass of the open string tachyon computed in Section 3.

The vertex operator for the photon is

$$V_{\text{photon}} \sim \int_{\partial \mathcal{M}} ds \, \zeta_a : \partial X^a e^{ip \cdot X} :$$

(5.10)

where the index $a = 0, \ldots, p$ now runs only over those directions with Neumann boundary conditions that lie parallel to the brane worldvolume. The requirement that this is a primary operator gives $p^a \zeta_a = 0$, while Weyl invariance tells us that $p^2 = 0$. This is the expected behaviour for the momentum and polarization of a photon.

### 5.4.3 More General CFTs

Let’s now consider a string propagating in four-dimensional Minkowski space $\mathcal{M}_4$, together with some internal CFT with $c = 22$. Then any primary operator of the internal CFT with weight $(h, h)$ can be assigned momentum $p^\mu$, for $\mu = 0, 1, 2, 3$ by dressing the operator with $e^{ip \cdot X}$. In order to get a primary operator of weight $(+1, +1)$ as required, we must have

$$\frac{\alpha' p^2}{4} = 1 - h$$

We see that the mass spectrum of closed string states is given by

$$M^2 = \frac{4}{\alpha'} (h - 1)$$

where $h$ runs over the spectrum of primary operators of the internal CFT. Some comments:

- Relevant operators in the internal CFT have $h < 1$ and give rise to tachyons in the spectrum. Marginal operators, with $h = 1$, give massless particles. And irrelevant operators result in massive states.

- Notice that requiring the vertex operators to be Weyl invariant determines the mass formula for the state. We say that the vertex operators are “on-shell”, in the same sense that external legs of Feynman diagrams are on-shell. We will have more to say about this in the next section.
6. String Interactions

So far, despite considerable effort, we’ve only discussed the free string. We now wish to consider interactions. If we take the analogy with quantum field theory as our guide, then we might be led to think that interactions require us to add various non-linear terms to the action. However, this isn’t the case. Any attempt to add extra non-linear terms for the string won’t be consistent with our precious gauge symmetries. Instead, rather remarkably, all the information about interacting strings is already contained in the free theory described by the Polyakov action. (Actually, this statement is almost true).

To see that this is at least feasible, try to draw a cartoon picture of two strings interacting. It looks something like the worldsheet shown in the figure. The worldsheet is smooth. In Feynman diagrams in quantum field theory, information about interactions is inserted at vertices, where different lines meet. Here there are no such points. Locally, every part of the diagram looks like a free propagating string. Only globally do we see that the diagram describes interactions.

6.1 What to Compute?

If the information about string interactions is already contained in the Polyakov action, let’s go ahead and compute something! But what should we compute? One obvious thing to try is the probability for a particular configuration of strings at an early time to evolve into a new configuration at some later time. For example, we could try to compute the amplitude associated to the diagram above, stipulating fixed curves for the string ends.

No one knows how to do this. Moreover, there are words that we can drape around this failure that suggests this isn’t really a sensible thing to compute. I’ll now try to explain these words. Let’s start by returning to the familiar framework of quantum field theory in a fixed background. There the basic objects that we can compute are correlation functions,

\[ \langle \phi(x_1) \cdots \phi(x_n) \rangle \quad (6.1) \]

After a Fourier transform, these describe Feynman diagrams in which the external legs carry arbitrary momenta. For this reason, they are referred to as off-shell. To get the scattering amplitudes, we simply need to put the external legs on-shell (and perform a few other little tricks captured in the LSZ reduction formula).
The discussion above needs amendment if we turn on gravity. Gravity is a gauge theory and the gauge symmetries are diffeomorphisms. In a gauge theory, only gauge invariant observables make sense. But the correlation function (6.1) is not gauge invariant because its value changes under a diffeomorphism which maps the points $x_i$ to another point. This emphasizes an important fact: there are no local off-shell gauge invariant observables in a theory of gravity.

There is another way to say this. We know, by causality, that space-like separated operators should commute in a quantum field theory. But in gravity the question of whether operators are space-like separated becomes a dynamical issue and the causal structure can fluctuate due to quantum effects. This provides another reason why we are unable to define local gauge invariant observables in any theory of quantum gravity.

Let’s now return to string theory. Computing the evolution of string configurations for a finite time is analogous to computing off-shell correlation functions in QFT. But string theory is a theory of gravity so such things probably don’t make sense. For this reason, we retreat from attempting to compute correlation functions, back to the S-matrix.

**The String S-Matrix**

The object that we can compute in string theory is the S-matrix. This is obtained by taking the points in the correlation function to infinity: $x_i \to \infty$. This is acceptable because, just like in the case of QED, the redundancy of the system consists of those gauge transformations which die off asymptotically. Said another way, points on the boundary don’t fluctuate in quantum gravity. (Such fluctuations would be over an infinite volume of space and are suppressed due to their infinite action).

So what we’re really going to calculate is a diagram of the type shown in the figure, where all external legs are taken to infinity. Each of these legs can be placed in a different state of the free string and assigned some spacetime momentum $p_i$. The resulting expression is the string $S$-matrix.

Using the state-operator map, we know that each of these states at infinity is equivalent to the insertion of an appropriate vertex operator on the worldsheet. Therefore, to compute this S-matrix element we use a conformal transformation to bring each of these infinite legs to a finite distance. The end result is a worldsheet with the topology
of the sphere, dotted with vertex operators where the legs used to be. However, we already saw in the previous section that the constraint of Weyl invariance meant that vertex operators are necessarily on-shell. Technically, this is the reason that we can only compute on-shell correlation functions in string theory.

6.1.1 Summing Over Topologies

The Polyakov path integral instructs us to sum over all metrics. But what about worldsheets of different topologies? In fact, we should also sum over these. It is this sum that gives the perturbative expansion of string theory. The scattering of two strings receives contributions from worldsheets of the form

\[
\begin{align*}
\text{Figure 33:} \\
\quad + \\
\quad + \\
\quad + \\
\end{align*}
\]

The only thing that we need to know is how to weight these different worldsheets. Thankfully, there is a very natural coupling on the string that we have yet to consider and this will do the job. We augment the Polyakov action by

\[
S_{\text{string}} = S_{\text{Poly}} + \lambda \chi 
\]

Here \( \lambda \) is simply a real number, while \( \chi \) is given by an integral over the (Euclidean) worldsheet

\[
\chi = \frac{1}{4\pi} \int d^2 \sigma \sqrt{g} R 
\]

where \( R \) is the Ricci scalar of the worldsheet metric. This looks like the Einstein-Hilbert term for gravity on the worldsheet. It is simple to check that it is invariant under reparameterizations and Weyl transformations.

In four-dimensions, the Einstein-Hilbert term makes gravity dynamical. But life is very different in 2d. Indeed, we’ve already seen that all the components of the metric can be gauged away so there are no propagating degrees of freedom associated to \( g_{\alpha\beta} \). So, in two-dimensions, the term (6.4) doesn’t make gravity dynamical: in fact, classically, it doesn’t do anything at all!
The reason for this is that $\chi$ is a topological invariant. This means that it doesn’t actually depend on the metric $g_{\alpha\beta}$ at all – it depends only on the topology of the worldsheet. (More precisely, $\chi$ only depends on those global properties of the metric which themselves depend on the topology of the worldsheet). This is the content of the Gauss-Bonnet theorem: the integral of the Ricci scalar $R$ over the worldsheet gives an integer, $\chi$, known as the Euler number of the worldsheet. For a worldsheet without boundary (i.e. for the closed string) $\chi$ counts the number of handles $h$ on the worldsheet. It is given by,

$$\chi = 2 - 2h = 2(1 - g)$$

(6.5)

where $g$ is called the genus of the surface. The simplest examples are shown in the figure. The sphere has $g = 0$ and $\chi = 2$; the torus has $g = 1$ and $\chi = 0$. For higher $g > 1$, the Euler character $\chi$ is negative.

![Figure 34: Examples of increasingly poorly drawn Riemann surfaces with $\chi = 2$, 0 and −2.](image)

Now we see that the number $\lambda$ — or, more precisely, $e^\lambda$ — plays the role of the string coupling. The integral over world sheets is weighted by,

$$\sum_{\text{topologies}} e^{-S_{\text{string}}} \sim \sum_{\text{topologies}} e^{-2\lambda(1-g)} \int \mathcal{D}X \mathcal{D}g \ e^{-S_{\text{poly}}}$$

For $e^\lambda \ll 1$, we have a good perturbative expansion in which we sum over all topologies. (In fact, it is an asymptotic expansion, just as in quantum field theory). It is standard to define the string coupling constant as

$$g_s = e^\lambda$$

After a conformal map, tree-level scattering corresponds to a worldsheet with the topology of a sphere: the amplitudes are proportional to $1/g_s^2$. One-loop scattering corresponds to toroidal world sheets and, with our normalization, have no power of $g_s$. (Although, obviously, these are suppressed by $g_s^2$ relative to tree-level processes). The end
result is that the sum over worldsheets in (6.2) becomes a sum over Riemann surfaces of increasing genus, with vertex operators inserted for the initial and final states,

\[ \int \mathcal{D}X \mathcal{D}g \, e^{-S_{\text{poly}}} \prod_{i=1}^{m} V_{\Lambda_i}(p_i) \]

The Riemann surface of genus \( g \) is weighted by \( (g_s^2)^{g-1} \)

While it may look like we’ve introduced a new parameter \( g_s \) into the theory and added the coupling (6.3) by hand, we will later see why this coupling is a necessary part of the theory and provide an interpretation for \( g_s \).

**Scattering Amplitudes**

We now have all the information that we need to explain how to compute string scattering amplitudes. Suppose that we want to compute the S-matrix for \( m \) states: we will label them as \( \Lambda_i \) and assign them spacetime momenta \( p_i \). Each has a corresponding vertex operator \( V_{\Lambda_i}(p_i) \). The S-matrix element is then computed by evaluating the correlation function in the 2d conformal field theory, with insertions of the vertex operators.

\[ A^{(m)}(\Lambda_i, p_i) = \sum_{\text{topologies}} g_s^{-\chi} \frac{1}{\text{Vol}} \int \mathcal{D}X \mathcal{D}g \, e^{-S_{\text{poly}}} \prod_{i=1}^{m} V_{\Lambda_i}(p_i) \]

This is a rather peculiar equation. We are interpreting the correlation functions of a two-dimensional theory as the S-matrix for a theory in \( D = 26 \) dimensions!

To properly compute the correlation function, we should introduce the \( b \) and \( c \) ghosts that we saw in the last chapter and treat them carefully. However, if we’re only interested in tree-level amplitudes, then we can proceed naively and ignore the ghosts. The reason can be seen in the ghost action (5.5) where we see that the ghosts couple only to the worldsheet metric, not to the other worldsheet fields. This means that if our gauge fixing procedure fixes the worldsheet metric completely — which it does for worldsheets with the topology of a sphere — then we can forget about the ghosts. (At least, we can forget about them as soon as we’ve made sure that the Weyl anomaly cancels). However, as we’ll explain in 6.4, for higher genus worldsheets, the gauge fixing does not fix the metric completely and there are residual dynamical modes of the metric, known as moduli, which couple the ghosts and matter fields. This is analogous to the statement in field theory that we only need to worry about ghosts running in loops.
6.2 Closed String Amplitudes at Tree Level

The tree-level scattering amplitude is given by the correlation function of the 2d theory, evaluated on the sphere,

\[ A^{(m)} = \frac{1}{g_s^2} \frac{1}{\text{Vol}} \int D X D g \ e^{-S_{pol}} \prod_{i=1}^{m} V_{\Lambda_i}(p_i) \]

where \( V_{\Lambda_i}(p_i) \) are the vertex operators associated to the states.

We want to integrate over all metrics on the sphere. At first glance that sounds rather daunting but, of course, we have the gauge symmetries of diffeomorphisms and Weyl transformations at our disposal. Any metric on the sphere is conformally equivalent to the flat metric on the plane. For example, the round metric on the sphere of radius \( R \) can be written as

\[ ds^2 = \frac{4R^2}{(1 + |z|^2)^2} \, dzd\bar{z} \]

which is manifestly conformally equivalent to the plane, supplemented by the point at infinity. The conformal map from the sphere to the plane is the stereographic projection depicted in the diagram. The south pole of the sphere is mapped to the origin; the north pole is mapped to the point at infinity. Therefore, instead of integrating over all metrics, we may gauge fix diffeomorphisms and Weyl transformations to leave ourselves with the seemingly easier task of computing correlation functions on the plane.

6.2.1 Remnant Gauge Symmetry: \( \text{SL}(2,\mathbb{C}) \)

There’s a subtlety. And it’s a subtlety that we’ve seen before: there is a residual gauge symmetry. It is the conformal group, arising from diffeomorphisms which can be undone by Weyl transformations. As we saw in Section 4, there are an infinite number of such conformal transformations. It looks like we have a whole lot of gauge fixing still to do.

However, global issues actually mean that there’s less remnant gauge symmetry than you might think. In Section 4, we only looked at infinitesimal conformal transformations, generated by the Virasoro operators \( L_n, n \in \mathbb{Z} \). We did not examine whether these transformations are well-defined and invertible over all of space. Let’s take a
look at this. Recall that the coordinate changes associated to $L_n$ are generated by the vector fields (4.49),

$$l_n = z^{n+1} \partial_z$$

which result in the shift $\delta z = \epsilon z^{n+1}$. This is non-singular at $z = 0$ only for $n \geq -1$. If we restrict to smooth maps, that gets rid of half the transformations right away. But, since we’re ultimately interested in the sphere, we now also need to worry about the point at $z = \infty$ which, in stereographic projection, is just the north pole of the sphere. To do this, it’s useful to work with the coordinate

$$u = \frac{1}{z}$$

The generators of coordinate transformations for the $u$ coordinate are

$$l_n = z^{n+1} \partial_z = \frac{1}{u^{n+1}} \frac{\partial u}{\partial z} \partial_u = -u^{1-n} \partial_u$$

which is non-singular at $u = 0$ only for $n \leq 1$.

Combining these two results, the only generators of the conformal group that are non-singular over the whole Riemann sphere are $l_{-1}$, $l_0$ and $l_1$ which act infinitesimally as

$$l_{-1} : \quad z \rightarrow z + \epsilon$$
$$l_0 : \quad z \rightarrow (1 + \epsilon) z$$
$$l_1 : \quad z \rightarrow (1 + \epsilon z) z$$

The global version of these transformations is

$$l_{-1} : \quad z \rightarrow z + \alpha$$
$$l_0 : \quad z \rightarrow \lambda z$$
$$l_1 : \quad z \rightarrow \frac{z}{1 - \beta z}$$

which can be combined to give the general transformation

$$z \rightarrow \frac{az + b}{cz + d} \quad (6.6)$$

with $a, b, c$ and $d \in \mathbb{C}$. We have four complex parameters, but we’ve only got three transformations. What happened? Well, one transformation is fake because an overall
scaling of the parameters doesn’t change \( z \). By such a rescaling, we can always insist that the parameters obey

\[
ad - bc = 1
\]

The transformations (6.6) subject to this constraint have the group structure \( SL(2; C) \), which is the group of \( 2 \times 2 \) complex matrices with unit determinant. In fact, since the transformation is blind to a flip in sign of all the parameters, the actual group of global conformal transformations is \( SL(2; C)/\mathbb{Z}_2 \), which is sometimes written as \( PSL(2; C) \). (This \( \mathbb{Z}_2 \) subtlety won’t be important for us in what follows).

The remnant global transformations on the sphere are known as conformal Killing vectors and the group \( SL(2; C)/\mathbb{Z}_2 \) is the conformal Killing group. This group allows us to take any three points on the plane and move them to three other points of our choosing. We will shortly make use of this fact to gauge fix, but for now we leave the \( SL(2; C) \) symmetry intact.

6.2.2 The Virasoro-Shapiro Amplitude

We will now compute the S-matrix for closed string tachyons. You might think that this is the least interesting thing to compute: after all, we’re ultimately interested in the superstring which doesn’t have tachyons. This is true, but it turns out that tachyon scattering is much simpler than everything else, mainly because we don’t have a plethora of extra indices on the states to worry about. Moreover, the lessons that we will learn from tachyon scattering hold for the scattering of other states as well.

The \( m \)-point tachyon scattering amplitude is given by the flat space correlation function

\[
A^{(m)}(p_1, \ldots, p_m) = \frac{1}{g_s^2 \text{Vol}(SL(2; C))} \int D X \ e^{-S_{\text{poly}}} \prod_{i=1}^{m} V(p_i)
\]

where the tachyon vertex operator is given by,

\[
V(p_i) = g_s \int d^2 z \ e^{ip_i \cdot X} \equiv g_s \int d^2 z \ \hat{V}(z, p_i)
\]

Note that, in contrast to (5.8), we’ve added an appropriate normalization factor to the vertex operator. Heuristically, this reflects the fact that the operator is associated to the addition of a closed string mode. A rigorous derivation of this normalization can be found in Polchinski.
The amplitude can therefore be written as,

$$ \mathcal{A}^{(m)}(p_1, \ldots, p_m) = \frac{g^m}{\text{Vol}(SL(2; \mathbb{C}))} \int \prod_{i=1}^{m} d^2 z_i \langle \hat{V}(z_1, p_1) \ldots \hat{V}(z_m, p_m) \rangle $$

where the expectation value \langle \ldots \rangle is computed using the gauge fixed Polyakov action. But the gauge fixed Polyakov action is simply a free theory and our correlation function is something eminently computable: a Gaussian integral,

$$ \langle \hat{V}(z_1, p_1) \ldots \hat{V}(z_m, p_m) \rangle = \int D\mathcal{X} \exp \left( -\frac{1}{16\pi} \int d^2 z \partial \mathcal{X} \cdot \overline{\partial} \mathcal{X} + iJ \cdot \mathcal{X} \right) \exp \left( i \sum_{i=1}^{m} p_i \cdot X(z_i, \bar{z}_i) \right) $$

The normalization in front of the Polyakov action is now $1/2\pi\alpha'$ instead of $1/4\pi\alpha'$ because we’re working with complex coordinates and we need to remember that $\partial_\alpha \overline{\partial}^\alpha = 4\overline{\partial} \partial$ and $d^2 z = 2d^2 \sigma$.

**The Gaussian Integral**

We certainly know how to compute Gaussian integrals. Let’s go slow. Consider the following general integral,

$$ \int D\mathcal{X} \exp \left( \int d^2 z \frac{1}{2\pi \alpha'} X \cdot \overline{\partial} \partial X + iJ \cdot X \right) \sim \exp \left( \frac{\pi \alpha'}{2} \int d^2 zd^2 z' J(z, \bar{z}) \frac{1}{\partial \overline{\partial}} J(z', \bar{z}') \right) $$

Here the $\sim$ symbol reflects the fact that we’ve dropped a whole lot of irrelevant normalization terms, including $\text{det}^{-1/2}(-\partial \overline{\partial})$. The inverse operator $1/\partial \overline{\partial}$ on the right-hand-side of this equation is shorthand for the propagator $G(z, z')$ which solves

$$ \partial \overline{\partial} G(z, \bar{z}; z', \bar{z}') = \delta(z - z', \bar{z} - \bar{z}') $$

As we’ve seen several times before, in two dimensions this propagator is given by

$$ G(z, \bar{z}; z', \bar{z}') = \frac{1}{2\pi} \ln |z - z'|^2 $$

**Back to the Scattering Amplitude**

Comparing our scattering amplitude with this general expression, we need to take the source $J$ to be

$$ J(z, \bar{z}) = \sum_{i=1}^{m} p_i \delta(z - z_i, \bar{z} - \bar{z}_i) $$
Inserting this into the Gaussian integral gives us an expression for the amplitude
\[
A^{(m)} \sim \frac{g_s^{m-2}}{\text{Vol}(SL(2; \mathbb{C}))} \int \prod_{i=1}^{m} d^2z_i \exp \left( \frac{\alpha'}{2} \sum_{j,l} p_j \cdot p_l \ln |z_j - z_l| \right)
\]
The terms with \( j = l \) seem to be problematic. In fact, they should just be left out. This follows from correctly implementing normal ordering and leaves us with
\[
A^{(m)} \sim \frac{g_s^{m-2}}{\text{Vol}(SL(2; \mathbb{C}))} \int \prod_{i=1}^{m} d^2z_i \prod_{j<l} |z_j - z_l|^{\alpha' p_j \cdot p_l} \tag{6.8}
\]
Actually, there’s something that we missed. (Isn’t there always!). We certainly expect scattering in flat space to obey momentum conservation, so there should be a \( \delta^{26}(\sum_{i=1}^{m} p_i) \) in the amplitude. But where is it? We missed it because we were a little too quick in computing the Gaussian integral. The operator \( \partial \bar{\partial} \) annihilates the zero mode, \( x^\mu \), in the mode expansion. This means that its inverse, \( 1/\partial \bar{\partial} \), is not well-defined.

But it’s easy to deal with this by treating the zero mode separately. The derivatives \( \partial^2 \) don’t see \( x^\mu \), but the source \( J \) does. Integrating over the zero mode in the path integral gives us our delta function
\[
\int dx \exp(i \sum_{i=1}^{m} p_i \cdot x) \sim \delta^{26}(\sum_{i=1}^{m} p_i)
\]

So, our final result for the amplitude is
\[
A^{(m)} \sim \frac{g_s^{m-2}}{\text{Vol}(SL(2; \mathbb{C}))} \delta^{26}(\sum_{i=1}^{m} p_i) \int \prod_{i=1}^{m} d^2z_i \prod_{j<l} |z_j - z_l|^{\alpha' p_j \cdot p_l} \tag{6.9}
\]

The Four-Point Amplitude

We will compute only the four-point amplitude for two-to-two scattering of tachyons. The \( \text{Vol}(SL(2; \mathbb{C})) \) factor is there to remind us that we still have a remnant gauge symmetry floating around. Let’s now fix this. As we mentioned before, it provides enough freedom for us to take any three points on the plane and move them to any other three points. We will make use of this to set
\[
z_1 = \infty \quad , \quad z_2 = 0 \quad , \quad z_3 = z \quad , \quad z_4 = 1
\]

Inserting this into the amplitude (6.9), we find ourselves with just a single integral to evaluate,
\[
A^{(4)} \sim g_s^2 \delta^{26}(\sum_{i} p_i) \int d^2z \ |z|^{\alpha' p_2 \cdot p_1} |1 - z|^{\alpha' p_3 \cdot p_4} \tag{6.10}
\]
(There is also an overall factor of \(|z_1|^4\), but this just gets absorbed into an overall normalization constant). We still need to do the integral. It can be evaluated exactly in terms of gamma functions. We relegate the proof to Appendix 6.5, where we show that

\[ \int d^2 z \, |z|^{2a-2}|1-z|^{2b-2} = \frac{2\pi \Gamma(a)\Gamma(b)\Gamma(c)}{\Gamma(1-a)\Gamma(1-b)\Gamma(1-c)} \]  

(6.11)

where \(a + b + c = 1\).

Four-point scattering amplitudes are typically expressed in terms of Mandelstam variables. We choose \(p_1\) and \(p_2\) to be incoming momenta and \(p_3\) and \(p_4\) to be outgoing momenta, as shown in the figure. We then define

\[ s = -(p_1 + p_2)^2 \quad , \quad t = -(p_1 + p_3)^2 \quad , \quad u = -(p_1 + p_4)^2 \]

These obey

\[ s + t + u = -\sum_i p_i^2 = \sum_i M_i^2 = -\frac{16}{\alpha'} \]

where, in the last equality, we’ve inserted the value of the tachyon mass \((2.27)\). Writing the scattering amplitude (6.10) in terms of Mandelstam variables, we have our final answer

\[ A^{(4)} \sim g_s^2 \delta^{26}(\sum_i p_i) \frac{\Gamma(-1-\alpha's/4)\Gamma(-1-\alpha't/4)\Gamma(-1-\alpha'u/4)}{\Gamma(2+\alpha's/4)\Gamma(2+\alpha't/4)\Gamma(2+\alpha'u/4)} \]  

(6.12)

This is the Virasoro-Shapiro amplitude governing tachyon scattering in the closed bosonic string.

Remarkably, the Virasoro-Shapiro amplitude was almost the first equation of string theory! (That honour actually goes to the Veneziano amplitude which is the analogous expression for open string tachyons and will be derived in Section 6.3.1). These amplitudes were written down long before people knew that they had anything to do with strings: they simply exhibited some interesting and surprising properties. It took several years of work to realise that they actually describe the scattering of strings. We will now start to tease apart the Virasoro-Shapiro amplitude to see some of the properties that got people hooked many years ago.

### 6.2.3 Lessons to Learn

So what’s the physics lying behind the scattering amplitude (6.12)? Obviously it is symmetric in \(s\), \(t\) and \(u\). That is already surprising and we’ll return to it shortly. But we’ll start by fixing \(t\) and looking at the properties of the amplitude as we vary \(s\).
The first thing to notice is that $A^{(4)}$ has poles. Lots of poles. They come from the factor of $\Gamma(-1 - \alpha' s/4)$ in the numerator. The first of these poles appears when

$$-1 - \frac{\alpha' s}{4} = 0 \quad \Rightarrow \quad s = -\frac{4}{\alpha'}$$

But that’s the mass of the tachyon! It means that, for $s$ close to $-4/\alpha'$, the amplitude has the form of a familiar scattering amplitude in quantum field theory with a cubic vertex,

$$\sim \frac{1}{s - M^2}$$

where $M$ is the mass of the exchanged particle, in this case the tachyon.

Other poles in the amplitude occur at $s = 4(n - 1)/\alpha'$ with $n \in \mathbb{Z}^+$. This is precisely the mass formula for the higher states of the closed string. What we’re learning is that the string amplitude is summing up an infinite number of tree-level field theory diagrams,

$$= \sum_n M_n$$

where the exchanged particles are all the different states of the free string.

In fact, there’s more information about the spectrum of states hidden within these amplitudes. We can look at the residues of the poles at $s = 4(n - 1)/\alpha'$, for $n = 0, 1, \ldots$. These residues are rather complicated functions of $t$, but the highest power of momentum that appears for each pole is

$$A^{(4)} \sim \sum_{n=0}^{\infty} \frac{t^{2n}}{s - M_n^2}$$

(6.13)

The power of the momentum is telling us the highest spin of the particle states at level $n$. To see why this is, consider a field corresponding to a spin $J$ particle. It has a whole bunch of Lorentz indices, $\chi_{\mu_1 \ldots \mu_J}$. In a cubic interaction, each of these must be soaked up by derivatives. So we have $J$ derivatives at each vertex, contributing powers of $(\text{momentum})^{2J}$ to the numerator of the Feynman diagram. Comparing with the string scattering amplitude, we see that the highest spin particle at level $n$ has $J = 2n$. This is indeed the result that we saw from the canonical quantization of the string in Section 2.
Finally, the amplitude (6.12) has a property that is very different from amplitudes in field theory. Above, we framed our discussion by keeping $t$ fixed and expanding in $s$. We could just have well done the opposite: fix $s$ and look at poles in $t$. Now the string amplitude has the interpretation of an infinite number of $t$-channel scattering amplitudes, one for each state of the string

$$\sum_n M_n$$

Usually in field theory, we sum up both $s$-channel and $t$-channel scattering amplitudes. Not so in string theory. The sum over an infinite number of $s$-channel amplitudes can be reinterpreted as an infinite sum of $t$-channel amplitudes. We don’t include both: that would be overcounting. (Similar statements hold for $u$). The fact that the same amplitude can be written as a sum over $s$-channel poles or a sum over $t$-channel poles is sometimes referred to as “duality”. (A much overused word). In the early days, before it was known that string theory was a theory of strings, the subject inherited its name from this duality property of amplitudes: it was called the dual resonance model.

**High Energy Scattering**

Let’s use this amplitude to see what happens when we collide strings at high energies. There are different regimes that we could look at. The most illuminating is $s, t \to \infty$, with $s/t$ held fixed. In this limit, all the exchanged momenta become large. It corresponds to high-energy scattering with the angle $\theta$ between incoming and outgoing particles kept fixed. To see this consider, for example, massless particles (our amplitude is really for tachyons, but the same considerations hold). We take the incoming and outgoing momenta to be

$$p_1 = \frac{\sqrt{s}}{2} (1,1,0,\ldots) , \quad p_2 = \frac{\sqrt{s}}{2} (1,-1,0,\ldots)$$

$$p_3 = \frac{\sqrt{s}}{2} (1,\cos \theta, \sin \theta, \ldots) , \quad p_4 = \frac{\sqrt{s}}{2} (1, -\cos \theta, -\sin \theta, \ldots)$$

Then we see explicitly that $s \to \infty$ and $t \to \infty$ with the ratio $s/t$ fixed also keeps the scattering angle $\theta$ fixed.
We can evaluate the scattering amplitude $A^{(4)}$ in this limit by using $\Gamma(x) \sim \exp(x \ln x)$. We send $s \to \infty$ avoiding the poles. (We can achieve this by sending $s \to \infty$ in a slightly imaginary direction. Ultimately this is valid because all the higher string states are actually unstable in the interacting theory which will shift their poles off the real axis once taken into account). It is simple to check that the amplitude drops off exponentially quickly at high energies,

$$A^{(4)} \sim g_s^2 \delta^{26}(\sum_i p_i) \exp\left(-\frac{\alpha'}{2} (s \ln s + t \ln t + u \ln u)\right) \quad \text{as } s \to \infty \quad (6.14)$$

The exponential fall-off seen in (6.14) is much faster than the amplitude of any field theory which, at best, fall off with power-law decay at high energies and, at worse, diverge. For example, consider the individual terms (6.13) corresponding to the amplitude for $s$-channel processes involving the exchange of particles with spin $2n$. We see that the exchange of a spin 2 particle results in a divergence in this limit. This is reflecting something you already know about gravity: the dimensionless coupling is $G_N E^2$ (in four-dimensions) which becomes large for large energies. The exchange of higher spin particles gives rise to even worse divergences. If we were to truncate the infinite sum (6.13) at any finite $n$, the whole thing would diverge. But infinite sums can do things that finite sums can’t and the final behaviour of the amplitude (6.14) is much softer than any of the individual terms. The infinite number of particles in string theory conspire to render finite any divergence arising from an individual particle species.

Phrased in terms of the $s$-channel exchange of particles, the high-energy behaviour of string theory seems somewhat miraculous. But there is another viewpoint where it’s all very obvious. The power-law behaviour of scattering amplitudes is characteristic of point-like charges. But, of course, the string isn’t a point-like object. It is extended and fuzzy at length scales comparable to $\sqrt{\alpha'}$. This is the reason the amplitude has such soft high-energy behaviour. Indeed, this idea that smooth extended objects give rise to scattering amplitudes that decay exponentially at high energies is something that you’ve seen before in non-relativistic quantum mechanics. Consider, for example, the scattering of a particle off a Gaussian potential. In the Born approximation, the differential cross-section is just given by the Fourier transform which is again a Gaussian, now decaying exponentially for large momentum.

It’s often said that theories of quantum gravity should have a “minimum length”, sometimes taken to be the Planck scale. This is roughly true in string theory, although not in any crude simple manner. Rather, the minimum length reveals itself in different
ways depending on which question is being asked. The above discussion highlights
one example of this: strings can’t probe distance scales shorter than \( l_s = \sqrt{\alpha'} \) simply
because they are themselves fuzzy at this scale. It turns out that D-branes are much
better probes of sub-stringy physics and provide a different view on the short distance
structure of spacetime. We will also see another manifestation of the minimal length
scale of string theory in Section 8.3.

**Graviton Scattering**

Although we’ve derived the result (6.14) for tachyons, all tree-level amplitudes have this
soft fall-off at high-energies. Most notably, this includes graviton scattering. As we
noted above, this is in sharp contrast to general relativity for which tree-level scattering
amplitudes diverge at high-energies. This is the first place to see that UV problems of
general relativity might have a good chance of being cured in string theory.

Using the techniques described in this section, one can compute \( m \)-point tree-level
amplitudes for graviton scattering. If we restrict attention to low-energies (i.e. much
smaller than \( 1/\sqrt{\alpha'} \)), one can show that these coincide with the amplitudes derived
from the Einstein-Hilbert action in \( D = 26 \) dimensions

\[
S = \frac{1}{2\kappa^2} \int d^{26}X \sqrt{-G} \mathcal{R}
\]

where \( \mathcal{R} \) is the \( D = 26 \) Ricci scalar (not to be confused with the worldsheet Ricci scalar
which we call \( R \)). The gravitational coupling, \( \kappa^2 \) is related to Newton’s constant in
26 dimensions. It plays no role for pure gravity, but is important when we couple to
matter. We’ll see shortly that it’s given by

\[
\kappa^2 \approx g_s^2 (\alpha')^{12}
\]

We won’t explicitly compute graviton scattering amplitudes in this course, partly be-
cause they’re fairly messy and partly because building up the Einstein-Hilbert action
from \( m \)-particle scattering is hardly the best way to look at general relativity. Instead,
we shall derive the Einstein-Hilbert action in a much better fashion in Section 7.

**6.3 Open String Scattering**

So far our discussion has been entirely about closed strings. There is a very similar
story for open strings. We again compute S-matrix elements. Conformal symmetry now
maps tree-level scattering to the disc, with vertex operators inserted on the boundary
of the disc.
For the open string, the string coupling constant that we add to the Polyakov action requires the addition of a boundary term to make it well defined,

$$\chi = \frac{1}{4\pi} \int_M d^2\sigma \sqrt{g} R + \frac{1}{2\pi} \int_{\partial M} ds k$$  \hspace{1cm} (6.15)

where $k$ is the geodesic curvature of the boundary. To define it, we introduce two unit vectors on the worldsheet: $t^\alpha$ is tangential to the boundary, while $n^\alpha$ is normal and points outward from the boundary. The geodesic curvature is defined as

$$k = -t^\alpha n_\beta \nabla_\alpha t^\beta$$

Boundary terms of the type seen in (6.15) are also needed in general relativity for manifolds with boundaries: in that context, they are referred to as Gibbons-Hawking terms.

The Gauss-Bonnet theorem has an extension to surfaces with boundary. For surfaces with $h$ handles and $b$ boundaries, the Euler character is given by

$$\chi = 2 - 2h - b$$

Some examples are shown in Figure 38. The expansion for open-string scattering consists of adding consecutive boundaries to the worldsheet. The disc is weighted by $1/g_s$; the annulus has no factor of $g_s$ and so on. We see that the open string coupling is related to the closed string coupling by

$$g_{\text{open}}^2 = g_s$$  \hspace{1cm} (6.16)

One of the key steps in computing closed string scattering amplitudes was the implementation of the conformal Killing group, which was defined as the surviving gauge
symmetry with a global action on the sphere. For the open string, there is again a residual gauge symmetry. If we think in terms of the upper-half plane, the boundary is \( \text{Im}z = 0 \). The conformal Killing group is composed of transformations

\[
z \rightarrow \frac{az + b}{cz + d}
\]

again with the requirement that \( ad - bc = 1 \). This time there is one further condition: the boundary \( \text{Im}z = 0 \) must be mapped onto itself. This requires \( a, b, c, d \in \mathbb{R} \). The resulting conformal Killing group is \( SL(2; \mathbb{R})/\mathbb{Z}_2 \).

### 6.3.1 The Veneziano Amplitude

Since vertex operators now live on the boundary, they have a fixed ordering. In computing a scattering amplitude, we must sum over all orderings. Let’s look again at the 4-point amplitude for tachyon scattering. The vertex operator is

\[
V(p_i) = \sqrt{g_s} \int dx \, e^{ip_i \cdot X}
\]

where the integral \( \int dx \) is now over the boundary and \( p^2 = 1/\alpha' \) is the on-shell condition for an open-string tachyon. The normalization \( \sqrt{g_s} \) is that appropriate for the insertion of an open-string mode, reflecting (6.16).

Going through the same steps as for the closed string, we find that the amplitude is given by

\[
\mathcal{A}^{(4)} \sim \frac{g_s}{\text{Vol}(SL(2; \mathbb{R}))} \delta^{26} \left( \sum_i p_i \right) \int dx_i \prod_{i=1}^{4} \prod_{j<l} |x_j - x_l|^{2\alpha' p_j \cdot p_l}
\]

\hspace{1cm} (6.17)

Note that there’s a factor of 2 in the exponent, differing from the closed string expression (6.8). This comes about because the boundary propagator (4.57) has an extra factor of 2 due to the image charge.
We now use the $SL(2; \mathbb{R})$ residual gauge symmetry to fix three points on the boundary. We choose a particular ordering and set $x_1 = 0$, $x_2 = x$, $x_3 = 1$ and $x_4 \to \infty$. The only free insertion point is $x_2 = x$ but, because of the restriction of operator ordering, this must lie in the interval $x \in [0, 1]$. The interesting part of the integral is then given by

$$A^{(4)} \sim g_s \int_0^1 dx \ |x|^{2\alpha' p_1 \cdot p_2} |1 - x|^{2\alpha' p_2 \cdot p_3}$$

This integral is well known: as shown in Appendix 6.5, it is the Euler beta function

$$B(a, b) = \int_0^1 dx \ x^{a-1} (1 - x)^{b-1} = \frac{\Gamma(a) \Gamma(b)}{\Gamma(a + b)}$$

After summing over the different orderings of vertex operators, the end result for the amplitude for open string tachyon scattering is,

$$A^{(4)} \sim g_s \left[ B(-\alpha' s - 1, -\alpha' t - 1) + B(-\alpha' s - 1, -\alpha' u - 1) + B(-\alpha' t - 1, -\alpha' u - 1) \right]$$

This is the famous Veneziano Amplitude, first postulated in 1968 to capture some observed features of the strong interactions. This was before the advent of QCD and before it was realised that the amplitude arises from a string.

The open string scattering amplitude contains the same features that we saw for the closed string. For example, it has poles at

$$s = \frac{n - 1}{\alpha'} \quad n = 0, 1, 2, \ldots$$

which we recognize as the spectrum of the open string.

### 6.3.2 The Tension of D-Branes

Recall that we introduced D-branes as surfaces in space on which strings can end. At the time, I promised that we would eventually discover that these D-branes are dynamical objects in their own right. We’ll look at this more closely in the next section, but for now we can do a simple computation to determine the tension of D-branes.

The tension $T_p$ of a $D_p$-brane is defined as the energy per spatial volume. It has dimension $[T_p] = p + 1$. The tension is telling us the magnitude of the coupling between the brane and gravity. Or, in our new language, the strength of the interaction between a closed string state and an open string. The simplest such diagram is shown in the figure, with a graviton vertex operator inserted. Although we won’t compute this
diagram completely, we can figure out its most important property just by looking at it: it has the topology of a disc, so is proportional to $1/g_s$. Adding powers of $\alpha'$ to get the dimension right, the tension of a $D_p$-brane must scale as

$$T_p \sim \frac{1}{l_s^{p+1}} \frac{1}{g_s}$$

where the string length is defined as $l_s = \sqrt{\alpha'}$. The $1/g_s$ scaling of the tension is one of the key characteristic features of a D-brane.

I should confess that there’s a lot swept under the carpet in the above discussion, not least the question of the correct normalization of the vertex operators and the difference between the string frame and the Einstein frame (which we will discuss shortly). Nonetheless, the end result (6.18) is correct. For a fuller discussion, see Section 8.7 of Polchinski.

6.4 One-Loop Amplitudes

We now return to the closed string to discuss one-loop effects. As we saw above, this corresponds to a worldsheet with the topology of a torus. We need to integrate over all metrics on the torus.

For tree-level processes, we used diffeomorphisms and Weyl transformations to map an arbitrary metric on the sphere to the flat metric on the plane. This time, we use these transformations to map an arbitrary metric on the torus to the flat metric on the torus. But there’s a new subtlety that arises: not all flat metrics on the torus are equivalent.

6.4.1 The Moduli Space of the Torus

Let’s spell out what we mean by this. We can construct a torus by identifying a region in the complex $z$-plane as shown in the figure. In general, this identification depends on a single complex parameter, $\tau \in \mathbb{C}$.

$$z \equiv z + 2\pi \quad \text{and} \quad z \equiv z + 2\pi \tau$$

Do not confuse $\tau$ with the Minkowski worldsheet time: we left that behind way back in Section 3. Everything here is Euclidean worldsheet and $\tau$ is just a parameter telling us how skewed the torus is. The flat metric on the torus is now simply

$$ds^2 = dz d\bar{z}$$

subject to the identifications above.
A general metric on a torus can always be transformed to a flat metric for some value of \( \tau \). But the question that interests us is whether two tori, parameterized by different \( \tau \), are conformally equivalent. In general, the answer is no. The space of conformally inequivalent tori, parameterized by \( \tau \), is called the \textit{moduli space} \( \mathcal{M} \).

However, there are some values of \( \tau \) that do correspond to the same torus. In particular, there are a couple of obvious ways in which we can change \( \tau \) without changing the torus. They go by the names of the \( S \) and \( T \) transformations:

- \( T : \tau \to \tau + 1 \): This clearly gives rise to the same torus, because the identification is now

  \[ z \equiv z + 2\pi \quad \text{and} \quad z \equiv z + 2\pi(\tau + 1) \equiv z + 2\pi\tau \]

- \( S : \tau \to -1/\tau \): This simply flips the sides of the torus. For example, if \( \tau = ia \) is purely imaginary, then this transformation maps \( \tau \to i/a \), which can then be undone by a scaling.

It turns out that these two changes \( S \) and \( T \) are the only ones that keep the torus intact. They are sometimes called \textit{modular transformations}. A general modular transformations is constructed from combinations of \( S \) and \( T \) and takes the form,

\[
\tau \to \frac{a\tau + b}{c\tau + d} \quad \text{with} \quad ad - bc = 1 \tag{6.19}
\]

where \( a, b, c \) and \( d \in \mathbb{Z} \). This is the group \( SL(2,\mathbb{Z}) \). (In fact, we have our usual \( \mathbb{Z}_2 \) identification and the group is actually \( PSL(2,\mathbb{Z}) = SL(2,\mathbb{Z})/\mathbb{Z}_2 \)). The moduli space \( \mathcal{M} \) of the torus is given by

\[
\mathcal{M} \cong \mathbb{C}/SL(2,\mathbb{Z})
\]

What does this space look like? Using \( T : \tau \to \tau + 1 \), we can always shift \( \tau \) until it lies within the interval

\[
\text{Re} \, \tau \in \left[-\frac{1}{2}, +\frac{1}{2}\right]
\]

where the edges of the interval are identified. Meanwhile, \( S : \tau \to -1/\tau \) inverts the
modulus $|\tau|$, so we can use this to map a point inside the circle $|\tau| < 1$ to a point outside $|\tau| > 1$. One can show that by successive combinations of $S$ and $T$, it is possible to map any point to lie within the shaded region shown in the figure, defined by

$$|\tau| \geq 1 \quad \text{and} \quad \text{Re} \, \tau \in \left[-\frac{1}{2}, +\frac{1}{2}\right]$$

This is referred to as the fundamental domain of $SL(2; \mathbb{Z})$.

We could have just as easily chosen one of the other fundamental domains shown in the figure. But the shaded region is the standard one.

**Integrating over the Moduli Space**

In string theory we’re invited to sum over all metrics. After gauge fixing diffeomorphisms and Weyl invariance, we still need to integrate over all inequivalent tori. In other words, we integrate over the fundamental domain. The $SL(2; \mathbb{Z})$ invariant measure over the fundamental domain is

$$\int \frac{d^2 \tau}{(\text{Im} \, \tau)^2}$$

To see that this is $SL(2; \mathbb{Z})$ invariant, note that under a general transformation of the form (6.19) we have

$$d^2 \tau \rightarrow \frac{d^2 \tau}{|c\tau + d|^4} \quad \text{and} \quad \text{Im} \, \tau \rightarrow \frac{\text{Im} \, \tau}{|c\tau + d|^2}$$

**Figure 41:** The fundamental domain.
There's some physics lurking within these rather mathematical statements. The integration over the fundamental domain in string theory is analogous to the loop integral over momentum in quantum field theory. Consider the square tori defined by $\text{Re} \, \tau = 0$. The tori with $\text{Im} \, \tau \to \infty$ are squashed and chubby. They correspond to the infra-red region of loop momenta in a Feynman diagram. Those with $\text{Im} \, \tau \to 0$ are long and thin. Those correspond to the ultra-violet limit of loop momenta in a Feynman diagram. Yet, as we have seen, we should not integrate over these UV regions of the loop since the fundamental domain does not stretch down that far. Or, more precisely, the thin tori are mapped to chubby tori. This corresponds to the fact that any putative UV divergence of string theory can always be reinterpreted as an IR divergence. This is the second manifestation of the well-behaved UV nature of string theory. We will see this more explicitly in the example of Section 6.4.2.

Finally, when computing a loop amplitude in string theory, we still need to worry about the residual gauge symmetry that is left unfixed after the map to the flat torus. In the case of tree-level amplitudes on the sphere, this residual gauge symmetry was due to the conformal Killing group $SL(2; \mathbb{C})$. For the torus, the conformal Killing group is generated by the obvious generators $\partial_z$ and $\bar{\partial}_z$. It is $U(1) \times U(1)$.

**Higher Genus Surfaces**

The moduli space $\mathcal{M}_g$ of the Riemann surface of genus $g > 1$ can be shown to have dimension,

$$\dim \mathcal{M}_g = 3g - 3$$

There are no conformal Killing vectors when $g > 1$. These facts can be demonstrated as an application of the Riemann-Roch theorem. For more details, see section 5.2 of Polchinski, or sections 3.3 and 8.2 of Green, Schwarz and Witten.

**6.4.2 The One-Loop Partition Function**

We won’t compute any one-loop scattering amplitudes in string theory. Instead, we will look at something a little simpler: the one-loop vacuum to vacuum amplitude. A Euclidean worldsheet with periodic time has the interpretation of a finite temperature partition function for the theory defined on a cylinder. In $D = 26$ dimensional spacetime, it is related to the cosmological constant in bosonic string theory.

Consider firstly the partition function of a theory on a square torus, with $\text{Re} \, \tau = 0$. Compactifying Euclidean time, with period $(\text{Im} \, \tau)$ is equivalent to putting the theory at temperature $T = 1/(\text{Im} \, \tau)$,

$$Z[\tau] = \text{Tr} \, e^{-2\pi(\text{Im} \, \tau)H}$$
where the Tr is over all states in the theory. For any CFT defined on a cylinder, the Hamiltonian given by

\[ H = L_0 + \tilde{L}_0 - \frac{c + \tilde{c}}{24} \]

where the final term is the Casimir energy computed in Section 4.4.1.

What then is the interpretation of the vacuum amplitude computed on a torus with \( \text{Re} \tau \neq 0 \)? From the diagram, we see that the effect of such a skewed torus is to translate a given point around the cylinder by \( \text{Re} \tau \). But we know which operator implements such a translation: it is \( \exp(2\pi i (\text{Re} \tau) P) \), where \( P \) is the momentum operator on the cylinder. After the map to the plane, this becomes the rotation operator

\[ P = L_0 - \tilde{L}_0 \]

So the vacuum amplitude on the torus has the interpretation of the sum over all states in the theory, weighted by

\[ Z[\tau] = \text{Tr} e^{-2\pi i (\text{Im} \tau) (L_0 + \tilde{L}_0)} e^{-2\pi i (\text{Re} \tau) (L_0 - \tilde{L}_0)} e^{2\pi (\text{Im} \tau) (c + \tilde{c})/24} \]

We define

\[ q = e^{2\pi i \tau}, \quad \bar{q} = e^{-2\pi i \bar{\tau}} \]

The partition function can then be written in slick notation as

\[ Z[\tau] = \text{Tr} q^{L_0 - c/24} \bar{q}^{\tilde{L}_0 - \tilde{c}/24} \]

Let’s compute this for the free string. We know that each scalar field \( X \) decomposes into a zero mode and an infinite number harmonic oscillator modes \( \alpha_{-n} \) which create states of energy \( n \). We’ll deal with the zero mode shortly but, for now, we focus on the oscillators. Acting \( d \) times with the operator \( \alpha_{-n} \) creates states with energy \( dn \). This gives a contribution to \( \text{Tr} q^{L_0} \) of the form

\[ \sum_{d=0}^{\infty} q^{nd} = \frac{1}{1 - q^n} \]
But the Fock space of a single scalar field is built by acting with oscillator modes $n \in \mathbb{Z}^+$. Including the central charge, $c = 1$, the contribution from the oscillator modes of a single scalar field is therefore

$$\text{Tr} \; q^{L_0 - c/24} = \frac{1}{q^{1/24}} \prod_{n=1}^{\infty} \frac{1}{1 - q^n}$$

There is a similar expression from the $\bar{q}^{\bar{L}_0 - \bar{c}/24}$ sector. We’re still left with the contribution from the zero mode $p$ of the scalar field. The contribution to the energy $H$ of the state on the worldsheet is

$$\frac{1}{4\pi \alpha'} \int d\sigma \; (\alpha' p)^2 = \frac{1}{2} \alpha' p^2$$

The trace in the partition function requires us to sum over all states, which gives

$$\int \frac{dp}{2\pi} \; e^{-\pi \alpha' (\text{Im} \tau) p^2} \sim \frac{1}{\sqrt{\alpha' \text{Im} \tau}}$$

So, including both the zero mode and oscillators, we get the partition function for a single free scalar field,

$$Z_{\text{scalar}}[\tau] \sim \frac{1}{\sqrt{\alpha' \text{Im} \tau}} \frac{1}{(q\bar{q})^{1/24}} \prod_{n=1}^{\infty} \frac{1}{1 - q^n} \prod_{n=1}^{\infty} \frac{1}{1 - \bar{q}^n} \quad (6.20)$$

where I haven’t been careful to keep track of constant factors.

To build the string partition function, we should really work in covariant quantization and include the ghost fields. Here we’ll cheat and work in lightcone gauge. This is dodgy because, if we do it honestly, much of the physics gets pushed to the $p^+ = 0$ limit of the lightcone momentum where the gauge choice breaks down. So instead we’ll do it dishonestly.

In lightcone gauge, we have 24 oscillator modes. But we have 26 zero modes. (You may worry that we still have to impose level matching...this is the dishonest part of the calculation. We’ll see partly where it comes from shortly). Finally, there’s a couple of extra steps. We need to divide by the volume of the conformal Killing group. This is just $U(1) \times U(1)$, acting by translations along the cycles of the torus. The volume is just $\text{Vol} = 4\pi^2 \text{Im} \tau$. Finally, we also need to integrate over the moduli space of the torus. Our final result, neglecting all constant factors, is

$$Z_{\text{string}} = \int d^2 \tau \frac{1}{(\text{Im} \tau)} \frac{1}{(\alpha' \text{Im} \tau)^{13}} \frac{1}{q\bar{q}} \left( \prod_{n=1}^{\infty} \frac{1}{1 - q^n} \right)^{24} \left( \prod_{n=1}^{\infty} \frac{1}{1 - \bar{q}^n} \right)^{24} \quad (6.21)$$
Modular Invariance

The function appearing in the partition function for the scalar field has a name: it is the inverse of the Dedekind eta function

\[ \eta(q) = q^{1/24} \prod_{n=1}^{\infty} (1 - q^n) \]

It was studied in the 1800s by mathematicians interested in the properties of functions under modular transformations \( T : \tau \rightarrow \tau + 1 \) and \( S : \tau \rightarrow -1/\tau \). The eta-function satisfies the identities

\[ \eta(\tau + 1) = e^{2\pi i/24} \eta(\tau) \quad \text{and} \quad \eta(-1/\tau) = \sqrt{-i\tau} \eta(\tau) \]

These two statements ensure that the scalar partition function (6.20) is a modular invariant function. Of course, that kinda had to be true: it follows from the underlying physics.

Written in terms of \( \eta \), the string partition function (6.21) takes the form

\[ Z_{\text{string}} = \int d^2 \tau (\text{Im } \tau)^2 \left( \frac{1}{\sqrt{\text{Im } \tau}} \frac{1}{\eta(q)} \frac{1}{\eta(\bar{q})} \right)^{24} \]

Both the measure and the integrand, are individually modular invariant.

6.4.3 Interpreting the String Partition Function

It’s probably not immediately obvious what the string partition function (6.21) is telling us. Let’s spend some time trying to understand it in terms of some simpler concepts.

We know that the free string describes an infinite number of particles with mass \( m_n^2 = 4(n - 1)/\alpha' \), \( n = 0, 1, \ldots \). The string partition function should just be a sum over vacuum loops of each of these particles. We’ll now show that it almost has this interpretation.

Firstly, let’s figure out what the contribution from a single particle would be? We’ll consider a free massive scalar field \( \phi \) in \( D \) dimensions. The partition function is given by,

\[
Z = \int \mathcal{D}\phi \exp \left( -\frac{1}{2} \int d^D x \ \phi (-\partial^2 + m^2) \phi \right) \\
\sim \det^{-1/2}(-\partial^2 + m^2) \\
= \exp \left( \frac{1}{2} \int \frac{d^D p}{(2\pi)^D} \ln(p^2 + m^2) \right)
\]
This is the partition function of a field theory. It contains vacuum loops for all numbers of particles. To compare to the string partition function, we want the vacuum amplitude for just a single particle. But that’s easy to extract. We write the field theory partition function as,

$$Z = \exp(Z_1) = \sum_{n=0}^{\infty} \frac{Z_1^n}{n!}$$

Each term in the sum corresponds to \(n\) particles propagating in a vacuum loop, with the \(n!\) factor taking care of Bosonic statistics. So the vacuum amplitude for a single, free massive particle is simply

$$Z_1 = \frac{1}{2} \int \frac{d^Dp}{(2\pi)^D} \ln(p^2 + m^2)$$

Clearly this diverges in the UV range of the integral, \(p \rightarrow \infty\). There’s a nice way to rewrite this integral using something known as Schwinger parameterization. We make use of the identity

$$\int_0^{\infty} dl \ e^{-xl} = 1 \Rightarrow \int_0^{\infty} dl \ \frac{e^{-xl}}{l} = -\ln x$$

We then write the single particle partition function as

$$Z_1 = \int \frac{d^Dp}{(2\pi)^D} \int_0^{\infty} dl \ \frac{1}{2l} e^{-(p^2 + m^2)l} \tag{6.22}$$

It’s worth mentioning that there’s another way to see that this is the single particle partition function that is a little closer in spirit to the method we used in string theory. We could start with the einbein form of the relativistic particle action (1.8). After fixing the gauge to \(e = 1\), the exponent in (6.22) is the energy of the particle traversing a loop of length \(l\). The integration measure \(dl/l\) sums over all possible sizes of loops.

We can happily perform the \(\int d^Dp\) integral in (6.22). Ignoring numerical factors, we have

$$Z_1 = \int_0^{\infty} dl \ \frac{1}{l^{1+D/2}} e^{-m^2l} \tag{6.23}$$

Note that the UV divergence as \(p \rightarrow \infty\) has metamorphosised into a divergence associated to small loops as \(l \rightarrow 0\).
Equation (6.23) gives the answer for a single particle of mass $m$. In string theory, we expect contributions from an infinite number of species of particles of mass $m_n$. Specializing to $D = 26$, we expect the partition function to be

$$Z = \int_0^\infty dl \frac{1}{l^{14}} \sum_{n=0}^\infty e^{-n^2l}$$

But we know that the mass spectrum of the free string: it is given in terms of the $L_0$ and $\tilde{L}_0$ operators by

$$m^2 = \frac{4}{\alpha'}(L_0 - 1) = \frac{4}{\alpha'}(\tilde{L}_0 - 1) = \frac{2}{\alpha'}(L_0 + \tilde{L}_0 - 2)$$

subject to the constraint of level matching, $L_0 = \tilde{L}_0$. It’s easy to impose level matching: we simply throw in a Kronecker delta in its integral representation,

$$\frac{1}{2\pi} \int_{-1/2}^{+1/2} ds \ e^{2\pi is(L_0 - \tilde{L}_0)} = \delta_{L_0,\tilde{L}_0} \quad (6.24)$$

Replacing the sum over species, with the trace over the spectrum of states subject to level matching, the partition function becomes,

$$Z = \int_0^\infty dl \frac{1}{l^{14}} \int_{-1/2}^{+1/2} ds \ \text{Tr} e^{2\pi is(L_0 - \tilde{L}_0)} e^{-2(L_0 + \tilde{L}_0 - 2)l/\alpha'} \quad (6.25)$$

We again use the definition $q = \exp(2\pi i \tau)$, but this time the complex parameter $\tau$ is a combination of the length of the loop $l$ and the auxiliary variable that we introduced to impose level matching,

$$\tau = s + \frac{2li}{\alpha'}$$

The trace over the spectrum of the string once gives the eta-functions, just as it did before. We’re left with the result for the partition function,

$$Z_{\text{string}} = \int \frac{d^2\tau}{(\text{Im} \tau)^2} \left( \frac{1}{\sqrt{\text{Im} \tau}} \frac{1}{\eta(q)} \frac{1}{\bar{\eta}(\bar{q})} \right)^{24}$$

But this is exactly the same expression that we saw before. With a difference! In fact, the difference is hidden in the notation: it is the range of integration for $d^2\tau$ which can be found in the original expressions (6.23) and (6.24). Re $\tau$ runs over the same interval $[-\frac{1}{2}, +\frac{1}{2}]$ that we saw in string theory. As is clear from this discussion, it is this integral which implements level matching. The difference comes in the range of Im $\tau$ which, in this naive analysis, runs over $[0, \infty)$. This is in stark contrast to string theory where we only integrate over the fundamental domain.
This highlights our previous statement: the potential UV divergences in field theory are encountered in the region $\text{Im} \tau \sim l \to 0$. In the above analysis, this corresponds to particles traversing small loops. But this region is simply absent in the correct string theory computation. It is mapped, by modular invariance, to the infra-red region of large loops.

It is often said that in the $g_s \to 0$ limit string theory becomes a theory of an infinite number of free particles. This is true of the spectrum. But this calculation shows that it’s not really true when we compute loops because the modular invariance means that we integrate over a different range of momenta in string theory than in a naive field theory approach.

So what happens in the infra-red region of our partition function? The easiest place to see it is in the $l \to \infty$ limit of the integral (6.25). We see that the integral is dominated by the lightest state which, for the bosonic string is the tachyon. This has $m^2 = -4/\alpha'$, or $(L_0 + \tilde{L}_0 - 2) = -2$. This gives a contribution to the partition function of,

$$\int_{-\infty}^{\infty} \frac{dl}{l^{14}} e^{4l/\alpha'}$$

which clearly diverges. This IR divergence of the one-loop partition function is another manifestation of tachyonic trouble. In the superstring, there is no tachyon and the IR region is well-behaved.

### 6.4.4 So is String Theory Finite?

The honest answer is that we don’t know. The UV finiteness that we saw above holds for all one-loop amplitudes. This means, in particular, that we have a one-loop finite theory of gravity interacting with matter in higher dimensions. This is already remarkable.

There is more good news: One can show that UV finiteness continues to hold at the two-loops. And, for the superstring, state-of-the-art techniques using the “pure-spinor” formalism show that certain objects remain finite up to five-loops. Moreover, the exponential suppression (6.14) that we saw when all momentum exchanges are large continues to hold for all amplitudes.

However, no general statement of finiteness has been proven. The danger lurks in the singular points in the integration over Riemann surfaces of genus 3 and higher.
6.4.5 Beyond Perturbation Theory?

From the discussion in this section, it should be clear that string perturbation theory is entirely analogous to the Feynman diagram expansion in field theory. Just as in field theory, one can show that the expansion in $g_s$ is asymptotic. This means that the series does not converge, but we can nonetheless make sense of it.

However, we know that there are many phenomena in quantum field theory that aren’t captured by Feynman diagrams. These include confinement in the strongly coupled regime and instantons and solitons in the weakly coupled regime. Does this mean that we are missing similarly interesting phenomena in string theory? The answer is almost certainly yes! In this section, I’ll very briefly allude to a couple of more advanced topics which allow us to go beyond the perturbative expansion in string theory. The goal is not really to teach you these things, but merely to familiarize you with some words.

One way to proceed is to keep quantum field theory as our guide and try to build a non-perturbative definition of string theory in terms of a path integral. We’ve already seen that the Polyakov path integral over worldsheets is equivalent to Feynman diagrams. So we need to go one step further. What does this mean? Recall that in QFT, a field creates a particle. In string theory, we are now looking for a field which creates a loop of string. We should have a different field for each configuration of the string. In other words, our field should itself be a function of a function: $\Phi(X^\mu(\sigma))$. Needless to say, this is quite a complicated object. If we were brave, we could then consider the path integral for this field,

$$Z = \int D\Phi \, e^{iS[\Phi(X(\sigma))]}$$

for some suitable action $S[\Phi]$. The idea is that this path integral should reproduce the perturbative string expansion and, furthermore, defines a non-perturbative completion of the theory. This line of ideas is known as string field theory. It should be clear that this is one step further in the development: particles $\rightarrow$ fields $\rightarrow$ string fields. Or, in more historical language, if field theory is “second quantization”, then string field theory is “third quantization”.

String field theory has been fairly successful for the open string and some interesting non-perturbative results have been obtained in this manner. However, for the closed string this approach has been much less useful. It is usually thought that there are deep reasons behind the failure of closed string field theory, related to issues that we mentioned at the beginning of this section: there are no off-shell quantities in a theory...
of gravity. Moreover, we mentioned in Section 4 that a theory of interacting open strings necessarily includes closed strings, so somehow the open string field theory should already contain gravity and closed strings. Quite how this comes about is still poorly understood.

There are other ways to get a handle on non-perturbative aspects of string theory using the low-energy effective action (we will describe what the “low-energy effective action” is in the next section). Typically these techniques rely on supersymmetry to provide a window into the strongly coupled regime and so work only for the superstring. These methods have been extremely successful and any course on superstring theory would be devoted to explaining various aspects of such as dualities and M-theory.

Finally, in asymptotically AdS spacetimes, the AdS/CFT correspondence gives a non-perturbative definition of string theory and quantum gravity in the bulk in terms of Yang-Mills theory, or something similar, on the boundary. In some sense, the boundary field theory is a “string field theory”.

6.5 Appendix: Games with Integrals and Gamma Functions

The gamma function is defined by the integral representation

$$\Gamma(z) = \int_0^\infty dt \ t^{z-1} e^{-t}$$

(6.26)

which converges if $\Re z > 0$. It has a unique analytic expression to the whole $z$-plane. The absolute value of the gamma function over the $z$-plane is shown in the figure.

The gamma function has a couple of important properties. Firstly, it can be thought of as the analytic continuation of the factorial function for positive integers, meaning

$$\Gamma(n) = (n-1)! \quad n \in \mathbb{Z}^+$$

Secondly, $\Gamma(z)$ has poles at non-positive integers. More precisely when $z \approx -n$, with $n = 0, 1, \ldots$, there is the expansion

$$\Gamma(z) \approx \frac{1}{z+n} \frac{(-1)^n}{n!}$$
The Euler Beta Function

The Euler beta function is defined for \( x, y \in \mathbb{C} \) by

\[
B(x, y) = \frac{\Gamma(x)\Gamma(y)}{\Gamma(x + y)}
\]

It has the integral representation

\[
B(x, y) = \int_0^1 dt \ t^{x-1}(1 - t)^{y-1}
\]  \hspace{1cm} (6.27)

Let’s prove this statement. We start by looking at

\[
\Gamma(x)\Gamma(y) = \int_0^\infty du \int_0^\infty dv \ e^{-u}u^{x-1}e^{-v}v^{y-1}
\]

We write \( u = a^2 \) and \( v = b^2 \) so the integral becomes

\[
\Gamma(x)\Gamma(y) = 4 \int_0^\infty da \int_0^\infty db \ e^{-(a^2+b^2)a^{2x-1}b^{2y-1}}
\]

We now change coordinates once more, this time to polar \( a = r \cos \theta \) and \( b = r \sin \theta \). We get

\[
\Gamma(x)\Gamma(y) = \int_0^\infty rdr \ e^{-r^2r^{2x+2y-2}} \int_0^{2\pi} d\theta \ |\cos \theta|^{2x-1}|\sin \theta|^{2y-1}
\]

\[
= \frac{1}{2}\Gamma(x + y) \times 4 \int_0^{\pi/2} d\theta \ (\cos \theta)^{2x-1}(\sin \theta)^{2y-1}
\]

\[
= \Gamma(x + y) \int_0^1 dt \ (1 - t)^{y-1}t^{x-1}
\]

where, in the final line, we made the substitution \( t = \cos^2 \theta \). This completes the proof.

The Virasoro-Shapiro Amplitude

In the closed string computation, we came across the integral

\[
C(a, b) = \int d^2z \ |z|^{2a-2}|1-z|^{2b-2}
\]

We will now evaluate this and show that it is given by (6.11). We start by using a trick. We can write

\[
|z|^{2a-2} = \frac{1}{\Gamma(1 - a)} \int_0^\infty dt \ t^{-a}e^{-|z|^2t}
\]
which follows from the definition (6.26) of the gamma function. Similarly, we can write

\[ |1 - z|^{2b - 2} = \frac{1}{\Gamma(1 - b)} \int_0^\infty du \ u^{-b} e^{-|1 - z|^2 u} \]

We decompose the complex coordinate \( z = x + iy \), so that the measure of the integral is \( d^2 z = 2dx dy \). We can then write the integral \( C(a, b) \) as

\[
C(a, b) = \int \frac{d^2 z \ du \ dt}{\Gamma(1 - a)\Gamma(1 - b)} \ t^{-a}u^{-b} e^{-|z|^2 t} e^{-|1 - z|^2 u}
\]

\[
= 2 \int \frac{dx \ dy \ du \ dt}{\Gamma(1 - a)\Gamma(1 - b)} \ t^{-a}u^{-b} e^{-(t+u)(x^2+y^2)+2ux-u}
\]

\[
= 2 \int \frac{dx \ dy \ du \ dt}{\Gamma(1 - a)\Gamma(1 - b)} \ t^{-a}u^{-b} \exp \left( -(t + u) \left[ \left( x - \frac{u}{t+u} \right)^2 + y^2 \right] - u + \frac{u^2}{t+u} \right)
\]

Now we do the \( dx \ dy \) integral which is simply Gaussian. We find

\[
C(a, b) = \frac{2\pi}{\Gamma(1 - a)\Gamma(1 - b)} \int_0^\infty du \ dt \ \frac{t^{-a}u^{-b}}{t+u} e^{-tu/(t+u)}
\]

Finally, we make a change of variables. We write \( t = \alpha \beta \) and \( u = (1 - \beta)\alpha \). In order for \( t \) and \( u \) to take values in the range \([0, \infty)\), we require \( \alpha \in [0, \infty) \) and \( \beta \in [0, 1] \). Taking into account the Jacobian arising from this transformation, which is simply \( \alpha \), the integral becomes

\[
C(a, b) = \frac{2\pi}{\Gamma(1 - a)\Gamma(1 - b)} \int d\alpha \ d\beta \ \frac{\alpha^{1-a-b}}{\alpha} \ \beta^{-a}(1 - \beta)^{-b} e^{-\alpha \beta(1 - \beta)}
\]

But we recognize the integral over \( d\alpha \): it is simply

\[
\int_0^\infty d\alpha \ \alpha^{-a-b} e^{-\beta \alpha(1 - \beta)} = [\beta(1 - \beta)]^{a+b-1} \Gamma(1 - a - b)
\]

We write \( c = 1 - a - b \). Finally, we’re left with

\[
C(a, b) = \frac{2\pi \Gamma(c)}{\Gamma(1 - a)\Gamma(1 - b)} \int_0^1 d\beta \ (1 - \beta)^{a-1} \beta^{b-1}
\]

But the final integral is the Euler beta function (6.27). This gives us our promised result,

\[
C(a, b) = \frac{2\pi \Gamma(a)\Gamma(b)\Gamma(c)}{\Gamma(1 - a)\Gamma(1 - b)\Gamma(1 - c)}
\]
7. Low Energy Effective Actions

So far, we’ve only discussed strings propagating in flat spacetime. In this section we will consider strings propagating in different backgrounds. This is equivalent to having different CFTs on the worldsheet of the string.

There is an obvious generalization of the Polyakov action to describe a string moving in curved spacetime,

\[ S = \frac{1}{4\pi\alpha'} \int d^2\sigma \sqrt{g} \, g^{\alpha\beta} \partial_{\alpha} X^\mu \partial_{\beta} X^\nu \, G_{\mu\nu}(X) \]  

(7.1)

Here \( g_{\alpha\beta} \) is again the worldsheet metric. This action describes a map from the worldsheet of the string into a spacetime with metric \( G_{\mu\nu}(X) \). (Despite its name, this metric is not to be confused with the Einstein tensor which we won’t have need for in this lecture notes).

Actions of the form (7.1) are known as non-linear sigma models. (This strange name has its roots in the history of pions). In this context, the \( D \)-dimensional spacetime is sometimes called the target space. Theories of this type are important in many aspects of physics, from QCD to condensed matter.

Although it’s obvious that (7.1) describes strings moving in curved spacetime, there’s something a little fishy about just writing it down. The problem is that the quantization of the closed string already gave us a graviton. If we want to build up some background metric \( G_{\mu\nu}(X) \), it should be constructed from these gravitons, in much the same manner that a laser beam is made from the underlying photons. How do we see that the metric in (7.1) has anything to do with the gravitons that arise from the quantization of the string?

The answer lies in the use of vertex operators. Let’s expand the metric as a small fluctuation around flat space

\[ G_{\mu\nu}(X) = \delta_{\mu\nu} + h_{\mu\nu}(X) \]

Then the partition function that we build from the action (7.1) is related to the partition function for a string in flat space by

\[ Z = \int \mathcal{D}X \mathcal{D}g \, e^{-S_{\text{Poly}} - V} = \int \mathcal{D}X \mathcal{D}g \, e^{-S_{\text{Poly}}}(1 - V + \frac{1}{2}V^2 + \ldots) \]
where \( S_{\text{Poly}} \) is the action for the string in flat space given in (1.22) and \( V \) is the expression

\[
V = \frac{1}{4\pi\alpha'} \int d^2\sigma \sqrt{g} \, g^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu \, h_{\mu\nu}(X) \tag{7.2}
\]

But we’ve seen this before: it’s the vertex operator associated to the graviton state of the string! For a plane wave, corresponding to a graviton with polarization given by the symmetric, traceless tensor \( \zeta_{\mu\nu} \) and momentum \( p^\mu \), the fluctuation is given by

\[
h_{\mu\nu}(X) = \zeta_{\mu\nu} \, e^{ip \cdot X}
\]

With this choice, the expression (7.2) agrees with the vertex operator (5.9). But in general, we could take any linear superposition of plane waves to build up a general fluctuation \( h_{\mu\nu}(X) \).

We know that inserting a single copy of \( V \) in the path integral corresponds to the introduction of a single graviton state. Inserting \( e^V \) in the path integral corresponds to a coherent state of gravitons, changing the metric from \( \delta_{\mu\nu} \) to \( \delta_{\mu\nu} + h_{\mu\nu} \). In this way we see that the background curved metric of (7.1) is indeed built of the quantized gravitons that we first met back in Section 2.

### 7.1 Einstein’s Equations

In conformal gauge, the Polyakov action in flat space reduces to a free theory. This fact was extremely useful, allowing us to compute the spectrum of the theory. But on a curved background, it is no longer the case. In conformal gauge, the worldsheet theory is described by an interacting two-dimensional field theory,

\[
S = \frac{1}{4\pi\alpha'} \int d^2\sigma \, G_{\mu\nu}(X) \, \partial_\alpha X^\mu \partial^\alpha X^\nu \tag{7.3}
\]

To understand these interactions in more detail, let’s expand around a classical solution which we take to simply be a string sitting at a point \( \bar{x}^\mu \).

\[
X^\mu(\sigma) = \bar{x}^\mu + \sqrt{\alpha'} \, Y^\mu(\sigma)
\]

Here \( Y^\mu \) are the dynamical fluctuations about the point which we assume to be small. The factor of \( \sqrt{\alpha'} \) is there for dimensional reasons: since \([X] = -1\), we have \([Y] = 0\) and statements like \( Y \ll 1 \) make sense. Expanding the Lagrangian gives

\[
G_{\mu\nu}(X) \, \partial X^\mu \partial X^\nu = \alpha' \left[ G_{\mu\nu}(\bar{x}) + \sqrt{\alpha'} G_{\mu\nu,\omega}(\bar{x}) \, Y^\omega + \frac{\alpha'}{2} \, G_{\mu\nu,\omega\rho}(\bar{x}) \, Y^\omega Y^\rho + \ldots \right] \partial Y^\mu \partial Y^\nu
\]

Each of the coefficients \( G_{\mu\nu,...} \) in the Taylor expansion are coupling constants for the interactions of the fluctuations \( Y^\mu \). The theory has an infinite number of coupling constants and they are nicely packaged into the function \( G_{\mu\nu}(X) \).
We want to know when this field theory is weakly coupled. Obviously this requires the whole infinite set of coupling constants to be small. Let’s try to characterize this in a crude manner. Suppose that the target space has characteristic radius of curvature $r_c$, meaning schematically that

\[ \frac{\partial G}{\partial X} \sim \frac{1}{r_c} \]

The radius of curvature is a length scale, so $[r_c] = -1$. From the expansion of the metric, we see that the effective dimensionless coupling is given by

\[ \frac{\sqrt{\alpha'}}{r_c} \]

This means that we can use perturbation theory to study the CFT (7.3) if the spacetime metric only varies on scales much greater than $\sqrt{\alpha'}$. The perturbation series in $\sqrt{\alpha'}/r_c$ is usually called the $\alpha'$-expansion to distinguish it from the $g_s$ expansion that we saw in the previous section. Typically a quantity computed in string theory is given by a double perturbation expansion: one in $\alpha'$ and one in $g_s$.

If there are regions of spacetime where the radius of curvature becomes comparable to the string length scale, $r_c \sim \sqrt{\alpha'}$, then the worldsheet CFT is strongly coupled and we will need to develop new methods to solve it. Notice that strong coupling in $\alpha'$ is hard, but the problem is at least well-defined in terms of the worldsheet path integral. This is qualitatively different to the question of strong coupling in $g_s$ for which, as discussed in Section 6.4.5, we’re really lacking a good definition of what the problem even means.

### 7.1.1 The Beta Function

Classically, the theory defined by (7.3) is conformally invariant. But this is not necessarily true in the quantum theory. To regulate divergences we will have to introduce a UV cut-off and, typically, after renormalization, physical quantities depend on the scale of a given process $\mu$. If this is the case, the theory is no longer conformally invariant. There are plenty of theories which classically possess scale invariance which is broken quantum mechanically. The most famous of these is Yang-Mills.

As we’ve discussed several times, in string theory conformal invariance is a gauge symmetry and we can’t afford to lose it. Our goal in this section is to understand the circumstances under which (7.3) retains conformal invariance at the quantum level.
The object which describes how couplings depend on a scale \( \mu \) is called the \( \beta \)-function. Since we have a functions worth of couplings, we should really be talking about a \( \beta \)-functional, schematically of the form

\[
\beta_{\mu\nu}(G) \sim \mu \frac{\partial G_{\mu\nu}(X; \mu)}{\partial \mu}
\]

The quantum theory will be conformally invariant only if

\[
\beta_{\mu\nu}(G) = 0
\]

We now compute this for the non-linear sigma model at one-loop. Our strategy will be to isolate the UV divergence of the theory and figure out what kind of counterterm we should add. The beta-function will vanish if this counterterm vanishes.

The analysis is greatly simplified by a cunning choice of coordinates. Around any point \( \bar{x} \), we can always pick Riemann normal coordinates such that the expansion in \( X^\mu = \bar{x}^\mu + \sqrt{\alpha'} Y^\mu \) gives

\[
G_{\mu\nu}(X) = \delta_{\mu\nu} - \frac{\alpha'}{3} \mathcal{R}_{\mu\lambda\nu\kappa}(\bar{x}) Y^\lambda Y^\kappa + \mathcal{O}(Y^3)
\]

To quartic order in the fluctuations, the action becomes

\[
S = \frac{1}{4\pi} \int d^2 \sigma \, \partial Y^\mu \partial Y^\nu \delta_{\mu\nu} - \frac{\alpha'}{3} \mathcal{R}_{\mu\lambda\nu\kappa} Y^\lambda Y^\kappa \partial Y^\mu \partial Y^\nu
\]

We can now treat this as an interacting quantum field theory in two dimensions. The quartic interaction gives a vertex with the Feynman rule,

\[
\begin{array}{c}
\times
\end{array} \sim \mathcal{R}_{\mu\lambda\nu\kappa} (k^\mu \cdot k^\nu)
\]

where \( k^\mu_\alpha \) is the 2d momentum (\( \alpha = 1, 2 \) is a worldsheet index) for the scalar field \( Y^\mu \). It sits in the Feynman rules because we are talking about derivative interactions.

Now we’ve reduced the problem to a simple interacting quantum field theory, we can compute the \( \beta \)-function using whatever method we like. The divergence in the theory comes from the one-loop diagram

![Diagram](attachment:image.png)
It’s actually simplest to think about this diagram in position space. The propagator for a scalar particle is

$$\langle Y^\lambda(\sigma)Y^\kappa(\sigma') \rangle = -\frac{1}{2} \delta^{\lambda\kappa} \ln |\sigma - \sigma'|^2$$

For the scalar field running in the loop, the beginning and end point coincide. The propagator diverges as $\sigma \to \sigma'$, which is simply reflecting the UV divergence that we would see in the momentum integral around the loop.

To isolate this divergence, we choose to work with dimensional regularization, with $d = 2 + \epsilon$. The propagator then becomes,

$$\langle Y^\lambda(\sigma)Y^\kappa(\sigma') \rangle = 2\pi \delta^{\lambda\kappa} \int \frac{d^{2+\epsilon} k}{(2\pi)^{2+\epsilon}} \frac{e^{ik \cdot (\sigma - \sigma')}}{k^2} \to \frac{\delta^{\lambda\kappa}}{\epsilon} \text{ as } \sigma \to \sigma'$$

The necessary counterterm for this divergence can be determined simply by replacing $Y^\lambda Y^\kappa$ in the action with $\langle Y^\lambda Y^\kappa \rangle$. To subtract the $1/\epsilon$ term, we add the counterterm

$$\mathcal{R}_{\mu\nu} Y^\lambda Y^\kappa \partial Y^\mu \partial Y^\nu \to \mathcal{R}_{\mu\nu} Y^\lambda Y^\kappa \partial Y^\mu \partial Y^\nu - \frac{1}{\epsilon} \mathcal{R}_{\mu\nu} \partial Y^\mu \partial Y^\nu$$

One can check that this can be absorbed by a wavefunction renormalization $Y^\mu \to Y^\mu + (\alpha'/6\epsilon) \mathcal{R}_{\mu\nu} Y^\nu$, together with the renormalization of the coupling constant which, in our theory, is the metric $G_{\mu\nu}$. We require,

$$G_{\mu\nu} \to G_{\mu\nu} + \frac{\alpha'}{\epsilon} \mathcal{R}_{\mu\nu}$$

(7.5)

From this we learn the beta function of the theory and the condition for conformal invariance. It is

$$\beta_{\mu\nu}(G) = \alpha' \mathcal{R}_{\mu\nu} = 0$$

(7.6)

This is a magical result! The requirement for the sigma-model to be conformally invariant is that the target space must be Ricci flat: $\mathcal{R}_{\mu\nu} = 0$. Or, in other words, the background spacetime in which the string moves must obey the vacuum Einstein equations! We see that the equations of general relativity also describe the renormalization group flow of 2d sigma models.

There are several more magical things just around the corner, but it’s worth pausing to make a few diverse comments.
**Beta Functions and Weyl Invariance**

The above calculation effectively studies the breakdown of conformal invariance in the CFT (7.3) on a flat worldsheet. We know that this should be the same thing as the breakdown of Weyl invariance on a curved worldsheet. Since this is such an important result, let’s see how it works from this other perspective. We can consider the worldsheet metric

\[ g_{\alpha\beta} = e^{2\phi} \delta_{\alpha\beta} \]

Then, in dimensional regularization, the theory is not Weyl invariant in \( d = 2 + \epsilon \) dimensions because the contribution from \( \sqrt{g} \) does not quite cancel that from the inverse metric \( g^{\alpha\beta} \). The action is

\[ S = \frac{1}{4\pi\alpha'} \int d^{2+\epsilon} \sigma \ e^{\phi \epsilon} \partial_{\alpha} X^\mu \partial^\alpha X^\nu G_{\mu\nu}(X) \approx \frac{1}{4\pi\alpha'} \int d^{2+\epsilon} \sigma \ (1 + \phi \epsilon) \partial_{\alpha} X^\mu \partial^\alpha X^\nu G_{\mu\nu}(X) \]

where, in this expression, the \( \alpha = 1, 2 \) index is now raised and lowered with \( \delta_{\alpha\beta} \). If we replace \( G_{\mu\nu} \) in this expression with the renormalized metric (7.5), we see that there’s a term involving \( \phi \) which remains even as \( \epsilon \to 0 \),

\[ S = \frac{1}{4\pi\alpha'} \int d^{2} \sigma \ \partial_{\alpha} X^\mu \partial^\alpha X^\nu \left[ G_{\mu\nu}(X) + \alpha' \phi R_{\mu\nu}(X) \right] \]

This indicates a breakdown of Weyl invariance. Indeed, we can look at our usual diagnostic for Weyl invariance, namely the vanishing of \( T^\alpha \). In conformal gauge, this is given by

\[ T_{\alpha\beta} = + \frac{4\pi}{\sqrt{g}} \frac{\partial S}{\partial g^{\alpha\beta}} = -2\pi \frac{\partial S}{\partial \phi} \delta_{\alpha\beta} \Rightarrow T^\alpha = -\frac{1}{2} R_{\mu\nu} \partial X^\mu \partial X^\nu \]

In this way of looking at things, we define the \( \beta \)-function to be the coefficient in front of \( \partial X \partial X \), namely

\[ T^\alpha = -\frac{1}{2\alpha'} \beta_{\mu\nu} \partial X^\mu \partial X^\nu \]

Again, we have the result

\[ \beta_{\mu\nu} = \alpha' R_{\mu\nu} \]
7.1.2 Ricci Flow

In string theory we only care about conformal theories with Ricci flat metrics. (And generalizations of this result that we will discuss shortly). However, in other areas of physics and mathematics, the RG flow itself is important. It is usually called Ricci flow,

\[ \mu \frac{\partial G_{\mu\nu}}{\partial \mu} = \alpha' R_{\mu\nu} \]

which dictates how the metric changes with scale \( \mu \).

As an illustrative and simple example, consider the target space \( S^2 \) with radius \( r \). This is an important model in condensed matter physics where it describes the low-energy limit of a one-dimensional Heisenberg spin chain. It is sometimes called the \( O(3) \) sigma-model. Because the sphere is a symmetric space, the only effect of the RG flow is to make the radius scale dependent: \( r = r(\mu) \). The beta function is given by

\[ \mu \frac{\partial r^2}{\partial \mu} = \frac{\alpha'}{2\pi} \]

Hence \( r \) gets large as we go towards the UV and small towards the IR. Since the coupling is \( 1/r \), this means that the non-linear sigma model with \( S^2 \) target space is asymptotically free. At low energies, the theory is strongly coupled and perturbative calculations — such as this one-loop beta function — are no longer trusted. In particular, one can show that the \( S^2 \) sigma-model develops a mass gap in the IR.

The idea of Ricci flow (7.7) was recently used by Perelman to prove the Poincaré conjecture. In fact, Perelman used a slightly generalized version of Ricci flow which we will see shortly. In the language of string theory, he introduced the dilaton field.

7.2 Other Couplings

We’ve understood how strings couple to a background spacetime metric. But what about the other modes of the string? In Section 2, we saw that a closed string has further massless states which are associated to the anti-symmetric tensor \( B_{\mu\nu} \) and the dilaton \( \Phi \). We will now see how the string reacts if these fields are turned on in spacetime.

7.2.1 Charged Strings and the \( B \) field

Let’s start by looking at how strings couple to the anti-symmetric field \( B_{\mu\nu} \). We discussed the vertex operator associated to this state in Section 5.4.1. It is given in
and takes the same form as the graviton vertex operator, but with $\zeta_{\mu\nu}$ anti-symmetric. It is a simple matter to exponentiate this, to get an expression for how strings propagate in background $B_{\mu\nu}$ field. We’ll keep the curved metric $G_{\mu\nu}$ as well to get the general action,

$$S = \frac{1}{4\pi\alpha'} \int d^2\sigma \sqrt{g} \left( G_{\mu\nu}(X) \partial_\alpha X^\mu \partial_\beta X^\nu g^{\alpha\beta} + i B_{\mu\nu}(X) \partial_\alpha X^\mu \partial_\beta X^\nu \epsilon^{\alpha\beta} \right)$$  (7.8)

Where $\epsilon^{\alpha\beta}$ is the anti-symmetric 2-tensor, normalized such that $\sqrt{g}\epsilon^{12} = +1$. (The factor of $i$ is there in the action because we’re in Euclidean space and this new term has a single “time” derivative). The action retains invariance under worldsheet reparameterizations and Weyl rescaling.

So what is the interpretation of this new term? We will now show that we should think of the field $B_{\mu\nu}$ as analogous to the gauge potential $A_\mu$ in electromagnetism. The action (7.8) is telling us that the string is “electrically charged” under $B_{\mu\nu}$.

**Gauge Potentials**

We’ll take a short detour to remind ourselves about some pertinent facts in electromagnetism. Let’s start by returning to a point particle. We know that a charged point particle couples to a background gauge potential $A_\mu$ through the addition of a worldline term to the action,

$$\int d\tau \ A_\mu(X) \dot{X}^\mu \ .$$  (7.9)

If this relativistic form looks a little unfamiliar, we can deconstruct it by working in static gauge with $X^0 \equiv t = \tau$, where it reads

$$\int dt \ A_0(X) + A_i(X) \dot{X}^i ,$$

which should now be recognizable as the Lagrangian that gives rise to the Coulomb and Lorentz force laws for a charged particle.

So what is the generalization of this kind of coupling for a string? First note that (7.9) has an interesting geometrical structure. It is the pull-back of the one-form $A = A_\mu dX^\mu$ in spacetime onto the worldline of the particle. This works because $A$ is a one-form and the worldline is one-dimensional. Since the worldsheet of the string is two-dimensional, the analogous coupling should be to a two-form in spacetime. This is an anti-symmetric
tensor field with two indices, $B_{\mu \nu}$. The pull-back of $B_{\mu \nu}$ onto the worldsheet gives the interaction,

$$
\int d^2 \sigma \ B_{\mu \nu}(X) \partial_\alpha X^\mu \partial_\beta X^\nu \epsilon^{\alpha \beta} .
$$

(7.10)

This is precisely the form of the interaction we found in (7.8).

The point particle coupling (7.9) is invariant under gauge transformations of the background field $A_\mu \rightarrow A_\mu + \partial_\mu \alpha$. This follows because the Lagrangian changes by a total derivative. There is a similar statement for the two-form $B_{\mu \nu}$. The spacetime gauge symmetry is,

$$
B_{\mu \nu} \rightarrow B_{\mu \nu} + \partial_\mu C_\nu - \partial_\nu C_\mu
$$

(7.11)

under which the Lagrangian (7.10) changes by a total derivative.

In electromagnetism, one can construct the gauge invariant electric and magnetic fields which are packaged in the two-form field strength $F = dA$. Similarly, for $B_{\mu \nu}$, the gauge invariant field strength $H = dB$ is a three-form,

$$
H_{\mu \nu \rho} = \partial_\mu B_{\nu \rho} - \partial_\nu B_{\rho \mu} + \partial_\rho B_{\mu \nu} .
$$

This 3-form $H$ is sometimes known as the torsion. It plays the same role as torsion in general relativity, providing an anti-symmetric component to the affine connection.

### 7.2.2 The Dilaton

Let’s now figure out how the string couples to a background dilaton field $\Phi(X)$. This is more subtle. A naive construction of the vertex operator is not primary and one must work a little harder. The correct derivation of the vertex operators can be found in Polchinski. Here I will simply give the coupling and explain some important features.

The action of a string moving in a background involving profiles for the massless fields $G_{\mu \nu}$, $B_{\mu \nu}$ and $\Phi(X)$ is given by

$$
S = \frac{1}{4 \pi \alpha'} \int d^2 \sigma \sqrt{g} \ (G_{\mu \nu}(X) \partial_\alpha X^\mu \partial_\beta X^\nu g^{\alpha \beta} + i B_{\mu \nu}(X) \partial_\alpha X^\mu \partial_\beta X^\nu \epsilon^{\alpha \beta} + \alpha' \Phi(X) R^{(2)})
$$

(7.12)

where $R^{(2)}$ is the two-dimensional Ricci scalar of the worldsheet. (Up until now, we’ve always denoted this simply as $R$ but we’ll introduce the superscript from hereon to distinguish the worldsheet Ricci scalar from the spacetime Ricci scalar).
The coupling to the dilaton is surprising for several reasons. Firstly, we see that the term in the action vanishes on a flat worldsheet, \( R^{(2)} = 0 \). This is one of the reasons that it’s a little trickier to determine this coupling using vertex operators.

However, the most surprising thing about the coupling to the dilaton is that it does not respect Weyl invariance! Since a large part of this course has been about understanding the implications of Weyl invariance, why on earth are we willing to throw it away now?! The answer, of course, is that we’re not. Although the dilaton coupling does violate Weyl invariance, there is a way to restore it. We will explain this shortly. But firstly, let’s discuss one crucially important implication of the dilaton coupling (7.12).

**The Dilaton and the String Coupling**

There is an exception to the statement that the classical coupling to the dilaton violates Weyl invariance. This arises when the dilaton is constant. For example, suppose

\[
\Phi(X) = \lambda , \text{ a constant}
\]

Then the dilaton coupling reduces to something that we’ve seen before: it is

\[
S_{\text{dilaton}} = \lambda \chi
\]

where \( \chi \) is the Euler character of the worldsheet that we introduced in (6.4). This tells us something important: the constant mode of the dilaton, \( \langle \Phi \rangle \) determines the string coupling constant. This constant mode is usually taken to be the asymptotic value of the dilaton,

\[
\Phi_0 = \lim_{X \to \infty} \Phi(X)
\]

(7.13)

The string coupling is then given by

\[
g_s = e^{\Phi_0}
\]

(7.14)

So the string coupling is not an independent parameter of string theory: it is the expectation value of a field. This means that, just like the spacetime metric \( G_{\mu \nu} \) (or, indeed, like the Higgs vev) it can be determined dynamically.

We’ve already seen that our perturbative expansion around flat space is valid as long as \( g_s \ll 1 \). But now we have a stronger requirement: we can only trust perturbation theory if the string is localized in regions of space where \( e^{\Phi(X)} \ll 1 \) for all \( X \). If the string ventures into regions where \( e^{\Phi(X)} \) is of order 1, then we will need to use techniques that don’t rely on string perturbation theory as described in Section 6.4.5.
7.2.3 Beta Functions

We now return to understanding how we can get away with the violation of Weyl invariance in the dilaton coupling (7.12). The key to this is to notice the presence of $\alpha'$ in front of the dilaton coupling. It’s there simply on dimensional grounds. (The other two terms in the action both come with derivatives $[\partial X] = -1$, so don’t need any powers of $\alpha'$).

However, recall that $\alpha'$ also plays the role of the loop-expansion parameter (7.4) in the non-linear sigma model. This means that the classical lack of Weyl invariance in the dilaton coupling can be compensated by a one-loop contribution arising from the couplings to $G_{\mu\nu}$ and $B_{\mu\nu}$.

To see this explicitly, one can compute the beta-functions for the two-dimensional field theory (7.12). In the presence of the dilaton coupling, it’s best to look at the breakdown of Weyl invariance as seen by $\langle T_\alpha^\alpha \rangle$. There are three different kinds of contribution that the stress-tensor can receive, related to the three different spacetime fields. Correspondingly, we define three different beta functions,

$$\langle T_\alpha^\alpha \rangle = \frac{1}{2\alpha'} \beta_{\mu\nu}(G) g^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu - \frac{i}{2\alpha'} \beta_{\mu\nu}(B) \epsilon^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu - \frac{1}{2} \beta(\Phi) R^{(2)} \quad (7.15)$$

We will not provide the details of the one-loop beta function computations. We merely state the results\(^8\),

$$\beta_{\mu\nu}(G) = \alpha' R_{\mu\nu} + 2\alpha' \nabla_\mu \nabla_\nu \Phi - \frac{\alpha'}{4} H_{\mu\lambda\nu} H^{\lambda\kappa}$$

$$\beta_{\mu\nu}(B) = -\frac{\alpha'}{2} \nabla^\lambda H_{\lambda\mu\nu} + \alpha' \nabla^\lambda \Phi H_{\lambda\mu\nu}$$

$$\beta(\Phi) = -\frac{\alpha'}{2} \nabla^2 \Phi + \alpha' \nabla_\mu \Phi \nabla^\mu \Phi - \frac{\alpha'}{24} H_{\mu\nu\lambda\kappa} H^{\mu\nu\lambda\kappa}$$

A consistent background of string theory must preserve Weyl invariance, which now requires $\beta_{\mu\nu}(G) = \beta_{\mu\nu}(B) = \beta(\Phi) = 0$.

7.3 The Low-Energy Effective Action

The equations $\beta_{\mu\nu}(G) = \beta_{\mu\nu}(B) = \beta(\Phi) = 0$ can be viewed as the equations of motion for the background in which the string propagates. We now change our perspective: we

\(^8\)The relationship between the beta function and Einstein’s equations was first shown by Friedan in his 1980 PhD thesis. A readable account of the full beta functions can be found in the paper by Callan, Friedan, Martinec and Perry “Strings in Background Fields”, Nucl. Phys. B262 (1985) 593. The full calculational details can be found in TASI lecture notes by Callan and Thorlacius which can be downloaded from the course webpage.
look for a \( D = 26 \) dimensional spacetime action which reproduces these beta-function equations as the equations of motion. This is the \textit{low-energy effective action} of the bosonic string,

\[
S = \frac{1}{2\kappa_0^2} \int d^{26}X \sqrt{-G} e^{-2\Phi} \left( \mathcal{R} - \frac{1}{12} H_{\mu\nu\lambda} H^{\mu\nu\lambda} + 4\partial_\mu \Phi \partial^\mu \Phi \right) \tag{7.16}
\]

where we have taken the liberty of Wick rotating back to Minkowski space for this expression. Here the overall constant involving \( \kappa_0 \) is not fixed by the field equations but can be determined by coupling these equations to a suitable source as described, for example, in 7.4.2. On dimensional grounds alone, it scales as \( \kappa_0^2 \sim l_s^{24} \) where \( \alpha' = l_s^2 \).

Varying the action with respect to the three fields can be shown to yield the beta functions thus,

\[
\delta S = \frac{1}{2\kappa_0^2} \int d^{26}X \sqrt{-G} e^{-2\Phi} \left( \delta G_{\mu\nu} \beta^\mu_\nu(G) - \delta B_{\mu\nu} \beta_{\mu\nu}(B) 
\right.
\]

\[
\left. - (2\delta \Phi + \frac{1}{2} G_{\mu\nu} \delta G_{\mu\nu}(\beta^\lambda_\lambda(G) - 4\beta(\Phi)) \right)
\]

Equation (7.16) governs the low-energy dynamics of the spacetime fields. The caveat “low-energy” refers to the fact that we only worked with the one-loop beta functions which requires large spacetime curvature.

Something rather remarkable has happened here. We started, long ago, by looking at how a single string moves in flat space. Yet, on grounds of consistency alone, we’re led to the action (7.16) governing how spacetime and other fields fluctuate in \( D = 26 \) dimensions. It feels like the tail just wagged the dog. That tiny string is seriously high-maintenance: its requirements are so stringent that they govern the way the whole universe moves.

You may also have noticed that we now have two different methods to compute the scattering of gravitons in string theory. The first is in terms of scattering amplitudes that we discussed in Section 6. The second is by looking at the dynamics encoded in the low-energy effective action (7.16). Consistency requires that these two approaches agree. They do.

7.3.1 String Frame and Einstein Frame

The action (7.16) isn’t quite of the familiar Einstein-Hilbert form because of that strange factor of \( e^{-2\Phi} \) that’s sitting out front. This factor simply reflects the fact that the action has been computed at tree level in string perturbation theory and, as we saw in Section 6, such terms typically scale as \( 1/g_s^2 \).

\[ \text{– 170 –} \]
It’s also worth pointing out that the kinetic terms for \( \Phi \) in (7.16) seem to have the wrong sign. However, it’s not clear that we should be worried about this because, again, the factor of \( e^{-2\Phi} \) sits out front meaning that the kinetic terms are not canonically normalized anyway.

To put the action in more familiar form, we can make a field redefinition. Firstly, it’s useful to distinguish between the constant part of the dilaton, \( \Phi_0 \), and the part that varies which we call \( \tilde{\Phi} \). We defined the constant part in (7.13); it is related to the string coupling constant. The varying part is simply given by

\[
\tilde{\Phi} = \Phi - \Phi_0
\]  

(7.17)

In \( D \) dimensions, we define a new metric \( \tilde{G}_{\mu\nu} \) as a combination of the old metric and the dilaton,

\[
\tilde{G}_{\mu\nu}(X) = e^{-4\tilde{\Phi}/(D-2)} G_{\mu\nu}(X)
\]  

(7.18)

Note that this isn’t to be thought of as a coordinate transformation or symmetry of the action. It’s merely a relabeling, a mixing-up, of the fields in the theory. We could make such redefinitions in any field theory. Typically, we choose not to because the fields already have canonical kinetic terms. The point of the transformation (7.18) is to get the fields in (7.16) to have canonical kinetic terms as well.

The new metric (7.18) is related to the old by a conformal rescaling. One can check that two metrics related by a general conformal transformation \( \tilde{G}_{\mu\nu} = e^{2\omega} G_{\mu\nu} \), have Ricci scalars related by

\[
\tilde{\mathcal{R}} = e^{-2\omega} \left( \mathcal{R} - 2(D-1)\nabla^2 \omega - (D-2)(D-1)\partial_\mu \omega \partial^\mu \omega \right)
\]

(We used a particular version of this earlier in the course when considering \( D = 2 \) conformal transformations). With the choice \( \omega = -2\tilde{\Phi}/(D-2) \) in (7.18), and restricting back to \( D = 26 \), the action (7.16) becomes

\[
S = \frac{1}{2\kappa^2} \int d^{26} X \sqrt{-G} \left( \tilde{\mathcal{R}} - \frac{1}{12} e^{-\tilde{\Phi}/3} H_{\mu\nu\lambda} H^{\mu\nu\lambda} - \frac{1}{6} \partial_\mu \tilde{\Phi} \partial^\mu \tilde{\Phi} \right)
\]  

(7.19)

The kinetic terms for \( \tilde{\Phi} \) are now canonical and come with the right sign. Notice that there is no potential term for the dilaton and therefore nothing that dynamically sets its expectation value in the bosonic string. However, there do exist backgrounds of the superstring in which a potential for the dilaton develops, fixing the string coupling constant.
The gravitational part of the action takes the standard Einstein-Hilbert form. The gravitational coupling is given by
\[ \kappa^2 = \kappa_0^2 e^{2\Phi_0} \sim l_s^{24} g_s^2 \]
(7.20)

The coefficient in front of Einstein-Hilbert term is usually identified with Newton’s constant

\[ 8\pi G_N = \kappa^2 \]

Note, however, that this is Newton’s constant in \( D = 26 \) dimensions: it will differ from Newton’s constant measured in a four-dimensional world. From Newton’s constant, we define the \( D = 26 \) Planck length \( 8\pi G_N = l_p^{24} \) and Planck mass \( M_p = l_p^{-1} \). (With the factor of 8\( \pi \) sitting there, this is usually called the reduced Planck mass). Comparing to (7.20), we see that weak string coupling, \( g_s \ll 1 \), provides a parametric separation between the Planck scale and the string scale,

\[ g_s \ll 1 \Rightarrow l_p \ll l_s \]

Often the mysteries of gravitational physics are associated with the length scale \( l_p \). We understand string theory best when \( g_s \ll 1 \) where much of stringy physics occurs at \( l_s \gg l_p \) and can be disentangled from strong coupling effects in gravity.

The original metric \( G_{\mu\nu} \) is usually called the string metric or sigma-model metric. It is the metric that strings see, as reflected in the action (7.1). In contrast, \( \tilde{G}_{\mu\nu} \) is called the Einstein metric. Of course, the two actions (7.16) and (7.19) describe the same physics: we have simply chosen to package the fields in a different way in each. The choice of metric — \( G_{\mu\nu} \) or \( \tilde{G}_{\mu\nu} \) — is usually referred to as a choice of frame: string frame, or Einstein frame.

The possibility of defining two metrics really arises because we have a massless scalar field \( \Phi \) in the game. Whenever such a field exists, there’s nothing to stop us measuring distances in different ways by including \( \Phi \) in our ruler. Said another way, massless scalar fields give rise to long range attractive forces which can mix with gravitational forces and violate the principle of equivalence. Ultimately, if we want to connect to Nature, we need to find a way to make \( \Phi \) massive. Such mechanisms exist in the context of the superstring.

### 7.3.2 Corrections to Einstein’s Equations

Now that we know how Einstein’s equations arise from string theory, we can start to try to understand new physics. For example, what are the quantum corrections to Einstein’s equations?
On general grounds, we expect these corrections to kick in when the curvature \( r_c \) of spacetime becomes comparable to the string length scale \( \sqrt{\alpha'} \). But that dovetails very nicely with the discussion above where we saw that the perturbative expansion parameter for the non-linear sigma model is \( \alpha'/r_c^2 \). Computing the next loop correction to the beta function will result in corrections to Einstein’s equations!

If we ignore \( H \) and \( \Phi \), the 2-loop sigma-model beta function can be easily computed and results in the \( \alpha' \), correction to Einstein’s equations:

\[
\beta_{\mu\nu} = \alpha' \mathcal{R}_{\mu\nu} + \frac{1}{2} \alpha'^2 \mathcal{R}_{\mu\lambda\rho\sigma} \mathcal{R}^{\lambda\rho\sigma}_{\nu} + \ldots = 0
\]

Such two loop corrections also appear in the heterotic superstring. However, they are absent for the type II string theories, with the first corrections appearing at 4-loops from the perspective of the sigma-model.

**String Loop Corrections**

Perturbative string theory has an \( \alpha' \) expansion and \( g_s \) expansion. We still have to discuss the latter. Here an interesting subtlety arises. The sigma-model beta functions arise from regulating the UV divergences of the worldsheet. Yet the \( g_s \) expansion cares only about the topology of the string. How can the UV divergences care about the global nature of the worldsheet. Or, equivalently, how can the higher-loop corrections to the beta-functions give anything interesting?

The resolution to this puzzle is to remember that, when computing higher \( g_s \) corrections, we have to integrate over the moduli space of Riemann surfaces. But this moduli space will include some tricky points where the Riemann surface degenerates. (For example, one cycle of the torus may pinch off). At these points, the UV divergences suddenly do care about global topology and this results in the \( g_s \) corrections to the low-energy effective action.

**7.3.3 Nodding Once More to the Superstring**

In section 2.5, we described the massless bosonic content for the four superstring theories: Heterotic \( SO(32) \), Heterotic \( E_8 \times E_8 \), Type IIA and Type IIB. Each of them contains the fields \( G_{\mu\nu} \), \( B_{\mu\nu} \) and \( \Phi \) that appear in the bosonic string, together with a collection of further massless fields. For each, the low-energy effective action describes the dynamics of these fields in \( D = 10 \) dimensional spacetime. It naturally splits up into three pieces,

\[
S_{\text{superstring}} = S_1 + S_2 + S_{\text{fermi}}
\]
Here $S_{\text{fermi}}$ describes the interactions of the spacetime fermions. We won’t describe these here. But we will briefly describe the low-energy bosonic action $S_1 + S_2$ for each of these four superstring theories.

$S_1$ is essentially the same for all theories and is given by the action we found for the bosonic string in string frame (7.16). We’ll start to use form notation and denote $H_{\mu\nu\lambda}$ simply as $H_3$, where the subscript tells us the degree of the form. Then the action reads

$$S_1 = \frac{1}{2\kappa_0^2} \int d^{10}X \sqrt{-G} e^{-2\Phi} \left( R - \frac{1}{2} |\tilde{H}_3|^2 + 4\partial\mu \Phi \partial^\mu \Phi \right)$$

(7.21)

There is one small difference, which is that the field $\tilde{H}_3$ that appears here for the heterotic string is not quite the same as the original $H_3$; we’ll explain this further shortly.

The second part of the action, $S_2$, describes the dynamics of the extra fields which are specific to each different theory. We’ll now go through the four theories in turn, explaining $S_2$ in each case.

- **Type IIA**: For this theory, $\tilde{H}_3$ appearing in (7.21) is $H_3 = dB_2$, just as we saw in the bosonic string. In Section 2.5, we described the extra bosonic fields of the Type IIA theory: they consist of a 1-form $C_1$ and a 3-form $C_3$. The dynamics of these fields is governed by the so-called Ramond-Ramond part of the action and is written in form notation as,

$$S_2 = -\frac{1}{4\kappa_0^2} \int d^{10}X \left[ \sqrt{-G} \left( |F_2|^2 + |\tilde{F}_4|^2 \right) + B_2 \wedge F_4 \wedge F_4 \right]$$

Here the field strengths are given by $F_2 = dC_1$ and $F_4 = dC_3$, while the object that appears in the kinetic terms is $\tilde{F}_4 = F_4 - C_1 \wedge H_3$. Notice that the final term in the action does not depend on the metric: it is referred to as a *Chern-Simons* term.

- **Type IIB**: Again, $\tilde{H}_3 \equiv H_3$. The extra bosonic fields are now a scalar $C_0$, a 2-form $C_2$ and a 4-form $C_4$. Their action is given by

$$S_2 = -\frac{1}{4\kappa_0^2} \int d^{10}X \left[ \sqrt{-G} \left( |F_1|^2 + |\tilde{F}_3|^2 + \frac{1}{2} |\tilde{F}_5|^2 \right) + C_4 \wedge H_3 \wedge F_3 \right]$$

where $F_1 = dC_0$, $F_3 = dC_2$ and $F_5 = dC_4$. Once again, the kinetic terms involve more complicated combinations of the forms: they are $\tilde{F}_3 = F_3 - C_0 \wedge H_3$ and
\[ \tilde{F}_5 = F_5 - \frac{1}{2} C_2 \wedge H_3 + \frac{1}{2} B_2 \wedge F_3. \] However, for type IIB string theory, there is one extra requirement on these fields that cannot be implemented in any simple way in terms of a Lagrangian: \( \tilde{F}_5 \) must be self-dual

\[ \tilde{F}_5 = * \tilde{F}_5 \]

Strictly speaking, one should say that the low-energy dynamics of type IIB theory is governed by the equations of motion that we get from the action, supplemented with this self-duality requirement.

- **Heterotic**: Both heterotic theories have just one further massless bosonic ingredient: a non-Abelian gauge field strength \( F_2 \), with gauge group \( SO(32) \) or \( E_8 \times E_8 \). The dynamics of this field is simply the Yang-Mills action in ten dimensions,

\[ S_2 = \frac{\alpha'}{8\kappa_0^2} \int d^{10}X \sqrt{-G} \text{Tr} |F_2|^2 \]

The one remaining subtlety is to explain what \( \tilde{H}_3 \) means in (7.21): it is defined as \( \tilde{H}_3 = dB_2 - \alpha' \omega_3 / 4 \) where \( \omega_3 \) is the Chern-Simons three form constructed from the non-Abelian gauge field \( A_1 \)

\[ \omega_3 = \text{Tr} \left( A_1 \wedge dA_1 + \frac{2}{3} A_1 \wedge A_1 \wedge A_1 \right) \]

The presence of this strange looking combination of forms sitting in the kinetic terms is tied up with one of the most intricate and interesting aspects of the heterotic string, known as anomaly cancelation.

The actions that we have written down here probably look a little arbitrary. But they have very important properties. In particular, the full action \( S_{\text{superstring}} \) of each of the Type II theories is invariant under \( N = 2 \) spacetime supersymmetry. (That means 32 supercharges). They are the unique actions with this property. Similarly, the heterotic superstring actions are invariant under \( N = 1 \) supersymmetry and, crucially, do not suffer from anomalies. The second book by Polchinski is a good place to start learning more about these ideas.

### 7.4 Some Simple Solutions

The spacetime equations of motion,

\[ \beta_{\mu\nu}(G) = \beta_{\mu\nu}(B) = \beta(\Phi) = 0 \]

have many solutions. This is part of the story of vacuum selection in string theory. What solution, if any, describes the world we see around us? Do we expect this putative
solution to have other special properties, or is it just a random choice from the many possibilities? The answer is that we don’t really know, but there is currently no known principle which uniquely selects a solution which looks like our world — with the gauge groups, matter content and values of fundamental constants that we observe — from the many other possibilities. Of course, these questions should really be asked in the context of the superstring where a greater understanding of various non-perturbative effects such as D-branes and fluxes leads to an even greater array of possible solutions.

Here we won’t discuss these problems. Instead, we’ll just discuss a few simple solutions that are well known. The first plays a role when trying to make contact with the real world, while the value of the others lies mostly in trying to better understand the structure of string theory.

7.4.1 Compactifications

We’ve seen that the bosonic string likes to live in $D = 26$ dimensions. But we don’t. Or, more precisely, we only observe three macroscopically large spatial dimensions. How do we reconcile these statements?

Since string theory is a theory of gravity, there’s nothing to stop extra dimensions of the universe from curling up. Indeed, under certain circumstances, this may be required dynamically. Here we exhibit some simple solutions of the low-energy effective action which have this property. We set $H_{\mu\nu\rho} = 0$ and $\Phi$ to a constant. Then we are simply searching for Ricci flat backgrounds obeying $R_{\mu\nu} = 0$. There are solutions where the metric is a direct product of metrics on the space

$$R^{1,3} \times X$$

where $X$ is a compact 22-dimensional Ricci-flat manifold.

The simplest such manifold is just $X = T^{22}$, the torus endowed with a flat metric. But there are a whole host of other possibilities. Compact, complex manifolds that admit such Ricci-flat metrics are called Calabi-Yau manifolds. (Strictly speaking, Calabi-Yau manifolds are complex manifolds with vanishing first Chern class. Yau’s theorem guarantees the existence of a unique Ricci flat metric on these spaces).

The idea that there may be extra, compact directions in the universe was considered long before string theory and goes by the name of Kaluza-Klein compactification. If the characteristic length scale $L$ of the space $X$ is small enough then the presence of these extra dimensions would not have been observed in experiment. The standard model of particle physics has been accurately tested to energies of a TeV or so, meaning that
if the standard model particles can roam around \( X \), then the length scale must be 
\[ L \lesssim (\text{TeV})^{-1} \sim 10^{-16} \text{ cm}. \]

However, one can cook up scenarios in which the standard model is stuck somewhere in these extra dimensions (for example, it may be localized on a D-brane). Under these circumstances, the constraints become much weaker because we would rely on gravitational experiments to detect extra dimensions. Present bounds require only 
\[ L \lesssim 10^{-5} \text{ cm}. \]

Consider the Einstein-Hilbert term in the low-energy effective action. If we are interested only in the dynamics of the 4d metric on \( \mathbb{R}^{1,3} \), this is given by
\[
S_{EH} = \frac{1}{2\kappa^2} \int d^{26}X \sqrt{-\tilde{G}} \tilde{\mathcal{R}} = \frac{\text{Vol}(X)}{2\kappa^2} \int d^4X \sqrt{-G_{4d}} R_{4d}
\]
(There are various moduli of the internal manifold \( X \) that are being neglected here). From this equation, we learn that effective 4d Newton constant is given in terms of 26d Newton constant by,
\[
8\pi G_{N}^{4d} = \frac{\kappa^2}{\text{Vol}(X)}
\]
Rewriting this in terms of the 4d Planck scale, we have \( l_{p}^{(4d)} \sim g_s l_s^{12} / \sqrt{\text{Vol}(X)} \). To trust this whole analysis, we require \( g_s \ll 1 \) and all length scales of the internal space to be bigger than \( l_s \). This ensures that \( l_{p}^{(4d)} < l_s \). Although the 4d Planck length is ludicrously small, \( l_{p}^{(4d)} \sim 10^{-33} \text{ cm} \), it may be that we don’t have to probe to this distance to uncover UV gravitational physics. The back-of-the-envelope calculation above shows that the string scale \( l_s \) could be much larger, enhanced by the volume of extra dimensions.

### 7.4.2 The String Itself

We’ve seen that quantizing small loops of string gives rise to the graviton and \( B_{\mu\nu} \) field. Yet, from the sigma model action (7.12), we also know that the string is charged under the \( B_{\mu\nu} \). Moreover, the string has tension, which ensures that it also acts as a source for the metric \( G_{\mu\nu} \). So what does the back-reaction of the string look like? Or, said another way: what is the sigma-model describing a string moving in the background of another string?

Consider an infinite, static, straight string stretched in the \( X^1 \) direction. We can solve for the background fields by coupling the equations of motion to a delta-function string
source. This is the same kind of calculation that we’re used to in electromagnetism. The resulting spacetime fields are given by

\[ ds^2 = f(r)^{-1} (-dt^2 + dX_1^2) + \sum_{i=2}^{25} dX_i^2 \]

\[ B = (f(r)^{-1} - 1) dt \wedge dX_1 \]

\[ e^{2\Phi} = f(r)^{-1} \quad (7.23) \]

The function \( f(r) \) depends only on the transverse direction \( r = \sum_{i=2}^{25} X_i^2 \) and is given by

\[ f(r) = 1 + \frac{g_s^2 N l_s^{22}}{r^{22}} \]

Here \( N \) is some constant which we will shortly demonstrate counts the number of strings which source the background. The string length scale in the solutions is \( l_s = \sqrt{\alpha'} \). The function \( f(r) \) has the property that it is harmonic in the space transverse to the string, meaning that it satisfies \( \nabla^2_{\mathbb{R}^{24}} f(r) = 0 \) except at \( r = 0 \).

Let’s compute the \( B \)-field charge of this solution. We do exactly what we do in electromagnetism: we integrate the total flux through a sphere which surrounds the object. The string lies along the \( X^1 \) direction so the transverse space is \( \mathbb{R}^{24} \). We can consider a sphere \( S^{23} \) at the boundary of this transverse space. We should be integrating the flux over this sphere. But what is the expression for the flux?

To see what we should do, let’s look at the action for \( H_{\mu\nu\rho} \) in the presence of a string source. We will use form notation since this is much cleaner and refer to \( H_{\mu\nu\rho} \) simply as \( H_3 \). Schematically, the action takes the form

\[ \frac{1}{g_s^2} \int_{\mathbb{R}^{26}} H_3 \wedge \ast H_3 + \int_{\mathbb{R}^{2}} B_2 = \frac{1}{g_s^2} \int_{\mathbb{R}^{26}} H_3 \wedge \ast H_3 + g_s^2 B_2 \wedge \delta(\omega) \]

Here \( \delta(\omega) \) is a delta-function source with support on the 2d worldsheet of the string. The equation of motion is

\[ d^{\ast} H_3 \sim g_s^2 \delta(\omega) \]

From this we learn that to compute the charge of a single string we need to integrate

\[ \frac{1}{g_s^2} \int_{S^{23}} \ast H_3 = 1 \]

After these general comments, we now return to our solution (7.23). The above discussion was schematic and no attention was paid to factors of 2 and \( \pi \). Keeping in this spirit, the flux of the solution (7.23) can be checked to be

\[ \frac{1}{g_s^2} \int_{S^{23}} \ast H_3 = N \]
This is telling us that the solution (7.23) describes the background sourced by $N$ coincident, parallel fundamental strings. Another way to check this is to compute the ADM mass per unit length of the solution: it is $NT \sim N/\alpha'$ as expected.

Note as far as the low-energy effective action is concerned, there is nothing that insists $N \in \mathbb{Z}$. This is analogous to the statement that nothing in classical Maxwell theory requires $e$ to be quantized. However, in string theory, as in QED, we know the underlying sources of the microscopic theory and $N$ must indeed take integer values.

Finally, notice that as $r \to 0$, the solution becomes singular. It is not to be trusted in this regime where higher order $\alpha'$ corrections become important.

### 7.4.3 Magnetic Branes

We’ve already seen that string theory is not just a theory of strings; there are also D-branes, defined as surfaces on which strings can end. We’ll have much more to say about D-branes in Section 7.5. Here, we will consider a third kind of object that exists in string theory. It is again a brane – meaning that it is extended in some number of spacetime directions — but it is not a D-brane because the open string cannot end there. In these lectures we will call it the magnetic brane.

#### Electric and Magnetic Charges

You’re probably not used to talking about magnetically charged objects in electromagnetism. Indeed, in undergraduate courses we usually don’t get much further than pointing out that $\nabla \cdot B = 0$ does not allow point-like magnetic charges. However, in the context of quantum field theory, much of the interesting behaviour often boils down to understanding how magnetic charges behave. And the same is true of string theory. Because this may be unfamiliar, let’s take a minute to discuss the basics.

In electromagnetism in $d = 3 + 1$ dimensions, we measure electric charge $q$ by integrating the electric field $\vec{E}$ over a sphere $S^2$ that surrounds the particle,

$$q = \int_{S^2} \vec{E} \cdot d\vec{S} = \int_{S^2} * F_2$$  

(7.24)

In the second equality we have introduced the notation of differential forms that we also used in the previous example to discuss the string solutions.

Suppose now that a particle carries magnetic charge $g$. This can be measured by integrating the magnetic field $\vec{B}$ over the same sphere. This means

$$g = \int_{S^2} \vec{B} \cdot d\vec{S} = \int_{S^2} F_2$$  

(7.25)
In $d = 3+1$ dimensions, both electrically and magnetically charged objects are particles. But this is not always true in any dimension! The reason that it holds in $4d$ is because both the field strength $F_2$ and the dual field strength $\star F_2$ are 2-forms. Clearly, this is rather special to four dimensions.

In general, suppose that we have a $p$-brane that is electrically charged under a suitable gauge field. As we discussed in Section 7.2.1, a $(p + 1)$-dimensional object naturally couples to a $(p + 1)$-form gauge potential $C_{p+1}$ through,

$$
\mu \int_W C_{p+1}
$$

where $\mu$ is the charge of the object, while $W$ is the worldvolume of the brane. The $(p + 1)$-form gauge potential has a $(p + 2)$-form field strength

$$
G_{p+2} = dC_{p+1}
$$

To measure the electric charge of the $p$-brane, we need to integrate the field strength over a sphere that completely surrounds the object. A $p$-brane in $D$-dimensions has a transverse space $\mathbb{R}^{D-p-1}$. We can integrate the flux over the sphere at infinity, which is $S^{D-p-2}$. And, indeed, the counting works out nicely because, in $D$ dimensions, the dual field strength is a $(D - p - 2)$-form, $\star G_{p+2} = \tilde{G}_{D-p-2}$, which we can happily integrate over the sphere to find the charge sitting inside,

$$
q = \int_{S^{D-p-2}} \star G_{p+2}
$$

This equation is the generalized version of (7.24)

Now let’s think about magnetic charges. The generalized version of (7.25) suggest that we should compute the magnetic charge by integrating $G_{p+2}$ over a sphere $S^{p+2}$. What kind of object sits inside this sphere to emit the magnetic charge? Doing the sums backwards, we see that it should be a $(D - p - 4)$-brane.

We can write down the coupling between the $(D - p - 4)$-brane and the field strength. To do so, we first need to introduce the magnetic gauge potential defined by

$$
\star G_{p+2} = \tilde{G}_{D-p-2} = d\tilde{C}_{D-p-3}
$$

(7.26)

We can then add the magnetic coupling to the worldvolume $\tilde{W}$ of a $(D - p - 4)$-brane simply by writing

$$
\tilde{\mu} \int_{\tilde{W}} \tilde{C}_{D-p-3}
$$

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where $\tilde{\mu}$ is the magnetic charge. Note that it’s typically not possible to write down a Lagrangian that includes both magnetically charged object and electrically charged objects at the same time. This would need us to include both $C_{p+1}$ and $\tilde{C}_{D-p-3}$ in the Lagrangian, but these are not independent fields: they’re related by the rather complicated differential equations (7.26).

**The Magnetic Brane in Bosonic String Theory**

After these generalities, let’s see what it means for the bosonic string. The fundamental string is a 1-brane and, as we saw in Section 7.2.1, carries electric charge under the 2-form $B$. The appropriate object carrying magnetic charge under $B$ is therefore a $(D - p - 4) = (26 - 1 - 4) = 21$-brane.

To stress a point: neither the fundamental string, nor the magnetic 21-brane are D-branes. They are not surfaces where strings can end. We are calling them *branes* only because they are extended objects.

The magnetic 21-brane of the bosonic string can be found as a solution to the low-energy equations of motion. The solution can be written in terms of the dual potential $\tilde{B}_{22}$ such that $d\tilde{B}_{22} = * dB$. It is

\[
d s^2 = \left(-dt^2 + \sum_{i=1}^{21} dX_i^2\right) + h(r) \left( dX_{22}^2 + \ldots + dX_{25}^2 \right)
\]

\[
The function $h(r)$ depends only on the radial direction in $\mathbb{R}^4$ transverse to the brane: $r^2 = \sum_{i=22}^{25} X_i^2$. It is a harmonic function in $\mathbb{R}^4$, given by

\[
h(r) = 1 + \frac{N l_s^2}{r^2}
\]

The role of this function in the metric (7.27) is to warp the transverse $\mathbb{R}^4$ directions. Distances get larger as you approach the brane and the origin, $r = 0$, is at infinite distance.

It can be checked that the solution carried $N$ units of magnetic charge and has tension

\[
T \sim \frac{N}{l_s^3} \frac{1}{g_s^2}
\]
Let’s summarize how the tension of different objects scale in string theory. The powers of $\alpha' = l_s^2$ are entirely fixed on dimensional grounds. (Recall that the tension is mass per spatial volume, so the tension of a $p$-brane has $[T_p] = p + 1$). More interesting is the dependence on the string coupling $g_s$. The tension of the fundamental string does not depend on $g_s$, while the magnetic brane scales as $1/g_s^2$. This kind of $1/g^2$ behaviour is typical of solitons in field theories. The D-branes sit between the two: their tension scales as $1/g_s$. Objects with this behaviour are somewhat rarer (although not unheard of) in field theory.

In the perturbative limit, $g_s \to 0$, both D-branes and magnetic branes are heavy. The coupling of an object with tension $T$ to gravity is governed by $T\kappa^2$ where the gravitational coupling scales as $\kappa \sim g_s^2$ (7.20). This means that in the weak coupling limit, the gravitational backreaction of the string and D-branes can be neglected. However, the coupling of the magnetic brane to gravity is always of order one.

The Magnetic Brane in Superstring Theory

Superstring theories also have a brane magnetically charged under $B$. It is a $(D - p - 4) = (10 - 1 - 4) = 5$-brane and is usually referred to as the NS5-brane. The solution in the transverse $\mathbb{R}^4$ again takes the form (7.27).

The NS5-brane exists in both type II and heterotic string. In many ways it is more mysterious than D-branes and its low-energy effective dynamics is still poorly understood. It is closely related to the 5-brane of M-theory.

7.4.4 Moving Away from the Critical Dimension

The beta function equations provide a new view on the critical dimension $D = 26$ of the bosonic string. To see this, let’s look more closely at the dilaton beta function $\beta(\Phi)$ defined in (7.15): it takes the same form as the Weyl anomaly that we discussed back in Section 4.4.2. This means that if we consider a string propagating in $D \neq 26$ then the Weyl anomaly simply arises as the leading order term in the dilaton beta function. So let’s relax the requirement of the critical dimension. The equations of motion arising from $\beta_{\mu \nu}(G)$ and $\beta_{\mu \nu}(B)$ are unchanged, while the dilaton beta function equation becomes

$$\beta(\Phi) = \frac{D - 26}{6} - \frac{\alpha'}{2} \nabla^2 \Phi + \alpha' \nabla_\mu \Phi \nabla^\mu \Phi - \frac{\alpha'}{24} H_{\mu \nu \lambda} H^{\mu \nu \lambda} = 0 \quad (7.28)$$

The low-energy effective action in string frame picks up an extra term which looks like a run-away potential for $\Phi$,

$$S = \frac{1}{2\kappa_0^2} \int d^D X \sqrt{-G} e^{-2\Phi} \left( R - \frac{1}{12} H_{\mu \nu \lambda} H^{\mu \nu \lambda} + 4 \partial_\mu \Phi \partial^\mu \Phi - \frac{2(D - 26)}{3\alpha'} \right)$$

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This sounds quite exciting. Can we really get string theory living in \( D = 4 \) dimensions so easily? Well, yes and no. Firstly, with this extra potential term, flat \( D \)-dimensional Minkowski space no longer solves the equations of motion. This is in agreement with the analysis in Section 2 where we showed that full Lorentz invariance was preserved only in \( D = 26 \).

Another, technical, problem with solving the string equations of motion this way is that we’re playing tree-level term off against a one-loop term. But if tree-level and one-loop terms are comparable, then typically all higher loop contributions will be as well and it is likely that we can’t trust our analysis.

**The Linear Dilaton CFT**

In fact, there is one simple solution to (7.28) which we can trust. It is the solution to

\[
\partial_\mu \Phi \partial^\mu \Phi = \frac{26 - D}{6\alpha'}
\]

Recall that we’re working in signature \((-,+,+,...)\), meaning that \( \Phi \) takes a spacelike profile if \( D < 26 \) and a timelike profile if \( D > 26 \),

\[
\Phi = \sqrt{\frac{26 - D}{6\alpha'}} X^1 \quad D < 26
\]

\[
\Phi = \sqrt{\frac{D - 26}{6\alpha'}} X^0 \quad D > 26
\]

This gives a dilaton which is linear in one direction. This can be compared to the study of the path integral for non-critical strings that we saw in 5.3.2. There are two ways of seeing the same physics.

The reason that we can trust this solution is that there is an exact CFT underlying it which we can analyze to all orders in \( \alpha' \). It’s called, for obvious reasons, the *linear dilaton CFT*. Let’s now look at this in more detail.

Firstly, consider the worldsheet action associated to the dilaton coupling. For now we’ll consider an arbitrary dilaton profile \( \Phi(X) \),

\[
S_{\text{dilaton}} = \frac{1}{4\pi} \int d^2 \sigma \sqrt{g} \Phi(X) R^{(2)}
\]

Although this term vanishes on a flat worldsheet, it nonetheless changes the stress-energy tensor \( T_{\alpha\beta} \) because this is defined as

\[
T_{\alpha\beta} = -4\pi \left. \frac{\partial S}{\partial g^{\alpha\beta}} \right|_{g_{\alpha\beta} = \delta_{\alpha\beta}}
\]
The variation of (7.29) is straightforward. Indeed, the term is akin to the Einstein-Hilbert term in general relativity but things are simpler in 2d because, for example \( R_{\alpha\beta} = \frac{1}{2} g_{\alpha\beta} R \). We have

\[
\delta (\sqrt{g} g^{\alpha\beta} R_{\alpha\beta}) = \sqrt{g} g^{\alpha\beta} \delta R_{\alpha\beta} = \sqrt{g} \nabla^\alpha v_\alpha
\]

where

\[
v_\alpha = \nabla^\beta \delta g_{\alpha\beta} - g^{\gamma\delta} \nabla_\alpha \delta g_{\gamma\delta}
\]

Using this, the variation of the dilaton term in the action is given by

\[
\delta S_{\text{dilaton}} = \frac{1}{4\pi} \int d^2 \sigma \sqrt{g} \left( \nabla^\alpha \nabla^\beta \Phi - \nabla^2 \Phi g^{\alpha\beta} \right) \delta g_{\alpha\beta}
\]

which, restricting to flat space \( g_{\alpha\beta} = \delta_{\alpha\beta} \), finally gives us the stress-energy tensor of a theory with dilaton coupling

\[
T_{\alpha\beta}^{\text{dilaton}} = -\partial_\alpha \partial_\beta \Phi + \partial^2 \Phi \delta_{\alpha\beta}
\]

Note that this stress tensor is not traceless. This is to be expected because, as we described above, the dilaton coupling is not Weyl invariant at tree-level. In complex coordinates, the stress tensor is

\[
T_{\alpha\beta}^{\text{dilaton}} = -\partial_\alpha \partial_\beta \Phi , \quad \bar{T}_{\alpha\beta}^{\text{dilaton}} = -\bar{\partial}^2 \Phi
\]

**Linear Dilaton OPE**

The stress tensor above holds for any dilaton profile \( \Phi(X) \). Let’s now restrict to a linear dilaton profile for a single scalar field \( X \),

\[
\Phi = QX
\]

where \( Q \) is some constant. We also include the standard kinetic terms for \( D \) scalar fields, of which \( X \) is a chosen one, giving the stress tensor

\[
T = -\frac{1}{\alpha'} : \partial X \partial X : -Q \partial^2 X
\]

It is a simple matter to compute the \( TT \) OPE using the techniques described in Section 4. We find,

\[
T(z) T(w) = \frac{c/2}{(z-w)^4} + \frac{2T(w)}{(z-w)^2} + \frac{\partial T(w)}{z-w} + \ldots
\]

where the central charge of the theory is given by

\[
c = D + 6\alpha' Q^2
\]

Note that \( Q^2 \) can be positive or negative depending on the whether we have a timelike or spacelike linear dilaton. In this way, we see explicitly how a linear dilaton gradient can absorb central charge.
7.4.5 The Elephant in the Room: The Tachyon

We’ve been waxing lyrical about the details of solutions to the low-energy effective action, all the while ignoring the most important, relevant field of them all: the tachyon. Since our vacuum is unstable, this is a little like describing all the beautiful pictures we could paint if only that damn paintbrush would balance, unaided, on its tip.

Of course, the main reason for discussing these solutions is that they all carry directly over to the superstring where the tachyon is absent. Nonetheless, it’s interesting to ask what happens if the tachyon is turned on. Its vertex operator is simply

\[ V_{\text{tachyon}} \sim \int d^2 \sigma \sqrt{g} e^{ipX} \]

where \( p^2 = 4/\alpha' \). Piecing together a general tachyon profile \( V(X) \) from these Fourier modes and exponentiating, results in a potential on the worldsheet of the string

\[ S_{\text{potential}} = \int d^2 \sigma \sqrt{g} \alpha' V(X) \]

This is a relevant operator for the worldsheet CFT. Whenever such a relevant operator turns on, we should follow the RG flow to the infra-red until we land on another CFT. The c-theorem tells us that \( c_{\text{IR}} < c_{\text{UV}} \), but in string theory we always require \( c = 26 \). The deficit, at least initially, is soaked up by the dilaton in the manner described above. The end point of the tachyon RG flow for the bosonic string is not understood. It may be that there is no end point and the bosonic string simply doesn’t make sense once the tachyon is turned on. Or perhaps we haven’t yet understood the true ground state of the bosonic string.

7.5 D-Branes Revisited: Background Gauge Fields

Understanding the constraints of conformal invariance on the closed string backgrounds led us to Einstein’s equations and the low-energy effective action in spacetime. Now we would like to do the same for the open string. We want to understand the restrictions that consistency places on the dynamics of D-branes.

We saw in Section 3 that there are two types of massless modes that arise from the quantization of an open string: scalars, corresponding to the fluctuation of the D-brane, and a \( U(1) \) gauge field. We will ignore the scalar fluctuations for now, but will return to them later. We focus initially on the dynamics of a gauge field \( A_a, a = 0, \ldots, p \) living on a Dp-brane.
The first question that we ask is: how does the end of the string react to a background gauge field? To answer this, we need to look at the vertex operator associated to the photon. It was given in (5.10)

\[ V_{\text{photon}} \sim \int_{\partial M} d\tau \, \zeta_a \, \partial^\tau X^a \, e^{ip \cdot X} \]

which is Weyl invariant and primary only if \( p^2 = 0 \) and \( p^a \zeta_a = 0 \). Exponentiating this vertex operator, as described at the beginning of Section 7, gives the coupling of the open string to a general background gauge field \( A_a(X) \),

\[ S_{\text{end-point}} = \int_{\partial M} d\tau \, A_a(X) \frac{dX^a}{d\tau} \]

But this is a very familiar coupling — we’ve already mentioned it in (7.9). It is telling us that the end of the string is charged under the background gauge field \( A_a \) on the brane.

### 7.5.1 The Beta Function

We can now perform the same type of beta function calculation that we saw for the closed string\(^9\). To do this, it’s useful to first use conformal invariance to map the open string worldsheet to the Euclidean upper-half plane as we described in Section 4.7. The action describing an open string propagating in flat space, with its ends subject to a background gauge field on the D-brane splits up into two pieces

\[ S = S_{\text{Neumann}} + S_{\text{Dirichlet}} \]

where \( S_{\text{Neumann}} \) describes the fluctuations parallel to the D\( p \)-brane and is given by

\[ S_{\text{Neumann}} = \frac{1}{4\pi\alpha'} \int_M d^2\sigma \, \partial^a X^a \partial_\alpha X^b \delta_{ab} + i \int_{\partial M} d\tau \, A_a(X) \dot{X}^a \]  
\[ \text{(7.30)} \]

Here \( a, b = 0, \ldots, p \). The extra factor of \( i \) arises because we are in Euclidean space. Meanwhile, the fields transverse to the brane have Dirichlet boundary conditions and take range \( I = p + 1, \ldots, D - 1 \). Their dynamics is given by

\[ S_{\text{Dirichlet}} = \frac{1}{4\pi\alpha'} \int_M d^2\sigma \, \partial^a X^I \partial_\alpha X^J \delta_{IJ} \]

\[^9\text{We’ll be fairly explicit here, but if you want to see more details then the best place to look is the original paper by Abouelsaood, Callan, Nappi and Yost, “Open Strings in Background Gauge Fields”, Nucl. Phys. B280 (1987) 599.}\]
The action $S_{\text{Dirichlet}}$ describes free fields and doesn’t play any role in the computation of the beta-function. The interesting part is $S_{\text{Neumann}}$ which, for non-zero $A_a(X)$, is an interacting quantum field theory with boundary. Our task is to compute the beta-function associated to the coupling $A_a(X)$. We use the same kind of technique that we earlier applied to the closed string. We expand the fields $X^a(\sigma)$ as

$$X^a(\sigma) = \bar{x}^a(\sigma) + \sqrt{\alpha'} Y^a(\sigma)$$

where $\bar{x}^a(\sigma)$ is taken to be some fixed background which obeys the classical equations of motion,

$$\partial^2 \bar{x}^a = 0$$

(In the analogous calculation for the closed string we chose the special case of $\bar{x}^a$ constant. Here we are more general). However, we also need to impose boundary conditions for this classical solution. In the absence of the gauge field $A_a$, we require Neumann boundary conditions $\partial_\sigma X^a = 0$ at $\sigma = 0$. However, the presence of the gauge field changes this. Varying the full action (7.30) shows that the relevant boundary condition is supplemented by an extra term,

$$\partial_\sigma \bar{x}^a + 2\pi \alpha' i F^{ab} \partial_\tau \bar{x}_b = 0 \quad \text{at } \sigma = 0 \quad (7.31)$$

where the $F_{ab}$ is the field strength

$$F_{ab}(X) = \frac{\partial A_b}{\partial X^a} - \frac{\partial A_a}{\partial X^b} = \partial_a A_b - \partial_b A_a$$

The fields $Y^a(\sigma)$ are the fluctuations which are taken to be small. Again, the presence of $\sqrt{\alpha'}$ in the expansion ensures that $Y^a$ are dimensionless. Expanding the action $S_{\text{Neumann}}$ (which we’ll just call $S$ from now on) to second order in fluctuations gives,

$$S[\bar{x} + \sqrt{\alpha'} Y] = S[\bar{x}] + \frac{1}{4\pi} \int_M d^2 \sigma \partial Y^a \partial Y^b \delta_{ab}$$

$$+ i\alpha' \int_{\partial M} d\tau \left( \partial_a A_b Y^a \dot{Y}^b + \frac{1}{2} \partial_a \partial_b A_c Y^a Y^b \dot{x}^c \right) + \ldots$$

where all expressions involving the background gauge fields are now evaluated on the classical solution $\bar{x}$. We can rearrange the boundary terms by splitting the first term up into two halves and integrating one of these pieces by parts,

$$\int d\tau (\partial_a A_b) Y^a \dot{Y}^b = \frac{1}{2} \int d\tau \partial_a A_b Y^a \dot{Y}^b - \partial_a A_b Y^a Y^b - \partial_c \partial_a A_b Y^a Y^b \dot{x}^c$$
Combining this with the second term means that we can write all interactions in terms of the gauge invariant field strength $F_{ab}$,

$$S[\bar{x} + \sqrt{\alpha'}Y] = S[\bar{x}] + \frac{1}{4\pi} \int_{\mathcal{M}} d^2 \sigma \, \partial Y^a \partial Y^b \delta_{ab} + \frac{i\alpha'}{2} \int_{\partial\mathcal{M}} d\tau \left( F_{ab} Y^a \dot{Y}^b + \partial_b F_{ac} Y^a Y^b \dot{x}^c \right) + \ldots \quad (7.32)$$

where the +\ldots refer to the higher terms in the expansion which come with higher derivatives of $F_{ab}$, accompanied by powers of $\alpha'$. We can neglect them for the purposes of computing the one-loop beta function.

The Propagator

This Lagrangian describes our interacting boundary theory to leading order. We can now use this to compute the beta function. Firstly, we should determine where possible divergences arise. The offending term is the last one in (7.32). This will lead to a divergence when the fluctuation fields $Y^a$ are contracted with their propagator

$$\langle Y^a(z, \bar{z}) Y^b(w, \bar{w}) \rangle = G_{ab}(z, \bar{z}; w, \bar{w})$$

We should be used to these free field Green’s functions by now. The propagator satisfies

$$\partial \bar{\partial} G_{ab}(z, \bar{z}) = -2\pi \delta^{ab} \delta(z, \bar{z}) \quad (7.33)$$

in the upper half plane. But now there’s a subtlety. The $Y^a$ fields need to satisfy a boundary condition at $\text{Im} \, z = 0$ and this should be reflected in the boundary condition for the propagator. We discussed this briefly for Neumann boundary conditions in Section 4.7. But we’ve also seen that the background field strength shifts the Neumann boundary conditions to (7.31). Correspondingly, the propagator $G(z, \bar{z}; w, \bar{w})$ must now satisfy

$$\partial_{\sigma} G^{ab}(z, \bar{z}; w, \bar{w}) + 2\pi \alpha' i F^a_c \partial_{\sigma} G^{cb}(z, \bar{z}; w, \bar{w}) = 0 \quad \text{at} \quad \sigma = 0 \quad (7.34)$$

In Section 4.7, we showed how Neumann boundary conditions could be imposed by considering an image charge in the lower half plane. A similar method works here. We extend $G^{ab} \equiv G^{ab}(z, \bar{z}; w, \bar{w})$ to the entire complex plane. The solution to (7.33) subject to (7.34) is given by

$$G^{ab} = -\delta^{ab} \ln|z - w| - \frac{1}{2} \left( \frac{1 - 2\pi \alpha' F}{1 + 2\pi \alpha' F} \right)^{ab} \ln(z - \bar{w}) - \frac{1}{2} \left( \frac{1 + 2\pi \alpha' F}{1 - 2\pi \alpha' F} \right)^{ab} \ln(\bar{z} - w)$$
The Counterterm and Beta Function

Let’s now return to the interacting theory (7.32) and see what counterterm is needed to remove the divergence. Since all interactions take place on the boundary, we should evaluate our propagator on the boundary, which means \( z = \bar{z} \) and \( w = \bar{w} \). In this case, all the logarithms become the same and, in the limit that \( z \to w \), gives the leading divergence \( \ln |z - w| \to \epsilon^{-1} \). We learn that the UV divergence takes the form,

\[
-\frac{1}{\epsilon} \left[ \delta^{ab} + \frac{1}{2} \left( \frac{1 - 2\pi\alpha' F}{1 + 2\pi\alpha' F} \right)^{ab} + \frac{1}{2} \left( \frac{1 + 2\pi\alpha' F}{1 - 2\pi\alpha' F} \right)^{ab} \right] = -\frac{2}{\epsilon} \left( \frac{1}{1 - 4\pi^2\alpha'^2 F^2} \right)^{ab}
\]

It’s now easy to determine the necessary counterterm. We simply replace \( Y^a Y^b \) in the final term with \( \langle Y^a Y^b \rangle \). This yields

\[
-\frac{i2\pi\alpha'^2}{\epsilon} \int_{\partial M} d\tau \, \partial_b F_{ac} \left[ \frac{1}{1 - 4\pi^2\alpha'^2 F^2} \right]^{ab} \dot{x}^c
\]

For the open string theory to retain conformal invariance, we need the associated beta function to vanish. This gives us the condition on the field strength \( F_{ab} \): it must satisfy the equation

\[
\partial_b F_{ac} \left[ \frac{1}{1 - 4\pi^2\alpha'^2 F^2} \right]^{ab} = 0 \quad (7.35)
\]

This is our final equation governing the equations of motion that \( F_{ab} \) must satisfy to provide a consistent background for open string propagation.

7.5.2 The Born-Infeld Action

Equation (7.35) probably doesn’t look too familiar! Following the path we took for the closed string, we wish to write down an action whose equations of motion coincide with (7.35). The relevant action was actually constructed many decades ago as a non-linear alternative to Maxwell theory: it goes by the name of the Born-Infeld action:

\[
S = -T_p \int d^{p+1}\xi \sqrt{-\det (\eta_{ab} + 2\pi\alpha' F_{ab})} \quad (7.36)
\]

Here \( \xi \) are the worldvolume coordinates on the brane and \( T_p \) is the tension of the \( Dp \)-brane (which, since it multiplies the action, doesn’t affect the equations of motion). The gauge potential is to be thought of as a function of the worldvolume coordinates: \( A_a = A_a(\xi) \). It actually takes a little work to show that the equations of motion that we derive from this action coincide with the vanishing of the beta function (7.35). Some hints on how to proceed are provided on Example Sheet 4.
For small field strengths, $F_{ab} \ll 1/\alpha'$, the action (7.36) coincides with Maxwell’s action. To see this, we need simply expand to get

$$S = -T_p \int d^{p+1}\xi \left( 1 + \frac{(2\pi \alpha')^2}{4} F_{ab} F^{ab} + \ldots \right)$$

The leading order term, quadratic in field strengths, is the Maxwell action. Terms with higher powers of $F_{ab}$ are suppressed by powers of $\alpha'$.

So, for small field strengths, the dynamics of the gauge field on a D-brane is governed by Maxwell’s equations. However, as the electric and magnetic field strengths increase and become of order $1/\alpha'$, non-linear corrections to the dynamics kick in and are captured by the Born-Infeld action.

The Born-Infeld action arises from the one-loop beta function. It is the exact result for constant field strengths. If we want to understand the dynamics of gauge fields with large gradients, $\partial F$, then we will have determine the higher loop contributions to the beta function.

### 7.6 The DBI Action

We’ve understood that the dynamics of gauge fields on the brane is governed by the Born-Infeld action. But what about the fluctuations of the brane itself. We looked at this briefly in Section 3.2 and suggested, on general grounds, that the action should take the Dirac form (3.6). It would be nice to show this directly by considering the beta function equations for the scalar fields $\phi^I$ on the brane. Turning these on corresponds to considering boundary conditions where the brane is bent. It is indeed possible to compute something along the lines of beta-function equations and to show directly that the fluctuations of the brane are governed by the Dirac action\(^{10}\).

More generally, one could consider both the dynamics of the gauge field and the fluctuation of the brane. This is governed by a mixture of the Dirac action and the Born-Infeld action which is usually referred to as the DBI action,

$$S_{DBI} = -T_p \int d^{p+1}\xi \sqrt{-\det(\gamma_{ab} + 2\pi \alpha' F_{ab})}$$

As in Section (3.2), $\gamma_{ab}$ is the pull-back of the the spacetime metric onto the worldvolume,

$$\gamma_{ab} = \frac{\partial X^\mu}{\partial \xi^a} \frac{\partial X^\nu}{\partial \xi^b} \eta_{\mu\nu}$$

\(^{10}\)A readable discussion of this calculation can be found in the original paper by Leigh, *Dirac-Born-Infeld Action from Dirichlet Sigma Model*, Mod. Phys. Lett. A4: 2767 (1989).
The new dynamical fields in this action are the embedding coordinates $X^\mu(\xi)$, with $\mu = 0, \ldots, D - 1$. This appears to be $D$ new degrees of freedom while we expect only $D - p - 1$ transverse physical degrees of freedom. The resolution to this should be familiar by now: the DBI action enjoys a reparameterization invariance which removes the longitudinal fluctuations of the brane.

We can use this reparameterization invariance to work in static gauge. For an infinite, flat $Dp$-brane, it is useful to set

$$X^a = \xi^a \quad a = 0, \ldots, p$$

so that the pull-back metric depends only on the transverse fluctuations $X^I$,

$$\gamma_{ab} = \eta_{ab} + \frac{\partial X^I}{\partial \xi^a} \frac{\partial X^J}{\partial \xi^b} \delta_{IJ}$$

If we are interested in situations with small field strengths $F_{ab}$ and small derivatives $\partial_a X$, then we can expand the DBI action to leading order. We have

$$S = - (2\pi\alpha')^2 T_p \int d^{p+1}\xi \left( \frac{1}{4} F_{ab} F^{ab} + \frac{1}{2} \partial_a \phi^I \partial^a \phi^I + \ldots \right)$$

where we have rescaled the positions to define the scalar fields $\phi^I = X^I / 2\pi\alpha'$. We have also dropped an overall constant term in the action. This is simply free Maxwell theory coupled to free massless scalar fields $\phi^I$. The higher order terms that we have dropped are all suppressed by powers of $\alpha'$.

### 7.6.1 Coupling to Closed String Fields

The DBI action describes the low-energy dynamics of a $Dp$-brane in flat space. We could now ask how the motion of the D-brane is affected if it moves in a background created by closed string modes $G_{\mu\nu}$, $B_{\mu\nu}$ and $\Phi$. Rather than derive this, we’ll simply write down the answer and then justify each term in turn. The answer is:

$$S_{DBI} = - T_p \int d^{p+1}\xi \ e^{-\Phi} \sqrt{- \det(\gamma_{ab} + 2\pi\alpha' F_{ab} + B_{ab})}$$

Let’s start with the coupling to the background metric $G_{\mu\nu}$. It’s actually hidden in the notation in this expression: it appears in the pull-back metric $\gamma_{ab}$ which is now given by

$$\gamma_{ab} = \frac{\partial X^\mu}{\partial \xi^a} \frac{\partial X^\nu}{\partial \xi^b} G_{\mu\nu}$$

It should be clear that this is indeed the natural place for it to sit.
Next up is the dilaton. As in (7.17), we have decomposed the dilaton into a constant piece and a varying piece: \( \Phi = \Phi_0 + \tilde{\Phi} \). The constant piece governs the asymptotic string coupling, \( g_s = e^{\Phi_0} \), and is implicitly sitting in front of the action because the tension of the D-brane scales as

\[
T_p \sim 1/g_s
\]

This, then, explains the factor of \( e^{-\tilde{\Phi}} \) in front of the action: it simply reunites the varying part of the dilaton with the constant piece. Physically, it’s telling us that the tension of the D-brane depends on the local value of the dilaton field, rather than its asymptotic value. If the dilaton varies, the effective string coupling at a point \( X \) in spacetime is given by \( g_s^{eff} = e^{\Phi(X)} = g_s e^{\tilde{\Phi}(X)} \). This, in turn, changes the tension of the D-brane. It can lower its tension by moving to regions with larger \( g_s^{eff} \).

Finally, let’s turn to the \( B_{\mu\nu} \) field. This is a 2-form in spacetime. The function \( B_{ab} \) appearing in the DBI action is the pull-back to the worldvolume

\[
B_{ab} = \frac{\partial X^\mu}{\partial \xi^a} \frac{\partial X^\nu}{\partial \xi^b} B_{\mu\nu}
\]

Its appearance in the DBI action is actually required on grounds of gauge invariance alone. This can be seen by considering an open string, moving in the presence of both a background \( B_{\mu\nu}(X) \) in spacetime and a background \( A_a(X) \) on the worldvolume of a brane. The relevant terms on the string worldsheet are

\[
\frac{1}{4\pi\alpha'} \int_M d^2\sigma \ \epsilon^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu B_{\mu\nu} + \int_{\partial M} d\tau \ A_a \dot{X}^a
\]

Under a spacetime gauge transformation

\[
B_{\mu\nu} \rightarrow B_{\mu\nu} + \partial_\mu C_\nu - \partial_\nu C_\mu \quad (7.37)
\]

the first term changes by a total derivative. This is fine for a closed string, but it doesn’t leave the action invariant for an open string because we pick up the boundary term. Let’s quickly look at what we get in more detail. Under the gauge transformation (7.37), we have

\[
S_B = \frac{1}{4\pi\alpha'} \int_M d^2\sigma \ \epsilon^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu B_{\mu\nu}
\]

\[
\rightarrow S_B + \frac{1}{2\pi\alpha'} \int_M d\sigma d\tau \ \epsilon^{\alpha\beta} \partial_\alpha X^\mu \partial_\beta X^\nu \partial_\mu C_\nu
\]

\[
= S_B + \frac{1}{2\pi\alpha'} \int_M d\sigma d\tau \ \epsilon^{\alpha\beta} \partial_\alpha (\partial_\beta X^\nu C_\nu)
\]

\[
= S_B + \frac{1}{2\pi\alpha'} \int_{\partial M} d\tau \dot{X}^\nu C_\nu = S_B + \frac{1}{2\pi\alpha'} \int_{\partial M} d\tau \dot{X}^a C_a
\]
where, in the last line, we have replaced the sum over all directions \( X^\nu \) with the sum over those directions obeying Neumann boundary conditions \( X^\nu \), since \( X^I = 0 \) at the end-points for any directions with Dirichlet boundary conditions.

The result of this short calculation is to see that the string action is not invariant under (7.37). To restore this spacetime gauge invariance, this boundary contribution must be canceled by an appropriate shift of \( A_a \) in the second term,

\[
A_a \rightarrow A_a - \frac{1}{2\pi \alpha'} C_a
\]

Note that this is not the usual kind of gauge transformation that we consider in electrodynamics. In particular, the field strength \( F_{ab} \) is not invariant. Rather, the gauge invariant combination under (7.37) and (7.38) is

\[
B_{ab} + 2\pi \alpha' F_{ab}
\]

This is the reason that this combination must appear in the DBI action. This is also related to an important physical effect. We have already seen that the string in spacetime is charged under \( B_{\mu\nu} \). But we’ve also seen that the end of the string is charged under the gauge field \( A_a \) on the D-brane. This means that the open string deposits \( B \) charge on the brane, where it is converted into \( A \) charge. The fact that the gauge invariant field strength involves a combination of both \( F_{ab} \) and \( B_{ab} \) is related to this interplay of charges.

7.7 The Yang-Mills Action

Finally, let’s consider the case of \( N \) coincident D-branes. We discussed this in Section 3.3 where we showed that the massless fields on the brane could be naturally packaged as \( N \times N \) Hermitian matrices, with the element of the matrix telling us which brane the end points terminate on. The gauge field then takes the form

\[
(A_a)^m_n
\]

with \( a = 0, \ldots, p \) and \( m, n = 1, \ldots, N \). Written this way, it looks rather like a \( U(N) \) gauge connection. Indeed, this is the correct interpretation. But how do we see this? Why is the gauge field describing a \( U(N) \) gauge symmetry rather than, say, \( U(1)^N \)?

The quickest way to see that coincident branes give rise to a \( U(N) \) gauge symmetry is to recall that the end point of the string is charged under the \( U(1) \) gauge field that inhabits the brane it’s ending on. Let’s illustrate this with the simplest example. Suppose that we have two branes. The diagonal components \( (A_a)^1_1 \) and \( (A_a)^2_2 \) arise
from strings which begin and end on the same brane. Each is a $U(1)$ gauge field. What about the off-diagonal terms $(A_a)^{12}$ and $(A_a)^{21}$? These come from strings stretched between the two branes. They are again massless gauge bosons, but they are charged under the two original $U(1)$ symmetries; they carry charge $(+1,-1)$ and $(-1,+1)$ respectively. But this is precisely the structure of a $U(2)$ gauge theory, with the off-diagonal terms playing a role similar to W-bosons. In fact, the only way to make sense of massless, charged spin 1 particles is through non-Abelian gauge symmetry.

So the massless excitations of $N$ coincident branes are a $U(N)$ gauge field $(A_a)^m_n$, together with scalars $(\phi^I)^m_n$ which transform in the adjoint representation of the $U(N)$ gauge group. We saw in Section 3 that the diagonal components $(\phi^I)^m_m$ have the interpretation of the transverse fluctuations of the $m$th brane. Can we now write down an action describing the interactions of these fields?

In fact, there are several subtleties in writing down a non-Abelian generalization of the DBI action and such an action is not known (if, indeed, it makes sense at all). However, we can make progress by considering the low-energy limit, corresponding to small field strengths. The field strength in question is now the appropriate non-Abelian expression which, neglecting the matrix indices, reads

$$F_{ab} = \partial_a A_b - \partial_b A_a + i [A_a, A_b]$$

The low-energy action describing the dynamics of $N$ coincident D$p$-branes can be shown to be (neglecting an overall constant term),

$$S = -(2\pi\alpha')^2 T_p \int d^{p+1}\xi \ Tr \left( \frac{1}{2} F_{ab} F^{ab} + \frac{1}{2} D_a \phi^I D^a \phi^I - \frac{1}{4} \sum_{I \neq J} [\phi^I, \phi^J]^2 \right)$$

We recognize the first term as the $U(N)$ Yang-Mills action. The coefficient in front of the Yang-Mills action is the coupling constant $1/g_{YM}^2$. For a D$p$-brane, this is given by $\alpha'^2 T_p$, or

$$g_{YM}^2 \sim l_s^{p-3} g_s$$

The kinetic term for $\phi^I$ simply reflects the fact that these fields transform in the adjoint representation of the gauge group,

$$D_a \phi^I = \partial_a \phi^I + i [A_a, \phi^I]$$

We won’t derive this action in these lectures: the first two terms basically follow from gauge invariance alone. The potential term is harder to see directly: the quick ways to derive it use T-duality or, in the case of the superstring, supersymmetry.
A flat, infinite Dp-brane breaks the Lorentz group of spacetime to

\[ S(1, D - 1) \rightarrow SO(1, p) \times SO(D - p - 1) \quad (7.40) \]

This unbroken group descends to the worldvolume of the D-brane where it classifies all low-energy excitations of the D-brane. The \( SO(1, p) \) is simply the Lorentz group of the D-brane worldvolume. The \( SO(D - p - 1) \) is a global symmetry of the D-brane theory, rotating the scalar fields \( \phi^I \).

The potential term in (7.39) is particularly interesting,

\[ V = -\frac{1}{4} \sum_{I \neq J} \text{Tr} [\phi^I, \phi^J]^2 \]

The potential is positive semi-definite. We can look at the fields that can be turned on at no cost of energy, \( V = 0 \). This requires that all \( \phi^I \) commute which means that, after a suitable gauge transformation, they take the diagonal form,

\[ \phi^I = \begin{pmatrix} \phi^I_1 \\ \vdots \\ \phi^I_N \end{pmatrix} \quad (7.41) \]

The diagonal component \( \phi^I_n \) describes the position of the \( n^{\text{th}} \) brane in transverse space \( \mathbb{R}^{D-p-1} \). We still need to get the dimensions right. The scalar fields have dimension \([\phi] = 1\). The relationship to the position in space (which we mentioned before in 3.2) is

\[ \vec{X}_n = 2\pi \alpha' \vec{\phi}_n \quad (7.42) \]

where we’ve swapped to vector notation to replace the \( I \) index.

The eigenvalues \( \phi^I_n \) are not quite gauge invariant: there is a residual gauge symmetry — the Weyl group of \( U(N) \) — which leaves \( \phi^I \) in the form (7.41) but permutes the entries by \( S_N \), the permutation group of \( N \) elements. But this has a very natural interpretation: it is simply telling us that the D-branes are indistinguishable objects.

When all branes are separated, the vacuum expectation value (7.41) breaks the gauge group from \( U(N) \rightarrow U(1)^N \). The W-bosons gain a mass \( M_W \) through the Higgs mechanism. Let’s compute this mass. We’ll consider a \( U(2) \) theory and we’ll separate
the two D-branes in the direction $X^D \equiv X$. This means that we turn on a vacuum expectation value for $\phi^D = \phi$, which we write as

$$\phi = \begin{pmatrix} \phi_1 & 0 \\ 0 & \phi_2 \end{pmatrix}$$

(7.43)

The values of $\phi_1$ and $\phi_2$ are the positions of the first and second brane. Or, more precisely, we need to multiply by the conversion factor $2\pi\alpha'$ as in (7.42) to get the position $X_m$ of the $m = 1^{st}, 2^{nd}$ brane,

Let’s compute the mass of the W-boson from the Yang-Mills action (7.39). It comes from the covariant derivative terms $D\phi$. We expand out the gauge field as

$$A_a = \begin{pmatrix} A_{a1}^1 & W_a \\ W_a^\dagger & A_{a2}^2 \end{pmatrix}$$

with $A_{a1}$ and $A_{a2}$ describing the two $U(1)$ gauge fields and $W$ the W-boson. The mass of the W-boson comes from the $[A_a, \phi]$ term inside the covariant derivative which, using the expectation value (7.43), is given by

$$\frac{1}{2} \text{Tr} [A_a, \phi]^2 = -(\phi_2 - \phi_1)^2 |W_a|^2$$

This gives us the mass of the W-boson: it is

$$M_{W}^2 = (\phi_2 - \phi_1)^2 = T^2 |X_2 - X_1|^2$$

where $T = 1/2\pi\alpha'$ is the tension of the string. But this has a very natural interpretation. It is precisely the mass of a string stretched between the two D-branes as shown in the figure above. We see that D-branes provide a natural geometric interpretation of the Higgs mechanism using adjoint scalars.

Notice that when branes are well separated, and the strings that stretch between them are heavy, their positions are described by the diagonal elements of the matrix given in (7.41). However, as the branes come closer together, these stretched strings become light and are important for the dynamics of the branes. Now the positions of the branes should be described by the full $N \times N$ matrices, including the off-diagonal elements. In this manner, D-branes begin to see space as something non-commutative at short distances.
In general, we can consider $N$ D-branes located at positions $\vec{X}_m$, $m = 1, \ldots, N$ in transverse space. The string stretched between the $m^{th}$ and $n^{th}$ brane has mass

$$M_W = |\vec{\phi}_n - \vec{\phi}_m| = T|\vec{X}_n - \vec{X}_m|$$

which again coincides with the mass of the appropriate W-boson computed using (7.39).

### 7.7.1 D-Branes in Type II Superstring Theories

As we mentioned previously, D-branes are ingredients of the Type II superstring theories. Type IIA has $D_p$-branes with $p$ even, while Type IIB is home to $D_p$-branes with $p$ odd. The D-branes have a very important property in these theories: they preserve half the supersymmetries.

Let’s take a moment to explain what this means. We’ll start by returning to the Lorentz group $SO(1, D - 1)$ now, of course, with $D = 10$. We’ve already seen that an infinite, flat Dp-brane is not invariant under the full Lorentz group, but only the subgroup (7.40). If we act with either $SO(1, p)$ or $SO(D - p - 1)$ then the D-brane solution remains invariant. We say that these symmetries are preserved by the solution.

However, the role of the preserved symmetries doesn’t stop there. The next step is to consider small excitations of the D-brane. These must fit into representations of the preserved symmetry group (7.40). This ensures that the low-energy dynamics of the D-brane must be governed by a theory which is invariant under (7.40) and we have indeed seen that the Lagrangian (7.39) has $SO(1, p)$ as a Lorentz group and $SO(D - p - 1)$ as a global symmetry group which rotates the scalar fields.

Now let’s return to supersymmetry. The Type II string theories enjoy a lot of supersymmetry: 32 supercharges in total. The infinite, flat D-branes are invariant under half of these; if we act with one half of the supersymmetry generators, the D-brane solutions don’t change. Objects that have this property are often referred to as BPS states. Just as with the Lorentz group, these unbroken symmetries descend to the worldvolume of the D-brane. This means that the low-energy dynamics of the D-branes is described by a theory which is itself invariant under 16 supersymmetries.

There is a unique class of theories with 16 supersymmetries and a non-Abelian gauge field and matter in the adjoint representation. This class is known as maximally supersymmetric Yang-Mills theory and the bosonic part of the action is given by (7.39). Supersymmetry is realized only after the addition of fermionic fields which also live on the brane. These theories describe the low-energy dynamics of multiple D-branes.
As an illustrative example, consider D3-branes in the Type IIB theory. The theory describing $N$ D-branes is $U(N)$ Yang-Mills with 16 supercharges, usually referred to as $U(N)\mathcal{N} = 4$ super-Yang-Mills. The bosonic part of the action is given by (7.39), where there are $D - p - 1 = 6$ scalar fields $\phi^I$ in the adjoint representation of the gauge group. These are augmented with four Weyl fermions, also in the adjoint representation.
8. Compactification and T-Duality

In this section, we will consider the simplest compactification of the bosonic string: a background spacetime of the form

\[ R^{1,24} \times S^1 \]  

(8.1)

The circle is taken to have radius \( R \), so that the coordinate on \( S^1 \) has periodicity

\[ X^{25} \equiv X^{25} + 2\pi R \]

We will initially be interested in the physics at length scales \( \gg R \) where motion on the \( S^1 \) can be ignored. Our goal is to understand what physics looks like to an observer living in the non-compact \( R^{1,24} \) Minkowski space. This general idea goes by the name of Kaluza-Klein compactification. We will view this compactification in two ways: firstly from the perspective of the spacetime low-energy effective action and secondly from the perspective of the string worldsheet.

8.1 The View from Spacetime

Let’s start with the low-energy effective action. Looking at length scales \( \gg R \) means that we will take all fields to be independent of \( X^{25} \): they are instead functions only on the non-compact \( R^{1,24} \).

Consider the metric in Einstein frame. This decomposes into three different fields on \( R^{1,24} \), a metric \( \tilde{G}_{\mu\nu} \), a vector \( A_\mu \) and a scalar \( \sigma \) which we package into the \( D = 26 \) dimensional metric as

\[ ds^2 = \tilde{G}_{\mu\nu} \, dX^\mu \, dX^\nu + e^{2\sigma} \left( dX^{25} + A_\mu \, dX^\mu \right)^2 \]  

(8.2)

Here all the indices run over the non-compact directions \( \mu, \nu = 0, \ldots, 24 \) only.

The vector field \( A_\mu \) is an honest gauge field, with the gauge symmetry descending from diffeomorphisms in \( D = 26 \) dimensions. To see this recall that under the transformation \( \delta X^\mu = V^\mu(X) \), the metric transforms as

\[ \delta G_{\mu\nu} = \nabla_\mu \Lambda_\nu + \nabla_\nu \Lambda_\mu \]

This means that diffeomorphisms of the compact direction, \( \delta X^{25} = \Lambda(X^\mu) \), turn into gauge transformations of \( A_\mu \),

\[ \delta A_\mu = \partial_\mu \Lambda \]
We’d like to know how the fields $G_{\mu \nu}$, $A_{\mu}$ and $\sigma$ interact. To determine this, we simply insert the ansatz (8.2) into the $D = 26$ Einstein-Hilbert action. The $D = 26$ Ricci scalar $\mathcal{R}^{(26)}$ is given by

$$\mathcal{R}^{(26)} = \mathcal{R} - 2e^{-\sigma} \nabla^2 e^\sigma - \frac{1}{4} e^{2\sigma} F_{\mu \nu} F^{\mu \nu}$$

where $\mathcal{R}$ in this formula now refers to the $D = 25$ Ricci scalar. The action governing the dynamics becomes

$$S = \frac{1}{2\kappa^2} \int d^{26} X \sqrt{-G^{(26)}} \mathcal{R}^{(26)} = \frac{2\pi R}{2\kappa^2} \int d^{25} X \sqrt{-G} e^\sigma \left( \mathcal{R} - \frac{1}{4} e^{2\sigma} F_{\mu \nu} F^{\mu \nu} + \partial_\mu \sigma \partial^\mu \sigma \right)$$

The dimensional reduction of Einstein gravity in $D$ dimensions gives Einstein gravity in $D - 1$ dimensions, coupled to a $U(1)$ gauge theory and a single massless scalar. This illustrates the original idea of Kaluza and Klein, with Maxwell theory arising naturally from higher-dimensional gravity.

The gravitational action above is not quite of the Einstein-Hilbert form. We need to again change frames, absorbing the scalar $\sigma$ in the same manner as we absorbed the dilaton in Section 7.3.1. Moreover, just as for the dilaton, there is no potential dictating the vacuum expectation value of $\sigma$. Changing the vev of $\sigma$ corresponds to changing $R$, so this is telling us that nothing in the gravitational action fixes the radius $R$ of the compact circle. This is a problem common to all Kaluza-Klein compactifications: there are always massless scalar fields, corresponding to the volume of the internal space as well as other deformations. Massless scalar fields, such as the dilaton $\Phi$ or the volume $\sigma$, are usually referred to as moduli.

If we want this type of Kaluza-Klein compactification to describe our universe — where we don’t see massless scalar fields — we need to find a way to “fix the moduli”. This means that we need a mechanism which gives rise to a potential for the scalar fields, making them heavy and dynamically fixing their vacuum expectation value. Such mechanisms exist in the context of the superstring.

Let’s now also look at the Kaluza-Klein reduction of the other fields in the low-energy effective action. The dilaton is easy: a scalar in $D$ dimensions reduces to a scalar in $D - 1$ dimensions. The anti-symmetric 2-form has more structure: it reduces to a 2-form $B_{\mu \nu}$, together with a vector field $\tilde{A}_\mu = B_{\mu 25}$.

The description of compactification on more general manifolds is a beautiful story involving aspects of differential geometry and topology. This story is told in the second volume of Green, Schwarz and Witten.
In summary, the low-energy physics of the bosonic string in $D-1$ dimensions consists of a metric $G_{\mu\nu}$, two $U(1)$ gauge fields $A_{\mu}$ and $\tilde{A}_{\mu}$ and two massless scalars $\Phi$ and $\sigma$.

### 8.1.1 Moving around the Circle

In the above discussion, we assumed that all fields are independent of the periodic direction $X^{25}$. Let’s now look at what happens if we relax this constraint. It’s simplest to see the resulting physics if we look at the scalar field $\Phi$ where we don’t have to worry about cluttering equations with indices. In general, we can expand this field in Fourier modes around the circle

$$\Phi(X^\mu; X^{25}) = \sum_{n=-\infty}^{\infty} \Phi_n(X^\mu)e^{inX^{25}/R}$$

where reality requires $\Phi^*_n = \Phi_{-n}$. Ignoring the coupling to gravity for now, the kinetic terms for this scalar are

$$\int d^{26}X \partial_\mu \Phi \partial^\mu \Phi + (\partial_{25}\Phi)^2 = 2\pi R \int d^{25}X \sum_{n=-\infty}^{\infty} \left( \partial_\mu \Phi_n \partial^\mu \Phi_{-n} + \frac{n^2}{R^2} |\Phi_n|^2 \right)$$

This simple Fourier decomposition is telling us something very important: a single scalar field on $\mathbb{R}^{1,D-1} \times S^1$ splits into an infinite number of scalar fields on $\mathbb{R}^{1,D-2}$, indexed by the integer $n$. These have mass

$$M_n^2 = \frac{n^2}{R^2} \quad (8.3)$$

For $R$ small, all particles are heavy except for the massless zero mode $n = 0$. The heavy particles are typically called Kaluza-Klein (KK) modes and can be ignored if we’re probing energies $\ll 1/R$ or, equivalently, distance scales $\gg R$.

There is one further interesting property of the KK modes $\Phi_n$ with $n \neq 0$: they are charged under the gauge field $A_{\mu}$ arising from the metric. The simplest way to see this is to look at the appropriate gauge transformation which, from the spacetime perspective, is the diffeomorphism $X^{25} \rightarrow X^{25} + \Lambda(X^\mu)$. Clearly, this shifts the KK modes

$$\Phi_n \rightarrow \exp \left( \frac{i n \Lambda}{R} \right) \Phi_n$$

This tells us that the $n^{th}$ KK mode has charge $n/R$. In fact, one usually rescales the gauge field to $A_{\mu}' = A_{\mu}/R$, under which the charge of the KK mode $\Phi_n$ is simply $n \in \mathbb{Z}$.
8.2 The View from the Worldsheets

We now consider the Kaluza-Klein reduction from the perspective of the string. We want to study a string moving in the background $\mathbb{R}^{1,24} \times S^1$. There are two ways in which the compact circle changes the string dynamics.

The first effect of the circle is that the spatial momentum, $p$, of the string in the circle direction can no longer take any value, but is quantized in integer units

$$ p^{25} = \frac{n}{R}, \quad n \in \mathbb{Z} $$

The simplest way to see this is simply to require that the string wavefunction, which includes the factor $e^{i p \cdot X}$, is single valued.

The second effect is that we can allow more general boundary conditions for the mode expansion of $X$. As we move around the string, we no longer need $X(\sigma + 2\pi) = X(\sigma)$, but can relax this to

$$ X^{25}(\sigma + 2\pi) = X^{25}(\sigma) + 2\pi m R, \quad m \in \mathbb{Z} $$

The integer $m$ tells us how many times the string winds around $S^1$. It is usually simply called the winding number.

Let’s now follow the familiar path that we described in Section 2 to study the spectrum of the string on the spacetime (8.1). We start by considering only the periodic field $X^{25}$, highlighting the differences with our previous treatment. The mode expansion of $X^{25}$ is now given by

$$ X^{25}(\sigma, \tau) = x^{25} + \frac{\alpha'}{R} \tau + m R \sigma + \text{oscillator modes} $$

which incorporates both the quantized momentum and the possibility of a winding number. Before splitting $X^{25}(\sigma, \tau)$ into right-moving and left-moving parts, it will be useful to introduce the quantities

$$ p_L = \frac{n}{R} + \frac{m R}{\alpha'}, \quad p_R = \frac{n}{R} - \frac{m R}{\alpha'} \quad (8.4) $$

Then we have $X^{25}(\sigma, \tau) = X_L^{25}(\sigma^+) + X_R^{25}(\sigma^-)$, where

$$ X_L^{25}(\sigma^+) = \frac{1}{2} x^{25} + \frac{1}{2} \alpha' p_L \sigma^+ + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \tilde{\alpha}_n^{25} e^{-i n \sigma^+}, $$

$$ X_R^{25}(\sigma^-) = \frac{1}{2} x^{25} + \frac{1}{2} \alpha' p_R \sigma^- + i \sqrt{\frac{\alpha'}{2}} \sum_{n \neq 0} \frac{1}{n} \alpha_n^{25} e^{-i n \sigma^-} $$
This differs from the mode expansion (1.36) only in the terms $p_L$ and $p_R$. The mode expansion for all the other scalar fields on flat space $\mathbb{R}^{1,24}$ remains unchanged and we don’t write them explicitly.

Let’s think about what the spectrum of this theory looks like to an observer living in $D = 25$ non-compact directions. Each particle state will be described by a momentum $p^\mu$ with $\mu = 0, \ldots, 24$. The mass of the particle is

$$M^2 = -\sum_{\mu=0}^{24} p_\mu p^\mu$$

As before, the mass of these particles is fixed in terms of the oscillator modes of the string by the $L_0$ and $\tilde{L}_0$ equations. These now read

$$M^2 = p_L^2 + \frac{4}{\alpha'}(\tilde{N} - 1) = p_R^2 + \frac{4}{\alpha'}(N - 1)$$

where $N$ and $\tilde{N}$ are the levels, defined in lightcone quantization by (2.24). (One should take the lightcone coordinate inside $\mathbb{R}^{1,24}$ rather than along the $S^1$). The factors of $-1$ are the necessary normal ordering coefficients that we’ve seen in several guises in this course.

These equations differ from (2.25) by the presence of the momentum and winding terms around $S^1$ on the right-hand side. In particular, level matching no longer tells us that $N = \tilde{N}$, but instead

$$N - \tilde{N} = nm \quad (8.5)$$

Expanding out the mass formula, we have

$$M^2 = \frac{n^2}{R^2} + \frac{m^2 R^2}{\alpha'^2} + \frac{2}{\alpha'}(N + \tilde{N} - 2) \quad (8.6)$$

The new terms in this formula have a simple interpretation. The first term tells us that a string with $n > 0$ units of momentum around the circle gains a contribution to its mass of $n/R$. This agrees with the result (8.3) that we found from studying the KK reduction of the spacetime theory. The second term is even easier to understand: a string which winds $m > 0$ times around the circle picks up a contribution $2\pi m R T$ to its mass, where $T = 1/2\pi \alpha'$ is the tension of the string.
8.2.1 Massless States

We now restrict attention to the massless states in $\mathbf{R}^{1,24}$. This can be achieved in the mass formula (8.6) by looking at states with zero momentum $n = 0$ and zero winding $m = 0$, obeying the level matching condition $N = \tilde{N} = 1$. The possibilities are

- $\alpha_{-1}^{\mu} \tilde{\alpha}_{-1}^{\nu} |0; p\rangle$: Under the $SO(1, 24)$ Lorentz group, these states decompose into a metric $G_{\mu\nu}$, an anti-symmetric tensor $B_{\mu\nu}$ and a scalar $\Phi$.

- $\alpha_{-1}^{\mu} \tilde{\alpha}_{-1}^{25} |0; p\rangle$ and $\alpha_{-1}^{25} \tilde{\alpha}_{-1}^{\mu} |0; p\rangle$: These are two vector fields. We can identify the sum of these $(\alpha_{-1}^{\mu} \tilde{\alpha}_{-1}^{25} + \alpha_{-1}^{25} \tilde{\alpha}_{-1}^{\mu}) |0; p\rangle$ with the vector field $A_\mu$ coming from the metric and the difference $(\alpha_{-1}^{\mu} \tilde{\alpha}_{-1}^{25} - \alpha_{-1}^{25} \tilde{\alpha}_{-1}^{\mu}) |0; p\rangle$ with the vector field $\tilde{A}_\mu$ coming from the anti-symmetric field.

- $\alpha_{-1}^{25} \tilde{\alpha}_{-1}^{25} |0; p\rangle$: This is another scalar. It is identified with the scalar $\sigma$ associated to the radius of $S^1$.

We see that the massless spectrum of the string coincides with the massless spectrum associated with the Kaluza-Klein reduction of the previous section.

8.2.2 Charged Fields

One can also check that the KK modes with $n \neq 0$ have charge $n$ under the gauge field $A_\mu$. We can determine the charge of a state under a given $U(1)$ by computing the 3-point function in which two legs correspond to the state of interest, while the third is the appropriate photon. We have two photons, with vertex operators given by,

$$V_\pm(p) \sim \int d^2 z \, \zeta_\mu(\partial X^\mu \bar{\partial} \bar{X}^{25} \mp \partial X^{25} \bar{\partial} \bar{X}^\mu)e^{ip \cdot X}$$

where $+$ corresponds to $A_\mu$ and $-$ to $\tilde{A}_\mu$ and we haven’t been careful about the overall normalization. Meanwhile, any state can be assigned momentum $n$ and winding $m$ by dressing the operator with the factor $e^{ip_L X^{25}(z) + ip_R \bar{X}^{25}(\bar{z})}$. As always, it’s simplest to work with the momentum and winding modes of the tachyon, whose vertex operators are of the form

$$V_{m,n}(p) \sim \int d^2 z \, e^{ip \cdot X} e^{ip_L X^{25} + ip_R \bar{X}^{25}}$$

The charge of a state is the coefficient in front of the 3-point coupling of the field and the photon,

$$\langle V_\pm(p_1) V_{m,n}(p_2) V_{-m,-n}(p_3) \rangle \sim \delta^{25}(\sum_i p_i) \zeta_\mu(p_2^\mu - p_3^\mu)(p_L \pm p_R)$$
The first few factors are merely kinematical. The interesting information is in the last factor. It is telling us that under $A_{\mu}$, fields have charge $p_L + p_R \sim n/R$. This is in agreement with the Kaluza-Klein analysis that we saw before. However, it’s also telling us something new: under $\tilde{A}_{\mu}$, fields have charge $p_L - p_R \sim mR/\alpha'$. In other words, winding modes are charged under the gauge field that arises from the reduction of $B_{\mu\nu}$. This is not surprising: winding modes correspond to strings wrapping the circle and we saw in Section 7 that strings are electrically charged under $B_{\mu\nu}$.

### 8.2.3 Enhanced Gauge Symmetry

With a circle in the game, there are other ways to build massless states that don’t require us to work at level $N = \tilde{N} = 1$. For example, we can set $N = \tilde{N} = 0$ and look at winding modes $m \neq 0$. The level matching condition (8.5) requires $n = 0$ and the mass of the states is

$$
M^2 = \left( \frac{mR}{\alpha'} \right)^2 - \frac{4}{\alpha'}
$$

and states can be massless whenever the radius takes special values $R^2 = 4\alpha'/m^2$ with $m \in \mathbb{Z}$. Similarly, we can set the winding to zero $m = 0$ and consider the KK modes of the tachyon which have mass

$$
M^2 = \frac{n^2}{R^2} - \frac{4}{\alpha'}
$$

which become massless when $R^2 = n^2\alpha'/4$.

However, the richest spectrum of massless states occurs when the radius takes a very special value, namely

$$R = \sqrt{\alpha'}$$

Solutions to the level matching condition (8.5) with $M^2 = 0$ are now given by

- $N = \tilde{N} = 1$ with $m = n = 0$. These give the states described above: a metric, two $U(1)$ gauge fields and two neutral scalars.

- $N = \tilde{N} = 0$ with $n = \pm 2$ and $m = 0$. These are KK modes of the tachyon field. They are scalars in spacetime with charges $(\pm 2, 0)$ under the $U(1) \times U(1)$ gauge symmetry.

- $N = \tilde{N} = 0$ with $n = 0$ and $m = \pm 2$. This is a winding mode of the tachyon field. They are scalars in spacetime with charges $(0, \pm 2)$ under $U(1) \times U(1)$. 

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\( N = 1 \) and \( \tilde{N} = 0 \) with \( n = m = \pm 1 \). These are two new spin 1 fields, \( \alpha_{-1}^{\mu} |0;p\rangle \). They carry charge \((\pm 1, \pm 1)\) under the two \( U(1) \times U(1) \).

\( N = 1 \) and \( \tilde{N} = 0 \) with \( n = -m = \pm 1 \). These are a further two spin 1 fields, \( \tilde{\alpha}_{-1}^{\mu} |0;p\rangle \), with charge \((\pm 1, \mp 1)\) under \( U(1) \times U(1) \).

How do we interpret these new massless states? Let’s firstly look at the spin 1 fields. These are charged under \( U(1) \times U(1) \). As we mentioned in Section 7.7, the only way to make sense of charged massless spin 1 fields is in terms of a non-Abelian gauge symmetry. Looking at the charges, we see that at the critical radius \( R = \sqrt{\alpha'} \), the theory develops an enhanced gauge symmetry

\[
U(1) \times U(1) \rightarrow SU(2) \times SU(2)
\]

The massless scalars from the \( N = \tilde{N} = 0 \) now join with the previous scalars to form adjoint representations of this new symmetry. We move away from the critical radius by changing the vacuum expectation value for \( \sigma \). This breaks the gauge group back to the Cartan subalgebra by the Higgs mechanism.

From the discussion above, it’s clear that this mechanism for generating non-Abelian gauge symmetries relies on the existence of the tachyon. For this reason, this mechanism doesn’t work in Type II superstring theories. However, it turns out that it does work in the heterotic string, even though it has no tachyon in its spectrum.

### 8.3 Why Big Circles are the Same as Small Circles

The formula (8.6) has a rather remarkable property: it is invariant under the exchange

\[
R \leftrightarrow \frac{\alpha'}{R}
\]

if, at the same time, we swap the quantum numbers

\[
m \leftrightarrow n
\]

This means that a string moving on a circle of radius \( R \) has the same spectrum as a string moving on a circle of radius \( \alpha'/R \). It achieves this feat by exchanging what it means to wind with that it means to move.

As the radius of the circle becomes large, \( R \rightarrow \infty \), the winding modes become very heavy with mass \( \sim R/\alpha' \) and are irrelevant for the low-energy dynamics. But the momentum modes become very light, \( M \sim 1/R \), and, in the strict limit form a continuum. From the perspective of the energy spectrum, this continuum of energy states is exactly what we mean by the existence of a non-compact direction in space.
In the other limit, $R \to 0$, the momentum modes become heavy and can be ignored: it takes way too much energy to get anything to move on the $S^1$. In contrast, the winding modes become light and start to form a continuum. The resulting energy spectrum looks as if another dimension of space is opening up!

The equivalence of the string spectrum on circles of radii $R$ and $\alpha'/R$ extends to the full conformal field theory and hence to string interactions. Strings are unable to tell the difference between circles that are very large and circles that are very small. This striking statement has a rubbish name: it is called $T$-duality.

This provides another mechanism in which string theory exhibits a minimum length scale: as you shrink a circle to smaller and smaller sizes, at $R = \sqrt{\alpha'}$, the theory acts as if the circle is growing again, with winding modes playing the role of momentum modes.

**The New Direction in Spacetime**

So how do we describe this strange new spatial direction that opens up as $R \to 0$? Under the exchange (8.7) and (8.8), we see that $p_L$ and $p_R$ transform as

$$p_L \to p_L, \quad p_R \to -p_R$$

Motivated by this, we define a new scalar field,

$$Y^{25} = X^{25}_L(\sigma^+) - X^{25}_R(\sigma^-)$$

It is simple to check that in the CFT for a free, compact scalar field all OPEs of $Y^{25}$ coincide with the OPEs of $X^{25}$. This is sufficient to ensure that all interactions defined in the CFT are the same.

We can write the new spatial direction $Y$ directly in terms of the old field $X$, without first doing the split into left and right-moving pieces. From the definition of $Y$, one can check that $\partial_\tau X = \partial_\nu Y$ and $\partial_\sigma X = \partial_\tau Y$. We can write this in a unified way as

$$\partial_\alpha X = \epsilon_{\alpha\beta} \partial^\beta Y$$

(8.9)

where $\epsilon_{\alpha\beta}$ is the antisymmetric matrix with $\epsilon_{\tau\sigma} = -\epsilon_{\sigma\tau} = +1$. (The minus sign from $\epsilon_{\sigma\tau}$ in the above equation is canceled by another from the Minkowski worldsheet metric when we lower the index on $\partial^\beta$).
The Shift of the Dilaton

The dilaton, or string coupling, also transforms under T-duality. Here we won’t derive this in detail, but just give a plausible explanation for why it’s the case. The main idea is that a scientist living in a stringy world shouldn’t be able to do any experiments that distinguish between a compact circle of radius \( R \) and one of radius \( \alpha'/R \). But the first place you would look is simply the low-energy effective action which, working in Einstein frame, contains terms like

\[
\frac{2\pi R}{2l_s^2 g_s^2} \int d^{25}X \sqrt{-\tilde{G}} e^{R} \mathcal{R} + \ldots
\]

A scientist cannot tell the difference between \( R \) and \( \tilde{R} = \alpha'/R \) only if the value of the dilaton is also ambiguous so that the term in front of the action remains invariant: i.e. \( R/g_s^2 = \tilde{R}/\tilde{g}_s^2 \). This means that, under T-duality, the dilaton must shift so that the coupling constant becomes

\[
g_s \rightarrow \tilde{g}_s = \frac{\sqrt{\alpha'} g_s}{R} \quad (8.10)
\]

8.3.1 A Path Integral Derivation of T-Duality

There’s a simple way to see T-duality of the quantum theory using the path integral. We’ll consider just a single periodic scalar field \( X \equiv X + 2\pi R \) on the worldsheet. It’s useful to change normalization and write \( X = R\varphi \), so that the field \( \varphi \) has periodicity \( 2\pi \). The radius \( R \) of the circle now sits in front of the action,

\[
S[\varphi] = \frac{R^2}{4\pi \alpha'} \int d^2\sigma \partial_\alpha \varphi \partial^\alpha \varphi \quad (8.11)
\]

The Euclidean partition function for this theory is \( Z = \int \mathcal{D}\varphi \, e^{-S[\varphi]} \). We will now play around with this partition function and show that we can rewrite it in terms of new variables that describe the T-dual circle.

The theory (8.11) has a simple shift symmetry \( \varphi \rightarrow \varphi + \lambda \). The first step is to make this symmetry local by introducing a gauge field \( A_\alpha \) on the worldsheet which transforms as \( A_\alpha \rightarrow A_\alpha - \partial_\alpha \lambda \). We then replace the ordinary derivatives with covariant derivatives

\[
\partial_\alpha \varphi \rightarrow \mathcal{D}_\alpha \varphi = \partial_\alpha \varphi + A_\alpha
\]

This changes our theory. However, we can return to the original theory by adding a new field, \( \theta \) which couples as

\[
S[\varphi, \theta, A] = \frac{R^2}{4\pi \alpha'} \int d^2\sigma \, \mathcal{D}_\alpha \varphi \mathcal{D}^\alpha \varphi + \frac{i}{2\pi} \int d^2\sigma \, \theta \epsilon^{\alpha\beta} \partial_\alpha A_\beta \quad (8.12)
\]
The new field $\theta$ acts as a Lagrange multiplier. Integrating out $\theta$ sets $\epsilon^{\alpha\beta}\partial_\alpha A_\beta = 0$. If the worldsheet is topologically $\mathbb{R}^2$, then this condition ensures that $A_\alpha$ is pure gauge which, in turn, means that we can pick a gauge such that $A_\alpha = 0$. The quantum theory described by (8.12) is then equivalent to that given by (8.11).

Of course, if the worldsheet is topologically $\mathbb{R}^2$ then we’re missing the interesting physics associated to strings winding around $\varphi$. On a non-trivial worldsheet, the condition $\epsilon^{\alpha\beta}\partial_\alpha A_\beta = 0$ does not mean that $A_\alpha$ is pure gauge. Instead, the gauge field can have non-trivial holonomy around the cycles of the worldsheet. One can show that these holonomies are gauge trivial if $\theta$ has periodicity $2\pi$. In this case, the partition function defined by (8.12),

$$Z = \frac{1}{\text{Vol}} \int \mathcal{D}\varphi \mathcal{D}\theta \mathcal{D}A \ e^{-S[\varphi,\theta,A]}$$

is equivalent to the partition function constructed from (8.11) for worldsheets of any topology.

At this stage, we make use of a clever and ubiquitous trick: we reverse the order of integration. We start by integrating out $\varphi$ which we can do by simply fixing the gauge symmetry so that $\varphi = 0$. The path integral then becomes

$$Z = \int \mathcal{D}\theta \mathcal{D}A \ \exp\left(-\frac{R^2}{4\pi\alpha'} \int d^2\sigma \ A_\alpha A^\alpha + \frac{i}{2\pi} \int d^2\sigma \ \epsilon^{\alpha\beta} (\partial_\alpha \theta) A_\beta \right)$$

where we have also taken the opportunity to integrate the last term by parts. We can now complete the procedure and integrate out $A_\alpha$. We get

$$Z = \int \mathcal{D}\theta \ \exp\left(-\frac{\tilde{R}^2}{4\pi\alpha'} \int d^2\sigma \ \partial_\alpha \theta \partial^\alpha \theta \right)$$

with $\tilde{R} = \alpha'/R$ the radius of the T-dual circle. In the final integration, we threw away the overall factor in the path integral, which is proportional to $\sqrt{\alpha'/R}$. A more careful treatment shows that this gives rise to the appropriate shift in the dilaton (8.10).

### 8.3.2 T-Duality for Open Strings

What happens to open strings and D-branes under T-duality? Suppose firstly that we compactify a circle in direction $X$ transverse to the brane. This means that $X$ has Dirichlet boundary conditions

$$X = \text{const} \Rightarrow \partial_\tau X^{25} = 0 \quad \text{at } \sigma = 0, \pi$$
But what happens in the T-dual direction \( Y \)? From the definition (8.9) we learn that the new direction has Neumann boundary conditions,

\[
\partial_\sigma Y = 0 \quad \text{at} \; \sigma = 0, \pi
\]

We see that T-duality exchanges Neumann and Dirichlet boundary conditions. If we dualize a circle transverse to a D\( p \)-brane, then it turns into a D\( (p + 1) \)-brane.

The same argument also works in reverse. We can start with a D\( p \)-brane wrapped around the circle direction \( X \), so that the string has Neumann boundary conditions. After T-duality, (8.9) changes these to Dirichlet boundary conditions and the D\( p \)-brane turns into a D\( (p - 1) \)-brane, localized at some point on the circle \( Y \).

In fact, this was how D-branes were originally discovered: by following the fate of open strings under T-duality.

### 8.3.3 T-Duality for Superstrings
To finish, let’s nod one final time towards the superstring. It turns out that the ten-dimensional superstring theories are not invariant under T-duality. Instead, they map into each other. More precisely, Type IIA and IIB transform into each other under T-duality. This means that Type IIA string theory on a circle of radius \( R \) is equivalent to Type IIB string theory on a circle of radius \( \alpha'/R \). This dovetails with the transformation of D-branes, since type IIA has D\( p \)-branes with \( p \) even, while IIB has \( p \) odd. Similarly, the two heterotic strings transform into each other under T-duality.

### 8.3.4 Mirror Symmetry
The essence of T-duality is that strings get confused. Their extended nature means that they’re unable to tell the difference between big circles and small circles. We can ask whether this confusion extends to more complicated manifolds. The answer is yes. The fact that strings can see different manifolds as the same is known as mirror symmetry.

Mirror symmetry is cleanest to state in the context of the Type II superstring, although similar behaviour also holds for the heterotic strings. The simplest example is when the worldsheet of the string is governed by a superconformal non-linear sigma-model with target space given by some Calabi-Yau manifold \( X \). The claim of mirror symmetry is that this CFT is identical to the CFT describing the string moving on a different Calabi-Yau manifold \( Y \). The topology of \( X \) and \( Y \) is not the same. Their Hodge diamonds are the mirror of each other; hence the name. The subject of mirror symmetry is an active area of research in geometry and provides a good example of the impact of string theory on mathematics.
8.4 Epilogue

We are now at the end of this introductory course on string theory. We began by trying to make sense of the quantum theory of a relativistic string moving in flat space. It is, admittedly, an odd place to start. But from then on we had no choices to make. The relativistic string leads us ineluctably to conformal field theory, to higher dimensions of spacetime, to Einstein’s theory of gravity at low-energies, to good UV behaviour at high-energies and to Yang-Mills theories living on branes. There are few stories in theoretical physics where such meagre input gives rise to such a rich structure.

This journey continues. There is one further ingredient that it is necessary to add: supersymmetry. Even this is in some sense not a choice, but is necessary to remove the troublesome tachyon that plagued these lectures. From there we may again blindly follow where the string leads, through anomalies (and the lack thereof) in ten dimensions, to dualities and M-theory in eleven dimensions, to mirror symmetry and moduli stabilization and black hole entropy counting and holography and the miraculous AdS/CFT correspondence.

However, the journey is far from complete. There is much about string theory that remains to be understood. This is true both of the mathematical structure of the theory and of its relationship to the world that we observe. The problems that we alluded to in Section 6.4.5 are real. Non-perturbative completions of string theory are only known in spacetimes which are asymptotically anti-de Sitter, but cosmological observations suggest that our home is not among these. In attempts to make contact with the standard models of particle physics and cosmology, we typically return to the old idea of Kaluza-Klein compactifications. Is this the right approach? Or are we missing some important and subtle conceptual ingredient? Or is the existence of this remarkable mathematical structure called string theory merely a red-herring that has nothing to do with the real world?

In the years immediately after its birth, no one knew that string theory was a theory of strings. It seems very possible that we’re currently in a similar situation. When the theory is better understood, it may have little to do with strings. We are certainly still some way from answering the simple question: what is string theory really?