# Analytic Aspects of Periodic Instantons 

by

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#### Abstract

The main result is a computation of the Nahm transform of a $\operatorname{SU}(2)$-instanton over $\mathbb{R} \times T^{3}$, called spatially-periodic instanton. It is a singular monopole over $T^{3}$, a solution to the Bogomolny equation, whose rank is computed and behavior at the singular points is understood under certain conditions. A full description of the Riemannian ADHMN construction of instantons on $\mathbb{R}^{4}$ is given, preceding a description of the heuristic behind the theory of instantons on quotients of $\mathbb{R}^{4}$. The Fredholm theory of twisted Dirac operators on cylindrical manifolds is derived, the spectra of spin Dirac operators on spheres and on product manifolds are computed. A brief discussion on the decay of spatially-periodic and doubly-periodic instantons is included.


Thesis Supervisor: Tomasz S. Mrowka
Title: Professor of Mathematics

## Acknowledgments

It all started on October 3, 1994, when I decided I would do math in life. But I already told that bit of the story in my masters thesis.

The readership of a thesis is usually quite small. This is in part because thesis are usually not well distributed. Soon indeed, this thesis will fall in the "microfilm zone," to be lost forever except to those of unfaltering determination. Thank you then, dear reader, for being here with me, and please accompany me in this section on a celebration of the many thanks I have to give.

Before continuing any further, it is imperative that I stop to warmly acknowledge my advisor and mentor, Tomasz Mrowka. By standing on the shoulders of a giant, it is much easier to see ahead. From him I learned many valuable lessons that I hope to retain and pass on to the next generation. One of these I included at the start of Chapter 8 His guidance - punctuated by his usual "So?" asked in the middle of a hallway - made this project a success. His patience, wisdom, and vision make him a very very good mentor to those that are willing to work hard without being told all the time what to do. The more I interact with Tom, the more I respect him, the more I understand how fortunate I am to be working with him.

The story of how he became my advisor is perhaps symptomatic of how easy the relationship can be with Tom. I remember that day very well: I was exiting the mail room as he walked by. We talked for a brief moment about an idea I had of working with one of his collaborator at Harvard. Tom told me that this collaborator had just had a baby and wasn't looking for more students at the moment, although I could still probably win him over should I push a little. And then he said, "Do you want a problem?" And as we walked to the subway, he gave me a problem. While I deviated a bit from this original problem, it is still on the back burner of my mind and permeates the research presented here.

Only a few months later I finally decided to choose Tom for advisor. We bumped into each other in front of the applied math common room, and I told him something along the lines of "I think it is about time for me to pick an advisor. I was thinking of you. What do I have to do?" He just replied, "Tell Linda." That was it! So easy!

I know that Tom's influence in my life doesn't stop with the completion of this thesis, it will persist, just as the influence of my earlier advisors, official or not, persisted. My style is greatly influenced by Tom and the four men I am about to describe.

Fernand Beaudet was my unofficial undergraduate advisor. He his the best ambassador for mathematics I have ever seen. He is the cause of my being a mathematician, and guided me through the early years. Back then, I never made a decision without consulting him.

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Pierre Bouchard, my masters thesis advisor, was always a great believer in me. His love of mathematics is vibrant, magnetic, and his generosity is exemplary.

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I wish to thank these four original mentors profusely. Thank you, thank you! I could feel you all with me all these years.

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My mathematical parents being acknowledged, I take this opportunity to recognize my real family, mom, dad, and my sisters. Although I have been away for a very long time, and only had a few chances to interact with my family, their unconditional love, their belief in who I am and who I can become, their trust in my ability to navigate through the hardship and to make good decisions most of the time mean a lot to me. I want to thank them for all they are.

Special thanks go to the special girls that shared my life, Anne-Marie when I came here and Caroline this past year. You've sure been a cause of a lot of non-productive time, as far as math goes, and thank you for that. Thank you for your great smiles and for believing in me. You've helped me become a better person.

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The mathematics department at MIT is a fantastic place, very effervescent. With about 120 students and 40 or more professors and postdocs, it is dynamic and I cherish the time I spent here. Major thanks go to the staff for running this department so smoothly.

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I arrived in office 2-251 at a time when it was a party zone and only the imminent deadlines of homework were keeping me working. Left to myself once the classes were over, I had to learn to work even with a far away deadline. I had to grow. My friend Leo Jesudian told me once to raise the bar in all I do. Thanks go to him and his wife Tiffany, as well as to my great friend Maneesh Bhatnagar for teaching me focus, discipline, goal setting and work ethic. When I arrived at MIT, I was not a person who could do research. But I was able to change and become one.

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The completion of a PhD is a very slow and gruelling process. It is very natural to feel dejected, frustrated, inadequate, braindead, desperate, miserable, close to tears. This is also a time when people's brain works overtime about all the different things they could be doing instead. Its the brain's attempt to cope by escapism. Thank you, thank you, thank you so much to Martin Pinsonnault and Suparna Sanyal for helping me stay sane.

While I never met Marcos Jardim, his papers had a great impact on my vision of this project and of further projects I want to tackle. I want to thank him and all the others whose papers made me explore over the last few years a gorgeous piece of mathematical landscape.

To all of those who paved the way: thank you! Perched at the top of the wall of knowledge, sitting on the shoulders of so many giants, you lead us in the quest for truth and beauty. Never stop pushing back the frontiers of ignorance. You are doing a great work, don't you dare come down.

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## Introduction

This introduction is composed of three parts. There is first a description of the main result of this thesis, then a historical account of the ideas leading to this thesis, and finally a road map composed of a brief description of each of the chapters in this thesis,
A Yang-Mills instanton on a Riemannian four-manifold is a vector bundle $E$ along with the gauge equivalence class of a connection $A$ whose curvature $F_{A}$ is anti-self dual and of finite $L^{2}$ norm.
The Nahm transform of an instanton $(E, A)$ on $\mathbb{R} \times T^{3}$ consists of a bundle $V$ over an open subset of $T^{3}$, a connection $B$ on $V$, and an element $\Phi$ of $\operatorname{End} V$. These objects are constructed as follows. Each point $z \in T^{3}$ correspond to a flat line bundle $L_{z}$ over $\mathbb{R} \times T^{3}$, and we consider the twisted spin Dirac operator

$$
\mathscr{P}_{A_{z}}^{*}: \Gamma\left(\mathbb{R} \times T^{3}, S^{-} \otimes E \otimes L_{z}\right) \rightarrow \Gamma\left(\mathbb{R} \times T^{3}, S^{+} \otimes E \otimes L_{z}\right)
$$

The bundle $V$ is defined by the equation

$$
V_{z}=L^{2} \cap \operatorname{ker}\left(\mathscr{D}_{A_{z}}^{*}\right)
$$

Let $t$ be the $\mathbb{R}$-coordinate in $\mathbb{R} \times T^{3}$, and $m_{t}$ denote multiplication by $t$. Let $P$ denote the $L^{2}$ projection on $V$, and $d^{z}$ the trivial connection for the trivial bundle with infinite dimensional fiber $L^{2}\left(\mathbb{R} \times T^{3}, S^{-} \otimes E\right)$. Then the connection $B$ and the Higgs field $\Phi$ are defined by the equations

$$
\begin{aligned}
B & =P d^{z} \\
\Phi & =-2 \pi i P m_{t}
\end{aligned}
$$

The main result of the present thesis is the following theorem.

Theorem 8.0-1 in the text on page 84. Outside of a set $W$ consisting of at most four points, the family of vector spaces $V$ described above defines a vector bundle of rank

$$
\frac{1}{8 \pi^{2}} \int_{\mathbb{R} \times T^{3}}\left|F_{A}\right|^{2}
$$

and the couple $(B, \Phi)$ satisfies the Bogomolny equation

$$
\nabla_{B} \Phi=* F_{B}
$$

For $w \in W$ and $z$ close enough to $w$, unless we are in the Scenario 2 of page 91 there are maps
$\Phi^{\perp}$ and $\Phi^{\lrcorner}$such that

$$
\Phi=\frac{-i}{2|z-w|} \Phi^{\perp}+\Phi^{\lrcorner}
$$

and $\Phi^{\perp}$ is the $L^{2}$-orthogonal projection on the orthogonal complement of a naturally defined subbundle $V_{\lrcorner}$of $V$.

We know heuristically, as shown in Chapter 2 that $(B, \Phi)$ must satisfies the Bogomolny equation, up to a limiting term coming out of some integration by parts on $\mathbb{R} \times T^{3}$. As it is shown in Chapter 7 harmonic spinors are exponentially decaying outside of $W$, and it then must be that the limiting term just mentioned is 0 .

The rank of $V$ is not really a surprise and follows for some relative index theorem. It is a sharp contrast with the $S^{1} \times \mathbb{R}^{3}$ case where the computation, and the formula itself, is slightly more involved; see [NS00].
The last part of the theorem follows from a careful analysis of some geometric splitting of $V$ coming from considering the kernel of the Dirac operator in some weighted $L^{2}$-space

$$
L_{\delta}^{2}:=e^{\delta t} L^{2}
$$

and variants. Taking $\epsilon>0$ small enough, we define the various spaces

$$
\begin{array}{rlrl}
V_{z} & :=e^{-\epsilon|t|} L^{2} \cap \operatorname{ker}\left(\mathscr{D}_{A_{z}}^{*}\right), & V_{z}:=e^{\epsilon|t|} L^{2} \cap \operatorname{ker}\left(\mathscr{D}_{A_{z}}^{*}\right), \\
K_{z}:=e^{\epsilon|t|} L^{2} \cap \operatorname{ker}\left(\mathscr{D}_{A_{z}}\right), & \mathcal{H}_{z}:=e^{\epsilon|t|} L^{2} \cap \operatorname{ker}\left(\mathscr{D}_{A_{z}}^{*} \mathscr{D}_{A_{z}}\right) .
\end{array}
$$

Then obviously $V \subset V \subset{ }^{\top} V$. But also, as shown in Section 8.2

$$
{ }^{\ulcorner } V_{z}=V_{z} \oplus \mathscr{D}_{A_{z}} \mathcal{H}_{z}
$$

## A progression of ideas

A concrete understanding of instantons played an important role in particle physics since their discovery in mid-' 70 by Belavin et al [BPST75]. More importantly for us, it played an important role in four-dimensional topology and geometry. For example, Donaldson has shown in [Don83] how to extract information about the intersection form on a given manifold from its moduli space of instantons; see [FU84] and [DK90] for more details.

Finding a complete description of all instantons on a given space is not an easy task and we have a description for a limited number of spaces. In particular, we do not completely understand the moduli spaces for quotients of $\mathbb{R}^{4}$ by lattices. In that picture, a non-linear analog of the Fourier transform, the "Nahm transform," appears.
This present thesis takes place in the quest for a unified understanding of moduli spaces of instantons on $\mathbb{R}^{4}$ invariant under the action of a group of translations via the Nahm transform heuristic.
The problem of describing all instantons on $\mathbb{R}^{4}$ was addressed by Atiyah, Drinfeld, Hitchin and Manin in 1978 in ADHM78. Their description became known as the ADHM construction. Using twistor methods, they were able to equate the moduli space of instanton on $\mathbb{R}^{4}$ to a finite dimensional space of algebraic data, called the "ADHM data." Still using twistor methods, and using the relationship between monopoles on $\mathbb{R}^{3}$ (solutions to Bogomolny equation) and time-invariant instantons on $\mathbb{R}^{4}$, Hitchin [Hit82] proved in 1982 that every monopole can be constructed from some algebraic geometry data, the "spectral curve."

Nahm in 1981-1982 proposed a simplification which he thought would be better understood by physicists. As it turned out, his idea was very fruitful. The main idea is to construct the ADHM data by considering the kernel of the Dirac operator coupled to the instanton connection. By twisting the connection by a flat connection parameterized by $t$, Nahm also explained how monopoles can arise from solutions to a set of differential equations on $\mathbb{R}$, which we now call the "Nahm equations."
These ideas were rapidly exploited by Corrigan and Goddard in [CG84] who formalized the $\mathbb{R}^{4}$ story, a complete proof of which with some algebraico-geometric flavor can be found in [DK90, Chap. 3], and by Hitchin in Hit83] who completed the $\mathrm{SU}(2)$-monopole story.
Around 1988, Braam noticed that Nahm's considerations can be used for instantons on flat tori. Exploiting Braam's observation, Schenk and Braam-van Baal in [Sch88] and [BvB89] proved independently a bijective correspondence between the moduli spaces of instantons over a flat torus and over its dual torus.
While the proofs of Corrigan-Goddard-Nahm and Schenk-Braam-van Baal are quite direct, it is not the case with Hitchin's construction, which sits in a triangle of equivalences:


In 1989, Hurtubise and Murray completed the monopole story for all classical groups, using in [HM89] a triangle of ideas similar to Hitchin's:


Note in both cases that not all arrows go both ways. While the "spectral curves" are interesting objects to study in themselves, it would be desirable to pass directly from monopoles to Nahm data, as we do for $\mathbb{R}^{4}$ and $T^{4}$. For $\mathrm{SU}(2)$-monopoles, this direct proof was accomplished in 1993 by Nakajima in [Nak93].
All those various correspondences fit in a more general framework. The Nahm transform takes an instanton over $\mathbb{R}^{4}$, invariant under the action of some group of translations $\Lambda$, and creates some Nahm data over $\mathbb{R}^{4 *}$, invariant under the action of

$$
\Lambda^{*}:=\left\{t \in \mathbb{R}^{4 *} \mid t(\Lambda) \subset \mathbb{Z}\right\}
$$

or equivalently, over $\mathbb{R}^{4 *} / \Lambda^{*}$.
More precisely, for each instanton $A$ on a bundle $E$ over $\mathbb{R}^{4} / \Lambda$, the Nahm transform creates a bundle $\hat{E}$ over $\mathbb{R}^{4 *} / \Lambda^{*}$ less a few points and a connection $\hat{A}$. The self-dual part of the curvature $F_{\hat{A}}$ encodes the behavior of solutions to the Dirac equation in the non-compact directions. The bundle $\hat{E}$ is assembled from kernels of twisted Dirac operators for perturbations of $A$ varying continuously over $\mathbb{R}^{4 *} / \Lambda^{*}$, less those points where the associated Dirac operator is not Fredholm.

For an expanded version of the Nahm transform, as well as for examples of non-flat Nahm transforms and a survey of the literature, read the survey paper [Jar].
This idea has been exploited quite successfully by Marcos Jardim in his doctoral thesis [Jar99] and a series of papers [Jar01 Jar02a, Jar02b. Some analytical details concerning asymptotics were tackled by Jardim and Biquard in [BJ01]. This work relates doubly-periodic instantons, or instantons on $T^{2} \times \mathbb{R}^{2}$, with singular Higgs pairs on $T^{2}$. It is worth noting that Jardim's construction does use Hitchin's approach and goes through the spectral curves realm.
And so do Cherkis and Kapustin in [CK98, CK99, CK01] where they relate monopoles on $\mathbb{R}^{2} \times S^{1}$ to solutions of Hitchin's equations on $S^{1} \times \mathbb{R}$ using the Nahm transform and Hitchin's approach.
While Nye's doctoral thesis's work [Nye01] on the Nahm story for calorons, which are instantons on $S^{1} \times \mathbb{R}^{3}$, does not directly use spectral curves, it relies on the construction of the Nahm data for monopoles of [HM89] which does use them. Nye's work, and the companion paper [NS00] with Singer, cover a lot of ground but bits and pieces are missing. As mentioned by Nye in his thesis, a direct proof of the $S U(n)$-monopole story through a careful analysis of the Dirac operator similar to Nakajima's proof for the $S U(2)$ case would help cover even more ground.
Of the four-dimensional quotients of $\mathbb{R}^{4}$, there remains only $\mathbb{R} \times T^{3}$. At this point in time, very little is known about instantons on $\mathbb{R} \times T^{3}$ : some comments about the Nahm transform heuristic, and numerical approximations are found in [vB96]. This current thesis is a step forward.

## Road Map

The heart of this present thesis is composed of Chapter 2 where the heuristic guiding our steps is presented, and Chapter $\boxed{\square}$ where the main result is described and proved. The experienced reader might want to pick and choose what he wants to read from the other chapters in order to get to the main result. To facilitate this approach, we now rapidly explore the whole thesis.
In Chapter $\square$ we explore the ADHM construction of instantons on $\mathbb{R}^{4}$, incorporating the idea of Nahm and using only Riemannian constructions and avoiding at all cost any use of the complex structure of $\mathbb{R}^{4}$. Acknowledging those facts, this chapter is called "The Riemannian ADHMN construction."
In Chapter 2] we explore in more details the Nahm Transform heuristic which guides the research in this field of study. The curvature computation presented in that chapter is the key ingredient in understanding why the pair $(B, \Phi)$ satisfies the Bogomolny equation on almost all of $T^{3}$.
In Chapter 3, we study the Dirac spectrum of product manifolds. Of particular interest is the Dirac Spectrum Formula given in Theorem 3.2-1, see page 41 This formula constructs the spectrum $\Sigma_{M \times N}$ of the Dirac operator on a spinor bundle of the product manifold $M^{m} \times N^{n}$ in terms of the spectra $\Sigma_{M}$ on $M^{m}$ and $\Sigma_{N}$ on $N$. More precisely, we get

$$
\Sigma_{M \times N}= \begin{cases} \pm\left|\Sigma_{M} \times \Sigma_{N}\right|, & \text { if } m \text { and } n \text { are odd; } \\ \pm\left|\Sigma_{M}^{>0} \times \Sigma_{N}\right| \cup\left(\Sigma_{N}\right)^{\# k_{M}^{+} \cup\left(-\Sigma_{N}\right)^{\# k_{M}^{-}},} & \text {if } m \text { is even. }\end{cases}
$$

This formula might not be present in the literature. As a corollary, we derived in Theorem 3.4-1 a formula for the spectrum of the Dirac operator on the spinor bundle of $T^{3}$ twisted by a flat line bundle.
In Chapter[4 we derive formulas for eigenvalues and multiplicities of the Dirac operator on spheres. Section 4.1 computes the spectrum for $S^{3}$. The proof presented here is quite similar to a proof of Hitchin of which the author was not aware at the time of the writing. Knowledge of this spectrum
is necessary to understand the asymptotic behavior of harmonic spinors on $\mathbb{R}^{4}$ proved in Chapter $\square$ and used in Chapter 1 Section 4.2 presents a construction of Trautman for the eigenvalues on all spheres and confirms to some extend the results of the other section.

In Chapter [5] we take note of certain results concerning the asymptotic decay of instantons on $\mathbb{R} \times T^{3}$. The proof exists elsewhere in the literature and is not included here. Should one be able to adapt the center manifold proof for instantons on cylindrical manifolds presented in [MMR94] to warped cylinders, one could use Theorem 5.2-2] on the decay of instantons on $T^{2} \times \mathbb{R}^{2}$ living in the gauge group translates of the zero Fourier mode to prove a conjecture of Jardim on finite energy and quadratic decay.
In Chapter we define weighted Sobolev spaces and study conditions on the weights for a Dirac operator twisted by an instanton to be Fredholm. An analysis of the time-independent case provides a formula for the difference of the indices for different weights. A short story of the concepts of weighted Sobolev spaces is presented to get the chapter off the ground.
In Chapter 7 we derive knowledge of the asymptotic behavior of harmonic spinor. To achieve that goal, the Fredholm theory of Chapter 6is extended to weighted Sobolev spaces on half-cylinders. Once this task is accomplished, a diagram chase gives the desired result. This chapter closes with an analysis of the asymptotic behavior on $\mathbb{R}^{4}$. The knowledge of this behavior is necessary for part of the algebraic data in the ADHMN construction of Chapter
In Chapter 8 we describe the Nahm transform of spatially periodic instantons. It is a singular monopole on $T^{3}$. The excision proof of Chapter 6 allows for a computation of the $L^{2}$-index of the Dirac operator, which is presented in Section 8.1. A geometric splitting of the bundle $V$ given in Section 8.2 allows for an understanding of the behavior of the Higgs field at the singular points, which is given in Section 8.3 A derivation of a precise formula for the Green's operator on $S^{ \pm} \otimes L_{z}$ presented in Section 8.4 constitutes some preliminary work on the behavior of the connection $B$ at the singular points.
Four appendices complete this thesis. In Appendix $\mathbb{A}$ we derive the various dimensional reductions of the anti-self-dual equation. In Appendix B we study an excision principle for the index of Fredholm operators. In Appendix C] we state and prove an algebraic lemma useful for simplifying the exposition in Chapter 8 In Appendix $\mathbb{D}$ we study how the Dirac operator changes under a conformal change of the metric. In Appendix E we visit the treatment of Bartnik of weighted Sobolev spaces on $\mathbb{R}^{n}$ and Fredholm properties for operators asymptotic to the ordinary Laplacian, merely cleaning up a part of his paper [Bar86] by adding proofs where needed. The results presented in this appendix are used in Chapter $\square$ and parallel to a certain extend our treatment of Dirac operator on cylindrical manifolds of Chapters 6 and 7

## Chapter 1

## The Riemannian ADHMN construction

On a four dimensional riemannian manifold $X$, a $G$-instanton is a $G$-bundle $E$ equipped with the gauge equivalence class of a connection $A$ which is such that its curvature $F_{A}$ is anti-self-dual (written ASD for short)

$$
* F_{A}=-F_{A}
$$

and has finite energy

$$
\left\|F_{A}\right\|_{L^{2}}<\infty .
$$

In the case where $X$ is compact, we can associate to the $\mathrm{SU}(n)$-instanton $(E, A)$ its instanton number $c_{2}(E)$. In fact, the equalities

$$
\begin{aligned}
c_{2}(E) & =\frac{1}{8 \pi^{2}} \int \operatorname{Tr}\left(F_{A}\right)^{2} \\
& =\frac{1}{8 \pi^{2}} \int-\left|F_{A}^{+}\right|^{2}+\left|F_{A}^{-}\right|^{2} d \mu
\end{aligned}
$$

and

$$
\left\|F_{A}\right\|_{L^{2}}^{2}=\int\left|F_{A}^{+}\right|^{2}+\left|F_{A}^{-}\right|^{2} d \mu
$$

indicate that

$$
\left\|F_{A}\right\|_{L^{2}}^{2}=8 \pi^{2} c_{2}(E) \text { if and only if } A \text { is ASD. }
$$

Hence not every bundle admit a ASD connection: an obstruction to the existence of a ASD connection on $E$ is $c_{2}(E) \geq 0$.
In this chapter, we explore the ADHM construction of instantons on $\mathbb{R}^{4}$ from a strictly riemannian viewpoint. Most treatments found in the literature exploit the holomorphic possibilities stemming from the ASD condition. Nahm's [Nah84] and Corrigan-Goddard's [CG84] papers are unlike those, but provide more of a backbone than a complete construction.

### 1.1 The setting

In this chapter, we consider only instantons for the group $\operatorname{SU}(n)$ on the space $\mathbb{R}^{4}$.
Let $S=S^{+} \oplus S^{-}$be the spinor bundle of $\mathbb{R}^{4}$. Recall that $S^{+}$and $S^{-}$are trivial bundles with quarternionic fiber $\mathbb{H}$. Let's denote the Clifford multiplication by $\rho$.
Let $E$ be a complex vector bundle with structure group $\mathrm{SU}(n)$. Let $E$ be equipped with a connection
$A$. We denote by $D_{A}$ the Dirac operator $\Gamma\left(S^{+} \otimes E\right) \rightarrow \Gamma\left(S^{-} \otimes E\right)$ and $D_{A}^{*}$ its adjoint. The Laplacian $\nabla_{A}^{*} \nabla_{A}$ we denote $\Delta_{A}$. Thus

$$
\Delta_{A} f=-\sum_{i=1}^{4}\left(\partial_{i}\right)^{2}-2 \sum_{i=1}^{4} A_{i} \partial_{i} f-\sum_{i=1}^{4}\left(\left(\partial_{i} A_{i}\right)+A_{i}^{2}\right) f .
$$

The main object we are studying are instantons on $\mathbb{R}^{4}$. An instanton connection is
a $\mathrm{SU}(n)$ bundle $E$, and a connection $A$ on $E$ such that $F_{A}^{+}=0($ ASD condition), and

$$
\left\|F_{A}\right\|_{L^{2}}<\infty .
$$

An instanton is the gauge equivalence class of an instanton connection. It must be in fact that $k:=\left\|F_{A}\right\|_{L^{2}} / 8 \pi^{2}$ is an integer that we call the charge. Let $\mathcal{M}_{k, n}^{\text {ASD }}$ denote the moduli space of instantons of charge $k$ and rank $n$.
Equally important are the ADHM data. They are
a hermitian vector space $V$ of rank $k$,
a hermitian vector space $W$ of rank $n$,
a 1-form $a$ with values in hermitian endomorphisms of $V$, and
a map $\Psi: V \rightarrow S^{+} \otimes W$.
In the more general framework of Chapter 2, ADHM data are called Nahm data.
There is a natural notion of isomorphism of ADHM data. Of course, any two hermitian vector spaces of same rank are isomorphic, so $a$ can be thought as a 1 -form with values in hermitian $k \times k$ matrices, and $\Psi$ as a $2 n \times k$ matrix. The ADHM data $(V, W, a, \Psi)$ and $\left(V^{\prime}, W^{\prime}, a^{\prime}, \Psi^{\prime}\right)$ are to be considered equivalent if there exist $u \in \mathrm{SU}(n)$ and $v \in \mathrm{U}(k)$ for which

$$
\begin{equation*}
u a^{\prime} u^{-1}=a, \text { and }\left(1 \otimes u^{-1}\right) \Psi^{\prime} v=\Psi . \tag{1.1}
\end{equation*}
$$

The aim of the ADHM construction is to place in correspondence the space of instantons and the space of equivalence classes of ADHM data satisfying Conditions (1.3) and (1.4) described below. We identify $S^{+}$to its dual using a complex skewform $\omega$ on $S^{+}$:

$$
\begin{aligned}
S^{+} & \rightarrow\left(S^{+}\right)^{*} \\
s & \mapsto \omega(\cdot, s)
\end{aligned}
$$

Hence we can associate to the map $\Psi: V \rightarrow S^{+} \otimes W$ the map

$$
\Phi=(\omega \otimes 1) \circ(1 \otimes \Psi): S^{+} \otimes V \rightarrow W
$$

We use $a$ and $\Phi$ to define the map

$$
\begin{gather*}
Q_{x}: S^{+} \otimes V \rightarrow S^{-} \otimes V \oplus W \\
Q_{x}=\left[\begin{array}{c}
\sum_{i=1}^{4} \rho\left(\partial_{i}\right) \otimes\left(a_{i}+x_{i}\right) \\
\Phi
\end{array}\right] . \tag{1.2}
\end{gather*}
$$

The conditions referred to above are the

$$
\begin{equation*}
-\rho([a, a])+2 \Phi^{*} \Phi=1 \otimes \Psi^{*} \Psi \quad \text { (ADHM equation) } \tag{1.3}
\end{equation*}
$$

$Q_{x}$ is everywhere injective. (non-degeneracy condition)
Let $\mathcal{M}_{n, k}^{\mathrm{ADHM}}$ denote the space of ADHM data satisfying the ADHM and non-degeneracy conditions, modulo the equivalence relation of Equation (1.1).
The goal of this chapter is to prove the following theorem.

Theorem 1.1-1 (ADHM construction). The map

$$
\mathfrak{N}: \mathcal{M}_{k, n}^{\mathrm{ASD}} \rightarrow \mathcal{M}_{n, k}^{\mathrm{ADHM}}
$$

constructed in Section 1.2 and the map

$$
\mathfrak{F}: \mathcal{M}_{n, k}^{\mathrm{ADHM}} \rightarrow \mathcal{M}_{k, n}^{\mathrm{ASD}},
$$

constructed in Section 1.3 are inverses of each other.

### 1.2 From instanton to ADHM data

We build up the ADHM data bit by bit.
The Weitzenbock formula

$$
D_{A}^{*} D_{A}=\Delta_{A}+\rho\left(F_{A}^{+}\right)+\frac{1}{4} \text { scalar curvature }
$$

tells us that the ASD (anti-self-dual) condition for the connection $A$ is equivalent to the condition that

$$
D_{A}^{*} D_{A} \text { commutes with quaternion multiplication. }
$$

It also tells us that for an $A$ connection whose curvature is ASD, $\operatorname{ker}\left(D_{A}\right) \cap L^{2}=\{0\}$. Indeed, because of the Weitzenbock formula, when $D_{A} \phi=0$ it must be that $\phi$ is parallel. But to be $L^{2}$ on $\mathbb{R}^{4}$, a parallel section must then be 0 .
Set

$$
V_{E}:=L^{2} \cap \operatorname{ker}\left(D_{A}^{*}\right),
$$

and

$$
W_{E}:=\text { bounded harmonic sections of } E \text {. }
$$

Elements of $W_{E}$ are in natural bijection with sections parallel at infinity. Set the scalar product on $W_{E}$ to be

$$
\begin{equation*}
\left\langle w_{1}, w_{2}\right\rangle=4 \pi^{2}\left(w_{1}^{\infty}, w_{2}^{\infty}\right) . \tag{1.5}
\end{equation*}
$$

Let $\left\{\psi_{1}, \ldots, \psi_{k}\right\}$ be a $L^{2}$-orthonormal basis of $V_{E}$. We use the $L^{2}$ scalar product, which we also denote $\langle$,$\rangle , and define the projection \Pi$ by the formula

$$
\Pi:=\sum_{j=1}^{k}\left\langle\psi_{j}, \cdot\right\rangle \psi_{j} .
$$

Let $m_{\mu}$ denote multiplication by $x_{\mu}$. Consider the linear map

$$
\begin{align*}
a_{\mu}: V_{E} & \rightarrow V_{E} \\
\psi & \mapsto-\Pi m_{\mu} \psi . \tag{1.6}
\end{align*}
$$

The endomorphism $a_{\mu}$ has matrix

$$
\left[-\left\langle\psi_{i}, x_{\mu} \psi_{j}\right\rangle_{L^{2}}\right]_{1 \leq i, j \leq k}
$$

This matrix is clearly hermitian, that is $a_{\mu}^{*}=a_{\mu}$.
The $L^{2}$ condition imposes a particular asymptotic behavior to elements of $V_{E}$. We study in Chapter 7 how harmonic spinors decay on cylindrical manifolds. Since $\mathbb{R}^{4} \backslash\{0\}=\mathbb{R} \times S^{3}$ conformally, we reprove in Section 7.3 the classical result that any element of $V_{E}$ has an asymptotic expansion of the type

$$
\begin{equation*}
|x|^{-4} \rho(x) \hat{\phi}+O\left(|x|^{-4}\right) \tag{1.7}
\end{equation*}
$$

for a parallel section $\hat{\phi}$ of $S^{+} \otimes E$.
We define the map

$$
\begin{gathered}
\Psi: V_{E} \rightarrow S^{+} \otimes W_{E} \\
\phi \mapsto \hat{\phi} / 2 .
\end{gathered}
$$

We package the obtained ADHM data as

$$
\begin{equation*}
\mathfrak{N}(E, A)=\left(V_{E}, W_{E}, a, \Psi\right) \tag{1.8}
\end{equation*}
$$

The rest of this section is devoted to the analysis justifying the given description of $\mathfrak{N}(E, A)$ and preparing the way for the proof that $\mathfrak{N}(E, A)$ satisfies the conditions (1.3) and (1.4).

## The Green's function $G$ and Projections

As multiplication by $x_{i}$ could potentially kick an element of $L^{2}$ out of it, we have to prove that on $\operatorname{ker}\left(\mathscr{P}_{A}^{*}\right)$ it doesn't. To do so, we observe that the $L^{2}$-condition on $V_{E}$ is actually too weak. This observation is best described in the realm of weighted Sobolev spaces see Appendix Eor [Bar86] for conventions and results.
Again using the conformal identification of $\mathbb{R}^{4} \backslash\{0\}$ with $\mathbb{R} \times S^{3}$, and using the fact that the interval ( $-3 / 2,3 / 2$ ) contains no eigenvalue of the Dirac operator on $S^{3}$ (see Section 4.1), we can use the technology of Chapter 6 or of Appendix Eto prove that

$$
V_{E}=W_{\delta}^{1,2} \cap \operatorname{ker}\left(D_{A}^{*}\right)
$$

for $\delta \in(-3,0)$. In that range, the kernel is constant.
For $-2<\delta<0$ and $1<p<\infty$, the operator

$$
\Delta_{A}: W_{\delta}^{k+2, p}(S \otimes E) \rightarrow W_{\delta-2}^{k, p}(S \otimes E)
$$

is invertible; see for example [KN90, lemma 5.1, p. 279]. Let $G_{A}$ denote its inverse, the so called Green's operator. Observe that as $\Delta_{A}$ is defined independently of $\delta$, so is $G_{A}$ for $\delta<0$.

Set

$$
P_{A}:=I d-D_{A} G_{A} D_{A}^{*} .
$$

In a finite dimensional setting, it is obvious that $\Pi=P_{A}$. The next lemma tells us for which weighted Sobolev spaces these projections are indeed the same.

Lemma 1.2-1. When $\delta \in(-3,-1)$, the projection $P_{A}: W_{\delta}^{1,2} \rightarrow V_{E}$ is a well-defined continuous map. When $\delta<-1$, the projection $\Pi$ : $L_{\delta}^{2} \rightarrow V_{E}$ is a well-defined continuous map. On the spaces $W_{\delta}^{1,2}$ for $\delta \in(-3,-1)$, we have $\Pi=P_{A}$.

Proof: All the maps in the sequence

$$
W_{\delta}^{1,2} \xrightarrow{D_{A}^{*}} L_{\delta-1}^{2} \xrightarrow{G_{A}} W_{\delta+1}^{2,2} \xrightarrow{D_{A}} W_{\delta}^{1,2}
$$

are continuous when $\delta+1 \in[-2,0]$, thus when $\delta \in[-3,-1]$. Since $D_{A}^{*} P_{A}=0$ in the interior of that range, $P_{A}$ maps into $V_{E}$ for $\delta \in(-3,-1)$.
We have

$$
\left|\left\langle\psi_{j}, \phi\right\rangle\right| \leq\left\|\psi_{j}\right\|_{2, \delta_{1}}\|\phi\|_{2, \delta}
$$

with $\delta+\delta_{1}=-4$. Since $\psi_{j} \in V_{E} \subset W_{\delta_{1}}^{1,2}$ for $\delta_{1}>-3$, we have that the scalar product is finite when $\phi \in L_{\delta}^{2}$ for $\delta<-1$. We can clearly see that $\Pi$ is continuous and maps into $V_{E}$ in that range. Now, suppose $\phi \in V_{E}^{\perp} \cap W_{\delta}^{1,2}$. Then

$$
\left\langle P_{A} \phi, \psi_{j}\right\rangle=\left\langle\phi, \psi_{j}\right\rangle-\left\langle D_{A} G_{A} D_{A}^{*} \phi, \psi_{j}\right\rangle .
$$

The first term of the right hand side is clearly 0 since $\psi_{j} \in V_{E}$. For $\phi \in W_{\delta}^{1,2}$ with $\delta \in(-3,-1)$, we have the equality

$$
\left\langle D_{A} G_{A} D_{A}^{*} \phi, \psi_{j}\right\rangle=\left\langle G_{A} D_{A}^{*} \phi, D_{A}^{*} \psi_{j}\right\rangle=0
$$

Hence $\left\langle P_{A} \phi, \psi_{j}\right\rangle=0$ for all $j$ and $P_{A} \phi \in V_{E}^{\perp}$. Since we already know that $P_{A} \phi \in V_{E}$, we must have $P_{A} \phi=0$ and $P_{A}=\Pi$.

## Asymptotic for $G \phi$

We also need to know the asymptotic behavior of $G_{A} \phi$ for $\phi \in V_{E}$.
Lemma 1.2-2. For $\phi \in V_{E}$, we have

$$
\begin{equation*}
G_{A} \phi=r^{2} \frac{\phi}{4}+O\left(r^{-2}\right) \tag{1.9}
\end{equation*}
$$

Proof: Notice first that

$$
\begin{equation*}
\nabla_{A}^{*} \nabla_{A} r^{2} \phi=-8 \phi-4 \nabla_{A}^{x} \phi+r^{2} \nabla_{A}^{*} \nabla_{A} \phi . \tag{1.10}
\end{equation*}
$$

Since $\phi \in V_{E}$, we have from Equation (1.7) that

$$
\phi=|x|^{-4} \rho(x) \hat{\phi}+O\left(r^{-4}\right)
$$

In fact, by doing the decomposition in some higher order Sobolev spaces, we see that

$$
\begin{aligned}
& \nabla_{A}^{x} O\left(r^{-4}\right)=O\left(r^{-4}\right), \text { and } \\
& \nabla_{A}^{*} \nabla_{A} O\left(r^{-4}\right)=O\left(r^{-6}\right) .
\end{aligned}
$$

Note now that

$$
\begin{aligned}
\nabla_{A}^{x} \rho(\nu) r^{-3} \hat{\phi} & =r \rho\left(\nabla_{A}^{\nu} \nu\right) r^{-3} \hat{\phi}+r \rho(\nu)\left(\nu \cdot r^{-3}\right) \hat{\phi}+\rho(\nu) r^{-3} r \nabla_{A}^{\nu} \hat{\phi} \\
& =0+r \rho(\nu)\left(-3 r^{-4}\right) \hat{\phi}+0 \\
& =-3 \rho(\nu) r^{-3} \hat{\phi}
\end{aligned}
$$

Thus, we have

$$
\begin{align*}
\nabla_{A}^{x} \phi & =-3 \phi+O\left(r^{-4}\right)+\nabla_{A}^{x} O\left(r^{-4}\right) \\
& =-3 \phi+O\left(r^{-4}\right) . \tag{1.11}
\end{align*}
$$

Similarly, since

$$
\begin{aligned}
\sum_{i=1}^{4} \nabla_{A}^{i} \nabla_{A}^{i}\left(r^{-4} \sum_{j=1}^{4} x_{j} \rho\left(\partial_{j}\right) \hat{\phi}\right)= & \sum_{i=1}^{4} \nabla_{A}^{i}\left(r^{-4} \sum_{j=1}^{4} \nabla_{A}^{i} x_{j} \rho\left(\partial_{j}\right) \hat{\phi}+\partial_{i}\left(r^{-4}\right) \sum_{j=1}^{4} x_{j} \rho\left(\partial_{j}\right) \hat{\phi}\right) \\
= & \sum_{i=1}^{4} \nabla_{A}^{i}\left(r^{-4} \rho\left(\partial_{i}\right) \hat{\phi}\right)+\sum_{i=1}^{4} \partial_{i}^{2}\left(r^{-4}\right) \rho(x) \hat{\phi} \\
& +\sum_{i=1}^{4} \partial_{i}\left(r^{-4}\right) \nabla_{A}^{i}(\rho(x) \hat{\phi}) \\
= & 2 \sum_{i=1}^{4} \partial_{i}\left(r^{-4}\right) \rho\left(\partial_{i}\right) \hat{\phi}+8 r^{-6} \rho(x) \hat{\phi} \\
= & 2 \sum_{i=1}^{4}(-4) r^{-6} x_{i} \rho\left(\partial_{i}\right) \hat{\phi}+8 r^{-6} \rho(x) \hat{\phi} \\
= & 0
\end{aligned}
$$

we have

$$
\begin{align*}
\nabla_{A}^{*} \nabla_{A} \phi & =-\sum_{i=1}^{4} \nabla_{A}^{i} \nabla_{A}^{i}\left(r^{-4} \rho(x) \hat{\phi}\right)+\nabla_{A}^{*} \nabla_{A} O\left(r^{-4}\right) \\
& =O\left(r^{-6}\right) \tag{1.12}
\end{align*}
$$

Summing up what we know with Equations (1.10), (1.11), and (1.12), we find that

$$
\nabla_{A}^{*} \nabla_{A} r^{2} \phi=-8 \phi-4\left(-3 \phi+O\left(r^{-4}\right)\right)+r^{2} O\left(r^{-6}\right) .
$$

Applying $G_{A}$ on both sides, we get

$$
G_{A} \phi=r^{2} \frac{\phi}{4}+G_{A} O\left(r^{-4}\right) .
$$

Lemma 3.3.35 from [DK90, p. 105] tells us that

$$
G_{A} O\left(r^{-4}\right)=O\left(r^{-4}\right)+O\left(r^{-2}\right)+O\left(r^{4-(2+4)}\right)=O\left(r^{-4}\right) .
$$

Substituting in the previous equation, we complete the proof.

## The Curvature of $a$

In view of the heuristic we explore in Chapter 2] we choose to temporarily view $V_{E}$ as the fiber of a trivial bundle over $\mathbb{R}^{4}$. The endomorphisms $a_{\mu}$ then team up to produce the constant connection

$$
a=a_{1} d x^{1}+\cdots+a_{4} d x^{4}
$$

Lemma 1.2-3. The curvature $F_{a}=\frac{1}{2}[a, a]$ of the connection $a$ on the trivial bundle with fiber $V_{E}$ over $\mathbb{R}^{4}$, seen as an element of $\operatorname{End}\left(V_{E}\right) \otimes \bigwedge^{2} \mathbb{R}^{4}$ is given by

$$
F_{a}=\Pi G_{A} \sum_{i=1}^{3} \rho\left(\bar{\epsilon}_{i}\right) \otimes \bar{\epsilon}_{i}-\frac{1}{8} \sum_{i=1}^{3} \Psi^{*} \rho\left(\epsilon_{i}\right) \Psi \otimes \epsilon_{i} .
$$

Proof: Component-wise, the curvature is

$$
\left(F_{a}\right)_{i j}=\frac{1}{2}\left[a_{i}, a_{j}\right] .
$$

To compute this curvature, we need to to see that

$$
\begin{equation*}
\left[m_{i}, D_{A}\right]=-\rho\left(\partial_{i}\right) \tag{1.13}
\end{equation*}
$$

and that

$$
\begin{equation*}
\left[\rho\left(\partial_{i}\right), G_{A}\right]=0 \tag{1.14}
\end{equation*}
$$

These two results are independent of the ASD condition on $A$. The second has to do with the fact that the $\partial_{i}$ are parallel.
Let $\phi \in V_{E}$. We have

$$
\begin{equation*}
\left[a_{i}, a_{j}\right](\phi)=\Pi\left(m_{i} \Pi m_{j} \phi\right)-\Pi\left(m_{j} \Pi m_{i} \phi\right) . \tag{1.15}
\end{equation*}
$$

Since at this point we have two formulas for $\Pi$, let's use both and use $P_{A}$ to denote the usage of the $1-D_{A} G_{A} D_{A}^{*}$ formula, and $\Pi$ for the scalar product type formula. We then compute

$$
\begin{equation*}
\Pi\left(m_{i} P_{A} m_{j} \phi\right)=\sum_{l=1}^{k} \lim _{r \rightarrow \infty}\left(\int_{B^{4}(r)}\left(m_{i}\left(1-D_{A} G_{A} D_{A}^{*}\right) m_{j} \phi, \psi_{l}\right)\right) \psi_{l} \tag{1.16}
\end{equation*}
$$

but

$$
\begin{equation*}
\left(m_{i}\left(1-D_{A} G_{A} D_{A}^{*}\right) m_{j} \phi, \psi_{l}\right)=\left(m_{i} m_{j} \phi, \psi_{l}\right)-\left(m_{i} D_{A} G_{A} D_{A}^{*} m_{j} \phi, \psi_{l}\right) \tag{1.17}
\end{equation*}
$$

The first term gets killed when we antisymmetrize with respect to $i$ and $j$. We thus compute only the second term. Equation (1.13) tells us that

$$
\begin{align*}
\left(m_{i} D_{A} G_{A} D_{A}^{*} m_{j} \phi, \psi_{l}\right) & =\left(m_{i} D_{A} G_{A} m_{j} D_{A}^{*} \phi, \psi_{l}\right)+\left(m_{i} D_{A} G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right) \\
& =\left(D_{A} m_{i} G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right)-\left(\rho\left(\partial_{i}\right) G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right) . \tag{1.18}
\end{align*}
$$

The second term gives the $\Pi G_{A} \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right)$ part of the curvature once we integrate, take the limit, substitute, antisymmetrize as asked by Equation (1.15) and divide by 2 to get its part in $F_{a}$.
As proved in Roe98, p. 46], we have

$$
\left(D_{A} s, \psi_{l}\right)-\left(s, D_{A}^{*} \psi_{l}\right)=\sum_{h=1}^{4} \partial_{h}\left(\rho\left(\partial_{h}\right) s, \psi_{l}\right) .
$$

Since $D_{A}^{*} \psi_{l}=0$, the first term of Equation (1.18) transforms:

$$
\left(D_{A} m_{i} G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right)=\sum_{h=1}^{4} \partial_{h}\left(\rho\left(\partial_{h}\right) G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right) .
$$

We now integrate by parts over the ball of radius $r$ to obtain

$$
\int_{B^{4}(r)}\left(D_{A} m_{i} G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right)=\int_{S^{3}(r)}\left(\rho(\nu) m_{i} G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right) .
$$

We now use the asymptotic given by Equations (1.7) and (1.9) to get

$$
\begin{gathered}
\psi_{l}=r^{-3} \rho(\nu) \hat{\psi}_{l}+O\left(r^{-4}\right), \text { and } \\
G_{A} \phi=\rho(\nu) \hat{\phi} / 4 r+O\left(r^{-2}\right)
\end{gathered}
$$

Then

$$
\begin{aligned}
\int_{S^{3}(r)}(\rho(\nu) & \left.m_{i} G_{A} \rho\left(\partial_{j}\right) \phi, \psi_{l}\right) \\
& =\sum_{h=1}^{4} \int_{S^{3}(r)} r^{-1} x_{i} x_{h}\left(\rho(\nu) \rho\left(\partial_{j}\right) \rho\left(\partial_{h}\right) \hat{\phi} / 4 r+O\left(r^{-2}\right), r^{-3} \rho(\nu) \hat{\psi}_{l}+O\left(r^{-4}\right)\right) \\
& =\sum_{h=1}^{4} \int_{S^{3}(r)}\left(r^{-5} / 4\right) x_{i} x_{h}\left(\rho\left(\partial_{j}\right) \rho\left(\partial_{h}\right) \hat{\phi}, \hat{\psi}_{l}\right)+\int_{S^{3}(r)}\left(O\left(r^{-5}\right)+O\left(r^{-4}\right)+O\left(r^{-4}\right)\right) \\
& =\left(r^{-5} / 4\right) \sum_{h=1}^{4}\left(\rho\left(\partial_{j}\right) \rho\left(\partial_{h}\right) \hat{\phi}, \hat{\psi}_{l}\right) \int_{S^{3}(r)} x_{i} x_{h}+O\left(r^{-4}\right) \operatorname{Vol}\left(S^{3}(r)\right) .
\end{aligned}
$$

As $r \rightarrow \infty$, the volume of $S^{3}(r)$ is $O\left(r^{3}\right)$ thus the last term vanish in the limit. The integral $\int_{S^{3}(r)} x_{i} x_{h}$ vanishes when $i \neq l$ and is otherwise

$$
r^{2} \operatorname{Vol}\left(S^{3}(r)\right) / 4=r^{5} \operatorname{Vol}\left(S^{3}(1)\right) / 4=r^{5} \pi^{2} / 2 .
$$

Thus we have from Equations (1.17), (1.16) and (1.15) that

$$
\left[a_{i}, a_{j}\right](\phi)=2 \Pi G_{A} \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right) \phi-\frac{\pi^{2}}{4} \sum_{l=1}^{k}\left(\rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right) \hat{\phi}, \hat{\psi}_{l}\right) \psi_{l} .
$$

Note that

$$
\sum_{1 \leq i, j \leq 4} \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right) d x^{i} \wedge d x^{j}=\sum_{i=1}^{3} \rho\left(\epsilon_{i}\right) \epsilon_{i}+\rho\left(\bar{\epsilon}_{i}\right) \bar{\epsilon}_{i}
$$

Remember that $\wedge^{+}$acts trivially on $V_{E} \subset \Gamma\left(S^{-} \otimes E\right)$ and $\wedge^{-}$acts trivially on $\Gamma\left(S^{+} \otimes E\right)$. Hence, the curvature is

$$
\begin{aligned}
F_{a}(\phi) & =\Pi G_{A} \sum_{i=1}^{3} \rho\left(\bar{\epsilon}_{i}\right) \phi \bar{\epsilon}_{i}-\frac{\pi^{2}}{8} \sum_{i=1}^{3} \sum_{l=1}^{k}\left(\rho\left(\epsilon_{i}\right) \hat{\phi}, \hat{\psi}_{l}\right) \psi_{l} \epsilon_{i} \\
& =\Pi G_{A} \sum_{i=1}^{3} \rho\left(\bar{\epsilon}_{i}\right) \phi \bar{\epsilon}_{i}-\frac{\pi^{2}}{2} \sum_{i=1}^{3} \sum_{l=1}^{k}\left(\rho\left(\epsilon_{i}\right) \Psi(\phi), \Psi\left(\psi_{l}\right)\right) \psi_{l} \epsilon_{i} .
\end{aligned}
$$

Using the scalar product given by Equation (1.5), we complete the proof.

## The ADHM data satisfies the conditions

Before going further, let's walk through the association between $\Phi$ and $\Psi$ in more details. Using the identification $\mathbb{R}^{4}=\mathbb{H}=S^{+}=S^{-}$, the Clifford multiplication $\rho(x): S^{+} \rightarrow S^{-}$is multiplication by $-\bar{x}$ and $\rho(x): S^{-} \rightarrow S^{+}$is multiplication by $x$.
Let $\epsilon_{i}$ and $\bar{\epsilon}_{i}$ denote the usual basis of $\Lambda^{+}$and $\Lambda^{-}$respectively. The action of self-dual forms on $S^{+}$is

$$
\rho\left(\epsilon_{1}\right)=2 i, \quad \rho\left(\epsilon_{2}\right)=2 j, \quad \rho\left(\epsilon_{3}\right)=2 k .
$$

We use the complex basis $s_{1}=1, s_{2}=j$ of $S^{+}$, with the identification

$$
\begin{align*}
\mathbb{C} \oplus \mathbb{C} & =S^{+}  \tag{1.19}\\
\left(z_{1}, z_{2}\right) & \mapsto z_{1}+j z_{2}
\end{align*}
$$

Then,

$$
\rho\left(\epsilon_{1}\right)=2\left[\begin{array}{cc}
i & 0  \tag{1.20}\\
0 & -i
\end{array}\right], \quad \rho\left(\epsilon_{2}\right)=2\left[\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right], \quad \rho\left(\epsilon_{3}\right)=2\left[\begin{array}{cc}
0 & -i \\
-i & 0
\end{array}\right] .
$$

We split the map $\Psi$ as

$$
\begin{gathered}
\Psi: V \rightarrow S^{+} \otimes W \\
\Psi=s_{1} \otimes \Psi_{1}+s_{2} \otimes \Psi_{2} .
\end{gathered}
$$

We identify $S^{+}$to its dual using the skewform $\omega=s^{1} \wedge s^{2}$ :

$$
\begin{aligned}
S^{+} & \rightarrow\left(S^{+}\right)^{*} \\
s & \mapsto \omega(\cdot, s) .
\end{aligned}
$$

Thus

$$
\begin{aligned}
& s_{1} \mapsto-s^{2} \\
& s_{2} \mapsto s^{1} .
\end{aligned}
$$

As mentioned before, in doing so we identify $\Psi$ with the map

$$
\begin{gathered}
\Phi=(\omega \otimes 1) \circ(1 \otimes \Psi): S^{+} \otimes V \rightarrow W \\
\Phi=-s^{2} \Psi_{1}+s^{1} \Psi_{2}
\end{gathered}
$$

The adjoints are

$$
\begin{gathered}
\Psi^{*}: S^{+} \otimes W \rightarrow V \\
\Psi^{*}=s^{1} \Psi_{1}^{*}+s^{2} \Psi_{2}^{*},
\end{gathered}
$$

and

$$
\begin{gathered}
\Phi^{*}: W \rightarrow S^{+} \otimes V \\
\Phi^{*}=-s_{2} \otimes \Psi_{1}^{*}+s_{1} \otimes \Psi_{2}^{*} .
\end{gathered}
$$

Thus

$$
\Psi^{*}=-(\omega \otimes 1) \circ\left(1 \otimes \Phi^{*}\right) .
$$

We have

$$
\begin{gather*}
\Psi^{*} \Psi: V \rightarrow V, \\
\Psi^{*} \Psi=\Psi_{1}^{*} \Psi_{1}+\Psi_{2}^{*} \Psi_{2} \tag{1.21}
\end{gather*}
$$

and

$$
\Phi^{*} \Phi: S^{+} \otimes W \rightarrow S^{+} \otimes W
$$

is given by $s_{2} \otimes s^{2} \otimes \Psi_{1}^{*} \Psi_{1}+s_{1} \otimes s^{1} \otimes \Psi_{2}^{*} \Psi_{2}-s_{1} \otimes s^{2} \otimes \Psi_{2}^{*} \Psi_{1}-s_{2} \otimes s^{1} \otimes \Psi_{1}^{*} \Psi_{2}$. In matrix form, this expression becomes

$$
\Phi^{*} \Phi=\left[\begin{array}{cc}
\Psi_{2}^{*} \Psi_{2} & -\Psi_{2}^{*} \Psi_{1}  \tag{1.22}\\
-\Psi_{1}^{*} \Psi_{2} & \Psi_{1}^{*} \Psi_{1}
\end{array}\right]
$$

Theorem 1.2-4. The ADHM data $\left(V_{E}, W_{E}, a, \Psi\right)$ obtained from the $S U(n)$-instanton connection $(E, A)$ satisfies the ADHM and nondegeneracy conditions (1.3) and (1.4).

Proof: Let's first consider the action of $[a, a]=2 F_{a}$ on $S^{+} \otimes V$. On that space, only the self-dual part matters. Recall from Lemma 1.2-3 that for $\phi \in V_{E}$, we have

$$
[a, a]^{+}(\phi)=2 F_{a}^{+}=\frac{1}{4} \sum_{i=1}^{3} \Psi^{*} \rho\left(\epsilon_{i}\right) \Psi(\phi) \epsilon_{i} .
$$

Let's break it down using the identification of Equation (1.19) and the matrices of Equation (1.20). We have

$$
\Psi^{*} \rho\left(\epsilon_{1}\right) \Psi=\left[\begin{array}{ll}
\Psi_{1}^{*} & \Psi_{2}^{*}
\end{array}\right]\left[\begin{array}{cc}
i & 0 \\
0 & -i
\end{array}\right]\left[\begin{array}{c}
\Psi_{1} \\
\Psi_{2}
\end{array}\right]=2 i\left(\Psi_{1}^{*} \Psi_{1}-\Psi_{2}^{*} \Psi_{2}\right),
$$

and similarly

$$
\begin{aligned}
& \Psi^{*} \rho\left(\epsilon_{1}\right) \Psi=2\left(\Psi_{2}^{*} \Psi_{1}-\Psi_{1}^{*} \Psi_{2}\right), \\
& \Psi^{*} \rho\left(\epsilon_{1}\right) \Psi=-2 i\left(\Psi_{2}^{*} \Psi_{1}+\Psi_{1}^{*} \Psi_{2}\right) .
\end{aligned}
$$

For the map $\rho([a, a]): S^{+} \otimes V \rightarrow S^{-} \otimes V$, only the self-dual part matters and in term of the
identification of Equation (1.19), it is given by

$$
\begin{aligned}
\rho([a, a]) & =\frac{1}{4}\left(2\left[\begin{array}{cc}
i & 0 \\
0 & -i
\end{array}\right] \Psi^{*} \rho\left(\epsilon_{1}\right) \Psi+2\left[\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right] \Psi^{*} \rho\left(\epsilon_{2}\right) \Psi+2\left[\begin{array}{cc}
0 & -i \\
-i & 0
\end{array}\right] \Psi^{*} \rho\left(\epsilon_{3}\right) \Psi\right) \\
& =\left[\begin{array}{cc}
\Psi_{2}^{*} \Psi_{2}-\Psi_{1}^{*} \Psi_{1} & -2 \Psi_{2}^{*} \Psi_{1} \\
-2 \Psi_{1}^{*} \Psi_{2} & \Psi_{1}^{*} \Psi_{1}-\Psi_{2}^{*} \Psi_{2}
\end{array}\right] .
\end{aligned}
$$

Using Equations (1.21) and (1.22), we see that

$$
-\rho([a, a])+2 \Phi^{*} \Phi=1 \otimes \Psi^{*} \Psi \quad \text { (ADHM equation). }
$$

Hence the ADHM equation (1.3) is satisfied.

### 1.3 From ADHM data to instanton

We now start with some ADHM data $(V, W, a, \Psi)$ satisfying the ADHM and nondegeneracy conditions (1.3) and (1.4), and want to construct an instanton connection

$$
\mathfrak{F}(V, W, a, \Psi)=(E, A)
$$

Recall from Equation (1.2) that we define $Q_{x}: S^{+} \otimes V \rightarrow S^{-} \otimes V \oplus W$ as

$$
Q_{x}=\left[\begin{array}{c}
\sum_{i=1}^{4} \rho\left(\partial_{i}\right) \otimes\left(a_{i}+x_{i}\right) \\
\Phi
\end{array}\right] .
$$

Since $Q_{x}$ is injective for every $x$, the map $Q_{x}^{*} Q_{x}: S^{+} \otimes V \rightarrow S^{+} \otimes V$ is an isomorphism at for every $x$. Let

$$
\begin{equation*}
G_{Q_{x}}=\left(Q_{x}^{*} Q_{x}\right)^{-1} \tag{1.23}
\end{equation*}
$$

be its inverse.
Let $E$ be the bundle with fiber $\operatorname{ker}\left(Q_{x}^{*}\right)$ at $x$. The bundle $E$ sits in the trivial bundle with fiber $S^{-} \otimes V \oplus W$. To simplify the notation, we drop the subscript $x$. The map

$$
\begin{equation*}
R:=1-Q G_{Q} Q^{*} \tag{1.24}
\end{equation*}
$$

is the orthogonal projection on $E$. We equip $E$ with the induced connection

$$
A:=R d .
$$

Theorem 1.3-1. The pair $(E, A)$ is an instanton connection on $\mathbb{R}^{4}$.

Proof: To compute the curvature, we first need a better grip on $G_{Q}$. We have

$$
Q^{*}=\left[\begin{array}{ll}
-\sum_{i=1}^{4} \rho\left(\partial_{i}\right) \otimes\left(a_{i}+x_{i}\right) & \Phi^{*}
\end{array}\right],
$$

thus

$$
\begin{align*}
Q^{*} Q & =-\sum_{1 \leq i, j \leq 4} \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right) \otimes\left(a_{i}+x_{i}\right)\left(a_{j}+x_{j}\right)+\Phi^{*} \Phi \\
& =1 \otimes \sum_{i=1}^{4}\left(a_{i}+x_{i}\right)^{2}-\frac{1}{2} \rho([a, a])+\Phi^{*} \Phi \\
& =1 \otimes\left(\sum_{i=1}^{4}\left(a_{i}+x_{i}\right)^{2}+\frac{1}{2} \Psi^{*} \Psi\right) . \tag{1.25}
\end{align*}
$$

Thus, $Q^{*} Q$ commutes with the quaternions and so does its inverse $G_{Q}$. We can then write

$$
G_{Q}=1 \otimes g_{Q}
$$

with $g_{Q}$ the inverse of the map $q: V \rightarrow V$ given by

$$
q=\sum_{i=1}^{4}\left(a_{i}+x_{i}\right)^{2}-(1 / 2) \Psi^{*} \Psi
$$

Notice that

$$
\left[\partial_{\mu}, Q\right]=\left[\begin{array}{c}
\rho\left(\partial_{\mu}\right) \\
0
\end{array}\right]
$$

The curvature acts on $\phi \in E$ as

$$
R d R d \phi=\sum_{1 \leq i, j \leq 4} R \partial_{i} R \partial_{j} \phi d x^{i} \wedge d x^{j}
$$

and

$$
\left.\left.\begin{array}{rl}
R \partial_{i} R \partial_{j} \phi & =R \partial_{i} \partial_{j} \phi-R \partial_{i} Q G_{Q} Q^{*} \partial_{j} \phi \\
& =R \partial_{i} \partial_{j} \phi-R \partial_{i} Q G_{Q} \partial_{j} Q^{*} \phi+R \partial_{i} Q F\left[\begin{array}{cc}
\rho\left(\partial_{j}\right) & 0
\end{array}\right] \phi \\
& =R \partial_{i} \partial_{j} \phi+0+R Q \partial_{i} G_{Q}\left[\rho\left(\partial_{j}\right)\right.
\end{array} 0\right] \phi-R\left[\begin{array}{c}
\rho\left(\partial_{i}\right) \\
0
\end{array}\right] G_{Q}\left[\begin{array}{ll}
\rho\left(\partial_{j}\right) & 0
\end{array}\right] \phi-\begin{array}{cc}
G_{Q} \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right) & 0 \\
0 & 0
\end{array}\right] .
$$

Thus, we find

$$
F_{\widehat{a}}=-R \sum_{i=1}^{3}\left[\begin{array}{cc}
G_{Q} \rho\left(\bar{\epsilon}_{i}\right) & 0 \\
0 & 0
\end{array}\right] \bar{\epsilon}_{i}
$$

and the connection has ASD curvature.
Since $q=r^{2}+O(r)$ as $r$ tends to $\infty$, we have that

$$
g_{Q}=r^{-2}+O\left(r^{-3}\right) \text { as } r \rightarrow \infty
$$

Then

$$
\begin{aligned}
& R=1-Q G_{Q} Q^{*} \\
& =1-\left[\begin{array}{c}
\rho(x)+O(1) \\
O(1)
\end{array}\right]\left(r^{-2}+O\left(r^{-3}\right)\right)\left[\begin{array}{ll}
-\rho(x)+O(1) & O(1)
\end{array}\right] \\
& =1-\left[\begin{array}{cc}
-\rho(x)^{2} r^{-2}+O\left(r^{-1}\right) & O\left(r^{-1}\right) \\
O\left(r^{-1}\right) & O\left(r^{-2}\right)
\end{array}\right] \\
& =\left[\begin{array}{cc}
O\left(r^{-1}\right) & O\left(r^{-1}\right) \\
O\left(r^{-1}\right) & 1+O\left(r^{-2}\right)
\end{array}\right],
\end{aligned}
$$

thus the curvature of the connection $R d$ on $E$ satisfies

$$
R\left[\begin{array}{cc}
G_{Q} & 0 \\
0 & 0
\end{array}\right]=\left[\begin{array}{cc}
O\left(r^{-3}\right) & 0 \\
0 & 0
\end{array}\right] \text { as } r \rightarrow \infty
$$

and is consequently in $L^{2}$.
The proof is now complete.
In fact, we can prove even a better asymptotic formula for $g_{Q}$.
Lemma 1.3-2. In fact, we even have

$$
\begin{equation*}
g_{Q}=r^{-2}-2 \sum x_{j} a_{j} r^{-4}-\left(\sum a_{j}^{2}-\frac{1}{2} \Psi^{*} \Psi\right) r^{-4}+4 \sum_{j, k} a_{j} a_{k} x_{j} x_{k} r^{-6}+O\left(r^{-5}\right) \tag{1.26}
\end{equation*}
$$

Proof: The proof is mechanical. We build up the asymptotics of $g_{Q}$ from the asymptotics of $\mathscr{G}$, term by term.

### 1.4 Uniqueness

We now wish to prove that starting from some ADHM data, creating the associated instanton and looking at the ADHM data this instanton produce, we come back to where we started. In other words, we prove in this section that the composition

$$
\begin{array}{cc}
\text { ADHM data } \\
(V, W, a, \Psi) \\
\text { Sect[1.3] } \\
\mathfrak{F} & \begin{array}{c}
\text { instanton } \\
(E, A)
\end{array} \\
\underset{\sim}{\text { Sect[1.2] }} & \text { ADHM data } \\
\left(V^{\prime}, W^{\prime}, a^{\prime}, \Psi^{\prime}\right)
\end{array}
$$

gives $\left(V^{\prime}, W^{\prime}, a^{\prime}, \Psi^{\prime}\right)=(V, W, a, \Psi)$.
We are thus searching for a proof that $V$ is isomorphic, in a somewhat canonical way, to $\operatorname{ker}\left(D^{*}\right)$ in sections of $S^{-} \otimes E$. As $E$ is a subbundle of $S^{-} \otimes V \oplus W$, we would be happy to find a map

$$
\psi: V \rightarrow S^{-} \otimes\left(S^{-} \otimes V \oplus W\right)
$$

or equivalently

$$
\tilde{\psi}: S^{+} \otimes V \rightarrow S^{-} \otimes V \oplus W
$$

such that

$$
\begin{gathered}
\operatorname{Im}(\tilde{\psi}) \subset E \\
D^{*} \psi=0
\end{gathered}
$$

$$
\psi \text { is injective. }
$$

Using the identification $S^{+} \equiv S^{-}$, define the map

$$
\begin{align*}
b: S^{+} \otimes V & \rightarrow S^{-} \otimes V \oplus W  \tag{1.27}\\
x & \mapsto(x, 0) .
\end{align*}
$$

Our candidate is

$$
\begin{equation*}
\tilde{\psi}=R b G_{Q} \tag{1.28}
\end{equation*}
$$

Obviously, $\operatorname{Im}(\tilde{\psi}) \subset E$ as $R$ projects on $E$. To prove that $D^{*} \psi=0$, we observe that

$$
\begin{aligned}
\widetilde{D_{A}^{*} \psi} & =\left\langle D_{A}^{*} \psi, \cdot\right\rangle \\
& =\sum_{j=1}^{4}\left\langle\rho\left(\partial_{j}\right) \nabla_{A}^{j} \psi, \cdot\right\rangle \\
& =-\sum_{j=1}^{4}\left\langle\nabla_{A}^{j} \psi, \rho\left(\partial_{j}\right) \cdot\right\rangle \\
& =-\sum_{j=1}^{4}\left(\nabla_{A}^{j} \tilde{\psi}\right) \circ \rho\left(\partial_{j}\right) .
\end{aligned}
$$

Hence

$$
\begin{aligned}
\widetilde{D^{*} \psi} & =-\sum_{j=1}^{4} R\left(\partial_{j} \tilde{\psi}\right) \rho\left(\partial_{j}\right) \\
& =-\sum_{j=1}^{4} R \partial_{j}\left(R b G_{Q}\right) \rho\left(\partial_{j}\right) \\
& =-\sum_{j=1}^{4} R\left(\partial_{j} R\right) b G_{Q} \rho\left(\partial_{j}\right)-\sum_{j=1}^{4} R b\left(\partial_{j} G_{Q}\right) \rho\left(\partial_{j}\right) \\
& =\sum_{j=1}^{4} R \partial_{j}\left(Q G_{Q} Q^{*}\right) b G_{Q} \rho\left(\partial_{j}\right)-\sum_{j=1}^{4} R b\left(\partial_{j} G_{Q}\right) \rho\left(\partial_{j}\right) \\
& =\sum_{j=1}^{4} R\left(\partial_{j} Q\right) G_{Q} Q^{*} b G_{Q} \rho\left(\partial_{j}\right)+\sum_{j=1}^{4} R Q \partial_{j}\left(G_{Q} Q^{*}\right) b G_{Q} \rho\left(\partial_{j}\right)-\sum_{j=1}^{4} R b\left(\partial_{j} G_{Q}\right) \rho\left(\partial_{j}\right)
\end{aligned}
$$

On that last line, the second sum is obviously null as $R Q=0$. As for the first sum, observe that $\partial_{j} Q=b \rho\left(\partial_{j}\right)$, while

$$
\begin{equation*}
G_{Q} Q^{*} b G_{Q}=\frac{1}{2} \sum_{i=1}^{4} \rho\left(\partial_{i}\right) \partial_{i} G_{Q} \tag{1.29}
\end{equation*}
$$

Indeed, from Equation (1.25), we derive that $\partial_{i} G_{Q}=-2 G_{Q}\left(a_{i}+x_{i}\right) G_{Q}$. Equation (1.29) follows
immediately once we recognize that $Q^{*} b=-\sum_{i=1}^{4} \rho\left(\partial_{i}\right)\left(a_{i}+x_{i}\right)$.
Going back to where we left, we have

$$
\begin{aligned}
\widetilde{D_{A}^{*} \psi} & =\frac{1}{2} \sum_{1 \leq i, j \leq 4} R b \rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right)\left(\partial_{i} G_{Q}\right) \rho\left(\partial_{j}\right)-\sum_{j=1}^{4} R b\left(\partial_{j} G_{Q}\right) \rho\left(\partial_{j}\right) \\
& =R b \sum_{i=1}^{4}\left(\frac{1}{2} \sum_{j=1}^{4} \rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right)-\rho\left(\partial_{i}\right)\right) \partial_{i} G_{Q} \\
& =0
\end{aligned}
$$

Now that we know that $\psi$ maps $V$ to sections of $S^{-} \otimes E$ satisfying the Dirac equation, we would like to see that $\psi$ is actually an isomorphism $V \rightarrow V^{\prime}$. To prove this result, we use the following analytic lemma. But first let $\partial^{2}$ denote the Laplacian $\sum_{i} \partial_{i}^{2}$.

Lemma 1.4-1. For $\psi$ defined by Equation (1.28), we have

$$
\begin{equation*}
\psi^{*} \psi=-\frac{1}{4} \partial^{2} g_{Q} \tag{1.30}
\end{equation*}
$$

Proof: Let $t r_{2}$ denote the trace along the spinor factor. Notice that

$$
\psi^{*} \psi=\operatorname{tr}_{2}\left(\tilde{\psi}^{*} \tilde{\psi}\right)
$$

Indeed, as we have

$$
\begin{array}{ll}
\psi=s_{1} \otimes \psi_{1}+s_{2} \otimes \psi_{2}, & \psi^{*}=s^{1} \otimes \psi_{1}^{*}+s^{2} \otimes \psi_{2}^{*} \\
\tilde{\psi}=s^{1} \otimes \psi_{1}+s^{2} \otimes \psi_{2}, & \tilde{\psi}^{*}=s_{1} \otimes \psi_{1}^{*}+s_{2} \otimes \psi_{2}^{*}
\end{array}
$$

then

$$
\begin{aligned}
\operatorname{tr}_{2}\left(\tilde{\psi}^{*} \tilde{\psi}\right) & =s_{1} \otimes s^{1} \psi_{1}^{*} \psi_{1}+s_{1} \otimes s^{2} \psi_{1}^{*} \psi_{2}+s_{2} \otimes s^{1} \psi_{2}^{*} \psi_{1}+s_{2} \otimes s^{2} \psi_{2}^{*} \psi_{2} \\
& =\psi_{1}^{*} \psi_{1}+\psi_{2}^{*} \psi_{2}=\psi^{*} \psi
\end{aligned}
$$

On one hand, multiplying $\tilde{\psi}^{*} \tilde{\psi}$ by $G_{Q}^{-1}$ yields

$$
\begin{aligned}
\tilde{\psi}^{*} \tilde{\psi} & =G_{Q} b^{*} R b G_{Q} \\
& =G_{Q}\left(b^{*} b G_{Q}-b^{*} Q G_{Q} Q^{*} b G_{Q}\right) \\
& =G_{Q}{ }^{2}-\frac{1}{2} \sum_{i=1}^{4} G_{Q} b^{*} Q \rho\left(\partial_{i}\right) \partial_{i} G_{Q} .
\end{aligned}
$$

On the other hand, for each $k$, we have

$$
\begin{aligned}
\partial_{i}^{2}\left(G_{Q}^{-1} G_{Q}\right) & =\partial_{i}\left(2\left(a_{i}+x_{i}\right) G_{Q}\right)+\partial_{i}\left(G_{Q}^{-1} \partial_{i} G_{Q}\right) \\
& =2 G_{Q}+4\left(a_{i}+x_{i}\right) \partial_{i} G_{Q}+G_{Q}^{-1} \partial_{i}^{2} G_{Q}
\end{aligned}
$$

As $\partial^{2}=\sum_{i=1}^{4} \partial_{i}^{2}$, those equalities sum up to

$$
\begin{equation*}
\partial^{2} G_{Q}=-8 G_{Q}^{2}-4 \sum_{i=1}^{4} G_{Q}\left(a_{i}+x_{i}\right) \partial_{i} G_{Q} \tag{1.31}
\end{equation*}
$$

We hence have

$$
\begin{aligned}
\psi^{*} \psi & =\operatorname{tr}_{2}\left(\tilde{\psi}^{*} \tilde{\psi}\right) \\
& =\operatorname{tr}_{2}\left(G_{Q}{ }^{2}-\frac{1}{2} \sum_{i=1}^{4} G_{Q} b^{*} Q \rho\left(\partial_{i}\right) \partial_{i} G_{Q}\right) \\
& =2 g_{Q}{ }^{2}-\frac{t r_{2}}{2}\left(\sum_{1 \leq i, j \leq 4} G_{Q}\left(a_{j}+x_{j}\right) \rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right) \partial_{i} G_{Q}\right) \\
& =2 g_{Q}{ }^{2}+\sum_{i=1}^{4} g_{Q}\left(a_{i}+x_{i}\right) \partial_{i} g_{Q}-\frac{t r_{2}}{2}\left(G_{Q} \sum_{i \neq j}\left(a_{j}+x_{j}\right) \rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right) \partial_{i} G_{Q}\right) .
\end{aligned}
$$

The $t r_{2}$ part of this last line cancels as for $j \neq k$ we have $\rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right)=-\rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right)$ while $\operatorname{tr}_{2}\left(\rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right)\right)=\operatorname{tr}_{2}\left(\rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right)\right)$. Equation (1.30) then follows from this computation and Equation (1.31). The proof of the lemma is now complete.

Using that lemma, we show that $\psi$ is an isomorphism. Recall from Lemma 1.3-2 that we have the asymptotic behavior $g_{Q}=r^{-2}+O\left(r^{-3}\right)$ as $r \rightarrow \infty$, and more to the point, it is so that $\partial_{r} g_{Q}=-2 r^{-3}+O\left(r^{-4}\right)$ as $r \rightarrow \infty$. We then have

$$
\begin{aligned}
\int_{\mathbb{R}^{4}} \psi^{*} \psi & =-\frac{1}{4} \int_{\mathbb{R}^{4}} \partial^{2} g_{Q} \\
& =-\lim _{r \rightarrow \infty} \frac{1}{4} \int_{S^{3}(r)} \partial_{r} g_{Q} \\
& =\frac{1}{2} \operatorname{Vol}\left(S^{3}\right) i d_{V} \\
& =\pi^{2} \dot{d} d_{V} .
\end{aligned}
$$

Thus an orthonormal basis $v_{1}, \ldots, v_{k}$ of $V$ gives an orthonormal basis $\pi^{-1} \psi\left(v_{1}\right), \ldots, \pi^{-1} \psi\left(v_{k}\right)$ of $V^{\prime}$. Indeed, as $\psi=\sum_{j} \psi\left(v_{j}\right) v^{j}$, we have $\psi^{*}=\sum_{j} v_{j} \otimes\left\langle\psi\left(v_{j}\right), \cdot\right\rangle$ hence pointwise we have $\psi^{*} \psi=\sum_{i, j}\left\langle\psi\left(v_{j}\right), \psi\left(v_{i}\right)\right\rangle v_{j} \otimes v^{i}$, or once we integrate,

$$
\int_{\mathbb{R}^{4}} \psi^{*} \psi=\sum_{1 \leq i, j \leq 4}\left\langle\psi\left(v_{j}\right), \psi\left(v_{i}\right)\right\rangle_{L^{2}} v_{j} \otimes v^{i}
$$

Remembering that the $\psi\left(v_{j}\right)$ do not have norm 1 , we compute the endomorphism $a_{\mu}^{\prime}$ of $V^{\prime}$ :

$$
a_{\mu}^{\prime}=\frac{1}{\pi^{2}} \sum_{1 \leq i, j \leq 4}\left\langle\psi\left(v_{j}\right),-x_{\mu} \psi\left(v_{i}\right)\right\rangle_{L^{2}} v_{j} \otimes v^{i}=-\frac{1}{\pi^{2}} \int_{\mathbb{R}^{4}} x_{\mu} \psi^{*} \psi .
$$

Using from Lemma 1.3-2 the asymptotic knowledge that

$$
\partial_{r} g_{Q}=-2 r^{-3}+6 \sum_{j=1}^{4} x_{j} a_{j} r^{-5}+O\left(r^{-5}\right)
$$

we find that

$$
\begin{aligned}
\int_{\mathbb{R}^{4}} x_{\mu} \psi^{*} \psi= & -\frac{1}{4} \int_{\mathbb{R}^{4}} x_{\mu} \partial^{2} g_{Q} \\
= & \frac{1}{4} \lim _{r \rightarrow \infty} \int_{S^{3}(r)} \partial_{r}\left(x_{\mu}\right) g_{Q}-x_{\mu} \partial_{r} g_{Q} \\
= & \frac{1}{4} \lim _{r \rightarrow \infty} \int_{S^{3}(r)} \frac{x_{\mu}}{r}\left(r^{-2}-2 \sum_{j=1}^{4} x_{j} a_{j} r^{-4}+O\left(r^{-4}\right)\right) \\
& \quad+2 x_{\mu} r^{-3}-6 \sum_{j=1}^{4} x_{\mu} x_{j} a_{j} r^{-5}+O\left(r^{-5}\right) \\
= & \frac{1}{4} \lim _{r \rightarrow \infty}\left(3 r \int_{S^{3}(1)} x_{j}-8 \sum_{j=1}^{4} a_{j} \int_{S^{3}(1)} x_{\mu} x_{j}+O\left(r^{-1}\right)\right) \\
= & -2 a_{\mu} \int_{S^{3}} x_{j}^{2} \\
= & -2 a_{\mu} \operatorname{Vol}\left(S^{3}\right) / 4 \\
= & -\pi^{2} a_{\mu} .
\end{aligned}
$$

Hence we obtain back the same maps, $a_{\mu}^{\prime}=a_{\mu}$.

### 1.5 Completeness

We now close this chapter by proving that every instanton arise from some ADHM data. In other words, we prove in this section that the composition

| instanton | Sect[1.2 | ADHM data | Sect[13] | n |
| :---: | :---: | :---: | :---: | :---: |
| $(E, A)$ | $\mathfrak{N}$ | $\left(V^{\prime}, W^{\prime}, a^{\prime}, \Psi^{\prime}\right)$ | $\mathfrak{F}$ | $\left(E^{\prime}, A^{\prime}\right)$ |

gives an instanton $\left(E^{\prime}, A^{\prime}\right)$ gauge equivalent to $(E, A)$.
This last fact establishes the validity of Theorem 1.1-1
Since $\left(E^{\prime}, A^{\prime}\right)=\mathfrak{F}(V, W, a, \Psi)$, then $E_{x}^{\prime}$ sits in $S^{-} \otimes V \oplus W$ as $\operatorname{ker}\left(Q_{x}^{*}\right)$. Elements of $S^{-} \otimes V$, once contracted using the skewform $\omega$, give sections of $E$. Hence the map

$$
\begin{aligned}
\alpha_{x}: S^{-} \otimes V \oplus W & \rightarrow E_{x} \\
{\left[\begin{array}{l}
\psi \\
\phi
\end{array}\right] } & \mapsto \omega G_{A} \psi(x)+\frac{1}{2} \phi(x)
\end{aligned}
$$

is well defined. Its adjoint $\alpha_{x}^{*}$ gives the map we want between $E$ and $E^{\prime}$. To prove that fact, we need to show that $Q_{x}^{*} \alpha_{x}^{*}=0$, or equivalently $\alpha_{x} Q_{x}=0$. For any $s \otimes \psi \in S^{+} \otimes V$, we have

$$
\begin{aligned}
\alpha_{x} Q_{x}(s \otimes \psi) & =\sum_{j=1}^{4} \omega\left(\rho\left(\partial_{j}\right) s \otimes G_{A}\left(a_{j}+x_{j}\right) \psi\right)(x)+\frac{1}{2} \Phi(s \otimes \psi)(x) \\
& \left.=\omega\left(s, \sum_{j=1}^{4} \rho\left(\partial_{j}\right) G_{A}\left(a_{j}+x_{j}\right) \psi\right)(x)+\frac{1}{2} \Psi(\psi)(x)\right),
\end{aligned}
$$

hence it suffices to prove that for every $\psi \in V$,

$$
\begin{equation*}
\left.\sum_{j=1}^{4} \rho\left(\partial_{j}\right) G_{A}\left(a_{j}+x_{j}\right) \psi\right)(x)+\frac{1}{2} \Psi(\psi)(x)=0 \tag{1.32}
\end{equation*}
$$

Since $V$ sits in $L_{-3+\epsilon}^{2}$ for all small $\epsilon$, it must be that $\rho\left(\partial_{j}\right) G_{A} a_{j} \psi \in L_{-1+\epsilon}^{2}$. Equation (1.9) guarantees that

$$
\begin{aligned}
G_{A} \psi & =\frac{r^{2} \psi}{4}+O\left(r^{-2}\right) \\
& =\frac{\rho(x)}{2 r^{2}} \Psi(\psi)+O\left(r^{-2}\right)
\end{aligned}
$$

hence $\rho(x) G_{A} \psi+(1 / 2) \Psi(\psi) \in L_{-1+\epsilon}^{2}$ as well. Hence the left-hand-side of Equation (1.32) is all in $L_{-1+\epsilon}^{2}$, and thus must be 0 if harmonic.
Applying $\Delta$ kills the $\Psi(\psi)$ term, and for the rest we obtain

$$
\begin{aligned}
\left.\Delta\left(\sum_{j=1}^{4} \rho\left(\partial_{j}\right) G_{A}\left(a_{j}+x_{j}\right) \psi\right)(x)\right) & =\sum_{j=1}^{4} \rho\left(\partial_{j}\right) a_{j} \psi+\rho\left(\partial_{j}\right) m_{j} \psi-2 \rho\left(\partial_{j}\right) \nabla_{A}^{j} G_{A} \psi \\
& =\left(\sum_{j=1}^{4} \rho\left(\partial_{j}\right) D_{A} G_{A} D_{A}^{*} m_{j} \psi\right)-2 D_{A}^{*} G_{A} \psi \\
& =\left(\sum_{j=1}^{4} \rho\left(\partial_{j}\right)\left(\sum_{k=1}^{4} \rho\left(\partial_{i}\right) \nabla_{A}^{i}\right) \rho\left(\partial_{j}\right) G_{A} \psi\right)-2 D_{A}^{*} G_{A} \psi
\end{aligned}
$$

which is 0 since

$$
\sum_{j=1}^{4} \rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right) \rho\left(\partial_{j}\right)=2 \rho\left(\partial_{i}\right) .
$$

The validity of Equation (1.32) is now established, and so is the fact that $\alpha^{*}$ maps $E$ to $E^{\prime}$. Proving that $\left(\alpha^{*}\right)^{*} A^{\prime}$ is gauge equivalent to $A$ is an exercise in rewriting [KN90, Section 6a]. It is left to the author and the reader for further study.

## Chapter 2

## The Nahm transform heuristic

Heuristically, the Nahm transform places instantons on $\mathbb{R}^{4} / \Lambda$ in reciprocity with certain data on the dual $\mathbb{R}^{4^{*}} / \Lambda^{*}$. This short chapter describes in more details this heuristic.
Let $\Lambda$ be a closed subgroup of $\mathbb{R}^{4}$. We associate two $\mathbb{R}$-vector spaces to $\Lambda$ :

$$
\begin{aligned}
\Lambda_{\mathbb{R}} & :=\text { maximal } \mathbb{R} \text {-linear subspace of } \Lambda, \\
\Lambda_{\mathbb{Z}} & :=\text { orthogonal complement of } \Lambda_{\mathbb{R}} \text { in } \Lambda, \\
\mathfrak{L}(\Lambda) & :=\mathbb{R} \text {-vector subspace of } \mathbb{R}^{4} \text { generated by } \Lambda .
\end{aligned}
$$

Obviously, $\Lambda$ is isomorphic to some $\mathbb{R}^{r} \times \mathbb{Z}^{s}$ with $r+s \leq 4$, and then $\operatorname{dim} \Lambda_{\mathbb{R}}:=r, \operatorname{dim} \Lambda_{\mathbb{Z}}:=s$, and $\operatorname{dim} \mathfrak{L}(\Lambda):=r+s$.
The dual $\Lambda^{*}$ is defined to be

$$
\Lambda^{*}:=\left\{z \in \mathbb{R}^{4^{*}} \mid z(\Lambda) \subset \mathbb{Z}\right\}
$$

Obviously, $\operatorname{dim} \Lambda_{\mathbb{R}}^{*}=4-r-s, \operatorname{dim} \Lambda_{\mathbb{Z}}^{*}=s$, and $\operatorname{dim} \mathfrak{L}\left(\Lambda^{*}\right)=4-r$.
We start with a $\operatorname{SU}(n)$-bundle $E$ over $\mathbb{R}^{4}$ and a $\operatorname{SU}(n)$-connection $A$ on $E$, both invariant under the action of $\Lambda$. We require the curvature of $A$, denoted $F_{A}$, to be anti-self-dual (ASD) and to have finite $L^{2}$-norm on the quotient $\mathbb{R}^{4} / \Lambda$. Equivalently, we start with a vector bundle $E$ on $\mathbb{R}^{4} / \Lambda$, and a connection $A$ and endomorphisms $a_{1}, \ldots, a_{r}$ of $E$, such that $\left(A, a_{1}, \ldots, a_{r}\right)$ satisfies the $(4-r)$ dimensional reduction of the ASD equation, as given in Appendix and such that its $L^{2}$-energy, to be defined appropriately, is finite.
It might happen that the finite $L^{2}$-norm condition is too strong to get any interesting solutions, in which case we need to search for a better condition. This need arises for example on $\mathbb{R}^{2} \times S^{1}$, and Cherkis-Kapustin give in [CK98] an appropriate logarithmic decay condition for the endomorphism $a_{1}$. Anyhow, we are exploring a heuristic for studying instantons or their various dimensional reductions, not a precise recipe, and adjustments need to be made in many cases.
Suppose now $z$ is an element of $\mathbb{R}^{4^{*}}$, the space of $\mathbb{R}$-valued linear functions on $\mathbb{R}^{4}$. We define the bundle $L_{z}$ over $\mathbb{R}^{4}$ to be a trivial $\mathbb{C}$-bundle with connection

$$
\omega_{z}:=2 \pi i z=2 \pi i \sum_{j=1}^{4} z_{j} d x^{j} .
$$

Notice that $L_{z}$ is invariant under the action of $\Lambda$, and that it is flat. Furthermore, whenever $z^{\prime} \in \Lambda^{*}$, the bundles with connections $L_{z}$ and $L_{z+z^{\prime}}$ are isomorphic over $\mathbb{R}^{4} / \Lambda$, or equivalently we can
parameterize flat connections over $\mathbb{R}^{4} / \Lambda$ by elements of $\mathbb{R}^{4^{*}} / \Lambda^{*}$.
Indeed, for $z \in \mathbb{R}^{4^{*}}$, define the function $g_{z}: \mathbb{R}^{4} \rightarrow \mathrm{U}(1)$ by

$$
g_{z}(x)=e^{-2 \pi i z(x)}
$$

and notice that $g_{z}$ is invariant under the action of $\Lambda$ for each $z \in \Lambda^{*}$. But more importantly, we have

$$
g_{z} \cdot \omega_{z^{\prime}}=\omega_{z^{\prime}-z}
$$

We write $A_{z}$ for the connection $A \otimes 1+1 \otimes \omega_{z}$ on $E \otimes L_{z}=E$. For $z \in \mathbb{R}^{4^{*}}$, consider the operator

$$
D_{A_{z}}^{*}: \Gamma\left(\mathbb{R}^{4}, S^{-} \otimes E \otimes L_{z}\right) \rightarrow \Gamma\left(\mathbb{R}^{4}, S^{+} \otimes E \otimes L_{z}\right)
$$

A section of the bundle $S^{-} \otimes E \otimes L_{z}$ is said to be in $L_{\Lambda}^{2}$ if it is invariant under the action of $\Lambda$ and if its $L^{2}$-norm over $\mathbb{R}^{4} / \Lambda$ is finite.
We set

$$
V_{z}:=L_{\Lambda}^{2} \cap \operatorname{ker}\left(D_{A_{z}}^{*}\right)
$$

By putting some restrictions on $A$, for example that $(E, A)$ has no flat factor, we ensure via the use of the Weitzenbock formula that $L_{\Lambda}^{2} \cap \operatorname{ker}\left(D_{A_{z}}\right)=\{0\}$. At this point, we need to prove that $D_{A_{z}}^{*}$ is Fredholm to prove that $V$ is a bundle, and compute the index of $D_{A_{z}}^{*}$ to find its rank.
It turns out in many cases that $D_{A_{z}}^{*}$ is not Fredholm for every $z$, which is a good thing. Without going into details, suppose for example that $D_{A_{z}}^{*}$ was Fredholm everywhere when $\Lambda=\mathbb{Z}^{3}$. As we explore in this present thesis, the object created by the Nahm transform would be a monopole over $T^{3}$. But as one can show (see [Pau98, Prop. 1]), monopoles over compact 3-manifolds are not very interesting.
Notice that for any section $\phi$ of $S^{-} \otimes E$, we have

$$
\begin{aligned}
D_{A_{z}}^{*}\left(g_{z} \phi\right) & =g_{z}\left(D_{A}^{*} \phi+2 \pi i c l(z) \phi\right)+\operatorname{cl}\left(\operatorname{grad} g_{z}\right) \phi \\
& =g_{z} D_{A}^{*} \phi
\end{aligned}
$$

Then for all $z^{\prime} \in \Lambda^{*}$, we have an isomorphism

$$
\begin{equation*}
g_{z^{\prime}}: V_{z} \rightarrow V_{z+z^{\prime}} \tag{2.1}
\end{equation*}
$$

hence $V$ is a bundle over $\mathbb{R}^{4^{*}} / \Lambda^{*}$.
It is important to keep two points of view in parallel, the full $\mathbb{R}^{4^{*}}$ view and the quotient $\mathbb{R}^{4^{*}} / \Lambda^{*}$ view. In the first view, we perform a curvature computation and observe how far the curvature of the connection $B$ we introduce on $V$ is from being anti-self-dual. In the second view, we can sometime reduce dimension, as in the $\mathbb{R}^{4}$-ADHM case.
Let's stick to the $\mathbb{R}^{4^{*}}$ point of view for now. We define a connection $B$ on $V$. Each fiber $V_{z}$ is in fact contained in $L_{\Lambda}^{2}\left(S^{-} \otimes E\right)$. We can then consider the trivial connection $d^{z}$ in the trivial bundle of fiber $L_{\Lambda}^{2}\left(S^{-} \otimes E\right)$, and its projection $P d^{z}$ to $V$.
The operator $D_{A_{z}}^{*} D_{A_{z}}$ should be invertible, and we use the Green's operator $G_{A_{z}}=\left(D_{A_{z}}^{*} D_{A_{z}}\right)^{-1}$ to define the projection $P$ by the formula

$$
P=1-D_{A_{z}} G_{A_{z}} D_{A_{z}}^{*}
$$

To help ease the notation, set $\Omega:=2 \pi i \sum_{j=1}^{4} c l\left(d x^{j}\right) d z^{j}$.
Let's now compute the curvature of $B$. Notice that

$$
\left[d^{z}, D_{A_{z}}\right]=\left[d^{z}, D_{A}+2 \pi i \sum_{j=1}^{4} z_{j} c l\left(d x^{j}\right)\right]=\Omega
$$

and similarly for $D_{A_{z}}^{*}$.
The Leibnitz's rule tells us that

$$
d^{z}\left\langle P d^{z} \phi, \psi\right\rangle=\left\langle d^{z} P d^{z} \phi, \psi\right\rangle-\left\langle P d^{z} \phi, d^{z} \psi\right\rangle,
$$

but as $\left\langle P d^{z} \phi, \psi\right\rangle=\left\langle d^{z} \phi, \psi\right\rangle$, we also have

$$
\begin{aligned}
d^{z}\left\langle P d^{z} \phi, \psi\right\rangle & =d^{z}\left\langle d^{z} \phi, \psi\right\rangle \\
& =\left\langle d^{z 2} \phi, \psi\right\rangle-\left\langle d^{z} \phi, d^{z} \psi\right\rangle \\
& =-\left\langle d^{z} \phi, d^{z} \psi\right\rangle,
\end{aligned}
$$

hence the curvature $F_{B}$ can be computed as follows:

$$
\begin{aligned}
\left\langle\left(P d^{z}\right)^{2} \phi, \psi\right\rangle & =\left\langle d^{z} P d^{z} \phi, \psi\right\rangle \\
& =\left\langle P d^{z} \phi, d^{z} \psi\right\rangle-\left\langle d^{z} \phi, d^{z} \psi\right\rangle \\
& =-\left\langle D_{A_{z}} G_{A_{z}} D_{A_{z}} d^{z} \phi, d^{z} \psi\right\rangle \\
& =\left\langle D_{A_{z}} G_{A_{z}} \Omega \phi, d^{z} \psi\right\rangle .
\end{aligned}
$$

Let $\nu$ be the normal vector field to $S^{r-1}(R) \times T^{s}$. The integration by parts necessary to bring $D$ on the right-hand-side of the scalar product introduces a boundary term

$$
\begin{equation*}
\lim _{\partial}:=\lim _{R \rightarrow \infty} \int_{S^{r-1}(R) \times T^{s}}\left\langle c l(\nu) G \Omega \phi, d^{z} \psi\right\rangle . \tag{2.2}
\end{equation*}
$$

Performing the said integration by parts, we obtain

$$
\begin{aligned}
\left\langle F_{B} \phi, \psi\right\rangle & =\left\langle G_{A_{z}} \Omega \phi, D_{A_{z}} d^{z} \psi\right\rangle+\lim _{\partial} \\
& =-\left\langle G_{A_{z}} \Omega \phi, \Omega \psi\right\rangle+\lim _{\partial} \\
& =\left\langle G_{A_{z}} \phi, \Omega \wedge \Omega \psi\right\rangle+\lim _{\partial} .
\end{aligned}
$$

The first term is ASD since

$$
\begin{aligned}
\Omega \wedge \Omega & =-4 \pi^{2} \sum_{1 \leq i, j \leq 4} c l\left(d x^{i} \wedge d x^{j}\right) d z^{i} \wedge d z^{j} \\
& =-4 \pi^{2} \sum_{j=1}^{3}\left(c l\left(\epsilon_{j}\right) \epsilon_{j}+\operatorname{cl}\left(\bar{\epsilon}_{j}\right) \bar{\epsilon}_{j}\right) .
\end{aligned}
$$

and $\Lambda^{-}$acts on $S^{-} \otimes E$. (Remember from page 25 that the $\epsilon_{j}$ and $\bar{\epsilon}_{j}$ are the usual basis of $\Lambda^{+}$and $\Lambda^{-}$respectively.)
The key to the ADHMN construction of instantons on $\mathbb{R}^{4}$ as portrayed in Chapter $\square$ is to get a good grip of the boundary term $\lim _{\partial}$, which is precisely the role of $\Phi$ and the ADHM condition (1.3).

Let's now see how the heuristic described so far can lead to the ADHMN construction for $\mathbb{R}^{4}$, and explore what happen in the second point of view, where we look at everything on the quotient $\mathbb{R}^{4^{*}} / \Lambda^{*}$.
Of course, our first task is now to interpret the connection $B$ in that new setting. It passes really well to the quotient by $\Lambda_{\mathbb{Z}}^{*}$, the difficulty lies in quotienting by the remaining $\Lambda_{\mathbb{R}}^{*}$.
Suppose we have coordinates $\left(x_{1}, \ldots, x_{4}\right)$ on $\mathbb{R}^{4}$ and associated coordinates $\left(z_{1}, \ldots, z_{4}\right)$ on $\mathbb{R}^{4^{*}}$ such that

$$
\Lambda_{\mathbb{R}}^{*}=\left\{z_{1}=\cdots=z_{r+s}=0\right\} .
$$

Let's introduce new coordinates $\left(u_{1}, \ldots, u_{r+s}, v_{1}, \ldots, v_{4-r-s}\right)=\left(z_{1}, \ldots, z_{4}\right)$. Using the gauge transformation of Equation (2.1), we regard the space $V_{u}$ as $V_{(u, 0)}$. We go back to the other point of view using the isomorphism

$$
g_{v}: V_{u} \rightarrow V_{(u, v)} .
$$

Suppose $\phi_{u}$ is a section of $V$ on $\mathbb{R}^{4^{*}} / \Lambda^{*}$. Then

$$
\begin{aligned}
g_{v}^{-1} B\left(g_{v} \phi_{u}\right) & =g_{v}^{-1} P\left(\left(d^{u}+d^{v}\right) g_{v} \phi_{u}\right) \\
& =P d^{u} \phi_{u}-2 \pi i \sum_{j=r+s+1}^{4} P m_{x_{j}}\left(\phi_{u}\right) .
\end{aligned}
$$

Hence the connection matrices $B_{j}$ for $r+s+1 \leq j \leq 4$ get replaced by $-2 \pi i P m_{x_{j}}$. Of course, $m_{x_{j}}$ is multiplication by $x_{j}$, as in Equation (1.6).
The heuristic presented in this chapter therefore allows us to say that, to any connection $A$ on a bundle $E$ over $\mathbb{R}^{4} / \Lambda$ satisfying the appropriate dimensional reduction (see Appendix of the ASD equation, we can associate the following objects:

1. a family $V$ of vector spaces over $\mathbb{R}^{4^{*}} / \Lambda^{*}$ defined by the kernel $V_{z}=\operatorname{ker}\left(D_{A_{z}}^{*}\right)$ of the Dirac operator lifted to $\mathbb{R}^{4}$, family which forms a bundle over the open set on which $D_{A_{z}}^{*}$ is Fredholm;
2. a connection $B$ on $V$ defined by the projection $P d^{z}$ of the trivial connection $d^{z}$ on $\mathbb{R}^{4^{*}} / \Lambda^{*}$;
3. a family of endomorphisms $a_{\alpha}$, as many as $\operatorname{dim}\left(\mathfrak{L}(\Lambda)^{\perp}\right)=\operatorname{dim} \Lambda_{\mathbb{R}}^{*}$, of $V$, defined by the projection $-2 \pi i P m_{x_{\alpha}}$ of the multiplication $m_{x_{\alpha}}$ by the coordinate $x_{\alpha}$ on $\mathbb{R}^{4}$.

When the boundary term of Equation (2.2) is 0 , the $\left(B,\left\{a_{\alpha}\right\}\right)$ satisfies the appropriate dimensional reduction of the ASD equation on $\mathbb{R}^{4^{*}} / \Lambda^{*}$. When it is not, further ad hoc analysis is required.
The various dimensional reductions of the ASD equations are presented in Appendix A

## Chapter 3

## Dirac Spectrum of Product Manifolds

In this chapter, we compute the spectrum of the Dirac operator of manifolds which are products, using the spectrum of the Dirac operator on each factor.
The spectrum of the Dirac operator has been computed in many cases: spheres in [Tra93], threedimensional Berger spheres in [Hit74], odd-dimensional Berger spheres in [Bär96], flat manifolds in [Pfä00], tori in [Fri84], simply connected Lie groups in [Feg87], fibrations over $S^{1}$ in [Kra01], etc. To learn more about the current state of affairs related to spectra and eigenvalue estimates for Dirac operator, see the survey paper [Bär00].
The explicit formula of Theorem 3.2-1 for the spectrum of the Dirac operator on a general product manifold seems to be missing in the literature.
As a special case, we compute in theorem 3.3-1 the spectrum of the Dirac operator on the $n$ dimensional torus. Our computation confirms the result in [Fri84].
In Section 3.1 we construct the spinor bundles on $M \times N$ from the spinor bundles on the factor. This section is somewhat inspired by [Kli]. In Section 3.2] we compute the Dirac spectrum of product manifolds $M \times N$. In Section 3.3, we use the acquired knowledge to compute the Dirac spectrum for a torus. In Section 3.4 we consider the special case where $N=S^{1}$ and we twist the spinor bundle by a flat line bundle on $S^{1}$. We extend the result to the torus.
For the reader interested in a shortcut to the main result of this thesis, most of this chapter can be skipped, only Theorem 3.4-1 and the remarks that follow it are necessary.

### 3.1 Complex Spinor Bundles of $M \times N$

Let $M$ and $N$ be two manifolds equipped with spinor bundles $S_{M}$ and $S_{N}$ respectively. Let $p_{1}, p_{2}$ be the projections on the first and second factor of $M \times N$. In this section, we construct a spinor bundle $S$ for $M \times N$ using $S_{M}$ and $S_{N}$.
For a vector $v$ tangent to $M$, let $\rho_{v}$ and $c_{v}$ denote the Clifford multiplication on $S_{M}$ and $S$ respectively. For a vector $w$ tangent to $N$, we similarly use $\rho_{w}$ and $c_{w}$.
Suppose at least one of the manifolds, say $M$, is even-dimensional. Thus $S_{M}$ splits as $S_{M}^{+} \oplus S_{M}^{-}$. In that case, we set

$$
S:=p_{1}^{*} S_{M} \otimes p_{2}^{*} S_{N}
$$

and set $S_{ \pm}:=p_{1}^{*} S_{M}^{ \pm} \otimes p_{2}^{*} S_{N}$.

Suppose on the contrary that both manifolds are odd-dimensional. In that case, we set

$$
S:=\left(p_{1}^{*} S_{M} \otimes p_{2}^{*} S_{N}\right) \oplus\left(p_{1}^{*} S_{M} \otimes p_{2}^{*} S_{N}\right)
$$

Let $S_{+}$denote the first factor, and $S_{-}$the second.
In both cases, with respect to the decomposition $S=S_{+} \oplus S_{-}$, set

$$
c_{v}:=\left[\begin{array}{ll} 
& \rho_{v} \otimes 1 \\
\rho_{v} \otimes 1 &
\end{array}\right], \text { and } c_{w}:=\left[\begin{array}{ll}
1 \otimes \rho_{w} & \\
& -1 \otimes \rho_{w}
\end{array}\right] .
$$

Proposition 3.1-1. The map $c$ is a Clifford multiplication. The bundle $S$ is a spinor bundle.
Proof: First, for the Clifford multiplication. We have

$$
c_{w} c_{v}=\left[\begin{array}{ll} 
& -\rho_{v} \otimes \rho_{w} \\
\rho_{v} \otimes \rho_{w} &
\end{array}\right]=-c_{v} c_{w}
$$

hence $c_{v} c_{w}+c_{w} c_{v}=0$, as wanted.
Let $d_{m}$ be the dimension of a complex irreducible representation of $\mathbb{C} l_{m}$. From LM89, Thm 5.8, p. 33], we have

$$
d_{2 k+1}=d_{2 k}=2^{k} .
$$

Since $S$ has the required dimension and is a Clifford bundle, it must be a spin bundle.

## Positive and Negative Spinors

At this point, a warning is necessary: the splitting $S=S_{+} \oplus S_{-}$on $M^{m} \times N^{n}$ when both $m$ and $n$ are odd is not the same as the splitting $S=S^{+} \oplus S^{-}$in terms of positive and negative spinors. The second splitting appears through the isomorphism

$$
\begin{align*}
S^{+} \oplus S^{-} & \rightarrow S_{+} \oplus S_{-} \\
(a, b) & \mapsto(i a, a)+(b, i b), \tag{3.1}
\end{align*}
$$

and its inverse

$$
\begin{align*}
S_{+} \oplus S_{-} & \rightarrow S^{+} \oplus S^{-} \\
(a, b) & \mapsto \frac{1}{2}(b-a i, a-b i) . \tag{3.2}
\end{align*}
$$

Let's verify the accuracy of this last statement. The orientation class in $\mathbb{C l} l_{m}$ is

$$
\omega= \begin{cases}i^{n} v o l & \text { for } m=2 n \\ i^{n+1} \text { vol } & \text { for } m=2 n+1\end{cases}
$$

It satisfies $\omega^{2}=1$. We define $S^{ \pm}=(1 \pm c(\omega)) S$; see [LM89] Prop 5.15, p. 36], where his $\omega_{\mathbb{C}}$ is our $\omega$. Note that $c(\omega)$ acts as $\pm 1$ on $S^{ \pm}$.
Since $m$ is odd, we know from [LM89, Prop 5.9, p. 34] that $\rho\left(\omega_{M}\right)$ can be either $\pm 1$ and that the corresponding representations are inequivalent. However, the definition of the complex spin representation is independent of which irreducible representation of $\mathbb{C} l_{m}$ is used; see [LM89, Prop 5.15]. For simplicity, let's fix a sign and always choose the spinor bundle for which the orientation class acts with that sign.

The orientation class is $\omega=-i \omega_{M} \omega_{N}$, and

$$
c(\omega)=\left[\begin{array}{ll}
-i \rho\left(\omega_{M}\right) \otimes \rho\left(\omega_{N}\right) & i \rho\left(\omega_{M}\right) \otimes \rho\left(\omega_{N}\right) \\
-i &
\end{array}\right]
$$

The decomposition $S=S^{+} \oplus S^{-}$given by $S^{+}=\{(i a, a)\}$ and $S^{-}=\{(b, i b)\}$ is thus accurate. When $m$ and $n$ are both even, the orientation class is $\omega=\omega_{M} \omega_{N}$, and

$$
c(\omega)=\left[\begin{array}{ll}
\rho\left(\omega_{M}\right) \otimes \rho\left(\omega_{N}\right) & \\
& \rho\left(\omega_{M}\right) \otimes \rho\left(\omega_{N}\right)
\end{array}\right]
$$

The splitting $S_{N}=S_{N}^{+} \oplus S_{N}^{-}$induces splittings for $S_{+}$and $S_{-}$. Using the very evident notation coming from those splittings, we have

$$
\begin{aligned}
& S^{+}=S_{++} \oplus S_{--}, \text {and } \\
& S^{-}=S_{+-} \oplus S_{-+}
\end{aligned}
$$

### 3.2 Dirac Spectrum Formula

To describe the spectrum of the Dirac operator on $M \times N$, we need to work with multisets, as most eigenvalues appear with high multiplicity. For the multiset $A$, let $A^{\# a}$ be the union of $a$ copies of $A$. Of course, the kernel is always a very special set and we need some notation for the multiplicity of 0 in the spectrum when $m$ is even. For that purpose, set

$$
k_{M}^{ \pm}:=\operatorname{dim} \operatorname{ker} D_{M}^{ \pm}
$$

Theorem 3.2-1. The spectrum of the Dirac operator on $M^{m} \times N^{n}$ on the spinor bundle constructed in Section 3.1 from chosen spinor bundles on $M$ and $N$ is given as a multiset in terms of the respective spectrum $\Sigma_{M}$ and $\Sigma_{N}$ by the formula

$$
\Sigma_{M \times N}= \begin{cases} \pm\left|\Sigma_{M} \times \Sigma_{N}\right|, & \text { if } m \text { and } n \text { are odd; } \\ \pm\left|\Sigma_{M}^{00} \times \Sigma_{N}\right| \cup\left(\Sigma_{N}\right)^{\# k_{M}^{+}} \cup\left(-\Sigma_{N}\right)^{\# k-\bar{M}}, & \text { if } m \text { is even. }\end{cases}
$$

Proof: For the decomposition $S=S_{+} \oplus S_{-}$, the Dirac operator on the spinor bundle of $M \times N$ is

$$
D=\left[\begin{array}{cc}
D_{N} & D_{M} \\
D_{M} & -D_{N}
\end{array}\right]
$$

Suppose first that $m$ and $n$ are both odd.
Let $\left\{\psi_{\mu}\right\}_{\mu \in \Sigma_{M}}$ be a basis of eigenvectors of $D_{M}$ on $L^{2}\left(S_{M}\right)$, with $D_{M} \psi_{\mu}=\mu \psi_{\mu}$. There might be more than one function called $\psi_{\mu}$, but this abuse of notation should not cause any problems. Similarly, let $\left\{\phi_{\nu}\right\}_{\nu \in \Sigma_{N}}$ be a basis of eigenvectors of $D_{N}$ for $L^{2}\left(S_{N}\right)$.
We have $L^{2}(S)=\mathbb{C}^{2} \otimes L^{2}\left(S_{M}\right) \otimes L^{2}\left(S_{N}\right)$, and on $\mathbb{C}^{2} \otimes \psi_{\mu} \otimes \phi_{\nu}$, we have

$$
D=\left[\begin{array}{cc}
\nu & \mu \\
\mu & -\nu
\end{array}\right]
$$

This matrix has two eigenvectors, of respective eigenvalue $\pm \sqrt{\mu^{2}+\nu^{2}}$. The corresponding eigenvectors are respectively $\left[\nu \pm \sqrt{\mu^{2}+\nu^{2}} \mu\right]^{T}$ when $\mu \neq 0$ while for $\mu=0$, the vector $\left[\begin{array}{ll}1 & 0\end{array}\right]^{T}$ and $\left[\begin{array}{ll}0 & 1\end{array}\right]^{T}$ correspond respectively to the eigenvalues $\nu$ and $-\nu$.
We just proved the theorem for $m$ and $n$ both odd.
Suppose now that $m$ is even. As before, let $\left\{\psi_{\mu}\right\}$ and $\left\{\phi_{\nu}\right\}$ be eigenbasis for $D_{M}$ and $D_{N}$. Recall that the positive and negative eigenspaces of any Dirac operator $D_{M}$ (not just the spin one) are isomorphic via

$$
\psi_{\mu}=\psi_{\mu}^{+}+\psi_{\mu}^{-} \mapsto \psi_{-\mu}=\psi_{\mu}^{+}-\psi_{\mu}^{-}
$$

For $\mu=0$, we have the positive spinors $\psi_{0}^{+} \in L^{2}\left(S_{M}^{+}\right)$, and the negative spinors $\psi_{0}^{-} \in L^{2}\left(S_{M}^{-}\right)$.
We thus have a unique decomposition

$$
f=f_{0}^{+} \psi_{0}^{+}+f_{0}^{-} \psi_{0}^{-}+\sum_{\mu>0}\left(f_{\mu}+f_{-\mu}\right) \psi_{\mu}^{+}+\sum_{\mu>0}\left(f_{\mu}-f_{-\mu}\right) \psi_{\mu}^{-}
$$

Note that

$$
\begin{aligned}
& D\left(\psi_{0}^{+} \otimes \phi_{\nu}\right)=\nu \psi_{0}^{+} \otimes \phi_{\nu}, \text { and } \\
& D\left(\psi_{0}^{-} \otimes \phi_{\nu}\right)=-\nu \psi_{0}^{-} \otimes \phi_{\nu}
\end{aligned}
$$

Then ker $D_{M}$ contributes

$$
\left(\Sigma_{N}\right)^{\# k_{M}^{+}} \cup\left(-\Sigma_{N}\right)^{\# k_{M}^{-}}
$$

Suppose now that $\mu \neq 0$. Then

$$
D\left(\psi_{\mu} \otimes \phi_{\nu}\right)=\mu \psi_{\mu} \otimes \phi_{\nu}+\nu \psi_{-\mu} \otimes \phi_{\nu}
$$

Thus, for $\mu>0$, the Dirac operator acts on the span of $\psi_{\mu} \otimes \phi_{\nu}$ and $\psi_{-\mu} \otimes \phi_{\nu}$ as

$$
\left[\begin{array}{cc}
\mu & \nu \\
\nu & -\mu
\end{array}\right]
$$

This matrix has two eigenvectors, of respective eigenvalue $\pm \sqrt{\mu^{2}+\nu^{2}}$.
We just proved the theorem for $m$ even. The proof is now complete.
Corollary 3.2-2. When both $m$ and $n$ are odd, we have

$$
\operatorname{ker}\left(D^{+}\right) \text {is isomorphic to } \operatorname{ker}\left(D^{-}\right)
$$

When both $m$ and $n$ are even, we have

$$
\begin{aligned}
k_{M \times N}^{+} & =k_{M}^{+} k_{N}^{+}+k_{M}^{-} k_{N}^{-}, \text {and } \\
k_{M \times N}^{-} & =k_{M}^{+} k_{N}^{-}+k_{M}^{-} k_{N}^{+}
\end{aligned}
$$

Proof: When $m$ and $n$ are both odd, a basis of $\operatorname{ker}(D)$ is given by all the

$$
\left[\begin{array}{l}
i \\
1
\end{array}\right] \otimes \psi_{0} \otimes \phi_{0}\left(\text { sections of } S^{+}\right), \text {and }\left[\begin{array}{l}
1 \\
i
\end{array}\right] \otimes \psi_{0} \otimes \phi_{0}\left(\text { sections of } S^{-}\right)
$$

When $m$ and $n$ are both even, a basis of $\operatorname{ker}(D)$ is given by all the

$$
\begin{aligned}
& \psi_{0}^{+} \otimes \phi_{0}^{+}, \psi_{0}^{-} \otimes \phi_{0}^{-}\left(\text {sections of } S^{+}\right), \\
& \psi_{0}^{+} \otimes \phi_{0}^{-}, \psi_{0}^{-} \otimes \phi_{0}^{+}\left(\text {sections of } S^{-}\right) .
\end{aligned}
$$

The proof is now complete.
Corollary 3.2-3. The index of $D^{+}$on an even-dimensional product $M^{m} \times N^{n}$ is given by

$$
\operatorname{ind}(M \times N)= \begin{cases}0, & \text { if } m \text { and } n \text { are both odd; } \\ \operatorname{ind}(M) \cdot \operatorname{ind}(N), & \text { if } m \text { and } n \text { are both even. } .\end{cases}
$$

### 3.3 Dirac Spectrum of $T^{n}$

The work we have done so far allows us to compute the spectrum of the Dirac operator on the $n$-torus $T^{n}=\mathbb{R}^{n} / \Lambda$.
Recall first that there are two possibilities for the spinor bundle $S_{1}$, the trivial, denoted $\mathbb{S}_{0}$, and the nontrivial, denoted $\mathbb{S}_{1}$. In either case, we have $\rho_{\theta}=i$.

Let $\mathbb{S}_{\epsilon_{1} \cdots \epsilon_{n}}$ denote the spin structure constructed inductively using the procedure given in Section 3.1 Let $b_{1}^{*}, \ldots, b_{n}^{*}$ be a basis for the lattice

$$
\Lambda^{*}:=\left\{\lambda^{*} \in \mathbb{R}^{*} \mid \lambda^{*}(\Lambda) \subset \mathbb{Z}\right\}
$$

dual to the lattice $\Lambda$ defining $T^{n}$.
Theorem 3.3-1. The Dirac Spectrum for the spin structure $\mathbb{S}_{\epsilon_{1} \cdots \epsilon_{n}}$ on the torus $T^{n}=\mathbb{R}^{n} / \Lambda$ is the multiset of all the

$$
\pm 2 \pi\left|b^{*}+\sum \epsilon_{j} b_{j}^{*} / 2\right|
$$

for $b^{*} \in \Lambda^{*}$ given with multiplicity $2^{\lfloor n / 2\rfloor-1}$.
Proof: Note that we can rewrite this theorem as saying that the spectrum is

$$
\pm\left|\Sigma_{S^{1}} \times \cdots \times \Sigma_{S^{1}}\right|^{\# 2^{\lfloor n / 2\rfloor-1}}
$$

When $n=1$ and $S^{1}$ has length $\ell$, we have

$$
\Sigma_{S^{1}}=(2 \pi / \ell)(\epsilon / 2+\mathbb{Z}) .
$$

The factor $1 / 2$ counterbalances the $\pm$ as $-\Sigma_{S^{1}}=\Sigma_{S^{1}}$.
In fact, we will use the more general fact that $-\Sigma_{T^{n}}=\Sigma_{T^{n}}$, and that $k_{T^{2 k}}^{+}=k_{T^{2 k}}^{-}$.
Suppose now that $n=2 k+1$. Then $T^{n}=T^{2 k} \times S^{1}$. From Theorem 3.2-1] we know the spectrum is

$$
\pm\left|\Sigma_{T^{2 k}}^{>0} \times \Sigma_{S^{1}}\right| \cup\left(\Sigma_{S^{1}}\right)^{\# k_{T^{2 k}}^{+}} \cup\left(-\Sigma_{S^{1}}\right)^{\# k_{T^{2 k}}^{-}},
$$

which is

$$
\pm\left|\Sigma_{T^{2 k}} \times \Sigma_{S^{1}}\right| \not{ }^{\# \frac{1}{2}}
$$

By induction, this multiset is
as wanted.
Suppose now that $n=2 k$. Then $T^{n}=T^{2 k-1} \times S^{1}$, and using Theorem 3.2-1 we find that this spectrum is
as wanted.
The proof is now complete.

### 3.4 Tensoring by $L_{z}$

Suppose now we change the Clifford bundle, tensoring it by the flat bundle $L_{z}$ on $S^{1}$, which is trivial with connection $2 \pi i z d \theta$. Since it is constant in the $M$ direction, it doesn't affect $\Sigma_{M}$.
Whether $m$ is odd or even, the new Dirac operator is

$$
D=\left[\begin{array}{cc}
D_{S^{1}}-2 \pi z & D_{M} \\
D_{M} & -\left(D_{S^{1}}-2 \pi z\right)
\end{array}\right],
$$

hence we just need to shift the eigenvalues of $D_{S^{1}}$ by $-2 \pi z$.
For $m$ odd and even, the new spectrum is respectively

$$
\begin{gathered}
\pm\left|\Sigma_{M} \times\left(\Sigma_{S^{1}}-2 \pi z\right)\right|, \text { and } \\
\pm\left|\Sigma_{M}^{>0} \times\left(\Sigma_{S^{1}}-2 \pi z\right)\right| \cup\left(\Sigma_{S^{1}}-2 \pi z\right)^{\# k_{M}^{+}} \cup\left(-\Sigma_{S^{1}}+2 \pi z\right)^{\# k_{M}^{-}} .
\end{gathered}
$$

Suppose now we tensor the spinor bundle on $T^{n}$ by the flat bundle $L_{z}$ with connection

$$
2 \pi i \sum_{j=1}^{n} z_{j} d x^{j}
$$

with $z \in \Lambda^{*}$. A modification of the proof used in Section 3.3 works to prove by induction that the spectrum of the Dirac operator $D_{z}$ on $T^{n}$ for $n \geq 2$ is

$$
\left\{ \pm 2 \pi\left|b^{*}-z+\sum_{i=1}^{n} \epsilon_{i} b_{i}^{*} / 2\right| \mid b^{*} \in \Lambda^{*}\right\}^{\# 2^{\lfloor n / 2\rfloor-1}}
$$

This result implies that 0 is in the spectrum if and only if $z \in \sum_{i=1}^{n} \epsilon_{i} b_{i}^{*} / 2+\Lambda^{*}$.
As we need it later on, let's summarize the situation for the three-dimensional torus $T^{3}$.
Theorem 3.4-1. Choose the trivial spin structure $S$ on $T^{3}=\mathbb{R}^{3} / \Lambda$. Pick $z \in \mathbb{R}^{3 *}$. The spectrum of the Dirac operator on $S \otimes L_{z}$ is given by the multiset

$$
\pm 2 \pi\left|\Lambda^{*}-z\right| .
$$

The eigenspaces coming from $0 \in \Lambda^{*}$ are found as follows. Notice first that $S=\mathbb{C}^{2}$ and that $S^{+}$ and $S^{-}$on $T^{2}$ are found using Maps (3.1) and (3.2) by the projections

$$
P_{+}=\frac{1}{2}\left[\begin{array}{cc}
1 & i \\
-i & 1
\end{array}\right], \text { and } P_{-}=\frac{1}{2}\left[\begin{array}{cc}
1 & -i \\
i & 1
\end{array}\right] .
$$

Then

$$
\begin{aligned}
D & =D_{T^{2}, z}-2 \pi z_{3} P_{+}+2 \pi z_{3} P_{-} \\
& =2 \pi\left[\begin{array}{cc}
-z_{2} & -\left(z_{1}+z_{3} i\right) \\
-\left(z_{1}-z_{3} i\right) & z_{2}
\end{array}\right] .
\end{aligned}
$$

When $\left(z_{1}, z_{3}\right)=(0,0)$, the eigenspaces are the

$$
\begin{aligned}
& \mathbb{C}\left[\begin{array}{l}
1 \\
0
\end{array}\right] \quad \text { of eigenvalue }-2 \pi z_{2}, \\
& \mathbb{C}\left[\begin{array}{l}
0 \\
1
\end{array}\right] \quad \text { of eigenvalue } 2 \pi z_{2},
\end{aligned}
$$

while when $\left(z_{1}, z_{3}\right) \neq(0,0)$, they are the

$$
\mathbb{C}\left[\begin{array}{c}
-\left(z_{1}+z_{3} i\right) \\
z_{2} \pm|z|
\end{array}\right] \quad \text { of eigenvalue } \pm 2 \pi z_{2} .
$$

We can consider the bundle $V_{2 \pi \epsilon}$ created from the eigenspace of $D_{z}$ of eigenvalue $2 \pi \epsilon$ on the sphere $|z|=\epsilon$. A non-zero section is obviously given by $\left[-\left(z_{1}+z_{3} i\right) \quad z_{2}+|z|\right]^{T}$. This section vanishes only at $(0, \epsilon, 0)$ and its multiplicity is obviously 1 , hence $c_{1}\left(V_{2 \pi \epsilon}\right)=e\left(V_{2 \pi \epsilon}\right)= \pm 1$, depending on the choice of orientation class.

## Chapter 4

## Dirac Spectrum of $S^{n}$

Because of the splitting relation

$$
D_{\mathbb{R} \times S^{3}}^{ \pm}= \pm \frac{\partial}{\partial t}+D_{S^{3}}
$$

and because of the conformal equivalence $\mathbb{R} \times S^{3} \equiv \mathbb{R}^{4} \backslash\{0\}$, the eigenvalues of the Dirac operator $D$ on $S^{3}$ and the kernel of the Dirac operator on $\mathbb{R}^{4}$ are intimately related. So we first proceed to study the eigenvalues of $D$ on $S^{3}$. We then exploit these results in Section 7.3 to prove the asymptotic behavior of Equation (1.7).
In Section 4.1] we compute the Spectrum of $D_{S^{3}}$ in a way quite similar to Hitchin's [Hit74. In Section 4.2, we confirm the results of Section 4.1, using a construction of Trautman for the spectrum of the Dirac operators of spheres. The drawback of Trautman's method is that it does not give easily the multiplicities, which is why we need the computations.

## 4.1 $S^{3}$ : Spherical harmonics and representations

Let's start by writing down an explicit formula for $D_{S^{3}}$ and $D_{S^{3}}^{2}$. Consider the left-invariant orthonormal frame on $S^{3}$ given by

$$
\begin{aligned}
e_{1}(x) & :=x \cdot i, \\
e_{2}(x) & :=x \cdot j, \text { and } \\
e_{3}(x) & :=x \cdot k .
\end{aligned}
$$

As derivations, the $e_{i}$ satisfy the commuting relations obtained by cyclicly permuting $\{1,2,3\}$ in the expression

$$
\left[e_{1}, e_{2}\right]=2 e_{3} .
$$

The Levi-Civita connection matrix in that orthonormal frame is

$$
\left[\omega_{b}^{a}\right]=\left[\begin{array}{ccc}
0 & -e^{3} & e^{2} \\
e^{3} & 0 & -e^{1} \\
-e^{2} & e^{1} & 0
\end{array}\right] .
$$

The spinor bundle $S\left(S^{3}\right)$ of $S^{3}$ is a trivial $\mathbb{H}$-bundle. The vectors $e_{1}, e_{2}$ and $e_{3}$ act by Clifford multiplication on $S\left(S^{3}\right)$ simply by left-multiplication by $i, j$ and $-k$ respectively, so that the volume
element acts as +1 . Thus the spin connection endomorphism is

$$
\begin{aligned}
\Omega & =\frac{1}{2} \sum_{1 \leq a<b \leq 3} \omega_{b}^{a} \otimes c l\left(e_{b}\right) c l\left(e_{a}\right) \\
& =\frac{1}{2}\left(e^{1} \otimes i+e^{2} \otimes j+e^{3} \otimes k\right),
\end{aligned}
$$

and the spin connection is $d+\Omega$. The Dirac operator hence, obeys the rule

$$
D=i e_{1}+j e_{2}-k e_{3}+\frac{3}{2} .
$$

Then, we have the formula

$$
D^{2}=-e_{1} e_{1}-e_{2} e_{2}-e_{3} e_{3}+i e_{1}+j e_{2}-k e_{3}+\frac{9}{4}
$$

In this formula, the part which looks second order is actually the standard Laplacian on $S^{3}$ :

$$
\Delta:=-e_{1} e_{1}-e_{2} e_{2}-e_{3} e_{3} .
$$

The eigenvalues of $D$ are distributed symmetrically with respect to 0 , as we now establish.
Theorem 4.1-1. Let $n \equiv 3 \bmod 4$ and $M$ be a riemannian manifold of dimension $n$. Let $\phi$ be an orientation-reversing isometry. Then the spectrum $\Sigma$ of the Dirac operator $D$ on the spinor bundle of $M$ is symmetric: $\Sigma=-\Sigma$.
Proof: In this case, the square of the Clifford volume element $\omega$ is $\omega^{2}=1$. Thus $C l_{n}$ splits as

$$
C l_{n}=V_{+} \oplus V_{-} .
$$

It turns out that the $V_{ \pm}$are invariant for the action of $C l_{n}$.
The algebra map $\alpha: C l_{n} \rightarrow C l_{n}$, generated by $\alpha(v)=-v$ for $v \in \mathbb{R}^{n}$, exchange $V_{+}$and $V_{-}$. It is an isomorphism of $\operatorname{Spin}(n)$-representations since $\operatorname{Spin}(n) \subset C l_{n}^{0}$.
Note now that there is a canonical isomorphism $P_{S O}(M) \equiv \phi^{*} P_{S O}(M)$ : on the fibers over a given point, it is given by $\left(e_{1}, \ldots, e_{n}\right) \mapsto\left(d \phi\left(e_{n}\right), \ldots, d \phi\left(e_{1}\right)\right)$.
This isomorphism induces an isomorphism $P_{\text {Spin }}(M) \equiv \phi^{*} P_{\text {Spin }}(M)$.
Let $\ell: \operatorname{Spin}(n) \rightarrow V_{ \pm}$be left-multiplication and set

$$
S_{ \pm}:=P_{S p i n}(M) \times_{\ell} V_{ \pm} .
$$

Both are "the" spinor bundle on $M$. They are isomorphic via $\alpha$. Let's choose $S_{+}$to work with.
Set $C l_{\text {Spin }}(M):=P_{\text {Spin }}(M) \times C l_{n}$. Then, as Clifford-modules, we have the isomorphism $\phi^{*} C l_{\text {Spin }}(M) \equiv C l_{\text {Spin }}(M)$. This isomorphism exchange $S_{ \pm}$and $S_{\mp}$.
Suppose now that $s \in \Gamma\left(S_{+}\right)$and consider $\alpha\left(\phi^{*} s\right) \in \Gamma\left(S_{+}\right)$. At the point $x$, we have that $\alpha\left(\phi^{*} s\right)(x)=\alpha(s(\phi(x))) \in\left(S_{-}\right)_{\phi(x)} \equiv\left(S_{+}\right)_{x}$.
The connection on $C l_{\text {Spin }}(M)$, being $1 / 4 \cdot \sum_{1 \leq a, b \leq n} \omega_{b a} e_{a} e_{b}$, is preserved by the canonical isomorphism. Thus

$$
D\left(\alpha\left(\phi^{*} s\right)\right)=-\alpha\left(D\left(\phi^{*} s\right)\right)=-\alpha\left(\phi^{*}(D s)\right) .
$$

Hence, if $D s=\lambda s$, then $D\left(\alpha\left(\phi^{*} s\right)\right)=-\lambda \cdot \alpha\left(\phi^{*} s\right)$. The proof is now complete.

Recall that $L^{2}\left(S^{3}\right)$ has a decomposition in eigenspaces of $D$ and $\Delta$; in fact, they are linked since the Laplacian commutes with our canonical basis as it is parallel:

$$
\left[\Delta, e_{a}\right]=0 \text { for } a=1,2,3
$$

Indeed, we have, for example, that

$$
\begin{aligned}
{\left[\Delta, e_{1}\right] } & =-e_{1} e_{1} e_{1}-e_{2} e_{2} e_{1}-e_{3} e_{3} e_{1}+e_{1} e_{1} e_{1}+e_{1} e_{2} e_{2}+e_{1} e_{3} e_{3} \\
& =-e_{2}\left[e_{2}, e_{1}\right]-e_{2} e_{1} e_{2}-e_{3}\left[e_{3}, e_{1}\right]-e_{3} e_{1} e_{3} \\
& \quad+\left[e_{1}, e_{2}\right] e_{2}+e_{2} e_{1} e_{2}+\left[e_{1}, e_{3}\right] e_{3}+e_{3} e_{1} e_{3} \\
& =2 e_{2} e_{3}-2 e_{3} e_{2}+2 e_{3} e_{2}-2 e_{2} e_{3} \\
& =0 .
\end{aligned}
$$

Thus the eigenspaces are invariant under the action of $s p(1)$.
Let's review now some classical theory. Let $f$ be a function on $S^{3}$ and $F$ an extension of $f$ to $\mathbb{R}^{4}$. Then

$$
\begin{equation*}
\Delta(f)=\Delta_{\mathbb{R}^{4}}(F)+3 \frac{\partial F}{\partial r}+\frac{\partial^{2} F}{\partial r^{2}} \tag{4.1}
\end{equation*}
$$

This decomposition is fantastically simple and allows for a complete description of the eigenvalues of $\Delta$. Let $H_{m}\left(\mathbb{R}^{4}\right)$ denote the set of harmonic homogeneous polynomials of degree $m$ and denote by $H_{m}\left(S^{3}\right)$ its restriction to $S^{3}$. It follows from Equation (4.1) that $H_{m}\left(S^{3}\right)$ consists of eigenvectors for the Laplacian $\Delta$ on function on $S^{3}$. The corresponding eigenvalues are $m(m+2)$. In fact, these are all the eigenvalues.
In fact, the eigenvectors of the Laplacian on $S^{n}$ are always the corresponding $H_{m}\left(S^{n}\right)$ and the eigenvalues are correspondingly the $m(m+n-1)$; see [GHL90, Cor 4.49].
Since they correspond to different eigenvalues, the spaces $H_{m}\left(S^{3}\right)$ are invariant under the action of $s p(1)$. So we reduce our study of eigenvalues of $D$ to the study of its eigenvalues on eigenspaces of $\Delta$. For those we have the beautiful decomposition theorem that follows.

Theorem 4.1-2. We have the following isomorphism of complex representation of the Lie algebra $s p(1)$ :

$$
H_{m}\left(S^{3} ; \mathbb{C}\right) \cong(m+1) \text { Sym }^{m} \mathbb{C}^{2}
$$

Proof: The left-invariant vector fields $e_{1}, e_{2}, e_{3}$ satisfy the same commuting relation as $i, j$, and $k$. Thus, we can view them in $s u(2)$ as

$$
\begin{aligned}
e_{1} & \equiv\left[\begin{array}{cc}
i & 0 \\
0 & -i
\end{array}\right], \\
e_{2} & \equiv\left[\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right], \text { and } \\
e_{3} & \equiv\left[\begin{array}{cc}
0 & -i \\
-i & 0
\end{array}\right] .
\end{aligned}
$$

To study the representation theory of $s u(2)$, it is convenient to use the standard basis $H, X$, and $Y$ of $s l_{2}$ since $s u(2)$ and $s l_{2}$ have the same irreducible representations. In terms of the $e_{a}$, we have

$$
\begin{aligned}
H & =-i e_{1} \\
X & =\frac{1}{2}\left(-e_{2}+i e_{3}\right), \text { and } \\
Y & =\frac{1}{2}\left(e_{2}+i e_{3}\right)
\end{aligned}
$$

We set $z_{1}=x_{1}+i x_{2}$, and $z_{2}=x_{3}+i x_{4}$. In these coordinates, we have

$$
\begin{aligned}
H & =-\frac{\partial}{\partial r}+2\left(z_{1} \frac{\partial}{\partial z_{1}}+\bar{z}_{2} \frac{\partial}{\partial \bar{z}_{2}}\right) \\
X & =\bar{z}_{2} \frac{\partial}{\partial \bar{z}_{1}}-z_{1} \frac{\partial}{\partial z_{2}} \\
Y & =-z_{2} \frac{\partial}{\partial z_{1}}+\bar{z}_{1} \frac{\partial}{\partial \bar{z}_{2}}, \text { and } \\
\Delta_{\mathbb{R}^{n}} & =4\left(\frac{\partial}{\partial z_{1}} \frac{\partial}{\partial \bar{z}_{1}}+\frac{\partial}{\partial z_{2}} \frac{\partial}{\partial \bar{z}_{2}}\right)
\end{aligned}
$$

Consider now the $m+1$ homogeneous polynomials

$$
p_{a}:=z_{1}^{a} \bar{z}_{2}^{m-a}, \quad a=0, \ldots, m
$$

They are obviously killed by $\Delta_{\mathbb{R}^{n}}$ and thus are in $H_{m}\left(S^{3}, \mathbb{C}\right)$.
Equally obvious is the fact that they are killed by $X$. Hence, each $p_{a}$ generates an irreducible submodule of $H_{m}\left(S^{3}\right)$. Since

$$
\begin{aligned}
H\left(p_{a}\right) & =-m p_{a}+2\left(a p_{a}+(m-a) p_{a}\right) \\
& =m p_{a}
\end{aligned}
$$

this module is isomorphic to $S y m^{m} \mathbb{C}^{2}$ as a representation of $s u(2)$.
We have so far establish the presence of $(m+1) S y m^{m} \mathbb{C}^{2}$ inside $H_{m}\left(S^{3} ; \mathbb{C}\right)$. Since both spaces have dimension $(m+1)^{2}$ (see ABR01, p. 78, Prop. 5.8]) they must be equal.

Now consider the isomorphism between $\mathbb{H}$ and $\mathbb{C}^{2}$ given by the natural decomposition $z_{1}+j z_{2}$. This isomorphism induces an isomorphism between $S p(1)$ and $S U(2)$ as follows:

$$
\begin{aligned}
S p(1) & \equiv S U(2) \\
i & \mapsto\left[\begin{array}{cc}
i & 0 \\
0 & -i
\end{array}\right] \\
j & \mapsto\left[\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right] \\
k & \mapsto\left[\begin{array}{cc}
0 & -i \\
-i & 0
\end{array}\right] .
\end{aligned}
$$

Basically, the operator $D$ restricts to a set of operators, one for every $m$ :

$$
\begin{aligned}
H_{m}\left(S^{3} ; \mathbb{C}\right)^{2} & \rightarrow H_{m}\left(S^{3} ; \mathbb{C}\right)^{2} \\
{\left[\begin{array}{l}
F_{1} \\
F_{2}
\end{array}\right] } & \mapsto\left[\begin{array}{cc}
i e_{1} & -e_{2}+i e_{3} \\
e_{2}+i e_{3} & -i e_{1}
\end{array}\right]\left[\begin{array}{l}
F_{1} \\
F_{2}
\end{array}\right]+\frac{3}{2}\left[\begin{array}{l}
F_{1} \\
F_{2}
\end{array}\right] .
\end{aligned}
$$

Following Hitchin, we let $Q$ denote the operator

$$
Q:=\left[\begin{array}{cc}
-H & 2 X \\
2 Y & H
\end{array}\right] .
$$

The operator $D$ is thus the sum $Q+3 / 2$ on the invariant subspace $\left(S y m^{m} \mathbb{C}^{2}\right)^{2}$. Let's find the eigenvalues of $Q$ on this subspace.
Let $x, y$ be the standard basis of $\mathbb{C}^{2}$. Then $S y m^{m} \mathbb{C}^{2}$ is the set of homogeneous polynomials of degree $m$ in $x$ and $y$. As such, there is an obvious basis $h_{a}:=x^{a} y^{m-a}$ for $S y m^{m} \mathbb{C}^{2}$ and we have

$$
\begin{aligned}
& H\left(h_{a}\right)=(2 a-m) h_{a}, \\
& X\left(h_{a}\right)=(m-a) h_{a+1}, \text { and } \\
& Y\left(h_{a}\right)=a h_{a-1} .
\end{aligned}
$$

The vectors

$$
\left[\begin{array}{c}
h_{0} \\
0
\end{array}\right] \text {, and }\left[\begin{array}{c}
0 \\
h_{m}
\end{array}\right]
$$

are eigenvectors of $Q$, of eigenvalues $m$. Each of them appear with multiplicity $m+1$ in the space $H_{m}\left(S^{3} ; \mathbb{C}\right)^{2}$. Consider now the vectors

$$
\left[\begin{array}{c}
(1+m-a) h_{a} \\
a h_{a-1}
\end{array}\right] \text { and }\left[\begin{array}{c}
h_{a} \\
-h_{a-1}
\end{array}\right] .
$$

These $2 m$ vectors, along with the two others, span $\left(S y m^{m} \mathbb{C}^{2}\right)^{2}$. Furthermore, they are eigenvectors of $Q$, of respective eigenvalues $m$ and $-m-2$. The $m+1$ diagonal $\left(S y m^{m} \mathbb{C}^{2}\right)^{2}$ factors in $H_{m}\left(S^{3} ; \mathbb{C}\right)^{2}$ span the whole space. Hence these eigenvalues appear with multiplicities $m(m+1)$. We just proved the following theorem.

Theorem 4.1-3. The eigenvalues of the Dirac operator on the spinor bundle of $S^{3}$ are the

$$
\pm(k+3 / 2), \text { for } k \in \mathbb{N}
$$

each $\pm(k+3 / 2)$ appearing with multiplicity $(k+1)(k+2)$.
This result is confirmed by a similar computation in Hit74, Prop 3.2] and by a different method of Andrzej Trautman in [Tra93] which apply to all spheres, and which we describe in Section 4.2

### 4.2 Trautman's construction

We now confirm the results of the previous two sections by a method of Trautman which apply to all spheres. This method appeared as a first paper [Tra93] in a projected series of paper of Trautman
with E. Winkowska on the spectrum of the Dirac operator on hypersurface. The promised sequel Spinors and the Dirac operator on hypersurfaces. II. The spheres as an example was apparently never completed and does not appear in the literature.
Let $S$ be the spinor bundle on $\mathbb{R}^{n+1}$. Let $i: S^{n}(r) \rightarrow \mathbb{R}^{n+1}$ be the inclusion. Then $i^{*}(S)$ is a Clifford bundle on $S^{n}(r)$. On this bundle, we have a spin connection, which gives us a Dirac operator $D_{r}$. Let $D$ be the Dirac operator on $S$. We now look at the relationship between $D$ and $D_{r}$.
Let $e_{1}, \ldots, e_{n}$ be an orthonormal frame on a patch of $S^{n}(1)$. We can extend this frame by radial parallel transport to a cone of $\mathbb{R}^{n+1} \backslash\{0\}$. Let $e_{0}=\nu$ be the radial vector field. Let $\mathbb{I}_{r}$ be the second fundamental form of $S(r)$. Then

$$
\begin{aligned}
D & =\rho\left(e_{0}\right) \nabla_{e_{0}}+\sum_{0<i \leq n} \rho\left(e_{i}\right) \nabla_{e_{i}} \\
& =\rho(\nu) \frac{\partial}{\partial r}+\sum_{0<i \leq n} \rho\left(e_{i}\right) e_{i}+\frac{1}{2} \rho\left(e_{i}\right) \sum_{0 \leq j<k \leq n} \omega_{k}^{j}\left(e_{i}\right) \rho\left(e_{k} e_{j}\right) \\
& =\rho(\nu) \frac{\partial}{\partial r}+\frac{1}{2} \sum_{0<i \leq n} \sum_{0<k \leq n} \omega_{k}^{0}\left(e_{i}\right) \rho\left(e_{i} e_{k}\right) \rho(\nu)+D_{r} \\
& =\rho(\nu) \frac{\partial}{\partial r}-\frac{1}{2} \operatorname{tr}\left(I_{r}\right) \rho(\nu)+D_{r} .
\end{aligned}
$$

A simple computation shows that $\operatorname{tr}\left(I_{r}\right)=-n / r$. Thus

$$
\begin{equation*}
D=D_{r}+\rho(\nu) \frac{\partial}{\partial r}+\frac{n}{2 r} \rho(\nu) . \tag{4.2}
\end{equation*}
$$

Let $p: \mathbb{R}^{n+1} \rightarrow S$ be a spinor-valued homogeneous harmonic polynomial of degree $l+1$. The polynomial $D p$ has degree $l$ and is killed by $D$. Consider then

$$
s_{ \pm}=\frac{(1 \mp \rho(\nu))}{2} D p
$$

We have $D s_{ \pm}=0$. Since $D p$ is homogeneous of degree $l$, we have $\partial s_{ \pm} / \partial r=(l / r) s_{ \pm}$. Thus

$$
\begin{aligned}
D_{r} s_{ \pm} & =-\rho(\nu) \frac{\partial s_{ \pm}}{\partial r}-\frac{n}{2 r} \rho(\nu) s_{ \pm} \\
& =\frac{(l+n / 2)}{r} \rho(-\nu) s_{ \pm} \\
& = \pm \frac{(l+n / 2)}{r} s_{ \pm} .
\end{aligned}
$$

For $n=3$, we see the spectrum described by Theorem 4.1-3
Now, $i^{*}(S)$ splits as a direct sum of irreducible spinor bundles. Because of dimensional reason and because $n$-spheres ( $n>1$ ) have only one spin structure, for $n=3$, the bundle $i^{*}(S)$ splits as two copies of the spinor bundle of $S^{3}$, and for $n=2$, the bundle $i^{*}(S)$ is the spinor bundle of $S^{2}$.
As $i^{*}(S)$ splits, the eigenspaces of $D_{r}$ split too. So we are building genuine eigenspaces for the Dirac operator on the spinor bundle of $S^{n}$.

## Chapter 5

## Decay of instantons

On a cylindrical manifold $\mathbb{R} \times Y$ with any warped product metric, the ASD equation for a connection in temporal gauge is

$$
\begin{equation*}
\partial_{t} A=-*_{3} F_{A}^{(3)}, \tag{5.1}
\end{equation*}
$$

On $\mathbb{R} \times T^{3}$ with coordinates $(t, \theta)$, we can expand any connection according to its Fourier modes:

$$
A(\theta, t)=\sum_{\nu \in \mathbb{Z}^{3}} A^{\nu}(t) e^{i \theta \cdot \nu}
$$

While the ASD equation mixes terms from different Fourier modes, the zero-mode behaves particularly nicely. Set

$$
H:=\left\{A \mid A^{\nu}=0 \text { for } \nu \neq 0\right\} .
$$

For the product metric, Equation (5.1) is autonomous, and it turns out that $\mathcal{G} H$ is then a center manifold for the flow of that equation. Hence, every flow line with finite energy approaches exponentially a flow line in $\mathcal{G} H$.
Since finite energy correspond here to

$$
\int_{[1, \infty) \times T^{3}}\left|F_{A}\right|^{2}<\infty,
$$

it remains to study the flow lines in the finite dimensional space $H$ and in order to understand the decay of instantons.
This material is well known to [MMR94], where it is proved that every instanton converges to a flat connection, and the decay to that instanton is exponential if the flat limit is irreducible.
For the warped metric on $[1, \infty) \times T^{3}$ coming from $T^{2} \times \mathbb{R}^{2}$ by polar coordinate on the $\mathbb{R}^{2}$ factor, the Flow Equation (5.1) is not autonomous and consequently we cannot use the traditional center manifold theorem. It is worthwhile however to study the behavior of flow lines on $\mathcal{G H}$ as well.
Any element of $H$ can be expressed as

$$
A=a_{1} d \theta^{1}+a_{2} d \theta^{2}+a_{3} d \theta^{3},
$$

with $a_{j} \in \mathfrak{s u}(n)$.

The Flow Equation (5.1), once written with the $a_{j}$, gives rise to the equations

$$
\left.\begin{array}{rl}
a_{1}^{\prime} & =\left[a_{2}, a_{3}\right]  \tag{5.2}\\
a_{2}^{\prime} & =\left[a_{3}, a_{1}\right] \\
a_{3}^{\prime} & =\left[a_{1}, a_{2}\right]
\end{array}\right\} \text { for the product metric } \mathbb{R} \times T^{3}
$$

and

$$
\left.\begin{array}{l}
a_{1}^{\prime}=t^{-1}\left[a_{2}, a_{3}\right]  \tag{5.3}\\
a_{2}^{\prime}=t^{-1}\left[a_{3}, a_{1}\right] \\
a_{3}^{\prime}=t \quad\left[a_{1}, a_{2}\right]
\end{array}\right\} \text { for the warped metric } T^{2} \times \mathbb{R}^{2}
$$

These equations are quite symmetrical, and we can reduce the study of those systems to the study of their invariants.

### 5.1 On $\mathbb{R} \times T^{3}$

Let's restrict our attention to $\mathfrak{s u}(2)$, and let's first deal with System (5.2). Define the following real valued functions:

$$
\begin{align*}
f & :=\left|a_{1}^{\prime}\right|^{2}+\left|a_{2}^{\prime}\right|^{2}+\left|a_{3}^{\prime}\right|^{2}, \\
g & :=\left|a_{1}\right|^{2}+\left|a_{2}\right|^{2}+\left|a_{3}\right|^{2},  \tag{5.4}\\
d & :=\left\langle a_{1},\left[a_{2}, a_{3}\right]\right\rangle .
\end{align*}
$$

Lemma 5.1-1. For $a_{1}, a_{2}$, $a_{3}$ flowing according to System (5.2), and for $f, g$, $d$ defined by Equations (5.4), we have

$$
\begin{aligned}
f^{\prime} & =16 d g, \\
g^{\prime} & =6 d, \\
d^{\prime} & =f .
\end{aligned}
$$

Proof: We equip $\mathrm{SU}(2)$ with its bi-invariant metric that gives it the riemannian structure of $S^{3}$. For elements $X, Y, Z, W$ of the Lie algebra, the covariant derivative and Riemann tensor have the simple expressions

$$
\begin{gathered}
\nabla_{X} Y=\frac{1}{2}[X, Y], \text { and } \\
R(X, Y, Z, W)=\frac{1}{4}\langle[X, Y],[Z, W]\rangle ;
\end{gathered}
$$

see for example [GHL90 3.17].
Since the sectional curvature is 1 everywhere, we have

$$
R(X, Y, Z, W)=\langle X, Z\rangle\langle Y, W\rangle-\langle Y, Z\rangle\langle X, W\rangle
$$

Notice, then, that for any permutation $i j k$ of 123 , we have

$$
\left\langle a_{i}, a_{k}^{\prime}\right\rangle= \pm\left\langle a_{i},\left[a_{i}, a_{j}\right]\right\rangle= \pm 2\left\langle a_{i}, \nabla_{a_{i}} a_{j}\right\rangle=\mp 2\left\langle\nabla_{a_{i}} a_{i}, a_{j}\right\rangle=\mp\left\langle\left[a_{i}, a_{i}\right], a_{j}\right\rangle=0,
$$

while $d=\left\langle a_{1}, a_{1}^{\prime}\right\rangle=\left\langle a_{2}, a_{2}^{\prime}\right\rangle=\left\langle a_{3}, a_{3}^{\prime}\right\rangle$. It is then quite obvious that $g^{\prime}=6 d$.

We have

$$
\begin{aligned}
\left(\left|a_{1}^{\prime}\right|^{2}\right)^{\prime} & =2\left\langle\left[a_{2}, a_{3}\right],\left[a_{2}^{\prime}, a_{3}\right]\right\rangle+2\left\langle\left[a_{2}, a_{3}\right],\left[a_{2}, a_{3}^{\prime}\right]\right\rangle \\
& =8 R\left(a_{2}, a_{3}, a_{2}^{\prime}, a_{3}\right)+8 R\left(a_{2}, a_{3}, a_{2}, a_{3}^{\prime}\right) \\
& =8\left\langle a_{2}, a_{2}^{\prime}\right\rangle\left|a_{3}\right|^{2}+8\left|a_{2}\right|^{2}\left\langle a_{3}, a_{3}^{\prime}\right\rangle \\
& =8 d\left(\left|a_{3}\right|^{2}+\left|a_{2}\right|^{2}\right) .
\end{aligned}
$$

This equation together with similar equations for $\left(\left|a_{2}^{\prime}\right|^{2}\right)^{\prime}$ and $\left(\left|a_{3}^{\prime}\right|^{2}\right)^{\prime}$ yield $f^{\prime}=16 d g$.
As for $d$, we have

$$
\begin{aligned}
d^{\prime} & =\left\langle a_{1}^{\prime},\left[a_{2}, a_{3}\right]\right\rangle+\left\langle a_{1},\left[a_{2}^{\prime}, a_{3}\right]\right\rangle+\left\langle a_{1},\left[a_{2}, a_{3}^{\prime}\right]\right\rangle \\
& =\left|a_{1}^{\prime}\right|^{2}-\left\langle\left[a_{1}, a_{3}\right], a_{2}^{\prime}\right\rangle-\left\langle\left[a_{2}, a_{1}\right], a_{3}^{\prime}\right\rangle \\
& =\left|a_{1}^{\prime}\right|^{2}+\left|a_{2}^{\prime}\right|^{2}+\left|a_{3}^{\prime}\right|^{2}=f .
\end{aligned}
$$

The proof is now complete.
We are now ready to prove decay properties of instantons.
Theorem 5.1-2 (Decay of instantons). Let $A$ be an instanton on $\mathbb{R} \times T^{3}$. Then

$$
\left|F_{A}\right|=o\left(t^{-1}\right)
$$

as $t \rightarrow \infty$.
Proof: As we mentioned in the introduction to this chapter, we only have to study the flow lines in $H$ as any other is exponentially decaying to a flow line in $\mathcal{G} H$.
Notice first the $\left\|F_{A}\right\|^{2}=\int f=\int d^{\prime}$ hence $\lim d$ exist as $t \rightarrow \infty$. Suppose $\lim d=2 l>0$. For some $T$ and $t>T$, we have $d>l$. Hence $g^{\prime}>l$, or once we integrate, $g(t)>l t+C$. Thus $f^{\prime}=16 d g>16 l^{2} t+C$ and $\lim f^{\prime}=\infty$. But then surely $\lim f=\infty$ and $f$ cannot be integrable, which clearly contradicts the finite energy condition. Hence we proved

$$
\lim d \leq 0 .
$$

In fact, as $d^{\prime}=f \geq 0$, we have

$$
d \leq 0 \text { always. }
$$

Consequently, $f^{\prime}=16 d g \leq 0$. Thus $f$ must have a finite limit since $f \geq 0$. Since $\int f$ converges, we have

$$
\lim f=0
$$

As $g \geq 0$ and $g^{\prime}=6 d \leq 0$, we have

$$
\lim g \text { exists. }
$$

Then we judiciously apply l'Hospital's rule, denoted HR below. Since

$$
(\lim g)=\lim \frac{g / t}{1 / t} \stackrel{\mathrm{HR}}{=} \lim \frac{-g / t^{2}+6 d / t}{-1 / t^{2}}=(\lim g)-6 \lim t d
$$

we have

$$
d=o\left(t^{-1}\right)
$$

But then,

$$
0=\lim t d=\lim \frac{d}{1 / t} \stackrel{\mathrm{HR}}{=} \lim \frac{f}{-1 / t^{2}}=-\lim t^{2} f
$$

Since $f=\left|F_{A}\right|^{2}$, the conclusion follows.
We can actually pull out more decay from those equations, even exponential decay in [MMR94]. Here we prove polynomial decay to infinite order for non-zero limits.

Theorem 5.1-3 (Extra decay for non-zero limits). Let $A$ be an instanton on $\mathbb{R} \times T^{3}$. Suppose that the flat connection to which $A$ converges at infinity is not gauge equivalent to $0 \in H$. Then for all $k$,

$$
\left|F_{A}\right|=o\left(t^{-k}\right)
$$

as $t \rightarrow \infty$.

Proof: We work in $H$. We already proved in Theorem 5.1-2 and its proof that $f=o\left(t^{-2}\right)$ and $d=o\left(t^{-1}\right)$.
Once we suppose $d=o\left(t^{-k}\right)$, we find $0=\lim d / t^{-k}=-k \lim f / t^{-k-1}$ using l'Hospital's rule. Hence

$$
\begin{equation*}
d=o\left(t^{-k}\right) \text { implies } f=o\left(t^{-k-1}\right) \tag{5.5}
\end{equation*}
$$

Suppose now $\lim g \neq 0$, and $f=o\left(t^{-k}\right)$. Then

$$
0=\lim \frac{f}{t^{-k}}=-\frac{1}{k} \lim \frac{d g}{t^{-k-1}} \stackrel{\mathrm{HR}}{=}-\frac{1}{k}(\lim g) \lim \frac{d}{t^{-k-1}}
$$

Hence

$$
\begin{equation*}
f=o\left(t^{-k}\right) \text { implies } d=o\left(t^{-k-1}\right) \tag{5.6}
\end{equation*}
$$

under the condition $\lim g \neq 0$.
The conclusion follows by pumping up Equations (5.5) and (5.6).

### 5.2 On $T^{2} \times \mathbb{R}^{2}$

We keep our attention on $\mathfrak{s u}(2)$, and deal now with System (5.3). Define the following real valued functions:

$$
\begin{align*}
f & :=\left|a_{1}^{\prime}\right|^{2}+\left|a_{2}^{\prime}\right|^{2}+\frac{1}{t^{2}}\left|a_{3}^{\prime}\right|^{2} \\
u & :=\frac{1}{t^{2}}\left|a_{3}^{\prime}\right|^{2} \\
g_{1} & :=\left|a_{1}\right|^{2}+\left|a_{2}\right|^{2}  \tag{5.7}\\
g_{2} & :=\left|a_{3}\right|^{2} \\
d & :=\left\langle a_{1},\left[a_{2}, a_{3}\right]\right\rangle
\end{align*}
$$

Lemma 5.2-1. For $a_{1}, a_{2}, a_{3}$ flowing according to System (5.3), and for $f, u, g_{1}, g_{2}, d$ defined by

Equations (5.7), we have

$$
\begin{array}{rlrl}
g_{1}^{\prime} & =\frac{4}{t} d, & u^{\prime} & =\frac{8}{t} g_{1} d, \\
g_{2}^{\prime} & =2 t d, & d^{\prime}=t f
\end{array}
$$

Proof: We proceed as in the proof of Lemma 5.1-1
Using the Leibnitz rule, we find

$$
\begin{aligned}
d^{\prime} & =\left\langle a_{1}^{\prime},\left[a_{2}, a_{3}\right]\right\rangle+\left\langle a_{1},\left[a_{2}^{\prime}, a_{3}\right]\right\rangle+\left\langle a_{1},\left[a_{2}, a_{3}^{\prime}\right]\right\rangle \\
& =t\left|a_{1}^{\prime}\right|^{2}-\left\langle\left[a_{1}, a_{3}\right], a_{2}^{\prime}\right\rangle-\left\langle\left[a_{2}, a_{1}\right], a_{3}^{\prime}\right\rangle \\
& =t\left|a_{1}^{\prime}\right|^{2}+t\left|a_{2}^{\prime}\right|^{2}+\frac{1}{t}\left|a_{3}^{\prime}\right|^{2},
\end{aligned}
$$

hence proving $d^{\prime}=t f$.
While $g_{2}^{\prime}=2\left\langle a_{3}, a_{3}^{\prime}\right\rangle=2 t d$, we have

$$
\begin{aligned}
g_{1}^{\prime} & =2\left\langle a_{1}, a_{1}^{\prime}\right\rangle+2\left\langle a_{2}, a_{2}^{\prime}\right\rangle \\
& =\frac{2}{t}\left(\left\langle a_{1},\left[a_{2}, a_{3}\right]\right\rangle+\left\langle a_{2},\left[a_{3}, a_{1}\right]\right\rangle\right) \\
& =\frac{4}{t} d,
\end{aligned}
$$

thus proving the differential equations for $g_{1}$ and $g_{2}$.
We have

$$
\begin{aligned}
u^{\prime} & =\left(\left|\left[a_{1}, a_{2}\right]\right|^{2}\right)^{\prime} \\
& =2\left\langle\left[a_{1}, a_{2}\right],\left[a_{1}^{\prime}, a_{2}\right]\right\rangle+2\left\langle\left[a_{1}, a_{2}\right],\left[a_{1}, a_{2}^{\prime}\right]\right\rangle \\
& =\frac{8}{t}\left(R\left(a_{1}, a_{2},\left[a_{2}, a_{3}\right], a_{2}\right)+R\left(a_{1}, a_{2}, a_{1},\left[a_{3}, a_{1}\right]\right)\right) \\
& =\frac{8}{t}\left(\left\langle a_{1},\left[a_{2}, a_{3}\right]\right\rangle\left|a_{2}\right|^{2}-\left\langle a_{2},\left[a_{2}, a_{3}\right]\right\rangle\left\langle a_{1}, a_{2}\right\rangle+\left|a_{1}\right|^{2}\left\langle a_{2},\left[a_{3}, a_{1}\right]\right\rangle-\left\langle a_{2}, a_{1}\right\rangle\left\langle a_{1} \cdot\left[a_{3}, a_{1}\right]\right\rangle\right) \\
& =\frac{8}{t} g_{1} d .
\end{aligned}
$$

The proof is now complete.

Theorem 5.2-2. In the gauge group translates of the zero Fourier mode on $T^{2} \times \mathbb{R}$, finite energy instantons have quadratically decaying curvature.

Proof: We of course aim to prove that $f=o\left(t^{-4}\right)$.
Notice first the $\left\|F_{A}\right\|^{2}=\int t f=\int d^{\prime}$ hence $\lim d$ exist as $t \rightarrow \infty$.
Suppose $\lim d=2 l>0$. For some $T$ and $t>T$, we have $d>l$. Hence $g_{1}^{\prime}=4 t^{-1} d>4 t^{-1} l$, or once we integrate, $g_{1}(t)-g_{1}(T)>4 l \log (t / T)$, or $g_{1}(t)>4 l \log (\gamma t)$ for some $\gamma>0$.

Then $u^{\prime}=8 t^{-1} g_{1} d>32 l^{2} t^{-1} \log (\gamma t)$. Since $\left(\log ^{2}(\gamma t)\right)^{\prime}=t^{-1} \log (\gamma t)$, we have

$$
\begin{aligned}
u(t)-u(T) & =\int_{T}^{t} u^{\prime} \\
& >32 l^{2} \int_{T}^{t}\left(\log ^{2}(\gamma t)\right)^{\prime} \\
& =32 l^{2}\left(\log ^{2}(\gamma t)-\log ^{2}(\gamma T)\right),
\end{aligned}
$$

or for some constant $\epsilon$, we get $u(t)>32 l^{2} \log ^{2}(\gamma t)+\epsilon$. Hence $\lim u=\infty$, and then surely $\lim f=\infty$ and $f$ cannot be integrable, which clearly contradicts the finite energy condition. Hence we proved

$$
\lim d \leq 0 .
$$

In fact, as $d^{\prime}=t f \geq 0$, we have

$$
d \leq 0 \text { always. }
$$

Since $g_{a}^{\prime}=2 t d \leq 0$ and $g_{2} \geq 0$, the quantity $\lim g_{2}$ exist and is finite.
Then we judiciously apply l'Hospital's rule. Since

$$
\left(\lim g_{2}\right)=\lim \frac{g_{2} / t}{1 / t} \stackrel{\mathrm{HR}}{=} \lim \frac{-g_{2} / t^{2}+g_{2}^{\prime} / t^{\prime}}{-1 / t^{2}}=(\lim g)-2 \lim t^{2} d
$$

we have

$$
d=o\left(t^{-2}\right) .
$$

But then,

$$
0=\lim t^{2} d=\lim \frac{d}{1 / t^{2}} \stackrel{\mathrm{HR}}{=} \lim \frac{t f}{-2 / t^{3}}=-\frac{1}{2} \lim t^{4} f .
$$

Since $f=\left|F_{A}\right|^{2}$, the conclusion follows.
Theorem $[5.2-2$ is perhaps a first step in a dynamical system approach to proving a conjecture of Jardim that every doubly-periodic instanton connection of finite action $\left\|F_{A}\right\|_{L^{2}}<\infty$ has quadratic curvature decay. This conjecture appeared in [Jar02a] p. 433] and is supported in part by Nahm transform considerations in [BJ01].

### 5.3 Notes on different quotients

We proved in Section 5.1] that given an instanton $A$ on $\mathbb{R} \times T^{3}$, its curvature $F_{A}$ decays like $o\left(r^{-1}\right)$. It was first observed by Mrowka that there are instantons who decay like $r^{-1}$. For example, let $\mathbf{i}, \mathbf{j}, \mathbf{k}$ be the usual basis of $\mathfrak{s u}(2)$, and consider

$$
A=\frac{\mathbf{i} d x+\mathbf{j} d y+\mathbf{k} d z}{2 r}
$$

The curvature of that connection is

$$
F_{A}=\frac{-\mathbf{i} d r \wedge d x-\mathbf{j} d r \wedge d y-\mathbf{k} d r \wedge d z+\mathbf{k} d x \wedge d y+\mathbf{j} d z \wedge d x+\mathbf{i} d y \wedge d z}{2 r^{2}},
$$

which is quite stronger than $o\left(r^{-1}\right)$.

On $\mathbb{R}^{2} \times T^{2}$, Jardim gave the example of the connection

$$
A=\left[\begin{array}{cc}
-i & 0 \\
0 & i
\end{array}\right] \frac{d \theta}{\log r}+\frac{1}{r \log r}\left[\begin{array}{cc}
0 & e^{-i \theta}(d x-i d y) \\
-e^{i \theta}(d x+i d y) & 0
\end{array}\right]
$$

with

$$
F_{A}=O\left(\frac{1}{r^{2} \log r}\right)
$$

which again is a bit stronger than the conjectured $O\left(r^{-2}\right)$ of Section 5.2
For the classical $\mathbb{R}^{4}$ case, it was proved originally by Uhlenbeck in Uhl82, Cor. 4.2] that the condition $\left\|F_{A}\right\|_{L^{2}}<\infty$ implies that $\left|F_{A}\right|=O\left(r^{-4}\right)$; see also the appendix of the seminal work of Donaldson [Don83]. This decay is achieved by the connection

$$
A=\frac{1}{1+r^{2}}\left(\theta_{1} \mathbf{i}+\theta_{2} \mathbf{j}+\theta_{3} \mathbf{k}\right)
$$

with

$$
\begin{aligned}
\theta_{1} & =x_{1} d x^{2}-x_{2} d x^{1}-x_{3} d x^{4}+x_{4} d x^{3} \\
\theta_{2} & =x_{1} d x^{3}-x_{3} d x^{1}-x_{4} d x^{2}+x_{2} d x^{4} \\
\theta_{3} & =x_{1} d x^{4}-x_{4} d x^{1}-x_{2} d x^{3}+x_{3} d x^{2} .
\end{aligned}
$$

While this connection in this particular gauge is $O\left(r^{-1}\right)$, its curvature

$$
F_{A}=\frac{1}{1+r^{4}}\left(d \theta_{1} \mathbf{i}+d \theta_{2} \mathbf{j}+d \theta_{3} \mathbf{k}\right)
$$

is exactly of order $1 / r^{4}$.
As for $S^{1} \times \mathbb{R}^{3}$, by taking a monopole $(A, \Phi)$ on $\mathbb{R}^{3}$, we get an example of an instanton $\Phi d r+A$ whose curvature $\left(\nabla_{A} \Phi\right) \wedge d r+F_{A}$ is exactly of order $1 / r^{2}$.

## Chapter 6

## Fredholm theory on $\mathbb{R} \times Y$

Let $Y$ be a three-dimensional compact orientable manifold. Let $(E, A)$ be a $S U(2)$-instanton over $\mathbb{R} \times Y$. We suppose for this section that $A$ is in temporal gauge, that is it has no $d t$ term. This assumption allows us to consider the restriction $A(t)$ to a cross-section $\{t\} \times Y$. The Dirac operator on $\{t\} \times Y$ is denoted $D_{A}$. We consider the Dirac operator

$$
\mathscr{D}_{A}=\partial_{t}+D_{A}
$$

on sections of $S^{+} \otimes E$.
Our aim in this chapter is to find spaces on which $\mathscr{D}_{A}$ is a Fredholm operator, and on those spaces compute its index. It is quite natural for such problems to consider Sobolev spaces, as in the compact case. While it is quite natural, it is perhaps too restrictive and what is happening on cylindrical manifolds in terms of Fredholmness is better understood in the realm of weighted Sobolev spaces.
It was observed long ago that for the usual Laplacian on $\mathbb{R}^{n}$, the classical Sobolev spaces are the wrong spaces for domains: the Laplacian is not Fredholm for those domains. The same is true for other elliptic partial differential operators on $\mathbb{R}^{n}$.
As an attempt to remedy the situation, Homer Walker in Wal71 Wal72, Wal73 proves that for certain domains, first order elliptic partial differential operators obtained from constant coefficients operators by adding a perturbation on a compact set are practically Fredholm, in the sense that the dimension of the kernel is finite and that the range can be described by a finite number of orthogonality condition.
In [NW73], the results are extended to a broader class of elliptic operators, perturbed in a less restrictive way, and $L^{p}$-type domains replace the $L^{2}$-type domains presents in the papers just described. We are presented with a sort of Gårding inequality decorated with weights, the treatment of which is not fully systematized at this point, and finite dimensionality of the kernel is proved.
The use of weights to describe which type of behavior is allowed at infinity was systematized in [Can75] with the introduction of weighted Sobolev spaces. Around the same time, Atiyah, Patodi and Singer in [APS75] made the crucial observation that the condition for an operator to be Fredholm in $L^{2}$ on a cylindrical manifold is that the restriction to the slice at infinity must have an empty kernel.
In [Loc81], and independently in [McO79], we are given specific APS-like conditions on the weights for an elliptic partial differential operator of any order with a certain type of asymptotic behavior to be Fredholm on the given weighted Sobolev spaces. The result and proof of that paper are extensions of [NW73], and partial results along these lines can be found in [Can75], where
isomorphism properties of the Laplacian were derived.
Choquet-Bruhat and Christodoulou in [CBC81] remove restriction on $p$ from another work of Cantor [Can79] and prove semi-Fredholmness, finite dimensionality of kernel and isomorphism theorems for operators on non-compact manifolds while giving improvements on imbedding and multiplication results for weighted Sobolev spaces. Those two papers constitute partial results along the lines of the more advanced and complete joint work [LM83, LM84] of Lockhart and McOwen. Their work extends the results of [Loc81] and [McO79] for systems of partial differential operators which are elliptic in the sense of Douglis-Nirenberg, and similar conditions on the weights are described to ensure Fredholmness.
Very similarly to what we do in this chapter, [LM85] study a much larger class of elliptic operator $C^{\infty}(E) \rightarrow C^{\infty}(f)$ for bundles $E$ and $F$ over a manifold with cylindrical ends. Conditions on weights to obtain Fredholmness and wall crossing formulas are derived. The paper also treats boundary valued problems with Lopatinski-Shapiro boundary conditions; see [APS75] along those lines.
As a prelude to proving that the mass of an asymptotically flat manifold is a geometric invariant, [Bar86] reviews the theory of weighted Sobolev spaces with an emphasis on two basic ideas which underlie the subject: the use of scaling arguments to pass from local estimates to global estimates and the derivation of sharp estimate from explicit formulas for Greens functions. Bartnik's paper add a number of technical improvements and some new observations to the theory. A simple example is that the indexing chosen for the weights is different from the one used by his predecessors, but it clearly reflects the growth at infinity allowed. An expanded version of his presentation, with complete proofs, can be found in Appendix E ,
These results can be put in a geometric form following Melrose; see [Mel93]. Melrose's technique involves adding a boundary at infinity to the underlying non-compact complete riemannian manifold. Later work of Mazzeo-Melrose [MM98] was used by Singer-Nye in [NS00] for computing the index of the Dirac operator twisted by a caloron on $S^{1} \times \mathbb{R}^{3}$.

### 6.1 Fredholmness

Suppose first that $A$ is independent of $t$.
The space $W^{k, p}(X, F)$ is the space of $L^{p}$ sections of the vector bundle $F$ over $X$, whose derivatives up to order $k$ for a given reference connection are also in $L^{p}$ on $X$. Because $\mathscr{D}_{A}$ is obviously defined on sections of $S^{+} \otimes E$ and $\mathscr{P}_{A}^{*}$ is obviously defined on sections of $S^{-} \otimes E$, we lighten up the notation by omiting the $F$, and most of time we omit the $X$ as well, in which case it is assumed to be $\mathbb{R} \times Y$, or $\mathbb{R} \times T^{3}$ when appropriate.

Lemma 6.1-1. Suppose the connection $A$ does not depend of $t$. Then

$$
\mathscr{D}_{A}: W^{1,2} \rightarrow L^{2}
$$

is Fredholm if and only if $0 \notin \operatorname{Spec}\left(D_{A}\right)$.
We actually prove a stronger version of the lemma, as this more powerful version is useful later on.
Lemma 6.1-2. Let $T$ be a formally self-adjoint elliptic 1 st order operator on a compact manifold $Y$. The operator $\hat{T}:=\partial_{t}+T$ seen as

$$
\hat{T}: W^{1,2} \rightarrow L^{2}
$$

is Fredholm if and only if $0 \notin \operatorname{Spec}(T)$. Furthermore, it is an isomorphism when Fredholm.

Proof: Suppose $0 \notin \operatorname{Spec}(T)$. Then we have the estimate $\|T \phi\|_{L^{2}(Y)} \geq C\|\phi\|_{W^{1,2}(Y)}$. Suppose $\phi$ is compactly supported. Then

$$
\begin{align*}
\|\hat{T} \phi\|^{2} & =\left\|\partial_{t} \phi\right\|^{2}+\|T \phi\|^{2}+2 \int_{\mathbb{R}} \partial_{t}\langle\phi, T \phi\rangle_{L^{2}(Y)} \\
& \geq\left(1+C^{2}\right)\left(\left\|\partial_{t} \phi\right\|^{2}+\| \| \phi\left\|_{W^{1,2}(Y)}\right\|_{L^{2}(\mathbb{R})}^{2}\right) \\
& \geq\left(1+C^{2}\right)\|\phi\|_{W^{1,2}}^{2} \tag{6.1}
\end{align*}
$$

thus $\hat{T}$ has closed range and no kernel.
If $\phi \in L^{2}$ is orthogonal to the image of $\hat{T}$, then it is a weak solution to $\hat{T} \phi=0$. Elliptic theory, for example in [LM89, Thm III.5.2(i), p. 193], implies that $\phi$ is $C^{\infty}$. But the only $C^{\infty}$ solution to be $L^{2}$ is 0 . Thus $\hat{T}$ is an isomorphism.
Suppose now $0 \in \operatorname{Spec}(T)$, and let $\phi_{0}$ be in $\operatorname{ker}(T)$. We show that $\hat{T}$ does not have closed range. It suffices to do so to prove that $\partial_{t}: W^{1,2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$ does not have closed image. Set

$$
\begin{aligned}
& f(x):= \begin{cases}1 / x, & \text { for }|x| \geq 1 \\
x, & \text { for }|x| \leq 1\end{cases} \\
& F(x):= \begin{cases}1 / 2-\log |x|, & \text { for }|x| \geq 1 \\
x^{2} / 2, & \text { for }|x| \leq 1\end{cases}
\end{aligned}
$$

The function $f$ clearly belongs to $L^{2}$. In fact, $\|f\|_{L^{2}}=2 \sqrt{2 / 3}$. We have $\partial_{t} F=f$, but $c+F \notin L^{2}$ for all $c \in \mathbb{R}$. Hence $f$ is not in the image of $\partial_{t}$.
Choose $\chi_{R}: \mathbb{R} \rightarrow[0, \infty)$ with $\chi_{R}(x)=\chi_{R}(-x)$, and

$$
\chi_{R}(x)= \begin{cases}0, & \text { when }|x| \geq 2 R \\ 1, & \text { when }|x| \leq R\end{cases}
$$

Set $f_{R}:=\chi_{R} f$. It is obvious that $f_{R} \rightarrow f$ in $L^{2}$, and $f_{R}(x)=-f_{R}(-x)$. This last property ensures that the function

$$
F_{R}:= \begin{cases}\int_{-3 R}^{x} f_{R}, & \text { for } x \leq 0 \\ -\int_{x}^{3 R} f_{R}, & \text { for } x \geq 0\end{cases}
$$

is well-defined at $x=0$. The function $F_{R}$ satisfies $\partial_{t} F_{R}=f_{R}$, and thus, since $F_{R}$ is compactly supported, $f_{R}$ is in the image of $\partial_{t}$. The image is therefore not closed. The proof is now complete.

We now add a number of weighted Sobolev spaces to our arsenal. The weight function we use here, denoted $\sigma_{\delta}$, depends only on $t$, and its definition depends on whether $\delta \in \mathbb{R}$ or $\delta \in \mathbb{R}^{2}$. For $\delta=\left(\delta_{-}, \delta_{+}\right) \in \mathbb{R}^{2}$, we want $\sigma_{\delta}>0$ with

$$
\sigma_{\delta}= \begin{cases}e^{-\delta_{-} t}, & \text { for } t<-1 \\ e^{-\delta_{+} t}, & \text { for } t>1\end{cases}
$$

To achieve such a weight function, choose $c$ smooth and positive with

$$
c(t)= \begin{cases}1, & \text { for } t \leq-1 \\ 0, & \text { for } t \geq 1\end{cases}
$$

Then set

$$
\sigma_{\delta}:=e^{-\left(c \delta_{-}+(1-c) \delta_{+}\right) t} .
$$

For $\delta \in \mathbb{R}$, set

$$
\sigma_{\delta}:=e^{-\delta t}
$$

The weighted Sobolev spaces are defined by the equation

$$
W_{\delta}^{k, p}:=\left\{f \mid\left\|\sigma_{\delta} f\right\|_{W^{k, p}}<\infty\right\}=\sigma_{-\delta} W^{k, p} .
$$

As usual, $L_{\delta}^{p}=W_{\delta}^{0, p}$. For $\delta \in \mathbb{R}$, notice that $\sigma_{\delta}=\sigma_{(\delta, \delta)}$ hence $W_{\delta}^{k, p}=W_{(\delta, \delta)}^{k, p}$.
Theorem 6.1-3. Suppose $A$ is translation invariant (it does not depend on $t$ ). Then

$$
\mathscr{D}_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}
$$

is Fredholm if and only if $\delta \notin \operatorname{Spec}\left(D_{A}\right)$. Moreover, it is an isomorphism if Fredholm.
Proof: The following diagram is commutative:


Because the columns are isomorphisms, the top row is Fredholm if and only if the bottom row is. But

$$
\sigma_{\delta} \mathscr{D}_{A} \sigma_{\delta}^{-1}=\partial_{t}+\left(\mathscr{D}_{A}+\delta\right) .
$$

Using Lemma 6.1-2 we see it is Fredholm if and only if $0 \notin \operatorname{Spec}\left(\mathscr{D}_{A}+\delta\right)$, or equivalently when $-\delta \notin \operatorname{Spec}\left(\mathscr{D}_{A}\right)$. Since $\operatorname{Spec}\left(\mathscr{D}_{A}\right)=-\operatorname{Spec}\left(\mathscr{D}_{A}\right)$, the conclusion follows.
Our ultimate goal is to find Fredholmness conditions for $\mathscr{D}_{A}$, with the only hypothesis that $A$ is an instanton. As we know, being an instanton forces $A$ to have flat limits at $\pm \infty$, say $\Gamma_{-}$and $\Gamma_{+}$. As a notational convenience, we define the grid

$$
\begin{equation*}
\mathfrak{G}_{A}:=\left(\operatorname{Sppc}\left(\Gamma_{-}\right) \times \mathbb{R}\right) \cup\left(\mathbb{R} \times \operatorname{Sppc}\left(\Gamma_{+}\right)\right) . \tag{6.2}
\end{equation*}
$$

Theorem 6.1-4. Let $\Gamma_{+}$and $\Gamma_{-}$be two flat connections on $Y$. Suppose $A$ is a connection on $\mathbb{R} \times Y$ such that

$$
A= \begin{cases}\Gamma_{-}, & \text {on }(-\infty,-R) \times Y \\ \Gamma_{+}, & \text {on }(R, \infty) \times Y\end{cases}
$$

Then for a weight $\delta \in \mathbb{R}^{2}$,

$$
\mathscr{D}_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}
$$

is Fredholm if and only if $\delta \notin \mathfrak{G}_{A}$.

Proof: Consider the three following manifolds:

$$
\begin{aligned}
& X_{1}=\mathbb{R} \times Y, \\
& X_{2}=\left([-R-2, R+2] /_{(R+2) \sim(-R-2)}\right) \times Y, \\
& X_{3}=\mathbb{R} \times Y .
\end{aligned}
$$

Using a path from $\Gamma_{-}$to $\Gamma_{+}$, we can find $\tilde{A}$ on $X_{2}$ such that $\tilde{A}=A$ on $[-R-1, R+1] \times Y$, and we can also find a function $\tilde{\sigma}$ defined on $X_{2}$ which restrict to $\sigma_{\delta}$ on that same subspace.
Choose a square root of a partition of unity

$$
\phi_{1}^{2}+\phi_{2}^{2}+\phi_{3}^{2}=1
$$

subordinate to the covering

$$
((-\infty,-R) \times Y,(-R-1, R+1) \times Y,(R, \infty) \times Y)
$$

Consider the operators

$$
\begin{aligned}
& D_{1}:=\sigma_{\delta_{-}} \mathscr{D}_{\Gamma_{-}} \sigma_{\delta_{-}}^{-1}, \\
& D_{2}:=\tilde{\sigma} \mathscr{D}_{\tilde{A}} \tilde{\sigma}^{-1}, \text { and } \\
& D_{3}:=\sigma_{\delta_{+}} \mathscr{D}_{\Gamma_{+}} \sigma_{\delta_{+}}^{-1}
\end{aligned}
$$

defined on the spaces $X_{1}, X_{2}$, and $X_{3}$ respectively.
When $\delta \notin \mathfrak{G}_{A}$, all the $D_{i}$ are Fredholm. In fact, $D_{1}$ and $D_{3}$ are even isomorphisms. Hence there exist

$$
\begin{aligned}
P_{i}: L^{2}\left(X_{i}\right) & \rightarrow W^{1,2}\left(X_{i}\right), i=1,2,3, \text { and } \\
K_{2}: L^{2}\left(X_{2}\right) & \rightarrow L^{2}\left(X_{2}\right)
\end{aligned}
$$

with $K_{2}$ compact such that

$$
\begin{gathered}
D_{1} P_{1}=1, D_{3} P_{3}=1, \\
D_{2} P_{2}=1+K_{2} .
\end{gathered}
$$

Set

$$
P:=\phi_{1} P_{1} \phi_{1}+\phi_{2} P_{2} \phi_{2}+\phi_{3} P_{3} \phi_{3} .
$$

Notice that $P$ is a well defined operator $L^{2}(\mathbb{R} \times Y) \rightarrow W^{1,2}(\mathbb{R} \times Y)$. Then

$$
\begin{aligned}
\sigma_{\delta} \mathscr{D}_{A} \sigma_{\delta}^{-1}(P f) & =\sum_{i} D_{i} \phi_{i} P_{i} \phi_{i} f \\
& =\left(\sum_{i} \phi_{i} D_{i} P_{i} \phi_{i} f\right)+\left(\sum_{i}\left[D_{i}, \phi_{i}\right] P_{i} \phi_{i} f\right) \\
& =\left(\sum_{i} \phi_{i}^{2} f\right)+\left(\phi_{2} K_{2} \phi_{2} f+\sum_{i}\left[D_{i}, \phi_{i}\right] P_{i} \phi_{i} f\right) \\
& =f+K f
\end{aligned}
$$

with $K$ compact.

Similarly, we can find left-parametrices for the $D_{i}$ and construct a left parametrix for $\sigma_{\delta} \mathscr{D}_{A} \sigma_{\delta}^{-1}$ using them. Hence $\delta \notin \mathfrak{G}_{A}$ implies $\mathscr{P}_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}$ is Fredholm.
The converse is a corollary of Theorem 6.4-1 It should be noted that we do not use this part of the result to establish Theorem 6.4-1

This last theorem now allows us to prove at last what we are really after.
Theorem 6.1-5. Let $(E, A)$ be a $S U(2)$-instanton on $\mathbb{R} \times Y$. Suppose that $A$ is in temporal gauge and that it converges to flat connections $\Gamma_{+}$at $+\infty$ and $\Gamma_{-}$at $-\infty$. Then the operator

$$
\mathscr{D}_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}
$$

is Fredholm if and only if $\delta \notin \mathfrak{G}_{A}$.
Proof: Let

$$
\left(\chi_{R}^{+}, \chi_{R}^{-}, \chi_{R}^{0}\right)
$$

be a partition of unity subordinate to the covering

$$
((R, \infty) \times Y,(-\infty,-R) \times Y,(-R-1, R+1) \times Y)
$$

Suppose $\Gamma_{ \pm}=d+\gamma_{ \pm}$and $A=d+a$. Then $a$ tends to $\gamma_{+}$and $\gamma_{-}$when $t$ tends to $+\infty$ and $-\infty$ respectively. Set

$$
\begin{equation*}
a_{R}=\chi_{R}^{+} \gamma_{+}+\chi_{R}^{-} \gamma_{-}+\chi_{R}^{0} a . \tag{6.3}
\end{equation*}
$$

For simplicity, we bring the discussion back to the classical Sobolev spaces $W^{1,2}$ and $L^{2}$ as we did in the proof of Theorem 6.1-3 Set

$$
\begin{aligned}
E_{n} & :=\sigma_{\delta} \mathscr{D}_{a_{n}} \sigma_{\delta}^{-1}, \text { and } \\
E & :=\sigma_{\delta} \mathscr{D}_{A} \sigma_{\delta}^{-1} .
\end{aligned}
$$

All the $E$ and $E_{n}$ are operators from $W^{1,2}$ to $L^{2}$. Our aim is to show that $E$ is Fredholm if and only if $\delta \notin \mathfrak{G}_{A}$. By virtue of Theorem 6.1-4 it is precisely out of that grid that $E_{1}$ is Fredholm. We now prove that $E-E_{1}$ is compact, whence the result.
Define the operator $K_{n}:=\mathscr{D}_{a_{n}}-\mathscr{P}_{a_{1}}$. Then

$$
\begin{aligned}
K_{n} & =\operatorname{cl}\left(a_{n}-a_{1}\right) \\
& =\left(\chi_{n}^{+}-\chi_{1}^{+}\right) \operatorname{cl}\left(\gamma_{+}\right)+\left(\chi_{n}^{-}-\chi_{1}^{-}\right) \operatorname{cl}\left(\gamma_{-}\right)+\left(\chi_{n}^{0}-\chi_{1}^{0}\right) \operatorname{cl}(a) .
\end{aligned}
$$

As it is a zeroth order operator, $K_{n}$ is continuous $W^{1,2} \rightarrow W^{1,2}$. Observe that the coefficients in $K_{n}$ have compact support:

$$
\begin{aligned}
\operatorname{supp}\left(\chi_{n}^{-}-\chi_{1}^{-}\right) & =[-n-1,-1] \times Y \\
\operatorname{supp}\left(\chi_{n}^{+}-\chi_{1}^{+}\right) & =[1, n+1] \times Y \\
\operatorname{supp}\left(\chi_{n}^{0}-\chi_{1}^{0}\right) & =\operatorname{supp}\left(\chi_{n}^{-}-\chi_{1}^{-}\right) \cup \operatorname{supp}\left(\chi_{n}^{+}-\chi_{1}^{+}\right)
\end{aligned}
$$

Hence $K_{n}$ factorizes through the compact inclusion

$$
W^{1,2}([-n-1, n+1] \times Y) \subset L^{2}([-n-1, n+1] \times Y)
$$

Thus $K_{n}$ is compact.
Suppose without loss of generality that $n<m$. Since

$$
K_{n}-K_{m}=\left(\chi_{n}^{+}-\chi_{m}^{+}\right) c l\left(\gamma_{+}\right)+\left(\chi_{n}^{-}-\chi_{m}^{-}\right) c l\left(\gamma_{-}\right)+\left(\chi_{n}^{0}-\chi_{m}^{0}\right) c l(a)
$$

we have

$$
\left\|\left(K_{n}-K_{m}\right) \phi\right\|^{2}=\left(\int_{n}^{m+1}\right)+\left(\int_{-m-1}^{-n}\right)
$$

The first integral involves only $\gamma_{+}$and $a$. On that domain, $\chi_{j}^{0}+\chi_{j}^{+}=1$ for any $j$. Hence on $[n, m+1]$,

$$
\left(K_{n}-K_{m}\right) \phi=\left(\chi_{n}^{0}-\chi_{m}^{0}\right) c l\left(a-\gamma_{+}\right) \phi
$$

Since $a=\gamma_{+}+O(1 / t)$, we have

$$
\left\|\left(K_{n}-K_{m}\right) \phi\right\| \leq C \frac{1}{n}\|\phi\|
$$

for any Sobolev norm. Hence the sequence of compact operator $K_{n}$ is Cauchy and its limit $K$ is compact. Now obviously $\mathscr{D}_{A}-\mathscr{D}_{a_{1}}=K$ hence $\mathscr{D}_{A}$ is Fredholm if and only if $\mathscr{P}_{a_{1}}$ is Fredholm. The proof is now complete.

### 6.2 Elliptic estimates

As for the compact case, we do have elliptic estimates but those are not sufficient to prove the Fredholmness of $\mathscr{D}$, which is why we need the more involved proofs of the previous section. We do however need those inequalities for finding the asymptotic behavior of harmonic spinors in the Chapter 7 Let's derive them.

Theorem 6.2-1 (Gårding Inequality). Let $A$ be an instanton on $\mathbb{R} \times Y$. If $s \in L^{2}$ and $\mathscr{D}_{A} s \in L^{2}$, then $s \in W^{1,2}$ and

$$
\begin{equation*}
\|s\|_{W^{1,2}} \leq C\left(\left\|\mathscr{P}_{A} s\right\|_{L^{2}}+\|s\|_{L^{2}}\right) \tag{6.4}
\end{equation*}
$$

Proof: Let $s_{c}$ denote the scalar curvature of $Y$. We start with the Weitzenbock formula:

$$
\mathscr{D}_{A}^{*} \mathscr{D}_{A} s=\nabla_{A}^{*} \nabla_{A} s+\left(c l\left(F_{A}\right)+\frac{s_{c}}{4}\right) s
$$

Suppose $s$ has compact support. Then $\left\|\mathscr{P}_{A} s\right\|_{L^{2}}^{2}=\left\|\nabla_{A} s\right\|_{L^{2}}^{2}+\left\langle\left(c l\left(F_{A}\right)+s_{c} / 4\right) s, s\right\rangle_{L^{2}}$, thus

$$
\begin{aligned}
\left\|\nabla_{A} s\right\|_{L^{2}}^{2} & \leq\left\|\mathscr{D}_{A} s\right\|_{L^{2}}^{2}+\sup \left(\left|F_{A}+s_{c} / 4\right|\right)\|s\|_{L^{2}}^{2} \\
& \leq \max \left(\sup \left(\left|F_{A}+s_{c} / 4\right|\right), 1\right)\left(\left\|\mathscr{P}_{A} s\right\|_{L^{2}}+\|s\|_{L^{2}}\right)^{2}
\end{aligned}
$$

While this inequality is good, we must not forget that the $W^{1,2}$-norm is defined using the trivial connection $\nabla$. Fortunately, for $C$ being, say, $2+\sup \left(\sqrt{\left|F_{A}+s_{c} / 4\right|}\right)+\sup (|A|)$, we have

$$
\begin{aligned}
\|\nabla s\|_{L^{2}} & \leq\left\|\nabla_{A} s\right\|_{L^{2}}+\|A s\|_{L^{2}} \\
& \leq\left\|\nabla_{A} s\right\|_{L^{2}}+\sup (|A|)\|s\|_{L^{2}} \\
& \leq C\left(\left\|\mathfrak{D}_{A} s\right\|_{L^{2}}+\|s\|_{L^{2}}\right)
\end{aligned}
$$

Since $A$ is in radial gauge, the ASD and $L^{2}$ conditions on the curvature imply that $C$ is finite. Thus Equation (6.4) is proved for $s$ with compact support.
Suppose now that $s$ does not have compact support. We use now a trick used also in [LM89] p. 117] to show that $\operatorname{ker}(\not \mathscr{D})=\operatorname{ker}\left(\mathscr{D}^{2}\right)$ on a complete manifold.
Choose $\chi \in C^{\infty}(\mathbb{R})$ such that

$$
\begin{gathered}
0 \leq \chi \leq 1 \\
\chi(t)=1 \text { for }|t| \leq 1, \\
\chi(t)=0 \text { for }|t| \geq 2, \text { and }\left|\chi^{\prime}\right| \leq 2 .
\end{gathered}
$$

We set $\chi_{n}(x, t):=\chi(t / n)$ on $\mathbb{R} \times Y$, and set $s_{n}:=\chi_{n} s$. The sequence $s_{n}$ has compact support and

$$
\begin{equation*}
s_{n} \rightarrow s \text { in } L^{2} . \tag{6.5}
\end{equation*}
$$

We know that

$$
\mathscr{D}_{A} s_{n}=\operatorname{cl}\left(\text { gmad } \chi_{n}\right) s+\chi_{n} \mathscr{D}_{A} s .
$$

Obviously $\chi_{n} \mathscr{D}_{A} s$ converges to $\mathscr{D}_{A} s$ in $L^{2}$, and $\left\|c l\left(g r o d ~ \chi_{n}\right) s\right\|_{L^{2}} \leq(2 / n)\|s\|_{L^{2}}$, hence

$$
\begin{equation*}
\mathscr{D}_{A} s_{n} \rightarrow \mathscr{D}_{A} s \text { in } L^{2} . \tag{6.6}
\end{equation*}
$$

Consequently, because of (6.5) and (6.6), and because Equation (6.4) is true for the $s_{n}$, we see that $s_{n}$ is a Cauchy sequence in $W^{1,2}$. Hence $s_{n}$ converges to, say, $\tilde{s}$ in $W^{1,2}$, whence it converges to $\tilde{s}$ in $L^{2}$. Thus $\tilde{s}=s$, and $s \in W^{1,2}$ as wanted, with norm bounded as in Equation (6.4). The proof is now complete.

We push things up the scale a tiny bit with the next result.
Corollary 6.2-2 (Elliptic Estimate). Let A be an instanton on $\mathbb{R} \times Y$ If $s \in W^{k, 2}$ and $\mathscr{D}_{A} s \in W^{k, 2}$, then $s \in W^{k+1,2}$ and

$$
\begin{equation*}
\|s\|_{W^{k+1,2}} \leq C\left(\left\|\oiint_{A} s\right\|_{W^{k, 2}}+\|s\|_{W^{k, 2}}\right) . \tag{6.7}
\end{equation*}
$$

Proof: We prove it by induction, the first step being the result of Theorem 6.2-1 Suppose the result is true for $k-1$, and suppose $s \in W^{k, 2}$ and $\mathscr{D}_{A} s \in W^{k, 2}$. Then $\nabla s \in W^{k-1,2}$ and

$$
\mathscr{D}_{A} \nabla s=\nabla \mathscr{D}_{A} s+\left[\mathscr{D}_{a}, \nabla\right] s .
$$

The first term of the right hand side is in $W^{k-1,2}$, while the second,

$$
\left[\mathscr{D}_{a} s, \nabla\right] s=-\sum_{j} c l\left(\partial_{j}\right) \nabla\left(A\left(\partial_{j}\right)\right) s,
$$

is in $W^{k-1,2}$ if $A \in W^{k-1, \infty}$, which is the case for a good choice of gauge as $A$ is an instanton. We hence have, by induction, that $\nabla s \in W^{k-1,2}$, which means that $s \in W^{k, 2}$ and

$$
\begin{aligned}
\|s\|_{W^{k, 2}} & =\|s\|_{L^{2}}+\|\nabla s\|_{W^{k-1,2}} \\
& \leq\|s\|_{L^{2}}+C\left(\left\|\not \wp_{A} \nabla s\right\|_{W^{k-1,2}}+\|\nabla s\|_{W^{k-1,2}}\right) \\
& \leq C\left(\|s\|_{L^{2}}+\left\|\nabla \mathscr{D}_{A} s\right\|_{W^{k-1,2}}+\left\|\left[\oiint_{A}, \nabla\right] s\right\|_{W^{k-1,2}}+\|\nabla s\|_{W^{k-1,2}}\right) \\
& \leq C^{\prime}\left(\left\|\oiint_{A} s\right\|_{W^{k, 2}}+\|s\|_{W^{k, 2}}\right) .
\end{aligned}
$$

The proof is now complete.
Remark 6.2-3. Note that in Estimates (6.4) and 6.7), we can choose a uniform $C$ for any family of connections $A_{z}$ parameterized by $z$ in a compact set, independently of whether or not $\mathscr{D}_{A_{z}}$ is Fredholm everywhere.

We can even prove a better estimate, which also hints to Fredholmness properties for $\delta \notin \mathfrak{G}_{A}$.
Theorem 6.2-4 (Gårding plus). Suppose $\delta \notin \mathfrak{G}_{A}$. Then there exist a compact subcylinder $K$ large enough so that for any $s \in W^{1,2}$, we have

$$
\begin{equation*}
\|s\|_{W_{\delta}^{1,2}} \leq C\left(\left\|\oiint_{A} s\right\|_{L_{\delta}^{2}}+\|s\|_{L_{\delta}^{2}(K)}\right), \tag{6.8}
\end{equation*}
$$

with $K$ and $C$ depending only on $A$ and $\delta$.
Proof: First set $K_{R}:=(-\infty, R] \times T^{3}$, and suppose $\operatorname{supp}(s) \cap K_{R}=\varnothing$. Then

$$
\left\|\left(\mathfrak{D}_{A}-\mathfrak{D}_{\Gamma}\right) s\right\|_{L^{2}} \leq\left(\sup _{t>R}|A-\Gamma|\right)\|s\|_{L^{2}}
$$

hence as $R \rightarrow \infty$, the operator norm of the restriction of $\mathscr{D}_{A}-\mathscr{D}_{\Gamma}$ on elements with support out of $K_{R}$, denoted $\left\|\mathscr{D}_{A}-\mathscr{D}_{\Gamma}\right\|_{o p, R}$, decreases to 0 .
Let $\chi_{R}$ be a cut off function, with $\chi_{R}(t)=0$ for $t>R+1$, and $\chi_{R}(t)=1$ for $t \leq R$. Write $s=s_{0}+s_{\infty}$, with $s_{0}=\chi_{R} s$ and $s_{\infty}=\left(1-\chi_{R}\right) s$. Since $\mathscr{D}_{\Gamma}$ is an isomorphism, we have

$$
\begin{aligned}
\left\|s_{\infty}\right\|_{W^{1,2}} & \leq C\left\|\mathfrak{D}_{\Gamma} s_{\infty}\right\|_{L^{2}} \\
& \leq C\left(\left\|\mathfrak{D}_{A} s_{\infty}\right\|_{L^{2}}+\left\|\left(\mathfrak{D}_{A}-\mathfrak{D}_{\Gamma}\right) s_{\infty}\right\|_{L^{2}}\right) \\
& \leq C\left(\left\|\mathscr{D}_{A} s_{\infty}\right\|_{L^{2}}+\left\|\mathfrak{D}_{A}-\mathfrak{D}_{\Gamma}\right\|_{o p, R}\left\|s_{\infty}\right\|_{L^{2}}\right) \\
& \leq C\left(\left\|\mathfrak{D}_{A} s_{\infty}\right\|_{L^{2}}+\left\|\mathfrak{D}_{A}-\mathfrak{D}_{\Gamma}\right\|_{o p, R}\left\|s_{\infty}\right\|_{W^{1,2}}\right) .
\end{aligned}
$$

But now,

$$
\begin{aligned}
&\left\|\mathfrak{D}_{A} s_{\infty}\right\|_{L^{2}} \leq\left\|\mathscr{D}_{A} s\right\|_{L^{2}}+\left\|\oiint_{A} s_{0}\right\|_{L^{2}} \\
& \leq\left\|\oiint_{A} s\right\|_{L^{2}}+\left\|\chi_{R} \mathfrak{D}_{A} s\right\|_{L^{2}}+\| c l(g r o d \\
&\left.\chi_{R}\right) s \|_{L^{2}} \\
& \leq C\left(\left\|\mathfrak{D}_{A} s\right\|_{L^{2}}+\|s\|_{L^{2}\left(K_{R+1}\right)}\right) .
\end{aligned}
$$

Hence

$$
\left\|s_{\infty}\right\|_{W^{1,2}} \leq C\left(\left\|\oiint_{A} s\right\|_{L^{2}}+\|s\|_{L^{2}\left(K_{R+1}\right)}+\left\|\mathscr{D}_{A}-\mathscr{D}_{\Gamma}\right\|_{o p, R}\left\|s_{\infty}\right\|_{W^{1,2}}\right) .
$$

For $R$ big enough, we can rearrange to obtain

$$
\left\|s_{\infty}\right\|_{W^{1,2}} \leq C\left(\left\|\oiint_{A} s\right\|_{L^{2}}+\|s\|_{L^{2}\left(K_{R+1}\right)}\right) .
$$

We can now play the same game at $-\infty$, splitting $s_{0}$ as $s_{-\infty}+\tilde{s}_{0}$ and we obtain a similar estimate. Once we patch all those estimates together, we find that there is a $R$ big enough so that for the compact subcylinder $K:=[-R, R] \times T^{3}$, we have the desired Inequality (6.8).
Note that the assumption $s \in W^{1,2}$ is important. This Theorem does not prove that $\mathscr{D}_{A} \phi \in L^{2}$ and $\left.\phi\right|_{K} \in L^{2}$ implies $\phi \in W^{1,2}$. If that implication were true, then in the language of Chapter $\mathbb{8}$ it
would rule out the possibility that $\mathfrak{D H} \cap V \neq\{0\}$, hence would imply that $V$ and $V$ are always equal.

### 6.3 Invariance of the kernels

We define the spaces

$$
\begin{align*}
\operatorname{ker}(\delta) & :=\operatorname{ker}\left(\mathscr{P}_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}\right), \\
\operatorname{ker}^{*}(\delta) & :=\operatorname{ker}\left(\mathfrak{P}_{A}^{*}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}\right), \tag{6.9}
\end{align*}
$$

and the integers

$$
\begin{align*}
\operatorname{ind}(\delta) & :=\operatorname{ind}\left(\oiint_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}\right) \\
N(\delta) & :=\operatorname{dim} \operatorname{ker}(\delta), \text { and }  \tag{6.10}\\
N^{*}(\delta) & :=\operatorname{dim} \operatorname{ker}^{*}(\delta) .
\end{align*}
$$

Since $\left(L_{\delta}^{2}\right)^{*}=L_{-\delta}^{2}$, Theorem 6.2-1 tells us that $\operatorname{dim} \operatorname{coker}\left(\mathscr{P}_{A}\right)=N^{*}(-\delta)$, hence

$$
\operatorname{ind}(\delta)=N(\delta)-N^{*}(-\delta)
$$

That the formal adjoint $\mathscr{P}_{A}^{*}$ on $W_{-\delta}^{1,2}$ is really the adjoint of $\mathscr{D}_{A}$ on $W_{\delta}^{1,2}$ is guaranteed by the following lemma.

Lemma 6.3-1. The subspace $\operatorname{ker}^{*}(-\delta)$ of $L_{-\delta}^{2}=\left(L_{\delta}^{2}\right)^{*}$ kills $\operatorname{Im}(\delta)$ in the $L^{2}$ natural pairing.
Proof: Suppose $\phi$ is a smooth function with compact support. Then for all $\psi \in \operatorname{ker}^{*}(-\delta)$, we have $\langle\psi, \mathscr{D} \phi\rangle=\left\langle\mathscr{D}^{*} \psi, \phi\right\rangle=0$. Since $C_{c}^{\infty}$ is dense in $W_{\delta}^{1,2}$, the lemma holds.

Let's say our instanton $A$ has limit $\Gamma_{ \pm}$as $t$ tends to $\pm \infty$. Recall from Equation (6.2) the definition of the grid

$$
\mathfrak{G}_{A}=\left(\operatorname{Spec}\left(\Gamma_{-}\right) \times \mathbb{R}\right) \cup\left(\mathbb{R} \times \operatorname{Spec}\left(\Gamma_{+}\right)\right)
$$

in $\mathbb{R}^{2}$. As we have shown in Theorem 6.1-5 the operator $\mathscr{D}_{A}: W_{\delta}^{1,2} \rightarrow L_{\delta}^{2}$ is Fredholm if and only if $\delta \notin \mathfrak{G}_{A}$. In fact, we have more, as in shown by the next theorem.

Theorem 6.3-2. In each open square of $\mathbb{R}^{2}$ delimited by the grid $\mathfrak{G}_{A}$, the quantities

$$
\operatorname{ind}(\delta), N(\delta), \text { and } N^{*}(\delta)
$$

are constant. In fact, for $\delta, \eta$ in a same square,

$$
\operatorname{ker}(\delta)=\operatorname{ker}(\eta), \text { and } \operatorname{ker}^{*}(\delta)=\operatorname{ker}^{*}(\eta)
$$

Proof: We use the family $D_{\delta}: W^{1,2} \rightarrow L^{2}$ of operators defined as

$$
D_{\delta}:=\sigma_{\delta} \not D_{A} \sigma_{\delta}^{-1} .
$$

The family is linear in $\delta$ as

$$
\begin{aligned}
D_{\delta} & =\bigoplus_{A}+\sigma_{\delta} c l\left(\text { grod } \sigma_{\delta}^{-1}\right) \\
& =\bigoplus_{A}+\left(c \delta_{-}+(1-c) \delta_{+}+t c^{\prime}\left(\delta_{-}-\delta_{+}\right)\right) .
\end{aligned}
$$

Hence the family depends continuously in the operator topology on the parameter $\delta \in \mathbb{R}^{2}$. Since

$$
\operatorname{ind}(\delta)=\operatorname{ind}\left(D_{\delta}\right),
$$

we see that indeed $\operatorname{ind}(\delta)$ is constant on each open square.
Let's define a partial ordering on $\mathbb{R}^{2}$ as follows

$$
\delta \leq \eta \Longleftrightarrow \delta_{-} \geq \eta_{-}, \text {and } \delta_{+} \leq \eta_{+}
$$

This ordering is designed so that

$$
\delta \leq \eta \Longrightarrow W_{\delta}^{k, p} \subset W_{\eta}^{k, p}
$$

Suppose for the moment that $\delta, \eta$ in the same open square are such that $\delta \leq \eta$. We then have

$$
\begin{gather*}
\operatorname{ker}(\delta) \subset \operatorname{ker}(\eta), \text { hence } \\
N(\delta) \leq N(\eta) . \tag{6.11}
\end{gather*}
$$

Similarly, as $-\delta \geq-\eta$, we have

$$
\begin{gather*}
\operatorname{ker}^{*}(-\delta) \supset \operatorname{ker}^{*}(-\eta), \text { hence } \\
N^{*}(-\delta) \geq N^{*}(-\eta) \tag{6.12}
\end{gather*}
$$

But then, $\operatorname{ind}(\delta)=\operatorname{ind}(\eta)$ implies

$$
N(\delta)-N(\eta)=N^{*}(-\delta)-N^{*}(-\eta)
$$

Inequality (6.11) shows that the left-hand-side is nonpositive, while Inequality (6.12) shows that the right-hand-side is nonnegative. Hence both sides must be zero, and moreover

$$
\begin{aligned}
\operatorname{ker}(\delta) & =\operatorname{ker}(\eta) \\
\operatorname{ker}^{*}(-\delta) & =\operatorname{ker}(-\eta)
\end{aligned}
$$

The proof is not complete yet, as $\delta$ and $\eta$ could be incomparable. In that case, we can find $\gamma$ in the same open square smaller than both. We then have

$$
\operatorname{ker}(\delta)=\operatorname{ker}(\gamma)=\operatorname{ker}(\eta)
$$

and similarly for ker $^{*}, N$ and $N^{*}$. The proof is now complete.

### 6.4 Wall crossing

The following theorem tells us how the index changes as we cross a wall to change square. This theorem is quite useful for our main purpose on $\mathbb{R} \times T^{3}$ especially once we know the index of $\mathscr{P}_{A}$ on weighted Sobolev spaces for weights contained in the open square around 0 , which we compute in Section 8.1

Theorem 6.4-1 (Wall Crossing). For an ASD connection $A$ on $\mathbb{R} \times Y$ converging to the flat con-
nections $\Gamma_{ \pm}$on $Y$ at $\pm \infty$, the index of $\mathscr{D}_{A}$ and $\mathscr{D}_{A}^{*}$ changes as follows:

$$
\begin{aligned}
\operatorname{ind}(\delta) & =\operatorname{ind}(\eta)+\operatorname{dim}\left\{D_{\Gamma_{+}} \phi=-\lambda \phi\right\}, \text { and } \\
\operatorname{ind}^{*}(\delta) & =\operatorname{ind}^{*}(\eta)+\operatorname{dim}\left\{D_{\Gamma_{+}} \phi=\lambda \phi\right\}
\end{aligned}
$$

when $\delta_{+}<\eta_{+}$, and $\delta$ and $\eta$ are in adjacent open squares separated by the wall $\mathbb{R} \times\{\lambda\} \subset \mathfrak{G}_{A}$;

$$
\begin{aligned}
\operatorname{ind}(\delta) & =\operatorname{ind}(\eta)+\operatorname{dim}\left\{D_{\Gamma_{-}} \phi=-\lambda \phi\right\}, \text { and } \\
\operatorname{ind}^{*}(\delta) & =\operatorname{ind}^{*}(\eta)+\operatorname{dim}\left\{D_{\Gamma_{-}} \phi=\lambda \phi\right\}
\end{aligned}
$$

when $\delta_{-}>\eta_{-}$, and $\delta$ and $\eta$ are in adjacent open squares separated by the wall $\{\lambda\} \times \mathbb{R} \subset \mathfrak{G}_{A}$.

Proof: We start by considering that $A$ is constant in $t$; say $A=\Gamma$. For simplicity, set

$$
\begin{gathered}
W_{\lambda}=\left\{D_{\Gamma} \phi=\lambda \phi\right\}, \text { and } \\
d_{\lambda}=\operatorname{dim} W_{\lambda}
\end{gathered}
$$

We have

$$
\begin{aligned}
\operatorname{ker}\left(\mathscr{D}_{\Gamma}\right) & =\bigoplus_{\lambda} e^{-\lambda t} W_{\lambda} \\
\operatorname{ker}\left(\mathscr{D}_{\Gamma}^{*}\right) & =\bigoplus_{\lambda} e^{\lambda t} W_{\lambda}
\end{aligned}
$$

Hence

$$
\begin{aligned}
& \operatorname{ker}\left(\mathscr{D}_{\Gamma}\right) \cap W_{\delta}^{1,2}=\bigoplus_{\delta_{-}<-\lambda<\delta_{+}} e^{-\lambda t} W_{\lambda}, \\
& \operatorname{ker}\left(\mathscr{D}_{\Gamma}^{*}\right) \cap W_{\delta}^{1,2}=\bigoplus_{\delta_{-}<\lambda<\delta_{+}} e^{\lambda t} W_{\lambda} .
\end{aligned}
$$

Now for $\delta \notin \mathfrak{G}_{\Gamma}$, we know $\mathscr{P}_{\Gamma}$ is Fredholm hence

$$
\begin{aligned}
\operatorname{ind}(\delta) & =N(\delta)-N^{*}(-\delta) \\
& =\sum_{\delta_{-}<-\lambda<\delta_{+}} d_{\lambda}-\sum_{-\delta_{-}<\lambda<-\delta_{+}} d_{\lambda} \\
& =\sum_{\delta_{-}<-\lambda<\delta_{+}} d_{\lambda}-\sum_{\delta_{+}<-\lambda<\delta_{-}} d_{-\lambda}
\end{aligned}
$$

Suppose $\delta$ and $\eta$ are in adjacent open squares delimited by $\mathbb{R}^{2} \backslash \mathfrak{G}_{\Gamma}$, say $\delta$ is in the square to the left of the square containing $\eta$, and both squares are separated by $\{a\} \times \mathbb{R} \subset \mathfrak{G}_{\Gamma}$.
Since the index is constant in each open square, we can pick $\delta$ and $\eta$ such that

$$
\begin{aligned}
& \delta=(a-\epsilon, b) \\
& \eta=(a+\epsilon, b)
\end{aligned}
$$

with $a+\epsilon<b$ or $b<a-\epsilon$.

Suppose $a+\epsilon<b$. Then $N^{*}(-\delta)=N^{*}(-\eta)=0$ and

$$
N(\delta)=\sum_{a-\epsilon<-\lambda<b} d_{\lambda}=d_{-a}+\sum_{a+\epsilon<-\lambda<b} d_{\lambda}=d_{-a}+N(\eta)
$$

Suppose on the contrary that $b<a-\epsilon$. Then $N(\delta)=N(\eta)=0$ and

$$
N^{*}(-\delta)=\sum_{b<-\lambda<a-\epsilon} d_{-\lambda}=-d_{a}+\sum_{b<-\lambda<a+\epsilon} d_{\lambda}=-d_{a}+N^{*}(-\eta)
$$

Hence in both cases, we find

$$
\operatorname{ind}(\delta)=d_{-a}+\operatorname{ind}(\eta)
$$

This formula also holds when $\delta$ is in the square above the one containing $\eta$.
Now suppose $A$ has limiting connections $\Gamma_{+}$and $\Gamma_{-}$at $+\infty$ and $-\infty$. We bring all the different operators we want to deal with on $W^{1,2}$ and $L^{2}$, and set

$$
\begin{aligned}
D_{1} & :=\sigma_{\delta} \mathscr{D}_{A} \sigma_{\delta}^{-1} \\
D_{2} & :=\sigma_{\eta} \mathscr{D}_{A} \sigma_{\eta}^{-1} \\
D_{3} & :=\sigma_{\delta} \mathscr{D}_{\Gamma_{+}} \sigma_{\delta}^{-1}, \\
D_{4} & :=\sigma_{\eta} \mathscr{D}_{\Gamma_{+}} \sigma_{\eta}^{-1} .
\end{aligned}
$$

Recall that $D_{1}=\mathscr{D}_{A}+\left(c \delta_{-}+(1-c) \delta_{+}+t c^{\prime}\left(\delta_{-}-\delta_{+}\right)\right)$, and similarly for the others.
Suppose that $\delta_{-}=\eta_{-}$. Notice that $D_{1}-D_{2}=D_{3}-D_{4}$. We can make up a compact operator $K$ so that $D_{1}-D_{2}=K$ for $t \leq 1$. Notice also that $D_{1}=D_{3}$ for $t>1$, and $D_{2}=D_{4}$ for $t>1$. Set

$$
\begin{aligned}
\tilde{D}_{2} & :=D_{2}+K \\
\tilde{D}_{4} & :=D_{4}+K
\end{aligned}
$$

Then we have

$$
\begin{aligned}
& D_{1}= \begin{cases}\tilde{D}_{2}, & \text { for } t \leq 1 \\
D_{3}, & \text { for } t>1\end{cases} \\
& \tilde{D}_{4}= \begin{cases}D_{3}, & \text { for } t \leq 1 \\
\tilde{D}_{2}, & \text { for } t>1\end{cases}
\end{aligned}
$$

The excision principles for indices (see Theorem B-1) tells us that

$$
\operatorname{ind}\left(D_{1}\right)-\operatorname{ind}\left(\tilde{D}_{2}\right)=\operatorname{ind}\left(D_{3}\right)-\operatorname{ind}\left(\tilde{D}_{4}\right)
$$

Since $\operatorname{ind}\left(\tilde{D}_{2}\right)=\operatorname{ind}\left(D_{2}\right)$ and $\operatorname{ind}\left(\tilde{D}_{4}\right)=\operatorname{ind}\left(D_{4}\right)$, we see that the index changes the same way for ASD connections and time-independent connections. The proof is now complete.

Remark 6.4-2. Notice that we proved something better for time-independent $\Gamma$. Indeed the analysis of $\operatorname{ker}(\delta)$ is such that we know its dimension $N(\delta)$ is lower semicontinuous: for

$$
\delta_{t}:=\left(\delta_{-}, \lambda-\epsilon+t\right) \quad \text { or } \quad \delta_{t}:=\left(\delta_{-}+\epsilon-t, \lambda\right),
$$

and for small $\epsilon$ and $t$, and some $a$ and $b$, we have that

$$
\begin{aligned}
& N\left(\delta_{t}\right)=a \text { for } t \leq 0, \\
& N\left(\delta_{t}\right)=b \text { for } t>0 .
\end{aligned}
$$

On $\mathbb{R} \times T^{3}$, the same is actually true for connections $A$ when $A$ decays exponentially to its limits. Suppose

$$
\lambda \in \operatorname{Spec}\left(D_{\Gamma_{-}}\right) \times \operatorname{Spec}\left(D_{\Gamma_{+}}\right),
$$

$\delta$ is in the upper left open square adjacent to $\lambda$, $\eta$ is in the lower right open square adjacent to $\lambda$.

We have $\operatorname{ker}(\lambda)=\operatorname{ker}(\eta)$.
Indeed, suppose now $\phi \in \operatorname{ker}(\lambda)$. Then $\phi \in \operatorname{ker}(\delta)$ hence by Theorem 7.2-1 we expand $\phi$ for $t>0$ as $\phi=e^{-\lambda_{+} t} \psi_{\lambda_{+}}+\bar{\phi}$, with $\bar{\phi} \in W_{\eta_{+}}^{1,2}\left([0, \infty) \times T^{3}\right)$. Since $\phi$ and $\bar{\phi}$ are both in $W_{\lambda_{+}}^{1,2}$, so is the term $e^{-\lambda_{+} t} \psi_{\lambda_{+}}$. This fact implies that $\psi_{\lambda_{+}}=0$. Using a similar proof at $-\infty$, we find $\phi \in W_{\eta}^{1,2}$.

## Chapter 7

## Asymptotic behavior of harmonic spinors

In this chapter, we study how $L^{2}$-harmonic spinors on various spaces decay with time. For a 3manifold $Y$ of scalar curvature $s_{c}$, the Weitzenbock formula (see [Roe98, Prop 3.18, p. 48]) says

$$
\mathscr{D}_{A}^{*} \mathfrak{D}_{A}=\nabla_{A}^{*} \nabla_{A}+\operatorname{cl}\left(F_{A}^{+}\right)+\frac{s_{c}}{4} .
$$

For an ASD connection $A$, we hence see that should $s_{c} \geq 0$, every positive harmonic spinor is parallel for the connection $A$, hence has constant norm. This conclusion certainly prevents it from being $L^{2}$ on the manifold $\mathbb{R} \times Y$ of infinite volume.
In view of theorem 6.3-2, any negative $L^{2}$ harmonic spinor can be seen in $W_{\delta}^{1,2}$ for any $\delta$ in the open square delimited by the grid $\mathfrak{G}_{A}$ (see Equation (6.2) and containing ( 0,0 ). The elliptic bootstrapping of Corollary 6.2-2 and the associated Sobolev embedding of $W^{3,2}$ in bounded $C^{0}$ functions (see Heb99, Thm 3.4, p. 68]) tells us that if

$$
\varphi \in W_{\delta}^{1,2}(\mathbb{R} \times Y), \text { and } \mathscr{P}_{A}^{*} \varphi=0
$$

then

$$
\varphi \leq C_{\alpha} e^{\alpha t}
$$

for all $\alpha$ shy of the first negative eigenvalue of $D_{\Gamma_{+}}$on $Y$ when $t \rightarrow \infty$ and shy of the first positive eigenvalue of $D_{\Gamma_{-}}$on $Y$ when $t \rightarrow-\infty$.
While this result is nice, we can get a better knowledge of the asymptotic behavior. We first build up the theory on half-cylinders, which we need, and then apply it to $T^{3}$ and $S^{3}$ in place of $Y$.

### 7.1 Translation invariant operators on half-cylinders

For a more compact notation, we introduce the following shorthands:

$$
\begin{aligned}
Y_{a} & :=\{a\} \times Y \\
Y_{a+} & :=[a, \infty) \times Y \\
Y^{+} & :=[0, \infty) \times Y
\end{aligned}
$$

We hope to construct the asymptotic expansion of harmonic spinors by comparing the $W_{\delta}^{1,2}$-kernel of $\mathscr{D}^{*}$ for varying $\delta$. To compare then, we need finite dimensionality, or better, finite index. Before studying what gives us those properties, let's first eliminate options that won't.
Given a Dirac operator $D$ on $Y$, with no zero eigenvalue, we have

$$
\|D \phi\|_{L^{2}(Y)} \geq C\|\phi\|_{W^{1,2}(Y)} .
$$

On the full cylinder $Y \times \mathbb{R}$, this estimate was enough to ensure that

$$
\not D: W^{1,2}(\mathbb{R} \times Y) \rightarrow L^{2}(\mathbb{R} \times Y)
$$

is an isomorphism. Working now on the half cylinder $[a, \infty) \times Y$, this estimate is not sufficient, as we now check. Define the $\psi_{\lambda}$ by the eigenvalue equation

$$
D \psi_{\lambda}=\lambda \psi_{\lambda}
$$

Then all the $e^{\lambda t} \psi_{\lambda}$ with $\lambda<0$ are in $W^{1,2}([a, \infty) \times Y) \cap \operatorname{ker}(\nsubseteq)$. So much for Fredholmness. Another option would be to consider the operator

$$
\nsupseteq: W_{\operatorname{tr}=0}^{1,2}([a, \infty) \times Y) \rightarrow L^{2}([a, \infty) \times Y)
$$

on the space of sections whose restriction to $Y_{a}$ is 0 The elliptic estimate (7.2) that we prove below still holds, but as $\phi(a)=0$, we have

$$
\|\phi\|_{W^{1,2}} \leq\|\nsubseteq \phi\|_{L^{2}} .
$$

Hence $\mathscr{D}$ is injective and has close range. It is therefore semi-Fredholm but it isn't Fredholm: its adjoint is the usual $\mathscr{D}^{*}$ with no boundary condition and it has infinite dimensional kernel on a halfspace.
The space $L^{2}(Y)$ splits according to the finite dimensional eigenspaces $W_{\lambda}$ for $D$. Let

$$
\begin{aligned}
\Pi_{\delta}^{+}: L^{2}(Y) & \rightarrow \bigoplus_{\lambda>\delta} W_{\lambda} \\
\Pi_{\delta}^{-}: L^{2}(Y) & \rightarrow \bigoplus_{\lambda<\delta} W_{\lambda} \\
\Pi_{\delta}: L^{2}(Y) & \rightarrow W_{\delta}
\end{aligned}
$$

be the projections. To simplify notation we omit $\delta$ when it is 0 and set $\phi^{ \pm}:=\Pi^{ \pm}(\phi)$.
For every $\phi \in L^{2}(Y)$, let $\phi_{\lambda}$ be its $W_{\lambda}$ component. Thus

$$
\phi=\sum \phi_{\lambda} .
$$

Using this decomposition, we can define the space $W^{\frac{1}{2}, 2}(Y)$ using the norm

$$
\begin{equation*}
\|\phi\|_{W^{\frac{1}{2}, 2}}^{2}=\sum(1+|\lambda|)\left\|\phi_{\lambda}\right\|_{L^{2}}^{2} . \tag{7.1}
\end{equation*}
$$

Because $Y$ is compact, the space $W^{\frac{1}{2}, 2}(Y)$ defined by two different Dirac operators are equal, with commensurate norms. The + and - part of $L^{2}$, however, depend highly on $D$.

Theorem 7.1-1. The operator

$$
\begin{aligned}
\not D: & W^{1,2}\left(Y_{a+}\right)
\end{aligned} \rightarrow L^{2}\left(Y_{a+}\right) \oplus \Pi^{+} W^{\frac{1}{2}, 2}\left(Y_{a}\right)
$$

is an isomorphism.

Keep in mind that $D$ has no kernel.
Proof: It all starts as in the full cylinder case:

$$
\begin{aligned}
\|\not \supset \phi\|_{L^{2}}^{2} & =\left\|\partial_{t} \phi\right\|_{L^{2}}^{2}+\|D \phi\|_{L^{2}}^{2}+\int_{a}^{\infty} \partial_{t}\langle\phi, D \phi\rangle_{L^{2}(Y)} \\
& \geq C\|\phi\|_{W^{1,2}}^{2}-\langle\phi(a), D \phi(a)\rangle_{L^{2}(Y)}
\end{aligned}
$$

Contrary to the full cylinder case, the boundary term here cannot be made to vanish and henceforth helps control the $W^{1,2}$-norm of $\phi$. For the decomposition $\phi=\sum \phi_{\lambda}$, we have

$$
\begin{align*}
\|\phi\|_{W^{1,2}}^{2} & \leq C\left(\|\not \supset \phi\|_{L^{2}}^{2}+\langle\phi(a), D \phi(a)\rangle_{L^{2}(Y)}\right) \\
& \leq C\left(\|\nsupseteq \phi\|_{L^{2}}^{2}+\sum \lambda\left\|\phi_{\lambda}(a)\right\|_{L^{2}(Y)}^{2}\right) \\
& \leq C\left(\|\nsupseteq \phi\|_{L^{2}}^{2}+\sum_{\lambda>0}|\lambda|\left\|\phi_{\lambda}(a)\right\|_{L^{2}(Y)}^{2}\right) \\
& \leq C\left(\|\not \supset \phi\|_{L^{2}}^{2}+\left\|\phi^{+}(a)\right\|_{W^{\frac{1}{2}, 2}(Y)}^{2}\right) \tag{7.2}
\end{align*}
$$

We just proved that $\|\phi\|_{W^{1,2}} \leq C\|\not D \phi\|$, hence $\not D D$ is semi-Fredholm and injective.
Suppose that $(\psi, \eta)$ is perpendicular to $\operatorname{Im}(\not D D)$. For all $\phi \in W^{1,2}([a, \infty) \times Y)$, we have

$$
\begin{aligned}
0 & =\langle\nsupseteq \phi, \psi\rangle+\left\langle\eta, \phi^{+}(a)\right\rangle \\
& =\left\langle\phi, \mathscr{D}^{*} \psi\right\rangle-\langle\phi(a), \psi(a)\rangle+\left\langle\eta, \phi^{+}(a)\right\rangle \\
& =\left\langle\phi, \mathfrak{D}^{*} \psi\right\rangle-\left\langle\phi^{-}(a), \psi^{-}(a)\right\rangle+\left\langle\eta-\psi^{+}(a), \phi^{+}(a)\right\rangle
\end{aligned}
$$

Going through all the $\phi$ with $\phi(a)=0$ in a first time, $\phi^{+}(a)=0$ then, and finally $\phi^{-}(a)=0$, we prove

$$
\begin{gathered}
\mathfrak{P}^{*} \psi=0 \\
\eta=\psi^{+}(a) \\
\psi^{-}(a)=0
\end{gathered}
$$

Thus we have $-\partial_{t} \psi+D \psi=0$, which means that $\psi$ is a linear combination of the $e^{\lambda t} \psi_{\lambda}$. The condition $\psi^{-}(a)=0$ forces out all the negative $\lambda$, while the positive ones are forced out by the $L^{2}$ condition. Hence $\psi=0$ and $I D$ is surjective. The proof is now complete.

While $\mathscr{D}: W^{1,2}([a, \infty) \times Y) \rightarrow L^{2}([a, \infty) \times Y)$ is not Fredholm, an easy corollary of Theorem $7.1-1$ is that it is surjective. Hence

$$
\mathfrak{D}: W^{1,2}([a, \infty) \times Y) \rightarrow L^{2}([a, \infty) \times Y)
$$

is semi-Fredholm with

$$
\operatorname{ind}(\mathfrak{D})=\infty .
$$

## Weighted version

As in the full cylinder case, we can look at weighted version of $\mathscr{D}$ and $\mathscr{D}$. For computing the asymptotic expansion of harmonic spinors, we actually need to consider the dual $\mathscr{D}^{*}$ and its counterpart

$$
\begin{aligned}
\mathbb{D}^{*}: W^{1,2}\left(Y_{a+}\right) & \rightarrow L^{2}\left(Y_{a+}\right) \oplus \Pi^{-} W^{\frac{1}{2}, 2}(Y) \\
\phi & \mapsto\left(\mathfrak{D}^{*} \phi, \Pi^{-} \phi(a)\right)
\end{aligned}
$$

Staring at the diagrams

shows that the top row $\mathscr{D}^{*}$ and $\mathscr{D}^{*}$ are respectively semi-Fredholm and Fredholm if and only if $\delta \notin \operatorname{Spec}(D)$. Moreover, when $\delta \notin \operatorname{Spec}(D)$ they are surjective and an isomorphism respectively.

## Independence of the norm

For any operator $T: W^{1,2}(Y) \rightarrow L^{2}(Y)$, the operator

$$
\hat{T}:=\partial_{t}+T
$$

has a norm independent of the half-cylinder on which that norm is taken. In other words,

$$
\begin{aligned}
& \hat{T}: W^{1,2}\left(Y_{a+}\right) \rightarrow L^{2}\left(Y_{a+}\right) \\
& \hat{T}: W^{1,2}\left(Y_{b+}\right) \rightarrow L^{2}\left(Y_{b+}\right)
\end{aligned}
$$

have the same operator norm. This fact is a manifestation of the translation invariance of $\hat{T}$. To prove this claim, consider the following characterization of the norm:

$$
\|\hat{T}\|=\sup \{\|\hat{T} f\|:\|f\|=1\} .
$$

Shifting a function in $t$ by $b-a$ doesn't change its $L^{2}$ or $W^{1,2}$ norm and shifts its value under $\hat{T}$. So if $f_{b-a}(y, t):=f(y, t+b-a)$, then

$$
\begin{aligned}
\|\hat{T}\|_{o p, a} & =\sup \left\{\|\hat{T} f\|_{L^{2}\left(Y_{a+}\right)}:\|f\|_{W^{1,2}\left(Y_{a+}\right)}=1\right\} \\
& =\sup \left\{\left\|(\hat{T} f)_{b-a}\right\|_{L^{2}\left(Y_{b+}\right)}:\|f\|_{W^{1,2}\left(Y_{a+}\right)}=1\right\} \\
& =\sup \left\{\left\|\hat{T}\left(f_{b-a}\right)\right\|_{L^{2}\left(Y_{b+}\right)}:\|f\|_{W^{1,2}\left(Y_{a+}\right)}=1\right\} \\
& =\|\hat{T}\|_{o p, b} .
\end{aligned}
$$

### 7.2 The commutative diagram

In this section, we derive the asymptotic behavior of harmonic spinors in the case where the connection $A$ decays exponentially to its limit $\Gamma$, with decay rate $\beta$,

$$
|A-\Gamma| \leq C e^{-\beta t}
$$

This feat is be achieved by some diagram chase. We first introduce maps to compose our diagram. Suppose $\eta<\delta$ and $\operatorname{Sppcc}(D) \cap[\eta, \delta]=\{\lambda\}$. Then the map

$$
\begin{aligned}
I: \Pi_{\eta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right) \oplus W_{\lambda} & \rightarrow \Pi_{\delta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right) \\
(\phi, \psi) & \mapsto \phi+e^{a \lambda} \psi
\end{aligned}
$$

is obviously an isomorphism.
Similarly, the map

$$
\begin{aligned}
J: W_{\eta}^{1,2}\left(Y_{a+}\right) \oplus W_{\lambda} & \rightarrow W_{\delta}^{1,2}\left(Y_{a+}\right) \\
(\phi, \psi) & \mapsto \phi+e^{\lambda t} \psi
\end{aligned}
$$

is obviously an injection.
Consider now the map

$$
\begin{aligned}
K: W_{\eta}^{1,2}\left(Y_{a+}\right) \oplus W_{\lambda} & \rightarrow L_{\eta}^{2}\left(Y_{a+}\right) \oplus \Pi_{\eta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right) \oplus W_{\lambda} \\
(\phi, \psi) & \mapsto\left(\mathscr{p}_{A}\left(\phi+e^{\lambda t} \psi\right), \Pi_{\eta}^{-} \phi, \psi+e^{-a \lambda} \Pi_{\lambda} \phi(a)\right) .
\end{aligned}
$$

Let's verify that this map is well-defined. As $\mathscr{D}_{A}^{*}\left(e^{\lambda t} \psi\right)=\mathscr{D}_{\Gamma}^{*}\left(e^{\lambda t} \psi\right)+c l(A-\Gamma) e^{\lambda t} \psi$, we have $\left|\mathscr{D}_{A}^{*}\left(e^{\lambda t} \psi\right)\right| \leq C e^{(\lambda-\beta) t}|\psi|$. Hence, if

$$
\begin{equation*}
\lambda-\beta<\eta, \tag{7.3}
\end{equation*}
$$

then $\mathscr{P}_{A}^{*}\left(e^{\lambda t} \psi\right) \in L_{\eta}^{2}\left(Y_{a+}\right)$, and $K$ is well-defined.
We put all these maps in a diagram

$$
\begin{array}{cl}
W_{\delta}^{1,2}\left(Y_{a+}\right) & \xrightarrow{\not D^{*}} \quad L_{\delta}^{2}\left(Y_{a+}\right) \oplus \\
{ }_{J} \uparrow & \Pi_{\delta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right)  \tag{7.4}\\
W_{\eta}^{1,2}\left(Y_{a+}\right) \oplus W_{\lambda} \xrightarrow[K]{ } & L_{\eta}^{2}\left(Y_{a+}\right) \oplus \Pi_{\eta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right) \oplus W_{\lambda}
\end{array}
$$

which is commutative as

$$
\begin{aligned}
\mathscr{D}^{*} J(\phi, \psi) & =\left(\mathfrak{D}_{A}^{*}\left(\phi+e^{\lambda t} \psi\right), \Pi_{\delta}^{-} \phi(a)+e^{a \lambda} \psi\right) \\
& =\left(\mathfrak{D}_{A}^{*}\left(\phi+e^{\lambda t} \psi\right), \Pi_{\eta}^{-} \phi(a)+\Pi_{\lambda} \phi(a)+e^{a \lambda} \psi\right) \\
& =\left(\mathfrak{D}_{A}^{*}\left(\phi+e^{\lambda t} \psi\right), \Pi_{\eta}^{-} \phi(a)+e^{a \lambda}\left(\psi+e^{-a \lambda} \Pi_{\lambda} \phi(a)\right)\right) \\
& =(\iota \oplus I) K(\phi, \psi) .
\end{aligned}
$$

Now that we know that the diagram is commutative, we want to exploit the fact that its rows are isomorphisms. While Theorem 7.1-1 assures us that $\mathscr{D}^{*}$ is an isomorphism, we still have to prove
that $K$ is one as well. Using the identification

$$
\mathscr{D}^{*}: W_{\eta}^{1,2}\left(Y_{a+}\right) \equiv L_{\eta}^{2}\left(Y_{a+}\right) \oplus \Pi_{\eta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right),
$$

we see that $K$ has the form

$$
\left[\begin{array}{ll}
1 & p \\
q & 1
\end{array}\right]
$$

for the splitting $W_{\eta}^{1,2}\left(Y_{a+}\right) \oplus W_{\lambda}$ of the domain and codomain. Hence $K-1$ is a compact operator, and $K$ is thus Fredholm of index 0 . If $K(x)=K(y)$, then $\not D^{*} J(x)=\not D^{*} J(y)$ as the diagram is commutative, hence $x=y$ and $K$ is injective. Being of index 0 , it henceforth must be an isomorphism.
Let's now exploit this fantastic diagram. Suppose

$$
\phi \in \operatorname{ker}\left(\mathfrak{P}_{A}^{*}\right) \cap W_{\delta}^{1,2}(\mathbb{R} \times Y) .
$$

Then for $a$ big enough, the diagram (7.4) has rows which are isomorphism for $\delta$ close to the first negative eigenvalue of $D_{\Gamma_{+}}$and $\eta$ past it, and satisfying condition (7.3). Theorem 6.3-2 guarantees us that $\phi \in W_{\delta}^{1,2}\left(Y_{a+}\right)$ for that particular $\delta$.
We now chase around the diagram. Since $I$ is an isomorphism, we know there exist $(\chi, \nu) \in$ $\Pi_{\eta}^{-} W^{\frac{1}{2}, 2}\left(Y_{a}\right) \oplus W_{\lambda}$ such that

$$
\iota \oplus I(0, \chi, \nu)=\not D^{*}(\phi) .
$$

But as $K$ is an isomorphism, there is $(\bar{\phi}, \bar{\psi}) \in W_{\eta}^{1,2}\left(Y_{a+}\right) \oplus W_{\lambda}$ such that

$$
K(\bar{\phi}, \bar{\psi})=(0, \chi, \nu) .
$$

By commutativity of the diagram, we have

$$
\mathscr{D}^{*} J(\bar{\phi}, \bar{\psi})=\not D^{*}(\phi)
$$

but $\mathscr{D D}^{*}$ is an isomorphism hence $\phi=e^{\lambda t} \bar{\psi}+\bar{\phi}$ for $t>a$.
Of course, at this point the choice of $a$ is artificial and we can choose $a=0$. We hence proved the following result.

Theorem 7.2-1. Suppose $\phi \in \operatorname{ker}\left(\mathfrak{D}_{A}^{*}\right) \cap W_{\delta}^{1,2}(\mathbb{R} \times Y)$. Suppose $\lambda-\beta<\eta<\delta$ and that $\lambda$ is the only eigenvalue of $D$ between $\eta$ and $\delta: \operatorname{Spec}(D) \cap[\eta, \delta]=\{\lambda\}$. Then there exist $\bar{\psi} \in W_{\lambda}$ and $\bar{\phi} \in W_{\eta}^{1,2}\left(Y^{+}\right)$such that

$$
\begin{equation*}
\phi=e^{\lambda t} \bar{\psi}+\bar{\phi} \text { for } t>0 \tag{7.5}
\end{equation*}
$$

Furthermore, $\bar{\phi}=O\left(e^{\eta t}\right)$ as $t \rightarrow \infty$.

### 7.3 Asymptotic on $\mathbb{R}^{4}$

In the spirit of this chapter, we want to study $\mathbb{R}^{4}$ as a cylindrical manifold. Let's then use the conformal equivalence

$$
\begin{aligned}
\mathbb{R} \times S^{3} & \rightarrow \mathbb{R}^{4} \backslash\{0\} \\
(t, x) & \mapsto e^{t} x
\end{aligned}
$$

The respective metrics of those spaces are related by the formula

$$
g_{\mathbb{R}^{4} \backslash\{0\}}=\left(e^{2 t}\right) g_{\mathbb{R} \times S^{3}} \quad \text { or } \quad g_{\mathbb{R}^{4} \backslash\{0\}}=|x|^{2} g_{\mathbb{R} \times S^{3}} .
$$

It then follows that $d v o l_{\mathbb{R}^{4} \backslash\{0\}}=|x|^{4} d v o l_{\mathbb{R} \times S^{3}}$, hence

$$
\|\phi\|_{L^{2}\left(\mathbb{R}^{4} \backslash\{0\}\right)}=\left\||x|^{2} \phi\right\|_{L^{2}\left(\mathbb{R} \times S^{3}\right)} .
$$

The spinor bundles of $\mathbb{R} \times S^{3}$ and $\mathbb{R}^{4} \backslash\{0\}$ are isomorphic. Once we fix a spinor bundle to work with, we can compare the Dirac operators given for the two metrics. The correct relation, as seen in Appendix (D is

$$
D_{\mathbb{R}^{4} \backslash\{0\}}=|x|^{-5 / 2} D_{\mathbb{R} \times S^{3}}|x|^{3 / 2} .
$$

Thus

$$
\begin{equation*}
D_{\mathbb{R}^{4} \backslash\{0\}} \phi=0 \quad \text { iff } \quad D_{\mathbb{R} \times S^{3}}\left(|x|^{3 / 2} \phi\right)=0 . \tag{7.6}
\end{equation*}
$$

Let $p$ be the projection $\mathbb{R} \times S^{3} \rightarrow S^{3}$. Let $S\left(S^{3}\right)$ be the spinor bundle of $S^{3}$. Set $S^{+}$and $S^{-}$to be $p^{*}\left(S\left(S^{3}\right)\right)$. The spinor bundle on $\mathbb{R} \times S^{3}$ is $S^{+} \oplus S^{-}$.
The Clifford multiplication exchanges $S^{+}$and $S^{-}$. For vectors tangent to $S^{3}$, the Clifford multiplication is already defined. The vector $\partial / \partial t$ acts as $i d: S^{+} \rightarrow S^{-}$and $-i d: S^{-} \rightarrow S^{+}$.
In this decomposition, the Dirac operator splits nicely:

$$
\begin{equation*}
D_{\mathbb{R} \times S^{3}}^{ \pm}= \pm \frac{\partial}{\partial t}+D_{S^{3}} . \tag{7.7}
\end{equation*}
$$

We use now the knowledge of the eigenvalues of the $D_{S^{3}}$ on $S^{3}$ obtained in Theorem 4.1-3 to understand the asymptotic expansion of solutions $\phi$ to the equation

$$
D_{\mathbb{R}^{4} \backslash\{0\}}^{-} \phi=0
$$

under the constraint of being $L^{2}$.
We have here a basis of the kernel of $D_{\mathbb{R}^{4} \backslash\{0\}}^{-}$. Indeed, if

$$
D_{S^{3}} \psi_{\lambda}=\lambda \psi_{\lambda}
$$

then, as suggested by Equation (7.7), we have

$$
D_{\mathbb{R} \times S^{3}}^{-}\left(e^{\lambda t} \psi_{\lambda}\right)=0
$$

and thus, because of the conformal relation 7.6, we have

$$
D_{\mathbb{R}^{4} \backslash\{0\}}^{-}\left(|x|^{\lambda-3 / 2} \psi_{\lambda}\right)=0 .
$$

Let's now use the notation of Chapter $\square$ Hence $A$ is an instanton connection on a bundle $E$ over $\mathbb{R}^{4}$ and $V_{E}=L^{2}\left(\mathbb{R}^{4}, S^{-} \otimes E\right) \cap \operatorname{ker}\left(D_{A}^{-}\right)$. Theorem 7.2-1 then tells us any $\phi \in V_{E}$ has an asymptotic behavior

$$
\phi=|x|^{-3} \psi_{-3 / 2}+O\left(|x|^{-4}\right) .
$$

Theorem 4.1-3 tells us the space of possible $\psi_{-3 / 2}$ has dimension $2 \operatorname{dim}(E)$. We build this space using parallel sections of $S^{+} \otimes E$ for the trivial connection on $E$. Let $a \in \Gamma\left(S^{+} \otimes E\right)$ be parallel.

Consider the section $\phi_{a}:=\rho(\nu) a / r^{3}$ of $S^{-} \otimes E$. In coordinates, we have

$$
\phi_{a}=\sum_{i=1}^{4} \frac{x_{i}}{r^{4}} \rho\left(\partial_{i}\right) a
$$

and we compute

$$
\begin{aligned}
D_{\mathbb{R}^{4} \backslash\{0\}} \phi_{a} & =\sum_{1 \leq i, j \leq 4} \rho\left(\partial_{j}\right) \partial_{j}\left(\frac{x_{i}}{r^{4}}\right) \rho\left(\partial_{i}\right) a \\
& =\sum_{1 \leq i, j \leq 4}\left(\frac{\delta_{i j} r^{4}-4 x_{i} x_{j} r^{2}}{r^{8}}\right) \rho\left(\partial_{j}\right) \rho\left(\partial_{i}\right) a \\
& =-\left(4 r^{-4}-4 r^{-4}\right) a \\
& =0 .
\end{aligned}
$$

But then, formula 4.2 implies that

$$
\begin{aligned}
D_{r}\left(\phi_{a}\right) & =-\rho(\nu) \frac{\partial \phi_{a}}{\partial r}-\frac{3}{2 r} \rho(\nu) \phi_{a} \\
& =-\frac{3}{r^{4}} a+\frac{3}{2 r^{4}} a \\
& =-\frac{3}{2} a / r^{4} \\
& =\rho(-\nu)\left(-\frac{3}{2} \phi / r\right) .
\end{aligned}
$$

Recall now that $S^{+}$and $S^{-}$are actually pullbacks of the spinor bundle of $S^{3}$. In this setting, $\rho(\nu): S^{-} \rightarrow S^{+}$is $-i d$, as explained in the beginning of this section. So restricting to $r=1$, we really find

$$
D_{S^{3}} \phi=-\frac{3}{2} \phi .
$$

So far we proved that for any $\phi \in V_{E}$, we have

$$
\phi=|x|^{4} \rho(x) a+O\left(|x|^{-4}\right)
$$

for a parallel section $a$ of $S^{+} \otimes E$, parallel that is for the trivial connection on $E$. This result is not exactly Equation (1.7), but leads to it. Indeed, the same analysis we did works for the Laplacian. Hence parallel sections of $S^{+} \otimes E$ for $A$ or for the trivial connection are the same to leading order, hence we have

$$
\phi=|x|^{4} \rho(x) \hat{\phi}+O\left(|x|^{-4}\right)
$$

for some $\hat{\phi} \in W_{E}$, and we proved Equation (1.7).

## Chapter 8

# Nahm Transform: Instantons to singular monopoles 

"It doesn't matter what you write as long as you write the truth. Then we can figure out what it means."<br>Tomasz S. Mrowka

Following the heuristic of Chapter 2, we show in this chapter that the Nahm transform

$$
\mathfrak{N}(E, A)=(V, B, \Phi)
$$

of a $S U(2)$-instanton $(E, A)$ on $\mathbb{R} \times T^{3}$ is a singular monopole $(V, B, \Phi)$ over over $T^{3}$.
As we found out in Chapter [5] once in a temporal gauge, the connection $A$ has limiting flat connections over the cross-section $T^{3}$ at $+\infty$ and $-\infty$, say

$$
\lim _{t \rightarrow \pm \infty} A=\Gamma_{ \pm}
$$

The flat connection $\Gamma_{ \pm}$gives a splitting $L_{w_{ \pm}} \oplus L_{-w_{ \pm}}$of the restriction of $E$ to $T^{3}$ at the infinities, for some $w_{ \pm} \in T^{3}$. Let $W$ denote the set

$$
W:=\left\{w_{+},-w_{+}, w_{-},-w_{-}\right\} .
$$

As before, we denote $A_{z}$ the connection on $E \otimes L_{z}$. We consider the Dirac operator

$$
\mathscr{P}_{A_{z}}^{*}: L^{2} \rightarrow L^{2} .
$$

Outside of $W$, Theorem 6.1-5 guarantees that $\mathscr{P}_{A_{z}}^{*}$ is Fredholm. Since $\operatorname{ker}\left(\mathscr{P}_{A_{z}}\right)=0$ as $F_{A_{z}}$ is ASD, we have a bundle $V$ over $T^{3} \backslash W$ whose fiber at $z$ is

$$
V_{z}:=\operatorname{ker}\left(\mathfrak{D}_{A_{z}}^{*}\right) \cap L^{2} .
$$

By a gauge transformation, we can make the connection $P d^{z}$ independent of the $\mathbb{R}$ factor. We can
thus see it as

$$
\begin{aligned}
& \text { a connection } B \text { on } T^{3} \backslash W, \\
& \text { a Higgs field } \Phi \in \Gamma\left(T^{3} \backslash W, \operatorname{End} V\right) \text {. }
\end{aligned}
$$

The main result of this present thesis is the following theorem.
Theorem 8.0-1. Outside of a set $W$ consisting of at most four points, the family of vector spaces $V$ described above defines a vector bundle of rank

$$
\frac{1}{8 \pi^{2}} \int\left|F_{A}\right|^{2},
$$

and the couple $(B, \Phi)$ satisfies the Bogomolny equation

$$
\nabla_{B} \Phi=* F_{B} .
$$

For $w \in W$ and $z$ close enough to $w$, unless we are in the Scenario 2 of page 91 there are maps $\Phi^{\perp}$ and $\Phi^{\lrcorner}$such that

$$
\Phi=\frac{-i}{2|z-w|} \Phi^{\perp}+\Phi^{\lrcorner},
$$

and $\Phi^{\perp}$ is the $L^{2}$-orthogonal projection on the orthogonal complement of a naturally defined subbundle $V$ of $V$.

The last part of the theorem is made clearer by the introduction of some notation in Section 8.2 The assumption that we are not in the Scenario 2 of page 91 can most probably be dropped.

Proof: The rank of $V$ is computed in Section 8.1
The limit term $\lim _{\partial}$ of Equation (2.2) is

$$
\lim _{\partial}=\left.\left\langle\nu \Omega G \phi, d^{z} \psi\right\rangle_{T^{3}}\right|_{-\infty} ^{\infty}
$$

For $z \notin W$, both $G \phi$ and $d^{z} \psi$ decay exponentially by Equation (7.5) hence

$$
\lim _{\partial}=0
$$

and the connection $P d^{z}$ on $\mathbb{R} \times\left(T^{3} \backslash W\right)$ is ASD. Thus, as explained in Chapter 2 the pair $(B, \Phi)$ satisfies outside of $W$ the appropriate dimensional reduction of the ASD equation, which is in this case the Bogomolny Equation (A.3):

$$
\nabla_{B} \Phi=* F_{B} .
$$

The last part of the theorem is the content of Section 8.3 and rest on the splitting of Section 8.2

### 8.1 An $L^{2}$-index theorem for $\mathbb{R} \times T^{3}$

The following theorem is reminiscent of the similar result for $\mathbb{R}^{4}$.
Theorem 8.1-1. For a $S U(2)$-instanton $(E, A)$ on $\mathbb{R} \times T^{3}$, the index of the Dirac operator

$$
\mathscr{D}_{A}: W^{1,2}\left(\mathbb{R} \times T^{3}\right) \rightarrow L^{2}\left(\mathbb{R} \times T^{3}\right)
$$

when $A$ has nonzero limits at $\pm \infty$ is given by the formula

$$
\operatorname{ind}\left(\mathscr{D}_{A}\right)=-\frac{1}{8 \pi^{2}} \int\left|F_{A}\right|^{2} .
$$

Proof: As seen earlier, the fact that $A$ has nonzero limits guarantees that the operator $\mathscr{D}_{A}$ is Fredholm on $W^{1,2}$. Moreover, $A$ decays exponentially to its limits.
Recall now that $\operatorname{ind}\left(\mathscr{D}_{A}\right)=\operatorname{ind}\left(\mathscr{D}_{a_{R}}\right)$ for all $R>0$. We now compute $\operatorname{ind}\left(\mathscr{D}_{a_{R}}\right)$ using the relative index theorem. It could be that $\Gamma_{-} \neq \Gamma_{+}$, but this case is easily converted to a situation where $\Gamma_{-}=\Gamma_{+}$, as we now see.
Choose a path $\Gamma_{s}$ in the space of flat connections on $T^{3}$ starting at $\Gamma_{+}$and ending at $\Gamma_{-}$, and avoiding the trivial connection. Hence $0 \notin \operatorname{Spec}\left(D_{\Gamma_{s}}\right)$ for all $s$. To define the family of connections $a_{R}^{s}$, replace $\Gamma_{+}$by $\Gamma_{s}$ in the definition of $a_{R}$ given by Equation (6.3).
The family $\mathscr{P}_{a_{R}^{s}}$ of Fredholm operator depends continuously on $s$. Hence

$$
\operatorname{ind}\left(\mathscr{D}_{A}\right)=\operatorname{ind}\left(\mathscr{D}_{a_{R}}\right)=\operatorname{ind}\left(\mathscr{P}_{a_{R}^{0}}\right)=\operatorname{ind}\left(\mathscr{D}_{a_{R}^{1}}\right) .
$$

Note now that the connection $a_{R}^{1}$ equals $\Gamma_{-}$outside $[-R-1, R+1] \times T^{3}$. Hence the relative index theorem tells us

$$
\begin{equation*}
\operatorname{ind}\left(\mathscr{P}_{a_{R}^{1}}\right)-\operatorname{ind}\left(\mathscr{P}_{\Gamma_{-}}\right)=\operatorname{ind}\left(\tilde{\mathscr{P}}_{a_{R}^{1}}\right)-\operatorname{ind}\left(\tilde{\mathscr{P}}_{\Gamma_{-}}\right), \tag{8.1}
\end{equation*}
$$

where the tilded operators are extensions to some compact manifold of the restriction of the operators $\mathscr{D}_{a_{R}^{1}}$ and $\mathscr{D}_{\Gamma_{-}}$to $[-R-1, R+1] \times T^{3}$.
Lemma 6.1-1 and Theorem 3.4-1 tell us that ind $\left(\mathscr{D}_{\Gamma_{-}}\right)=0$. Hence the left-hand-side of Equation (8.1) is equal to $\operatorname{ind}\left(\mathscr{D}_{A}\right)$.

To compute the right-hand-side, we embed $[-R-1, R+1] \times T^{3}$ in some flat $T^{4}$. The spinor bundles $S^{+}$and $S^{-}$on $[-R-1, R+1] \times T^{3}$ agree very nicely with those of $T^{4}$. We extend both $a_{R}^{1}$ and $\Gamma_{-}$by the trivial bundle with connection $\Gamma_{-}$.
The Atiyah-Singer index theorem (see [Roe98, Thm 12.27, p.164] or [LM89, Thm III.12.10, p. 256]) tells us that

$$
\begin{aligned}
\operatorname{ind}\left(\tilde{\mathscr{D}}_{\Gamma_{-}}\right) & =\left\{\operatorname{ch}\left(\Gamma_{-}\right) \cdot \hat{\mathbf{A}}\left(T^{4}\right)\right\}\left[T^{4}\right] \\
\operatorname{ind}\left(\tilde{\mathscr{D}}_{a_{R}^{1}}\right) & =\left\{\operatorname{ch}\left(a_{R}^{1}\right) \cdot \hat{\mathbf{A}}\left(T^{4}\right)\right\}\left[T^{4}\right] \\
& =\left(\frac{c_{1}^{2}}{2}-c_{2}\right)\left[T^{4}\right] .
\end{aligned}
$$

Since $a_{R}^{1}$ is in $S U(2)$, we have $c_{1}=0$, while

$$
c_{2}\left[T^{4}\right]=\frac{1}{8 \pi^{2}} \int_{T^{4}}\left|F_{a_{R}^{1}}\right|^{2} .
$$

Note that on the complement of $[-R-1, R+1] \times T^{3}$ in $T^{4}$, the connection $a_{R}^{1}$ equals $\Gamma_{-}$hence is flat there. Furthermore, on $[-R, R] \times T^{3}$, we have $a_{R}^{1}=A$. On $[R, R+1] \times T^{3}$ and $[-R-1,-R] \times T^{3}$, the curvature $F_{a_{R}^{1}}$ involves cut off functions, their derivatives and $\left(A-\Gamma_{-}\right)$terms. Since $A$ tends to $\Gamma_{-}$exponentially fast, we therefore have constant $C$ and $\beta$ such that

$$
\left.\left.\left|\operatorname{ind}\left(\mathscr{D}_{A}\right)+\frac{1}{8 \pi^{2}} \int_{[-R, R] \times T^{3}}\right| F_{A}\right|^{2} \right\rvert\, \leq C e^{-\beta R} .
$$

As $R \rightarrow \infty$, we have the wanted result.
Now suppose $(E, A)$ is a $S U(2)$-instanton on $\mathbb{R} \times T^{3}$. As mentioned before,

$$
\lim _{t \rightarrow \pm \infty}=\Gamma_{ \pm} .
$$

The flat connection $\Gamma_{ \pm}$gives a splitting $L_{w_{ \pm}} \oplus L_{-w_{ \pm}}$, for some $w_{ \pm} \in \Lambda^{*}$, of the bundle $E$ at $\pm \infty$ respectively.
We twist the connection $A$ by the flat connection parameterized by $z \in T^{3}$. Hence

$$
\begin{aligned}
& \operatorname{Sppcc}\left(\Gamma_{+} \otimes L_{z}\right)= \pm 2 \pi\left|\Lambda^{*}-z+w_{+}\right| \cup \pm 2 \pi\left|\Lambda^{*}-z-w_{+}\right|, \\
& \operatorname{Sppec}\left(\Gamma_{-} \otimes L_{z}\right)= \pm 2 \pi\left|\Lambda^{*}-z+w_{-}\right| \cup \pm 2 \pi\left|\Lambda^{*}-z-w_{-}\right| .
\end{aligned}
$$

Thus $\mathscr{D}_{A_{z}}$ is Fredholm as long as $z \pm w_{+} \notin \Lambda^{*}$ and $z \pm w_{-} \notin \Lambda^{*}$. Moreover, when it is Fredholm, the elements of its $L^{2}$-kernel decay exponentially.

### 8.2 A Geometric Splitting and Exact Sequences

In this section, we analyse a splitting of $V$ in a neighborhood of a point $w \in W$ where the solution $(B, \phi)$ to Bogomolny equation is singular. This point $w$ is associated to the limit $\Gamma=\Gamma_{+}$of $A$ at, say, $+\infty$, in the sense that $\Gamma$ splits $E$ as $L_{w} \oplus L_{-w}$ on $T^{3}$.
Suppose the connection $A$ decays at most with rate $\beta$, as in $\left|A-\Gamma_{+}\right| \leq C e^{-\beta t}$ for $t>0$ and $\left|A-\Gamma_{-}\right| \leq C e^{\beta t}$ for $t<0$. Set

$$
\epsilon:=\frac{1}{4} \min \left(\beta, \operatorname{dist}\left(w, \Lambda^{*}+W \backslash\{w\}\right)\right),
$$

and define the six weights

$$
\begin{array}{lll}
\ulcorner:=(-\epsilon, \epsilon) & \bar{\epsilon}:=(0, \epsilon) & \epsilon:=(\epsilon, \epsilon) \\
\epsilon:=(-\epsilon,-\epsilon) & \underline{\epsilon}:=(0,-\epsilon) & \epsilon\lrcorner:=(\epsilon,-\epsilon)
\end{array}
$$

displayed here in a way which is reminiscent of their position in $\mathbb{R}^{2}$.
Consider the ball $B^{3}(w)$ of radius $2 \epsilon$ around $w$. As $z$ varies in $B^{3}(w)$, and depending on whether $\Gamma_{+}=\Gamma_{-}$or not, there are two or one walls to cross to pass from 0 to ${ }^{\epsilon} \epsilon$ and from $\epsilon_{\lrcorner}$to 0 . In a picture, we have


As $z$ varies in $B^{3}(w)$, those walls move around without ever touching $\epsilon_{\lrcorner}$and ${ }_{\epsilon} \epsilon$. Hence for $L_{\epsilon}^{2}$ and $L_{\Gamma \in}^{2}$, the operators $\mathscr{D}_{A_{z}}, \mathscr{D}_{A_{z}}^{*}$ and $\mathscr{D}_{A_{z}}^{*} \mathscr{D}_{A_{z}}$ are Fredholm for all $z \in B^{3}(w)$.

Hence for $z \in B^{3}(w)$, the six vector spaces

$$
\begin{array}{rlrl}
V_{z}:=\operatorname{ker}\left(\mathscr{D}_{A_{z}}^{*}\right) \cap L_{\epsilon}^{2}, & \bar{K}_{z}:=\operatorname{ker}\left(\mathscr{D}_{A_{z}}\right) \cap L_{\Gamma_{\epsilon}}^{2}, & \mathcal{H}_{z}:=\operatorname{ker}\left(\nabla_{A_{z}}^{*} \nabla_{A_{z}}\right) \cap L_{\Gamma_{\epsilon}}^{2}, \\
V_{z}:=\operatorname{ker}\left(\mathscr{D}_{A_{z}}^{*}\right) \cap L_{\epsilon}^{2}, & K_{z z}:=\operatorname{ker}\left(\mathscr{D}_{A_{z}}\right) \cap L_{\epsilon}^{2}, & & K_{z}:=\operatorname{ker}\left(\mathscr{D}_{A_{z}}\right) \cap L^{2}
\end{array}
$$

are kernels of Fredholm operators. By contrast, the space $V_{z}$, already defined as $\operatorname{ker}\left(\mathscr{P}_{A_{z}}^{*}\right) \cap L^{2}$, is not the kernel of a Fredholm operator at $w$.
Notice that none of those vector space form a priori a bundle over $B^{3}(w)$ as the dimensions could jump at random. However, for $L_{\epsilon_{\epsilon}}^{2}$ and $L_{\epsilon}^{2}$, the operators $\mathscr{D}_{A_{z}}, \mathscr{P}_{A_{z}}^{*}$, and $\nabla_{A_{z}}^{*} \nabla_{A_{z}}$ are Fredholm operators for all $z \in B^{3}(w)$. The various indices are therefore constant and we have that, for example,

$$
\operatorname{dim} V_{z}-\operatorname{dim} K_{z} \text { is constant on } B^{3}(w) .
$$

We have the following obvious results:

$$
\begin{array}{lr}
V \subset V \subset V, & K \subset K \subset K, \\
\not \supset \mathcal{H} \subset V, & \Gamma \subset \mathcal{H}, \\
& K=K=\{0\} .
\end{array}
$$

It was remarked on page 73 that $V_{w}=V_{w}$. The following few lemmas describe in more detail the relationship between the various spaces.
We saw in Section 3.4 that the smallest eigenvalues of $D_{\Gamma_{z}}$ are $\pm 2 \pi|z-w|$. For simplicity, we set

$$
\lambda:=2 \pi|z-w|,
$$

and define

$$
W_{\lambda}:=\lambda \text { eigenspace of } D_{\Gamma_{z}} \text { on } T^{3} .
$$

The family $W_{\lambda}$ defines a bundle over the sphere $|z-w|=\lambda / 2 \pi$ around $w$. Its rank is given by

$$
\operatorname{rk} W_{\lambda}= \begin{cases}1, & \text { if } \lambda \neq 0 \text { and } 2 w \notin \Lambda^{*} ;  \tag{8.2}\\ 2, & \text { if } \lambda \neq 0 \text { and } 2 w \in \Lambda^{*}, \text { or } \lambda=0 \text { and } 2 w \notin \Lambda^{*} ; \\ 4, & \text { if } \lambda=0 \text { and } 2 w \in \Lambda^{*} .\end{cases}
$$

As suggested by Theorem 7.2-1 this $W_{\lambda}$ plays an important role in understanding the relations between the various spaces just introduced.
For any instanton connection $A^{\prime}$ on $\mathbb{R} \times T^{3}$, set

$$
\begin{aligned}
V(\delta) & :=\operatorname{ker}\left(\mathfrak{P}_{A^{\prime}}^{*}\right) \cap L_{\delta}^{2}, \\
K(\delta) & :=\operatorname{ker}\left(\mathfrak{P}_{A^{\prime}}\right) \cap L_{\delta}^{2},
\end{aligned}
$$

and let $[\delta]$ denote the open square in $\mathbb{R}^{2} \backslash \mathfrak{G}_{A^{\prime}}$ containing $\delta$.

Lemma 8.2-1 (one wall). Suppose $\delta, \eta \in \mathbb{R}^{2} \backslash \mathfrak{G}_{A^{\prime}}$ are weights for which $[\delta]$ and $[\eta]$ are adjacent and separated by the wall $\{\mu\} \times \mathbb{R}$ or $\mathbb{R} \times\{\mu\}$. Then the sequence

$$
\begin{equation*}
0 \longrightarrow V(\delta) \longrightarrow V(\eta) \xrightarrow{\lim \left(e^{-\mu t .}\right)} W_{\mu} \xrightarrow{\left(\lim \left(e^{\mu t} .\right)\right)^{*}} K(-\delta)^{*} \longrightarrow K(-\eta)^{*} \longrightarrow 0, \tag{8.3}
\end{equation*}
$$

where the limits are both evaluated at $+\infty$ when $[\eta]$ is above $[\delta]$ and at $-\infty$ when $[\eta$ ] is to the left of $[\delta]$, is exact.

Proof: Theorem 7.2-1 ensures that the limits give functions $\alpha$ and $\beta^{*}$ which are well defined, and that

$$
0 \longrightarrow V(\delta) \longrightarrow V(\eta) \longrightarrow W_{\mu} \quad \text { and } \quad 0 \longrightarrow K(-\eta) \longrightarrow K(-\delta) \longrightarrow W_{\mu}
$$

are exact.
It only remains to prove that Sequence (8.3) is exact at $W_{\mu}$. Suppose $\phi \in V(\eta)$ and $\psi \in K(-\eta)$. Then

$$
\begin{aligned}
0 & =\left\langle\mathscr{D}_{A^{\prime}}^{*} \phi, \psi\right\rangle-\left\langle\phi, \mathscr{D}_{A^{\prime}} \psi\right\rangle \\
& =\lim _{t \rightarrow \infty}\langle\phi, \nu \psi\rangle-\lim _{t \rightarrow-\infty}\langle\phi, \nu \psi\rangle \\
& =\lim _{t \rightarrow \infty}\left\langle e^{-\mu t} \phi, \nu e^{\mu t} \psi\right\rangle-\lim _{t \rightarrow-\infty}\left\langle e^{-\mu t} \phi, \nu e^{\mu t} \psi\right\rangle
\end{aligned}
$$

One of those limits is $\beta^{*} \alpha(\phi)(\psi)$ while the other one vanishes as we now see. Suppose $[\eta]$ is above $[\delta]$, and suppose $\left\{\mu^{\prime}\right\} \times \mathbb{R}$ is the wall to their right. Then $\phi=O\left(e^{\mu^{\prime} t}\right)$ as $t \rightarrow-\infty$ by Theorem 7.2-1 But for some $\mu^{\prime \prime}<\mu^{\prime}$, the wall $\left\{-\mu^{\prime \prime}\right\} \times \mathbb{R}$ is exactly to the right of $[-\eta]$ hence $\psi=O\left(e^{-\mu^{\prime \prime} t}\right)$ as $t \rightarrow-\infty$. But then

$$
\beta^{*} \alpha(\phi)(\psi)=\lim _{t \rightarrow-\infty} O\left(e^{\left(\mu^{\prime}-\mu^{\prime \prime}\right) t}\right)=0
$$

hence $\operatorname{Im}(\alpha) \operatorname{ker}\left(\beta^{*}\right)$. A similar argument establish the same fact when $[\eta]$ is to the left of $[\delta]$.
The sequence is then exact if $\operatorname{dim} \operatorname{Im}(\alpha)=\operatorname{dim} \operatorname{ker}\left(\beta^{*}\right)$. We have two short exact sequences:

$$
\begin{gathered}
0 \longrightarrow V(\delta) \longrightarrow V(\eta) \longrightarrow \operatorname{Im}(\alpha) \longrightarrow 0, \quad \text { and } \\
0 \longrightarrow W_{\mu} / \operatorname{ker}\left(\beta^{*}\right) \longrightarrow K(-\delta)^{*} \longrightarrow K(-\eta)^{*} \longrightarrow 0
\end{gathered}
$$

Using those short exact sequences and notation from Equations 6.10, we have

$$
\begin{aligned}
\operatorname{dim} \operatorname{Im}(\alpha)-\operatorname{dim} \operatorname{ker}\left(\beta^{*}\right) & =N^{*}(\eta)-N^{*}(\delta)-\operatorname{dim} W_{\mu}+N(-\delta)-N(-\eta) \\
& =\operatorname{ind}^{*}(\eta)-\operatorname{ind}^{*}(\delta)-\operatorname{dim} W_{\mu}
\end{aligned}
$$

The Wall Crossing Theorem 6.4-1 forces the last line to be 0 . The proof is thus complete.
Lemma 8.2-2. Suppose $\Gamma_{+} \neq \Gamma_{-}$. Then the sequences

$$
\begin{array}{r}
0 \longrightarrow V_{z} \longrightarrow V_{z} \longrightarrow W_{\lambda} \longrightarrow 0, \quad \text { for } \lambda \neq 0 \\
0 \longrightarrow V_{z} \longrightarrow V_{z} \longrightarrow W_{-\lambda} \longrightarrow K_{z} \longrightarrow 0, \quad \text { for } \lambda \neq 0 \\
0 \longrightarrow V_{w} \longrightarrow V_{w} \longrightarrow W_{0} \longrightarrow \Gamma_{w} \longrightarrow 0, \tag{8.6}
\end{array}
$$

are exact.
Proof: Apply Lemma 8.2-1 to the choice of weights $\left\{\left\ulcorner_{\epsilon}, 0\right\}\right.$ and $\left\{0, \epsilon_{\lrcorner}\right\}$for the connection $A^{\prime}=A_{z}$, and remember that $K_{\lrcorner}=K=\{0\}$.

Building up on that knowledge, we work out in Appendix Ca technology used to deal with the two walls involved in passing from $V$ to $V$ in the case $\Gamma_{+}=\Gamma_{-}$. The result of the beautiful abstract non-sense taking place there is summarized in the following lemma, which should be compared to Lemma8.2-2

Lemma 8.2-3. Suppose $\Gamma_{+}=\Gamma_{-}$. Then the sequences

$$
\begin{align*}
0 \longrightarrow V_{z} \longrightarrow V_{z} \longrightarrow W_{\lambda} \oplus W_{-\lambda} \longrightarrow 0, & \text { for } \lambda \neq 0,  \tag{8.7}\\
0 \longrightarrow V_{z} \longrightarrow V_{z} \longrightarrow W_{\lambda} \oplus W_{-\lambda} \longrightarrow K_{z} \longrightarrow 0, & \text { for } \lambda \neq 0,  \tag{8.8}\\
0 \longrightarrow V_{w} \longrightarrow V_{w} \longrightarrow W_{0} \oplus W_{0} \longrightarrow K_{w} \longrightarrow 0, & \tag{8.9}
\end{align*}
$$

are exact.

## Proof: See Appendix C

An analysis for $\nabla_{A_{z}}^{*} \nabla_{A_{z}}$ parallel to the one of Chapter for $\mathscr{D}_{A_{z}}$ brings a very similar wall crossing formula

$$
\operatorname{ind}\left(\nabla_{A_{z}}^{*} \nabla_{A_{z}}, \Gamma_{\epsilon}\right)-\operatorname{ind}\left(\nabla_{A_{z}}^{*} \nabla_{A_{z}}, \epsilon_{\lrcorner}\right)= \begin{cases}2 \operatorname{dim} W_{0}, & \text { for } \Gamma_{+} \neq \Gamma_{-} ; \\ 4 \operatorname{dim} W_{0}, & \text { for } \Gamma_{+}=\Gamma_{-} .\end{cases}
$$

However, since $\nabla_{A_{z}}^{*} \nabla_{A_{z}}$ is self-adjoint, $\operatorname{ind}\left(\nabla_{A_{z}}^{*} \nabla_{A_{z}}, r_{\epsilon}\right)=-\operatorname{ind}\left(\nabla_{A_{z}}^{*} \nabla_{A_{z}}, \epsilon_{\lrcorner}\right)$, whence

$$
\operatorname{rk} \mathcal{H}= \begin{cases}\operatorname{dim} W_{0}, & \text { for } \Gamma_{+} \neq \Gamma_{-} ; \\ 2 \operatorname{dim} W_{0}, & \text { for } \Gamma_{+}=\Gamma_{-}\end{cases}
$$

Using Equation (8.2), we can even say

$$
\operatorname{rk\mathcal {H}}= \begin{cases}2, & \text { for } \Gamma_{+} \neq \Gamma_{-} \text {and } 2 w \notin \Lambda^{*} ; \\ 4, & \text { for } \Gamma_{+} \neq \Gamma_{-} \text {and } 2 w \in \Lambda^{*}, \text { or } \Gamma_{+}=\Gamma_{-} \text {and } 2 w \notin \Lambda^{*} ; \\ 8, & \text { for } \Gamma_{+}=\Gamma_{-} \text {and } 2 w \in \Lambda^{*} .\end{cases}
$$

For $z \neq w$, an analysis parallel to the one of Chapter Zgives injective maps

$$
\begin{align*}
0 \longrightarrow \mathcal{H}_{z} \longrightarrow W_{\lambda} \oplus W_{-\lambda} \longrightarrow 0, & \text { for } z \neq w \text { and when } \Gamma_{+} \neq \Gamma_{-},  \tag{8.10}\\
0 \longrightarrow \mathcal{H}_{w} \longrightarrow W_{0} \longrightarrow 0, & \text { when } \Gamma_{+} \neq \Gamma_{-},  \tag{8.11}\\
0 \longrightarrow \mathcal{H}_{z} \longrightarrow\left(W_{\lambda} \oplus W_{-\lambda}\right)^{2} \longrightarrow 0, & \text { for } z \neq w \text { and when } \Gamma_{+}=\Gamma_{-},  \tag{8.12}\\
0 \longrightarrow \mathcal{H}_{w} \longrightarrow W_{0} \oplus W_{0} \longrightarrow 0, & \text { when } \Gamma_{+}=\Gamma_{-}, \tag{8.13}
\end{align*}
$$

which are surjective for dimensional reasons.
Bringing all of those sequences together allows us to conclude the following.

Theorem 8.2-4. On $B^{3}(w)$, we have

$$
\ulcorner=V \oplus \mathscr{D} \mathcal{H} .
$$

Proof: Denote $W^{\prime}{ }_{\lambda}$ the space

$$
W_{\lambda}^{\prime}:= \begin{cases}W_{\lambda} \oplus W_{-\lambda}, & \text { if } \Gamma_{+}=\Gamma_{-} ; \\ W_{\lambda}, & \text { if } \Gamma_{+} \neq \Gamma_{-} .\end{cases}
$$

Let $p: W_{\lambda}^{\prime} \oplus W^{\prime}{ }_{-\lambda} \rightarrow W_{\lambda}^{\prime}$ denote the map $p(a, b)=2 \lambda b$.

For $\lambda \neq 0$, we use the Snake Lemma on the diagram

coming from Sequences 8.4, 8.7, 8.10 , and 8.12 , to produce an exact sequence

$$
\begin{gather*}
\operatorname{ker}(0) \longrightarrow \operatorname{ker}(\not(D) \longrightarrow \operatorname{ker}(p) \longrightarrow \operatorname{coker}(0) \longrightarrow \operatorname{coker}(\nsupseteq) \longrightarrow \operatorname{coker}(p) \\
0 \longrightarrow \bar{K}_{z} \longrightarrow W^{\prime}{ }_{-\lambda} \longrightarrow \quad V_{z} \longrightarrow \operatorname{coker}(\nsupseteq) \longrightarrow 0 \tag{8.14}
\end{gather*}
$$

Note that the map $V \rightarrow \operatorname{coker}(\mathscr{D})$ being surjective forces $V$ to be spanned by $V$ and $\mathfrak{D H}$.
Sequences (8.5) and (8.8) imply

$$
\operatorname{dim} V_{z}=\operatorname{dim} V_{z}+\operatorname{dim} W_{\lambda}^{\prime}-\operatorname{dim}{ }_{K}{ }_{z}
$$

while sequences 8.4) and 8.7) imply

$$
\operatorname{dim} V_{z}=\operatorname{dim} V_{z}+\operatorname{dim} W_{\lambda}^{\prime} .
$$

Thus

$$
\operatorname{dim} V_{z}=\operatorname{dim} V_{z}+2 \operatorname{dim} W_{\lambda}^{\prime}-\operatorname{dim} K_{z}=\operatorname{dim} V_{z}+\operatorname{dim} \mathfrak{D} \mathcal{H} .
$$

Since Lemma 6.3-1 guarantees that $\langle\mathfrak{D H}, V\rangle=\{0\}$, we have $V \cap \nsubseteq \mathcal{H}$ perpendicular to $V$ for the $L^{2}$ inner product. Hence $\mathfrak{D H} \cap V=\{0\}$, and $V_{z}=V_{z} \oplus \not \mathscr{D} \mathcal{H}$.
It remains to prove the theorem for $z=w$. We already know $V_{w}=V_{w}$ and $\mathfrak{D H}_{w} \subset V_{w}$. We also know from Sequences (8.6) (8.9) that

$$
\begin{aligned}
\operatorname{dim} V_{w} & =\operatorname{dim} V_{w}+\operatorname{dim} W_{0}^{\prime}-\operatorname{dim} K_{w} \\
& =\operatorname{dim} V_{w}+\operatorname{dim} \not \mathfrak{D H}_{w} .
\end{aligned}
$$

We therefore only have to prove that the intersection $V_{w} \cap \mathscr{D}_{A_{w}} \mathcal{H}_{w}$ is $\{0\}$ to complete the proof. The asymptotic behavior of $\phi \in \mathcal{H}_{w}$ is

$$
\phi= \begin{cases}t \phi_{0}^{+}+\phi_{1}^{+}+o(1), & \text { as } t \rightarrow \infty ; \\ t \phi_{0}^{-}+\phi_{1}^{-}+o(1), & \text { as } t \rightarrow-\infty ;\end{cases}
$$

for some $\phi_{0}^{ \pm}, \phi_{1}^{ \pm} \in W_{0}$. If $\Gamma_{+} \neq \Gamma_{-}$, we must have $\phi_{0}^{-}=\phi_{1}^{-}=0$, as $w$ is associated to $\Gamma_{+}$. The asymptotic behavior of $\mathscr{P}_{A_{w}} \phi$ is

$$
\mathscr{D}_{A_{w}} \phi= \begin{cases}\phi_{0}^{+}+o(1), & \text { as } t \rightarrow \infty ; \\ \phi_{0}^{-}+o(1), & \text { as } t \rightarrow-\infty .\end{cases}
$$

Suppose $\mathscr{D}_{A_{w}} \phi \in L^{2}$. Then

$$
\begin{aligned}
\left\|\mathscr{D}_{A_{w}} \phi\right\|_{L^{2}}^{2} & =\left\langle\mathscr{D}_{A_{w}}^{*} \mathscr{D}_{A_{w}} \phi, \phi\right\rangle+\lim _{t \rightarrow \infty}\left\langle\mathscr{D}_{A_{w}} \phi, \nu \phi\right\rangle+\lim _{t \rightarrow-\infty}\left\langle\mathscr{D}_{A_{w}} \phi, \nu \phi\right\rangle \\
& =\left\langle\phi_{0}^{+}, \phi_{1}^{+}\right\rangle+\lim _{t \rightarrow \infty} t\left|\phi_{0}^{+}\right|^{2}-\left\langle\phi_{0}^{-}, \phi_{1}^{-}\right\rangle-\lim _{t \rightarrow-\infty} t\left|\phi_{0}^{-}\right|^{2} .
\end{aligned}
$$

For $\left\|\mathscr{P}_{A_{w}} \phi\right\|_{L^{2}}$ to be finite, we must get rid of the limits, thus forcing $\phi_{0}^{ \pm}=0$ and consequently we have $\mathscr{D}_{A_{w}} \phi=0$. The proof is now complete.
For a continuous family of Fredholm operators, like $\mathscr{D}_{A_{z}}$ on $L_{\text {F }}^{2}$ parameterized on $B^{3}(w)$, the dimension of the kernel can only drop in a small neighborhood of a given point, it cannot increase. However, not any random behavior is acceptable.

Lemma 8.2-5 (also found in Kat95, p. 241]). Let $T: X \rightarrow Y$ be Fredholm and $S: X \rightarrow Y$ a bounded operator. Then the operator $T+t S$ is Fredholm and $\operatorname{dim} \operatorname{ker}(T+t S)$ is constant for small $|t|>0$.

Before spelling out the proof of this lemma, which we obviously use with $T=\mathscr{D}_{A_{w}}, X=$ $W_{\tau \epsilon}^{1,2}, Y=L_{\Gamma_{\epsilon}}^{2}$, and $S=c l(e)$ for some direction $e \in \mathbb{R}^{3}$, let's note that three scenarios are possible.

1. $\operatorname{dim} \bar{K} K_{z}$ is constant on a neighborhood around $w$, say $B^{3}(w)$;
2. $\operatorname{dim} \bar{K} K_{z}$ is constant for $z \in B^{3}(w) \backslash\{w\}$, but is smaller than $\operatorname{dim} \bar{K} K_{w}$;
3. $\operatorname{dim} \bar{K} K_{w+\lambda e} \neq \operatorname{dim} K_{w+\lambda^{\prime} e^{\prime}}$ for small $\lambda, \lambda^{\prime}>0$ and some $e \neq e^{\prime}$.

We close this section with the proof of Lemma 8.2-5
Proof: The proof is a simplified proof of the one provided by Kato in [Kat95, p. 241] for more general $T$ and $S$.
Define the sequences $M_{n} \subset X$ and $R_{n} \subset Y$ by

$$
\begin{array}{lrl}
M_{0} & :=X, & R_{0} \\
M_{n} & :=Y \\
S^{-1} R_{n}, & R_{n+1} & :=T M_{n}
\end{array}
$$

All the $M_{n}$ and $R_{n}$ are imbricated as

$$
M_{0} \supset M_{1} \supset M_{2} \supset \cdots \quad \text { and } \quad R_{0} \supset R_{1} \supset R_{2} \supset \cdots
$$

That the $M_{n}$ are closed is a trivial fact once it is established that the $R_{n}$ are closed. But define $\tilde{X}:=X / \operatorname{ker}(T)$ and $\tilde{M}_{n}$ to be the set of corresponding $\operatorname{ker}(T)$-cosets. Then for the map $\tilde{T}$ defined as $\tilde{T}(x+\operatorname{ker}(T))=T(x)$, we have $T M_{n}=\tilde{T} \tilde{M}_{n}$. Since $\tilde{T}$ is injective and Fredholm, and since $\tilde{M}_{n}$ is closed in $\tilde{X}$, then $R_{n+1}=\tilde{T} \tilde{M}_{n}$ is closed as well.
Define now

$$
X^{\prime}:=\bigcap_{n} M_{n} \quad \text { and } \quad Y^{\prime}:=\bigcap_{n} R_{n}
$$

and let $T^{\prime}, S^{\prime}$ be the restriction to $X^{\prime}$.
If $x \in X^{\prime}$, then $x \in M_{n}$ for all $n$ hence $T^{\prime} x=T x \in T M_{n}=R_{n+1}$ for all $n$, and by definition $S x \in R_{n}$. So both $T^{\prime}$ and $S^{\prime}$ are bounded operators $X^{\prime} \rightarrow Y^{\prime}$.
We now prove $\operatorname{Im}\left(T^{\prime}\right)=Y^{\prime}$. Suppose $y \in Y^{\prime}$, then $y \in R_{n}=T M_{n-1}$ for all $n$, hence $T^{-1} y \cap M_{n-1} \neq \varnothing$. Since $T$ is Fredholm, $T^{-1} y$ is closed and finite dimensional. We hence have a descending sequence $T^{-1} y \cap M_{n} \supset T^{-1} y \cap M_{n+1} \supset \cdots$ of finite dimensional nonempty affine spaces, which must then be stationary after a finite number of steps. The limit, which is then nonempty, must be $T^{-1} y \cap X^{\prime}$, hence $y \in \operatorname{Im}\left(T^{\prime}\right)$, and $T^{\prime}$ is surjective.

Notice that trivially, $\operatorname{ker}\left(T^{\prime}+t S^{\prime}\right) \subset \operatorname{ker}(T+t S)$. But more interestingly, those kernels are equal for $t \neq 0$. Indeed, take $x \in \operatorname{ker}(T+t S)$. We know $x \in X=M_{0}$, and prove by induction that $x \in M_{n}$ for all $n$, hence proving that $x \in \operatorname{ker}\left(T^{\prime}+t S^{\prime}\right)$. The induction step is proved by staring at the definitions: being in the kernel forces $S(-t x)=T x \in R_{n+1}$ if $x \in M_{n}$; but then $-t x \in M_{n+1}$ and for $t \neq 0$, we then have $x \in M_{n+1}$. We thus established that

$$
\operatorname{ker}\left(T^{\prime}+t S^{\prime}\right)=\operatorname{ker}(T+t S) \text { for } t \neq 0
$$

Obviously, for $t$ small enough, $T^{\prime}+t S^{\prime}$ is Fredholm and surjective, hence for small enough $|t|>0$,

$$
\operatorname{dim} \operatorname{ker}(T+t S)=\operatorname{dim} \operatorname{ker}\left(T^{\prime}+t S^{\prime}\right)=\operatorname{ind}\left(T^{\prime}+t S^{\prime}\right)
$$

is constant.

### 8.3 Asymptotic of the Higgs field

We know study the behavior of the Higgs field $\Phi$ as $z$ approaches of a point in $w \in W$.
We know $w$ is associated to the limit $\Gamma$ of $A$ at $\infty$ or $-\infty$, in the sense that $\Gamma$ splits $E$ as $L_{w} \oplus L_{-w}$. Without loss of generality, we suppose

$$
\Gamma_{+}=\Gamma .
$$

We can break up the analysis depending on which scenario happens; see page 91
When $\Gamma_{+} \neq \Gamma_{-}$, and for $2 \pi|z-w|<\epsilon$, notice that

$$
\begin{aligned}
V_{z} & =L_{\epsilon \in}^{2} \cap \operatorname{ker}\left(\mathfrak{D}_{A_{z}}^{*}\right)=L_{\bar{\epsilon}}^{2} \cap \operatorname{ker}\left(\mathfrak{D}_{A_{z}}^{*}\right)=L_{\epsilon}^{2} \cap \operatorname{ker}\left(\mathfrak{D}_{A_{z}}^{*}\right), \text { and } \\
V_{z} & =L_{\epsilon \epsilon}^{2} \cap \operatorname{ker}\left(\mathfrak{D}_{A_{z}}^{*}\right)=L_{\underline{\epsilon}}^{2} \cap \operatorname{ker}\left(\mathscr{D}_{A_{z}^{*}}^{*}\right)=L_{\epsilon \in}^{2} \cap \operatorname{ker}\left(\mathfrak{D}_{A_{z}}^{*}\right) .
\end{aligned}
$$

When $\Gamma_{+}=\Gamma_{-}$, those spaces are a priori all different.
Theorem 8.3-1. Suppose $\operatorname{dim}{ }^{5} K_{z}$ is constant in a neighborhood of $w$. On a closed ball $B^{3}(w)$ around $w$, there exists families of operators $\Phi^{\perp}$ and $\Phi^{\lrcorner}$, bounded independently of $z$, such that

$$
\begin{equation*}
\Phi=-\frac{-i}{2|z-w|} \Phi^{\perp}+\Phi^{\lrcorner} \tag{8.15}
\end{equation*}
$$

Furthermore, $\Phi^{\perp}$ is the $L^{2}$-orthogonal projection on $\mathscr{D}_{A_{z}} \mathcal{H}_{z} \cap V_{z}$.
Proof: When $\operatorname{dim} K_{z}$ is constant on $B^{3}(w)$, so is $\operatorname{dim} V_{z}$. Hence $V$ is a bundle on $B^{3}(w)$, say of rank $l$. Obviously, $V_{\mathrm{J}}$ supports many different norms, and amongst those are the $L^{2}$ and $L_{\epsilon}^{2}$ norms, which must be equivalent since $V_{\checkmark}$ is finite rank on a compact space.
Hence, for $\phi \in V_{z}$, observe that

$$
\|t \phi\|_{L^{2}} \leq C_{\epsilon}\|\phi\|_{\epsilon} \leq C\|\phi\|_{L^{2}} .
$$

Denote $P^{\lrcorner}$the $L^{2}$-orthogonal projection of $V$ on $V$. We just proved that

$$
\Phi \circ P^{\lrcorner} \text {is bounded independently of } z \in B^{3}(w) \text {. }
$$

It is part of the map $\Phi^{\lrcorner}$announced in the statement of the theorem.

As suggested above, let $\Phi^{\perp}$ denote the $L^{2}$-orthogonal projection on $\mathscr{P}_{A_{z}} \mathcal{H}_{z} \cap V_{z}$. Then

$$
\begin{aligned}
\Phi=-2 \pi i P m_{t} & =\Phi P^{\lrcorner}-2 \pi i\left(P^{\lrcorner}+\Phi^{\perp}\right) m_{t} \Phi^{\perp} \\
& =\Phi P^{\lrcorner}+2 \pi i P^{\lrcorner} m_{t} \Phi^{\perp}-2 \pi i \Phi^{\perp} m_{t} \Phi^{\perp} .
\end{aligned}
$$

As it turns out, $P^{\lrcorner} m_{t} \Phi^{\perp}$ is also bounded independently of $z \in B^{3}(w)$. Indeed, suppose we have an $L^{2}$-orthonormal frame $\phi_{1}, \ldots, \phi_{l}$ of $V_{1}$ in some open subset of $B^{3}(w)$, then

$$
\begin{aligned}
\left\|P^{\lrcorner} m_{t} \Phi^{\perp}(\phi)\right\|_{L^{2}} & =\left\|\sum_{j=1}^{l}\left\langle\phi_{j}, t \Phi^{\perp}(\phi)\right\rangle \phi_{j}\right\|_{L^{2}} \\
& =\left\|\sum_{j=1}^{l}\left\langle t \phi_{j}, \Phi^{\perp}(\phi)\right\rangle \phi_{j}\right\|_{L^{2}} \\
& \leq \sum_{j=1}^{l} C\left\|\phi_{j}\right\|_{L^{2}}^{2}\left\|\Phi^{\perp}(\phi)\right\|_{L^{2}} \\
& \leq C\|\phi\|_{L^{2}} .
\end{aligned}
$$

It remains only to analyze $\Phi^{\perp} m_{t} \Phi^{\perp}$.
Pick a vector $e \in \mathbb{R}^{3}$ of length 1 . Let

$$
\text { Ray }=\left\{w+\frac{\lambda}{2 \pi} e\right\} \subset B^{3}(w)
$$

be a ray inside $B^{3}(w)$ emerging from $w$. As the notation suggests, we parameterize this ray by $\lambda=2 \pi|z-w|$. Pick a family $\phi_{z} \in \mathscr{D}_{A_{z}} \mathcal{H}_{z}$ for $z \in$ Ray, with

$$
\begin{gather*}
\phi_{z} \in V_{z} \text { for } \lambda>0, \\
\quad\left\|\phi_{z}\right\|_{L_{\epsilon \epsilon}^{2}}=1 . \tag{8.16}
\end{gather*}
$$

But then,

$$
\left\|\phi_{z}\right\|_{L^{2}} \rightarrow \infty \text { as } \lambda \rightarrow 0
$$

To prove this claim, suppose it is not true. Then there is a subsequence $\phi_{z_{j}} \rightharpoonup \tilde{\phi}_{w}$ weakly in $L^{2}$. Hence $\left\langle\phi_{z_{j}}, f\right\rangle \rightarrow\left\langle\tilde{\phi}_{w}, f\right\rangle$ for all $f \in L^{2}$, in particular for all $f \in L_{\epsilon_{\mathrm{j}}}^{2}=\left(L_{\tau_{\epsilon}}^{2}\right)^{*}$, whence $\phi_{z_{j}} \rightharpoonup \tilde{\phi}_{w}$ weakly in $L_{\digamma_{\epsilon}}^{2}$. Since $\phi_{z} \rightarrow \phi_{w}$ in $L_{\digamma}^{2}$, we have $\tilde{\phi}_{w}=\phi_{w}$, which is impossible as $\tilde{\phi}_{w}$ is in $L^{2}$ while $\phi_{w}$ is not.
Theorem 6.1-3 guarantees that the operator $\mathscr{D}_{\Gamma_{w}}^{*}$ is an isomorphism $W_{\leftarrow \in}^{1,2} \rightarrow L_{\leftarrow \in}^{2}$, and $W_{\epsilon}^{1,2} \rightarrow L_{\epsilon}^{2}$, hence there exist a constant $C$ such that

$$
\begin{array}{ll}
\|u\|_{W_{\leftarrow}^{1,2}} \leq C\left\|\mathfrak{P}_{\Gamma_{w}}^{*} u\right\|_{L_{\epsilon \epsilon}^{2}}, & \text { for } u \in W_{\leftarrow \epsilon}^{1,2}, \\
\|u\|_{W_{\epsilon}^{1,2}}^{1,2} \leq C\left\|\mathfrak{P}_{\Gamma_{w}}^{*} u\right\|_{L_{e}^{2}}^{2}, & \text { for } u \in W_{\epsilon}^{1,2} . \tag{8.18}
\end{array}
$$

Because $\phi_{z} \in V_{z}$ for $\lambda>0$, Theorem 7.2-1 tells us that for $t>0$, we can write $\phi_{z}=e^{-\lambda t} \psi_{-\lambda}+g_{z}$ for some eigenvector $\psi_{-\lambda}$ of eigenvalue $-\lambda$ of $D_{\Gamma_{z}}$ and some $g_{z} \in W_{-\epsilon}^{1,2}\left([0, \infty) \times T^{3}\right)$.
When $\Gamma_{-}=\Gamma_{+}$, we also have interest in understanding the asymptotic behavior at $-\infty$. Theorem
7.2-1 tells us that for $t<0$, we can write $\phi_{z}=e^{\lambda t} \psi_{\lambda}+j_{z}$ for some eigenvector $\psi_{\lambda}$ of eigenvalue $\lambda$ of $D_{\Gamma_{z}}$ and some $j_{z} \in W_{\epsilon}^{1,2}\left((-\infty, 0] \times T^{3}\right)$.
While $g_{z}$ and $j_{z}$ appear to be defined only for $t>0$ and $t<0$ respectively, let's define them globally on $\mathbb{R} \times T^{3}$ by $g_{z}=\phi_{z}-e^{-\lambda t} \psi_{-\lambda}$ and $j_{z}=\phi_{z}-e^{\lambda t} \psi_{\lambda}$.
Notice that

$$
\begin{equation*}
\mathscr{D}_{\Gamma_{z}}^{*} g_{z}=\mathscr{P}_{\Gamma_{z}}^{*} \phi_{z}=\left(\mathscr{D}_{\Gamma_{z}}^{*}-\mathscr{P}_{A_{z}}^{*}\right) \phi_{z}=c l(\Gamma-A) \phi_{z}, \tag{8.19}
\end{equation*}
$$

and similarly

$$
\begin{equation*}
\mathscr{P}_{\Gamma_{z}}^{*} j_{z}=\mathscr{D}_{\Gamma_{z}}^{*} \phi_{z}=\left(\mathscr{D}_{\Gamma_{z}}^{*}-\mathscr{D}_{A_{z}}^{*}\right) \phi_{z}=c l(\Gamma-A) \phi_{z}, \tag{8.20}
\end{equation*}
$$

Remember that we decided that $w$ is associated to $\Gamma=\Gamma_{+}$. Hence for $t>0$, we know that $|A-\Gamma| \leq C e^{-\beta t}$. For $t<0$, we have

$$
\begin{aligned}
|A-\Gamma| & \leq\left|A-\Gamma_{-}\right|+\left|\Gamma_{-}-\Gamma\right| \\
& \leq C e^{\beta t}+C^{\prime} .
\end{aligned}
$$

Hence overall, there is a constant such that $|c l(A-\Gamma)| \leq C \sigma_{(0, \beta)}$, and this estimate can be improved to $|c l(A-\Gamma)| \leq C \sigma_{(-\beta, \beta)}$ when $\Gamma_{-}=\Gamma_{+}$. Hence $c l(A-\Gamma)$ gives a bounded map $L_{\tau \epsilon}^{2} \rightarrow L_{\leftarrow \epsilon}^{2}$ in all cases and $L_{\bar{\epsilon}}^{2} \rightarrow L_{\epsilon}^{2}$ when $\Gamma_{-}=\Gamma_{+}$. Thus Equation (8.19) yields

$$
\begin{equation*}
\left\|\mathscr{P}_{\Gamma_{z}}^{*} g_{z}\right\|_{L_{\epsilon_{\epsilon}}^{2}} \leq C\left\|\phi_{z}\right\|_{L_{\epsilon \in}^{2}}, \tag{8.21}
\end{equation*}
$$

and for the special case $\Gamma_{-}=\Gamma_{+}$, Equation y.20 yields

$$
\begin{equation*}
\left\|\mathfrak{P}_{\Gamma_{z}}^{*} j_{z}\right\|_{L_{e}^{2}} \leq C\left\|\phi_{z}\right\|_{L_{\epsilon}^{2}} . \tag{8.22}
\end{equation*}
$$

From Equations 8.17, and 8.21, we derive

$$
\begin{aligned}
\left\|g_{z}\right\|_{W_{\leftarrow 乚}^{1,2}} & \leq C\left\|\mathfrak{p}_{\Gamma_{w}}^{*} g_{z}\right\|_{L_{\leftarrow \epsilon}^{2}} \\
& =C\left\|\mathfrak{p}_{\Gamma_{z}}^{*} g_{z}+\lambda c l(e) g_{z}\right\|_{L_{\llcorner\epsilon}^{2}} \\
& \leq C\left\|\phi_{z}\right\|_{L_{\epsilon \epsilon}^{2}}+C \lambda\left\|g_{z}\right\|_{L_{\llcorner\epsilon}^{2}},
\end{aligned}
$$

hence for $\lambda$ small enough, we can rearrange and obtain

$$
\begin{equation*}
\left\|g_{z}\right\|_{W_{\mathrm{L}^{e}}^{1,2}} \text { is bounded independently of small } z, \tag{8.23}
\end{equation*}
$$

while from Equations (8.18) and (8.22), we similarly obtain

$$
\begin{equation*}
\left\|j_{z}\right\|_{W_{e}^{1,2}} \text { is bounded independently of small } z . \tag{8.24}
\end{equation*}
$$

This last fact is also true for $\Gamma_{-} \neq \Gamma_{+}$, for in that case $j_{z}=\phi_{z}$ and its $L_{F_{\epsilon}}^{2}$-norm is equivalent to the $L_{\epsilon}^{2}$-norm, as both as defined on $V$ over $B^{3}(w)$.
While it is agreeable to work with a smooth splitting, nothing prevents us from considering the function

$$
h_{\lambda}= \begin{cases}e^{\lambda t} \psi_{\lambda}, & \text { for } t<0 \\ e^{-\lambda t} \psi_{-\lambda}, & \text { for } t>0\end{cases}
$$

and the associate splitting

$$
\phi_{z}=h_{\lambda}+r_{z},
$$

for which, obviously,

$$
r_{z}= \begin{cases}j_{z}, & \text { for } t<0  \tag{8.25}\\ g_{z}, & \text { for } t>0\end{cases}
$$

From the bounds of Equations 8.23) and 8.24, we have that

$$
\begin{equation*}
\left\|r_{z}\right\|_{L_{\epsilon^{2}}} \text { is bounded independently of small } z \tag{8.26}
\end{equation*}
$$

Consider the families

$$
\begin{aligned}
\bar{\phi}_{z} & :=\phi_{z} /\left\|\phi_{z}\right\|_{L^{2}}, \\
\bar{h}_{\lambda} & :=h_{\lambda} /\left\|\phi_{z}\right\|_{L^{2}}, \\
\bar{r}_{z} & :=r_{z} /\left\|\phi_{z}\right\|_{L^{2}} .
\end{aligned}
$$

The bound (8.26), and the fact that $\left\|\phi_{z}\right\|_{L^{2}} \rightarrow \infty$ imply that

$$
\left\|\bar{r}_{z}\right\|_{L_{\epsilon}^{2}} \rightarrow 0 \text { as } \lambda \rightarrow 0
$$

A fortiori, $\left\|\bar{r}_{z}\right\|_{L^{2}} \rightarrow 0$.
The triangle inequality guarantees

$$
\left|\left\|\bar{h}_{\lambda}\right\|_{L^{2}}-\left\|\bar{r}_{z}\right\|_{L^{2}}\right| \leq\left\|\bar{\phi}_{z}\right\|_{L^{2}} \leq\left\|\bar{h}_{\lambda}\right\|_{L^{2}}+\left\|\bar{r}_{z}\right\|_{L^{2}}
$$

Since $\left\|\bar{\phi}_{z}\right\|_{L^{2}}=1$, and $\left\|\bar{r}_{z}\right\|_{L^{2}} \rightarrow 0$, we must have

$$
\left\|\bar{h}_{\lambda}\right\|_{L^{2}} \rightarrow 1 \text { as } \lambda \rightarrow 0 .
$$

Let's now come back to our main worry. We study

$$
\left\langle t \bar{\phi}_{z}, \bar{\phi}_{z}\right\rangle=\left\langle t \bar{h}_{\lambda}, \bar{h}_{\lambda}\right\rangle+2\left\langle\bar{h}_{\lambda}, t \bar{r}_{z}\right\rangle+\left\langle t \bar{r}_{z}, \bar{r}_{z}\right\rangle .
$$

The last two terms are bounded by a multiple of $\left\|t \bar{r}_{z}\right\|_{L^{2}}$. But

$$
\left\|t \bar{r}_{z}\right\|_{L^{2}} \leq C\left\|\bar{r}_{z}\right\|_{L_{\epsilon}^{2}}=C\left\|\tilde{r}_{z}\right\|_{L_{\epsilon}^{2}} /\left\|\phi_{\lambda}\right\|_{L^{2}},
$$

hence it is going to 0 .
As for the first term, we have

$$
\begin{aligned}
\left\langle t \bar{h}_{\lambda}, \bar{h}_{\lambda}\right\rangle & =\frac{1}{\left\|\phi_{\lambda}\right\|_{L^{2}}^{2}}\left(\int_{0}^{\infty} t e^{-2 \lambda t}\left|\psi_{-\lambda}\right|^{2}+\int_{-\infty}^{0} t e^{2 \lambda t}\left|\psi_{\lambda}\right|^{2}\right) \\
& =\frac{1}{2 \lambda} \frac{1}{\left\|\phi_{\lambda}\right\|_{L^{2}}^{2}}\left(\int_{0}^{\infty} e^{-2 \lambda t}\left|\psi_{-\lambda}\right|^{2}+\int_{-\infty}^{0} e^{2 \lambda t}\left|\psi_{\lambda}\right|^{2}\right) \\
& =\frac{1}{2 \lambda}\left\|\bar{h}_{\lambda}\right\|_{L^{2}}^{2},
\end{aligned}
$$

hence

$$
\left\langle t \bar{\phi}_{\lambda}, \bar{\phi}_{\lambda}\right\rangle=\frac{1}{2 \lambda}+o(1) \text { as } \lambda \rightarrow 0 .
$$

Suppose now $\bar{\phi}_{z}^{1}$ and $\bar{\phi}_{z}^{2}$ are two such families, but so that

$$
\left\langle\bar{\phi}_{z}^{1}, \bar{\phi}_{z}^{2}\right\rangle_{L^{2}}=0 .
$$

Then

$$
\begin{aligned}
\left\langle t \bar{\phi}_{z}^{1}, \bar{\phi}_{z}^{2}\right\rangle & =\left\langle t \bar{h}_{\lambda}^{1}, \bar{h}_{\lambda}^{2}\right\rangle+\left\langle\bar{h}_{\lambda}^{1}, t \bar{r}_{z}^{2}\right\rangle+\left\langle t \bar{r}_{z}^{1}, \bar{h}_{\lambda}^{2}\right\rangle+\left\langle t \bar{r}_{z}^{1}, \bar{r}_{z}^{2}\right\rangle \\
& =\frac{1}{2 \lambda}\left\langle\bar{h}_{\lambda}^{1}, \bar{h}_{\lambda}^{2}\right\rangle+o(1),
\end{aligned}
$$

and of course $\left\langle\bar{h}_{\lambda}^{1}, \bar{h}_{\lambda}^{2}\right\rangle \rightarrow 0$, hence the result.
One of the crucial feature of this proof is our ability to find a uniform bound for $m_{t}$ on $V$. Such a bound exist in the case where $\operatorname{dim} K_{z}$ is constant precisely because this constant rank condition implies that $V$ is a bundle over $B^{3}(w)$, allowing us to say that the $L^{2}$-norm and $L_{\epsilon}^{2}$-norm are equivalent.
We can take the trace of $(B, \Phi)$ to obtain an abelian monopole $(a, \varphi)$ on $B^{3}(w) \backslash\{w\}$. The Bogomolny equation reduces to

$$
d \varphi=* d a,
$$

and thus $\Delta \varphi=0$. Since $\varphi$ is harmonic, not every possible behavior as $z \rightarrow w$ is acceptable. For one thing, there is a unique set of homogeneous harmonic polynomials $p_{m}$ and $q_{m}$ of degree $m$ which give a decomposition of $\varphi$ on $B^{3}(w) \backslash\{w\}$ as a Laurent series

$$
\varphi=\sum_{m=0}^{\infty} p_{m}(z-w)+\sum_{m=0}^{\infty} \frac{q_{m}(z-w)}{|z-w|^{2 m+1}} ;
$$

see ABR01 Thm 10.1, p. 209].
Whether or not the rank is constant, we can find for any sequence of points approaching $w$ a subsequence of points $z_{j} \rightarrow w$ for which the decomposition of Equation 8.15 is valid. We then have

$$
\lim _{j \rightarrow \infty} 2\left|z_{j}-w\right| \varphi_{z_{j}}=i \operatorname{dim} \mathscr{D}_{A_{z_{j}}} \mathcal{H}_{z_{j}}=i\left(\operatorname{rk} \mathcal{H}-\operatorname{dim}{ }^{〔} K_{z_{j}}\right) .
$$

By the Laurent series decomposition given above, this number must be the same in any way we approach $w$, hence $\operatorname{dim} K_{z}$ must be constant on $B^{3}(w) \backslash\{w\}$, thus eliminating Scenario 3 of page 91
Scenario 2 remains to be dealt with.

### 8.4 Preliminary work: Green's operator on $S^{ \pm} \otimes L_{z}$

This section consists of preliminary work on the study of the behavior of the connection $B$ at the singular points. Because of the formula $B=P d^{z}=\left(1-\mathscr{D}_{A_{z}} G_{A_{z}} \mathscr{D}_{A_{z}}^{*}\right) d^{z}$, getting an explicit formula for $G_{A_{z}}$ would greatly help in understanding the asymptotic behavior of $B$ at singular points. As a first step into achieving that goal, we compute the Green's operator for the Laplacian on $S^{ \pm} \otimes L_{z}$ on $\mathbb{R} \times T^{3}$.

Define the operator

$$
\begin{aligned}
T_{\lambda}: L^{2}(\mathbb{R}) & \rightarrow L^{2}(\mathbb{R}) \\
g & \mapsto \frac{1}{\lambda} \int_{-\infty}^{\infty} e^{-\lambda|t-s|} g(s) d s
\end{aligned}
$$

Let $\phi_{\nu}$ be a basis of the eigenspaces of the Laplacian on the spinor bundle of $T^{3}$. Then every section $\phi$ decomposes as

$$
\begin{equation*}
\phi=\sum_{\nu \in \mathbb{Z}^{3}} g_{\nu}(t) \phi_{\nu} \tag{8.27}
\end{equation*}
$$

Here multiplicities are hidden but keep them in mind.
Lemma 8.4-1. For all $z \in \hat{T}^{3}$, we have

$$
G_{L_{z}}(\phi)=\frac{1}{2} \sum_{\nu \in \mathbb{Z}^{3}} T_{2 \pi|\nu-z|}\left(g_{\nu}\right) \phi_{\nu}
$$

Proof: First notice that

$$
T_{\lambda} g=\frac{1}{\lambda}\left(\int_{-\infty}^{t} e^{-\lambda t} e^{\lambda s} g(s) d s+\int_{t}^{\infty} e^{\lambda t} e^{-\lambda s} g(s) d s\right)
$$

hence $\partial_{t}\left(T_{\lambda} g\right)=g(T) / \lambda-\left(\int_{-\infty}^{t}\right)+\left(\int_{t}^{\infty}\right)-g(t) / \lambda$, and

$$
\begin{aligned}
\partial_{t}^{2}\left(T_{\lambda} g\right) & =-g(t)+\lambda\left(\int_{-\infty}^{t}\right)+\lambda\left(\int_{t}^{\infty}\right)-g(t) / \lambda \\
& =-2 g+\lambda^{2} T_{\lambda} g
\end{aligned}
$$

Remember now that on $S^{ \pm} \otimes L_{z}$ on $\mathbb{R} \times T^{3}$, the Laplacian splits as

$$
\Delta_{L_{z}}=-\partial_{t}^{2}+\Delta_{T^{3}, L_{z}}
$$

Recall also that for $\nu \in \mathbb{Z}^{3}$, we have $\Delta_{T^{3}} \phi_{\nu}=(2 \pi|\nu|)^{2} \phi_{\nu}$ and $\Delta_{T^{3}, L_{z}} \phi_{\nu}=(2 \pi|\nu-z|)^{2} \phi_{\nu}$. Hence for the proposed $G$, we have

$$
\begin{aligned}
\Delta_{L_{z}} G_{L_{z}} \phi & =\frac{1}{2} \sum_{\nu \in \mathbb{Z}^{3}}\left(-\partial_{t}^{2} T_{2 \pi|\nu-z|}\left(g_{\nu}\right) \phi_{\nu}+(2 \pi|\nu-z|)^{2} g_{\nu} \phi_{\nu}\right) \\
& =\frac{1}{2} \sum_{\nu \in \mathbb{Z}^{3}}\left(2 g_{\nu}-(2 \pi|\nu-z|)^{2} g_{\nu}+(2 \pi|\nu-z|)^{2} g_{\nu}\right) \phi_{\nu} \\
& =\phi
\end{aligned}
$$

The proof is now complete.
It appears very important then to understand $T_{\lambda}$ carefully. Let $m_{\lambda}: \mathbb{R} \rightarrow \mathbb{R}$ be the multiplication by $\lambda$. We have the following identities.

Lemma 8.4-2. For different values of $\lambda, \eta>0$, we have

$$
\begin{equation*}
T_{\lambda}=\frac{\eta^{2}}{\lambda^{2}} m_{\lambda / \eta}^{*} T_{\eta} m_{\eta / \lambda}^{*} \tag{8.28}
\end{equation*}
$$

and in particular

$$
T_{\lambda}=\frac{1}{\lambda^{2}} m_{\lambda}^{*} T_{1} m_{\lambda^{-1}}^{*}
$$

Proof: We just compute

$$
\begin{aligned}
T_{\lambda}(g) & =\frac{1}{\lambda} \int_{-\infty}^{\infty} e^{-\lambda|t-s|} g(s) d s \\
& =\frac{1}{\lambda} \int_{-\infty}^{\infty} e^{-\eta\left|\frac{\lambda}{\eta} t-\frac{\lambda}{\eta} s\right|} g\left(\frac{\lambda s / \eta}{\lambda / \eta}\right) \frac{d(\lambda s / \eta)}{\lambda / \eta} \\
& =\frac{\eta^{2}}{\lambda^{2}} T_{\eta}\left(g \circ m_{\eta / \lambda}\right) \circ m_{\lambda / \eta},
\end{aligned}
$$

whence the conclusion. The proof is now complete.
Lemma 8.4-3. Viewed as operators $L^{2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R})$, the operators $m_{\lambda}^{*}$ and $T_{\lambda}$ have norm

$$
\begin{gathered}
\left\|m_{\lambda}^{*} f\right\|=\frac{1}{\sqrt{\lambda}}\|f\|, \forall f, \\
\left\|T_{\lambda}\right\|=\frac{1}{\lambda^{2}}\left\|T_{1}\right\|
\end{gathered}
$$

Furthermore for z close to 0 , the Green's operator has norm

$$
\left\|G_{L_{z}}\right\|=\frac{\left\|T_{1}\right\|}{2|z|^{2}}
$$

as an operator $L^{2} \rightarrow L^{2}$.
Proof: First notice

$$
\left\|m_{\lambda}^{*} f\right\|^{2}=\int\left|f(\lambda t)^{2}\right| d t=\int|f(s)|^{2} d s / \lambda=\frac{1}{\lambda}\|f\|^{2} .
$$

Then we compare. On one hand

$$
\begin{aligned}
\left\|T_{\lambda} g\right\| & =\frac{1}{\lambda^{2}}\left\|m_{\lambda}^{*} T_{1}\left(m_{\lambda^{-1}}^{*} g\right)\right\| \\
& =\frac{1}{\lambda^{2}} \frac{1}{\sqrt{\lambda}}\left\|T_{1}\left(m_{\lambda^{-1}}^{*} g\right)\right\| \\
& \leq \frac{1}{\lambda^{2}} \frac{1}{\sqrt{\lambda}}\left\|T_{1}\right\|\left\|m_{\lambda^{-1}}^{*} g\right\| \\
& =\frac{1}{\lambda^{2}}\left\|T_{1}\right\|\|g\|,
\end{aligned}
$$

hence

$$
\begin{equation*}
\left\|T_{\lambda}\right\| \leq \frac{1}{\lambda^{2}}\left\|T_{1}\right\| . \tag{8.29}
\end{equation*}
$$

On the other hand, we find in a similar fashion that

$$
\begin{equation*}
\left\|T_{1}\right\| \leq \lambda^{2}\left\|T_{\lambda}\right\| \tag{8.30}
\end{equation*}
$$

From Inequalities 8.29) and 8.30, we obtain the desired result for $T_{\lambda}$.

To find the norm of $G_{L_{z}}$, we use the decomposition given by Equation (8.27). We have

$$
\begin{aligned}
\left\|G_{L_{z}}(\phi)\right\|^{2} & =\frac{1}{4} \sum\left\|T_{2 \pi|\nu-z|}\left(g_{\nu}\right)\right\|^{2} \\
& \leq \frac{1}{4} \sum \frac{1}{(2 \pi)^{4}|\nu-z|^{4}}\left\|T_{1}\right\|^{2}\left\|g_{\nu}\right\|^{2} \\
& \leq \frac{\left\|T_{1}\right\|^{2}}{4} \sup \left(\frac{1}{(2 \pi)^{4}|\nu-z|^{4}}\right)\|\phi\|^{2},
\end{aligned}
$$

hence for $z$ close to 0 ,

$$
\left\|G_{L_{z}}\right\| \leq \frac{\left\|T_{1}\right\|}{8 \pi^{2}|z|^{2}}
$$

To prove equality, note that for $\phi=g_{0} \phi_{0}$, with $g_{0} \in L^{2}$, and $\left\|g_{0}\right\|=1$, we have

$$
G_{L_{z}}(\phi)=\frac{1}{2} T_{2 \pi|z|}\left(g_{0}\right) \phi_{0}
$$

But then

$$
\begin{aligned}
\left\|G_{z}\right\| & =\sup _{\|\phi\|=1}\left\|G_{z}(\phi)\right\| \\
& \geq \sup _{\left\|g_{0}\right\|=1}\left\|G_{z}\left(g_{0} \phi_{0}\right)\right\| \\
& =\frac{1}{2} \sup _{\left\|g_{0}\right\|=1}\left\|T_{2 \pi|z|}\left(g_{0}\right)\right\| \\
& =\frac{1}{8 \pi^{2}|z|^{2}}\left\|T_{1}\right\| .
\end{aligned}
$$

The proof is now complete.
Lemma 8.4-4. We have

$$
G_{L_{z}}=\frac{\mathfrak{L}_{z}}{|z|^{2}}+\mathfrak{M}_{z}
$$

with the $L^{2} \rightarrow L^{2}$ operator norms of $\mathfrak{L}_{z}$ and $\mathfrak{M}_{z}$ bounded independently of $z$ for $\operatorname{dist}\left(z, \mathbb{Z}^{3}\right)<1 / 2$.
Proof: Let $p_{0}$ be the projection $\phi \mapsto g_{0} \phi_{0}$ and $p_{1}=1-p_{0}$. Set

$$
\begin{aligned}
\mathfrak{L}_{z}(\phi) & :=|z|^{2} G_{z} p_{0}(\phi), \\
\mathfrak{M}_{z}(\phi) & =G_{z} p_{1}(\phi) .
\end{aligned}
$$

Obviously, $G_{L_{z}}=\mathfrak{L}_{z} /|z|^{2}+\mathfrak{M}_{z}$. It remains to show that $\mathfrak{L}_{z}$ and $\mathfrak{M}_{z}$ are uniformly bounded for $z$ close to $\Lambda^{*}$. For $\mathfrak{M}_{z}$ and $\operatorname{dist}\left(z, \mathbb{Z}^{3}\right)<1 / 2$, we have

$$
\begin{aligned}
\left\|G_{L_{z}} p_{1}(\phi)\right\|^{2} & =\frac{1}{4} \sum_{\nu \neq 0}\left\|T_{2 \pi|\nu-z|}\left(g_{\nu}\right)\right\|^{2} \\
& \leq \frac{\left\|T_{1}\right\|^{2}}{4} \sup _{\nu \neq 0} \frac{1}{(2 \pi|\nu-z|)^{4}}\|\phi\|^{2} \\
& \leq \frac{4}{\pi^{2}}\left\|T_{1}\right\|^{2}\|\phi\|^{2},
\end{aligned}
$$

proving the claim for $\mathfrak{M}_{z}$.
Obviously,

$$
G_{L_{z}} p_{0}(\phi)=\frac{1}{8 \pi^{2}|z|^{2}} m_{2 \pi|z|^{*}}^{*}\left(m_{(2 \pi|z|)^{-1}}^{*} g_{0}\right) \phi_{0}
$$

hence the claims for $\mathfrak{L}_{z}$ follows from Lemma 8.4-3 The proof is now complete.
While $\left\|\mathfrak{L}_{z}\right\|_{o p}$ is constant and not 0 , the family of operators $\mathfrak{L}_{z}$ in a very weak sense converges to 0 .
Lemma 8.4-5. Let $g \in C_{c}^{\infty}$. we have that

$$
\mathfrak{L}_{z}(g) \rightarrow 0 \text { in } L^{2} \text { norm as } \lambda \rightarrow 0 .
$$

Proof: Suppose the support of $g$ is $[m, M]$. Then

$$
\begin{aligned}
\left\|\mathfrak{L}_{z}(g)\right\|^{2} & =(2 \pi|z|)^{2} \int_{-\infty}^{\infty}\left|\int_{-\infty}^{\infty} e^{-2 \pi|z| \cdot|t-s|} g(s) d s\right|^{2} d t \\
& \leq(2 \pi|z|)^{2} \int_{-\infty}^{\infty}(M-m)^{2} \max (g)^{2} e^{-4 \pi|z| \text { dist }(t, s u p p(g))} \\
& =2 \pi|z|((M-m) \max (g))^{2} .
\end{aligned}
$$

The result follows.

## Appendix A

## Reduction of ASD equation to lower dimension

The curvature of the connection $A=A_{1} d x^{1}+\cdots A_{4} d x^{4}$ is given by

$$
\begin{aligned}
F & =d A+A \wedge A \\
& =\sum_{i, j} \partial_{j} A_{i} d x^{j} \wedge d x^{i}+\sum_{i, j} A_{i} A_{j} d x^{i} \wedge d x^{j} \\
& =\sum_{i<j} F_{i j} d x^{i} \wedge d x^{j}
\end{aligned}
$$

with $F_{i j}=\partial_{i} A_{j}-\partial_{j} A_{i}+\left[A_{i}, A_{j}\right]$.
To convert to the standard self-dual $\epsilon_{1}, \epsilon_{2}, \epsilon_{3}$, and anti-self-dual $\bar{\epsilon}_{1}, \bar{\epsilon}_{2}, \bar{\epsilon}_{3}$ basis of $\Lambda^{2}$, we collect terms. For example

$$
\begin{aligned}
F_{12} d x^{12}+F_{34} d x^{34}= & \left(\frac{F_{12}+F_{34}}{2}\right) d x^{12}+\left(\frac{F_{12}-F_{34}}{2}\right) d x^{12} \\
& +\left(\frac{F_{12}+F_{34}}{2}\right) d x^{34}+\left(\frac{F_{34}-F_{12}}{2}\right) d x^{34} \\
= & \left(\frac{F_{12}+F_{34}}{2}\right) \epsilon_{1}+\left(\frac{F_{12}-F_{34}}{2}\right) \bar{\epsilon}_{1}
\end{aligned}
$$

We keep collecting terms, and get

$$
\begin{aligned}
F= & \left(\frac{F_{12}+F_{34}}{2}\right) \epsilon_{1}+\left(\frac{F_{12}-F_{34}}{2}\right) \bar{\epsilon}_{1}+\left(\frac{F_{13}-F_{24}}{2}\right) \epsilon_{2} \\
& +\left(\frac{F_{13}+F_{24}}{2}\right) \bar{\epsilon}_{2}+\left(\frac{F_{14}+F_{23}}{2}\right) \epsilon_{3}+\left(\frac{F_{14}-F_{23}}{2}\right) \bar{\epsilon}_{3}
\end{aligned}
$$

So the ASD equations are

$$
\begin{align*}
& F_{12}+F_{34}=0, \\
& F_{13}-F_{24}=0,  \tag{A.1}\\
& F_{14}+F_{23}=0 .
\end{align*}
$$

We now peel off dimensions one at a time.

## Dimension 3: the Bogomolny equation

Let the $A_{i}$ be independent of $x_{1}$, and set $\Phi:=A_{1}$. Then the Equations A.1) reduce to

$$
\begin{align*}
& -\partial_{2} \Phi+\left[\Phi, A_{2}\right]+F_{34}=0, \\
& -\partial_{3} \Phi+\left[\Phi, A_{3}\right]-F_{24}=0,  \tag{A.2}\\
& -\partial_{4} \Phi+\left[\Phi, A_{4}\right]+F_{23}=0 .
\end{align*}
$$

Set $B:=A_{2} d x^{2}+A_{3} d x^{3}+A_{4} d x^{4}$. It is a connection on the $\left(x_{2}, x_{3}, x_{4}\right)$-space. On that space, the Hodge star works as follows:

$$
\begin{aligned}
* d x^{2} & =d x^{3} \wedge d x^{4}, \\
* d x^{3} & =-d x^{2} \wedge d x^{4}, \\
* d x^{4} & =d x^{2} \wedge d x^{3}, \text { and } \\
*^{2} & =1 \text { on } \wedge \wedge^{1} .
\end{aligned}
$$

Furthermore, the connection $B$ extends to endomorphisms by the formula

$$
\nabla_{B} \Phi=d \Phi+[B, \Phi] .
$$

Hence the Equations A.2) can be written as a single equation as

$$
\begin{equation*}
\nabla_{B} \Phi=* F_{B}, \tag{A.3}
\end{equation*}
$$

the Bogomolny equation.

## Dimension 2: the Hitchin equations

Let the $A_{i}$ be independent of $x_{3}$ and $x_{4}$, and set $\phi_{1}:=A_{3}, \phi_{2}:=A_{4}$. The the Equations A. 1 ) reduce to

$$
\begin{gathered}
F_{12}+\left[\phi_{1}, \phi_{2}\right]=0, \\
\partial_{1} \phi_{1}+\left[A_{1}, \phi_{1}\right]-\partial_{2} \phi_{2}-\left[A_{2}, \phi_{2}\right]=0, \\
\partial_{1} \phi_{2}+\left[A_{1}, \phi_{2}\right]+\partial_{2} \phi_{1}+\left[A_{2}, \phi_{1}\right]=0 .
\end{gathered}
$$

In other words, set $B:=A_{1} d x^{1}+A_{2} d x^{2}$, and we have

$$
\begin{gather*}
F_{12}=-\left[\phi_{1}, \phi_{2}\right], \\
\nabla_{B}^{1} \phi_{1}=\nabla_{B}^{2} \phi_{2},  \tag{A.4}\\
\nabla_{B}^{2} \phi_{1}=-\nabla_{B}^{1} \phi_{2} .
\end{gather*}
$$

Since all orientable 2-manifolds are complex, let $d z=d x_{1}+i d x_{2}$, and

$$
\Phi:=\frac{1}{2}\left(\phi_{1}+i \phi_{2}\right) d z
$$

Should the connection $A$ be on a bundle $E$, then $\Phi$ is a section of $\bigwedge^{1,0} \operatorname{End}(E)$ and is called a Higgs
field. On 1-forms, we consider the graded commutator

$$
\begin{aligned}
{\left[\Phi, \Phi^{*}\right] } & =\Phi \Phi^{*}+\Phi^{*} \Phi \\
& =\frac{1}{2}\left(\phi_{1}+i \phi_{2}\right)\left(\phi_{1}-i \phi_{2}\right) d z \wedge d \bar{z}+\frac{1}{2}\left(\phi_{1}-i \phi_{2}\right)\left(\phi_{1}+i \phi_{2}\right) d \bar{z} \wedge d z \\
& =-\frac{i}{2}\left[\phi_{1}, \phi_{2}\right] d z \wedge d \bar{z} \\
& =-\left[\phi_{1}, \phi_{2}\right] d x^{1} \wedge d x^{2} .
\end{aligned}
$$

Hence

$$
F_{B}=\left[\Phi, \Phi^{*}\right] .
$$

Consider the operator $\bar{\partial}_{B}=\frac{1}{2}\left(\nabla_{B}^{1}+i \nabla_{B}^{2}\right) d \bar{z}$, and we have

$$
\begin{aligned}
\bar{\partial}_{B} \Phi & =\left(\bar{\partial}_{B} \phi_{1}+i\left(\bar{\partial}_{B} \phi_{2}\right)\right) d \bar{z} \wedge d z \\
& =\frac{1}{2}\left(\nabla_{B}^{1} \phi_{1}+i \nabla_{B}^{2} \phi_{1}+i \nabla_{B}^{1} \phi_{2}-\nabla_{B}^{2} \phi_{2}\right) d \bar{z} \wedge d z \\
& =0 .
\end{aligned}
$$

Hence the Equations (A.4) can be written as two equations

$$
\begin{gather*}
F_{B}=\frac{1}{4}\left[\Phi, \Phi^{*}\right],  \tag{A.5}\\
\bar{\partial}_{B} \Phi=0,
\end{gather*}
$$

which we call the Hichin equations.

## Dimension 1: The Nahm Equations

Let the $A_{i}$ be independent of $x_{2}, x_{3}, x_{4}$, and set

$$
\begin{gathered}
t:=x_{1}, \\
B:=A_{1} d t \\
T_{1}:=A_{2}, T_{2}:=A_{3}, T_{3}:=A_{4} .
\end{gathered}
$$

Then the Equations (A.1) reduce to

$$
\begin{gathered}
\nabla_{B}^{t} T_{\sigma(1)}+\left[T_{\sigma(2)}, T_{\sigma(3)}\right]=0 \\
\text { for all even permutation } \sigma .
\end{gathered}
$$

We call those equations the Nahm equations. These equations first appeared in [Nah83].

## Appendix B

## Excision principle for the index of Fredholm operators

Let

$$
\begin{aligned}
& D_{1}: L^{2}\left(X_{1}\right) \rightarrow L^{2}\left(X_{1}\right) \\
& D_{2}: L^{2}\left(X_{2}\right) \rightarrow L^{2}\left(X_{2}\right)
\end{aligned}
$$

be unbounded Fredholm operators, defined locally.
Let $X_{1}=A_{1} \cup B_{1}$, and $X_{2}=A_{2} \cup B_{2}$, with compact intersections

$$
A_{1} \cap B_{1}=A_{2} \cap B_{2}
$$

and suppose $D_{1}=D_{2}$ on that intersection.
We construct $\tilde{X}_{1}=A_{1} \cup B_{2}$ and $\tilde{X}_{2}=A_{2} \cup B_{1}$. Let

$$
\begin{aligned}
& \tilde{D}_{1}: L^{2}\left(\tilde{X}_{1}\right) \rightarrow L^{2}\left(\tilde{X}_{1}\right) \\
& \tilde{D}_{2}: L^{2}\left(\tilde{X}_{2}\right) \rightarrow L^{2}\left(\tilde{X}_{2}\right)
\end{aligned}
$$

be unbounded Fredholm operators, defined locally, such that

$$
\begin{aligned}
& \tilde{D}_{1}= \begin{cases}D_{1}, & \text { on } A_{1} ; \\
D_{2}, & \text { on } B_{2} ;\end{cases} \\
& \tilde{D}_{2}= \begin{cases}D_{2}, & \text { on } A_{2} ; \\
D_{1}, & \text { on } B_{1} .\end{cases}
\end{aligned}
$$

Theorem B-1. Under the hypothesis just described, we have

$$
\operatorname{ind}\left(D_{1}\right)+\operatorname{ind}\left(D_{2}\right)=\operatorname{ind}\left(\tilde{D}_{1}\right)+\operatorname{ind}\left(\tilde{D}_{2}\right) .
$$

Proof: Choose square roots of partitions of unity

$$
\phi_{1}^{2}+\psi_{1}^{2}=1 \quad \phi_{2}^{2}+\psi_{2}^{2}=1
$$

subordinate to $\left(A_{1}, B_{1}\right)$ and $\left(A_{2}, B_{2}\right)$. Choose them so that

$$
\begin{equation*}
\phi_{1}=\phi_{2} \text { and } \psi_{1}=\psi_{2} \text { on } A_{1} \cap B_{1}=A_{2} \cap B_{2} . \tag{B.1}
\end{equation*}
$$

We define maps

$$
\begin{aligned}
& \Phi: L^{2}\left(X_{1}\right) \oplus L^{2}\left(X_{2}\right) \rightarrow L^{2}\left(\tilde{X}_{1}\right) \oplus L^{2}\left(\tilde{X}_{2}\right) \\
& \Psi: L^{2}\left(\tilde{X}_{1}\right) \oplus L^{2}\left(\tilde{X}_{2}\right) \rightarrow L^{2}\left(X_{1}\right) \oplus L^{2}\left(X_{2}\right)
\end{aligned}
$$

which in matrix form are written as

$$
\Phi=\left[\begin{array}{cc}
\phi_{1} & \psi_{2} \\
-\psi_{1} & \phi_{2}
\end{array}\right] \text { and } \Psi=\left[\begin{array}{cc}
\phi_{1} & -\psi_{1} \\
\psi_{2} & \phi_{2}
\end{array}\right] .
$$

Notice that outside of $A_{i} \cap B_{i}$, we clearly have $\psi_{1} \phi_{1}=\psi_{2} \phi_{2}$. Equation (B.1) shows that this equality is also true in the intersection. Hence

$$
\Phi \Psi=\left[\begin{array}{cc}
\phi_{1}^{2}+\psi_{2}^{2} & -\psi_{1} \phi_{1}+\psi_{2} \phi_{2} \\
-\psi_{1} \phi_{1}+\psi_{2} \phi_{2} & \psi_{1}^{2}+\phi_{2}^{2}
\end{array}\right]=\left[\begin{array}{cc}
1 & \\
& 1
\end{array}\right],
$$

and $\Phi$ and $\Psi$ are inverse of each other. They are in fact isometries. Indeed, we have

$$
\begin{aligned}
\left\|\Phi\left(f_{1}, f_{2}\right)\right\|^{2}= & \int_{\tilde{X}_{1}}\left|\phi_{1} f_{1}+\psi_{2} f_{2}\right|^{2}+\int_{\tilde{X}_{2}}\left|\phi_{2} f_{2}-\psi_{1} f_{1}\right|^{2} \\
= & \int_{\tilde{X}_{1}} \phi_{1}^{2}\left|f_{1}\right|^{2}+\int_{\tilde{X}_{1}} \psi_{2}^{2}\left|f_{2}\right|^{2}+2 \int_{\tilde{X}_{1}} \phi_{1} \psi_{2}\left\langle f_{1}, f_{2}\right\rangle \\
& +\int_{\tilde{X}_{2}} \phi_{2}^{2}\left|f_{2}\right|^{2}+\int_{\tilde{X}_{2}} \psi_{1}^{2}\left|f_{1}\right|^{2}-2 \int_{\tilde{X}_{2}} \phi_{2} \psi_{1}\left\langle f_{2}, f_{1}\right\rangle \\
= & \int_{X_{1}}\left|f_{1}\right|^{2}+\int_{X_{2}}\left|f_{2}\right|^{2}=\left\|\left(f_{1}, f_{2}\right)\right\|^{2} .
\end{aligned}
$$

Consider now $D=D_{1} \oplus D_{2}$ and $\tilde{D}=\tilde{D}_{1} \oplus \tilde{D}_{2}$. Then

$$
\operatorname{ind}(D)=\operatorname{ind}\left(D_{1}\right)+\operatorname{ind}\left(D_{2}\right) \quad \text { and } \quad \operatorname{ind}(\tilde{D})=\operatorname{ind}\left(\tilde{D}_{1}\right)+\operatorname{ind}\left(\tilde{D}_{2}\right)
$$

We pull back $\tilde{D}$ to $L^{2}\left(X_{1}\right) \oplus L^{2}\left(X_{2}\right)$ and compare it to $D$. Should the difference $\Psi \tilde{D} \Phi-D$ be compact, the theorem would be proved. We proceed:

$$
\begin{aligned}
\Psi \tilde{D} \Phi & =\left[\begin{array}{cc}
\phi_{1} & -\psi_{1} \\
\psi_{2} & \phi_{2}
\end{array}\right]\left[\begin{array}{cc}
\tilde{D}_{1} & \\
& \tilde{D}_{2}
\end{array}\right]\left[\begin{array}{cc}
\phi_{1} & \psi_{2} \\
-\psi_{1} & \phi_{2}
\end{array}\right] \\
& =\left[\begin{array}{cc}
\phi_{1} \tilde{D}_{1} \phi_{1}+\psi_{1} \tilde{D}_{2} \psi_{1} & \phi_{1} \tilde{D}_{1} \psi_{2}-\psi_{1} \tilde{D}_{2} \phi_{2} \\
\psi_{2} \tilde{D}_{1} \phi_{1}-\phi_{2} \tilde{D}_{2} \psi_{1} & \psi_{2} \tilde{D}_{1} \psi_{2}+\phi_{2} \tilde{D}_{2} \phi_{2}
\end{array}\right] \\
& =\left[\begin{array}{ll}
\phi_{1} D_{1} \phi_{1}+\psi_{1} D_{2} \psi_{1} & \phi_{1} D_{1} \psi_{2}-\psi_{1} D_{2} \phi_{2} \\
\psi_{2} D_{1} \phi_{1}-\phi_{2} D_{2} \psi_{1} & \psi_{2} D_{1} \psi_{2}+\phi_{2} D_{2} \phi_{2}
\end{array}\right] \\
& =D+K .
\end{aligned}
$$

Since each entry of $K$ is supported on the product $\left(A_{1} \cap B_{1}\right) \times\left(A_{2} \cap B_{2}\right)$, which is compact, the operator $K$ is compact and the proof is now complete.

## Appendix C

## An abstract non-sense lemma

Suppose we are given four exact sequences $\alpha, \beta, \gamma, \delta$ interlaced in the braided diagram

and such that all triangles and squares commute.
Lemma C-1. The sequence

$$
0 \longrightarrow A \xrightarrow{\epsilon_{1}} B \xrightarrow{\epsilon_{2}} V_{1} \oplus V_{2} \xrightarrow{\epsilon_{3}} C \xrightarrow{\epsilon_{4}} D \longrightarrow 0
$$

coming out of Diagram (C.I), with the obvious choice of maps

$$
\epsilon_{1}=\beta_{1} \alpha_{1}=\delta_{1} \gamma_{1} \quad \epsilon_{2}=\left[\begin{array}{l}
\delta_{2} \\
\beta_{2}
\end{array}\right] \quad \epsilon_{3}=\left[\begin{array}{ll}
\alpha_{3} & -\gamma_{3}
\end{array}\right] \quad \epsilon_{4}=\beta_{4} \alpha_{4}=\delta_{4} \gamma_{4},
$$

is exact.
Before proving this lemma, let's observe how it is used in the main text of this thesis on page 88 Lemma 8.2-2 tells us that for the weights $\delta, \eta$ situated in adjacent open squares of $\mathbb{R}^{2} \backslash \mathfrak{G}_{A^{\prime}}$, and separated by the wall $\{\mu\} \times \mathbb{R}$ or $\mathbb{R} \times\{\lambda\}$, we have an exact sequence

$$
\begin{equation*}
0 \longrightarrow V(\delta) \longrightarrow V(\eta) \longrightarrow W_{\mu} \longrightarrow K(-\delta)^{*} \longrightarrow K(-\eta)^{*} \longrightarrow 0 \tag{C.2}
\end{equation*}
$$

Suppose now we have the following choice of weights:

| $\delta$ | $\delta$ |
| :--- | :--- |
| $\delta$ | $\delta$ |
| $-\mu$ |  |

Denote $\iota$ any inclusion map, and $L_{\mu}^{ \pm}$the maps

$$
L_{\mu}^{+}(\phi)=\lim _{t \rightarrow \infty} e^{\mu t} \phi, \quad \text { and } \quad L_{\mu}^{-}=\lim _{t \rightarrow-\infty} e^{\mu t} \phi .
$$

Then sequences akin to Sequence (C.2) fit in a diagram

similar to Diagram (C.1).
Suppose $\phi V(\delta)$, and $\psi \in K\left(-\delta_{\jmath}\right)$. Then

$$
\begin{aligned}
0 & =\left\langle\mathscr{D}_{A^{\prime}} \phi, \psi\right\rangle-\left\langle\phi, \mathscr{D}_{A^{\prime}} \psi\right\rangle \\
& =\left.\langle\psi, \nu \psi\rangle\right|_{-\infty} ^{\infty} \\
& =\lim _{t \rightarrow \infty}\left\langle e^{-\mu t} \phi, \nu e^{\mu t} \psi\right\rangle-\lim _{t \rightarrow-\infty}\left\langle e^{\mu t} \phi, \nu e^{-\mu t} \psi\right\rangle \\
& =\left(L_{\mu}^{+*} L_{-\mu}^{+}(\phi)-L_{-\mu}^{-} L_{\mu}^{-}(\phi)\right)(\psi),
\end{aligned}
$$

hence the middle square commutes. It is quite obvious that all the other squares and triangles commute. The application of Lemma C -1 gives an exact sequence

$$
0 \longrightarrow V\left(\delta_{\lrcorner}\right) \longrightarrow V\left(\left) \longrightarrow W_{\mu} \oplus W_{-\mu} \longrightarrow K\left(-\delta_{\lrcorner}\right)^{*} \longrightarrow K(-\delta)^{*} \longrightarrow 0\right.\right.
$$

In particular, the sets of weights

$\begin{aligned} \text { at } z & \neq w \\ \Gamma_{+} & =\Gamma_{-}\end{aligned}$

at $z \neq w$
$\Gamma_{+} \neq \Gamma_{-}$

$\begin{aligned} \text { at } z & =w \\ \Gamma_{+} & =\Gamma_{-}\end{aligned}$
yield for $A^{\prime}=A_{z}$ the exact sequences

$$
\begin{gathered}
0 \longrightarrow V_{z} \longrightarrow V_{z} \longrightarrow W_{\lambda} \oplus W_{-\lambda} \longrightarrow K_{z}{ }^{*} \longrightarrow K_{z}{ }^{*} \longrightarrow 0, \\
0 \longrightarrow V_{z} \longrightarrow V_{z} \longrightarrow W_{\lambda} \oplus W_{-\lambda} \longrightarrow K_{z}{ }^{*} \longrightarrow K_{z}{ }^{*} \longrightarrow 0, \\
\text { for } \lambda \neq 0, \\
0 \longrightarrow V_{w} \longrightarrow V_{w} \longrightarrow W_{0} \oplus W_{0} \longrightarrow \bar{K}_{w}{ }^{*} \longrightarrow K_{w}{ }^{*} \longrightarrow 0 .
\end{gathered}
$$

We can now proceed to the postponed proof of Lemma C-1.

Proof: There is nothing deep in this proof: it is only a diagram chase. It is included here for completeness.
The sequence is obviously exact at $A$ and $D$ since any composition of injective maps is surjective, and any composition of injective maps is surjective.
The compositions

$$
\begin{aligned}
& \epsilon_{2} \epsilon_{1}=\left[\begin{array}{l}
\delta_{2} \\
\beta_{2}
\end{array}\right] \epsilon_{1}=\left[\begin{array}{l}
\delta_{2} \delta_{1} \gamma_{1} \\
\beta_{2} \beta_{1} \alpha_{1}
\end{array}\right]=0, \\
& \epsilon_{3} \epsilon_{2}=\left[\begin{array}{ll}
\alpha_{3} & -\gamma_{3}
\end{array}\right]\left[\begin{array}{l}
\delta_{2} \\
\beta_{2}
\end{array}\right]=\alpha_{3} \delta_{2}-\gamma_{3} \beta_{2}=0, \text { and } \\
& \epsilon_{4} \epsilon_{3}=\left[\begin{array}{ll}
\epsilon_{4} \alpha_{3} & -\epsilon_{4} \gamma_{3}
\end{array}\right]=\left[\begin{array}{ll}
\beta_{4} \alpha_{4} \alpha_{3} & -\delta_{4} \gamma_{4} \gamma_{3}
\end{array}\right]=0
\end{aligned}
$$

ensure that $\operatorname{Im}\left(\epsilon_{j}\right) \operatorname{ker}\left(\epsilon_{j}\right)$. We now prove the other inclusions.
To simplify notation, every element denoted by a small letter belongs to the space denoted by the corresponding capital letter. For example, $b \in B, x_{1} \in X_{1}, c_{1} \in C$.
Suppose $b \in \operatorname{ker}\left(\epsilon_{2}\right)$. Then $\delta_{2}(b)=\beta_{2}(b)=0$ hence $\beta_{1}\left(x_{1}\right)=b=\delta_{1}\left(y_{1}\right)$. But then $\alpha_{2}\left(x_{1}\right)=$ $\delta_{2} \beta_{1}\left(x_{1}\right)=0$ hence $x_{1}=\alpha_{1}(a)$, and thus $b=\beta_{1} \alpha_{1}(a)$ and $\operatorname{ker}\left(\epsilon_{2}\right) \subset \operatorname{Im}\left(\epsilon_{1}\right)$ and the sequence is exact at $B$.
Suppose $\left(v_{1}, v_{2}\right) \operatorname{ker}\left(\epsilon_{3}\right)$, or equivalently $\alpha_{3}\left(v_{1}\right)-\gamma_{3}\left(v_{2}\right)=0$. Then

$$
\begin{aligned}
& 0=\alpha_{4}\left(\alpha_{3}\left(v_{1}\right)-\gamma_{3}\left(v_{2}\right)\right)=-\beta_{3}\left(v_{2}\right), \\
& 0=\gamma_{4}\left(\alpha_{3}\left(v_{1}\right)-\gamma_{3}\left(v_{2}\right)\right)=-\delta_{3}\left(v_{1}\right) .
\end{aligned}
$$

But then, because the $\beta$ and $\delta$ sequences are exact, we have $v_{2}=\beta_{2}\left(b_{2}\right)$ and $v_{1}=\delta_{2}\left(b_{1}\right)$. As $\gamma_{3} \beta_{2}\left(b_{1}\right)=\alpha_{3} \delta_{2}\left(b_{1}\right)=\alpha_{3}\left(v_{1}\right)=0$, we have $\beta_{2}\left(b_{1}\right) \in \operatorname{ker}\left(\gamma_{3}\right)=\operatorname{Im}\left(\gamma_{2}\right)$ hence $\beta_{2}\left(b_{1}\right)=\gamma_{2}\left(y_{1}\right)$. Similarly, $\delta_{2}\left(b_{2}\right)=\alpha_{2}\left(x_{1}\right)$. But then

$$
\begin{aligned}
\epsilon_{2}\left(b_{1}-\delta_{1}\left(y_{1}\right)+b_{2}-\beta_{1}\left(x_{1}\right)\right) & =\left[\begin{array}{c}
\delta_{2}\left(b_{1}\right)-\delta_{2} \delta_{1}\left(y_{1}\right)+\delta_{2}\left(b_{2}\right)-\delta_{2} \beta_{1}\left(x_{1}\right) \\
\beta_{2}\left(b_{1}\right)-\beta_{2} \delta_{1}\left(y_{1}\right)+\beta_{2}\left(b_{2}\right)-\beta_{2} \beta_{1}\left(x_{1}\right)
\end{array}\right] \\
& =\left[\begin{array}{l}
v_{1}-0+\alpha_{2}\left(x_{1}\right)-\delta_{2} \beta_{1}\left(x_{1}\right) \\
\gamma_{2}\left(y_{1}\right)-\beta_{2} \delta_{1}\left(y_{1}\right)+v_{2}-0
\end{array}\right] \\
& =\left[\begin{array}{l}
v_{1} \\
v_{2}
\end{array}\right] .
\end{aligned}
$$

Hence $\operatorname{ker}\left(\epsilon_{3}\right) \subset \operatorname{Im}\left(\epsilon_{2}\right)$ and the sequence is exact at $V_{1} \oplus V_{2}$.
Suppose $\epsilon_{4}(c)=0$. Then $\delta_{4} \gamma_{4}(c)=0$ but $\operatorname{ker}\left(\delta_{4}\right)=\operatorname{Im}\left(\delta_{3}\right)$ hence $\gamma_{4}(c)=\delta_{3}\left(v_{1}\right)=\gamma_{4} \alpha_{3}\left(v_{1}\right)$ and thus $\alpha_{3}\left(v_{1}\right)-c \in \operatorname{ker}\left(\gamma_{4}\right)=\operatorname{Im}\left(\gamma_{3}\right)$. Hence $\alpha_{3}\left(v_{1}\right)-c=\gamma_{3}\left(v_{2}\right)$ and then $c=\epsilon_{3}\left(v_{1}, v_{2}\right)$. Hence $\operatorname{ker}\left(\epsilon_{4}\right) \subset \operatorname{Im}\left(\epsilon_{3}\right)$ and the sequence is exact at $D$.
The proof is now complete.

## Appendix D

## Dirac operators and conformal change of metric

Consider a spin manifold $M$ of dimension $n$. The Dirac operators $D$ and $D^{\prime}$ associated to the conformally related metrics $g$ and $g^{\prime}=e^{2 f} g$ on $M$ are, once a spinor bundle is chosen, linked by the formula

$$
\begin{equation*}
D^{\prime}=e^{-\frac{n+1}{2} f} \circ D \circ e^{\frac{n-1}{2} f} . \tag{D.1}
\end{equation*}
$$

In 1986, Bourguignon in [Bou86, p. 339, Prop. 10] had the above formula and claimed that Hitchin in Hit74] had it wrong in 1974. Then, Lawson and Michelsohn, in their wonderful book [LM89, p. 134], had a $n-1$ on the left hand side instead of a $n+1$. Finally, in 1990, in their inspiring book, Donaldson and Kronheimer, in the case $n=4$, had a factor of $-1 / 2$ on the left hand side instead of a $-5 / 2$; see [DK90, p. 102]. The formula was however only used in [DK90] and [LM89] to see how the kernels of $D$ and $D^{\prime}$ are related, so no harm was done.

## A proof of the Formula

We now prove the formula D. 1 Let's denote the spinor bundles for $g$ and $g^{\prime}$ by $S$ and $S^{\prime}$. Let $\mu: \operatorname{Spin}(n) \rightarrow \operatorname{Aut}(\Delta)$ be the spin representation. Then

$$
\begin{aligned}
S & =P_{\text {Spin }}(M, g) \times_{\mu} \Delta, \text { and } \\
S^{\prime} & =P_{\text {Spin }}\left(M, g^{\prime}\right) \times_{\mu} \Delta .
\end{aligned}
$$

We assume that they are the "same" spin structure. Hence $S$ and $S^{\prime}$ are isomorphic as vector bundles but not as Clifford bundles.
The bundles $S$ and $S^{\prime}$ come equipped with extra structure. The Clifford multiplication for $g$ is a map $\rho: T M \rightarrow \operatorname{Aut}(S)$ satisfying $\rho(v)^{2}=-g(v, v)$. Define a new Clifford multiplication $\rho^{\prime}: T M \rightarrow \operatorname{Aut}(S)$ for $g^{\prime}$ by the formula $\rho^{\prime}=e^{f} \rho$. It is still skew-adjoint.
Let $e_{i}$ be an orthonormal frame for $g$ over an open set $U$ and $\omega_{i j}$ be the Levi-Civita connection matrix for that frame. Over $U$, the spin connection (see [LM89] p. 110]) is $d+\Omega$ with

$$
\Omega(V)=\frac{1}{4} \sum_{i, j} \omega_{i j}(V) \rho\left(e_{j}\right) \rho\left(e_{i}\right) .
$$

Set $e_{j}^{\prime}=e^{-f} e_{j}$. The $e_{j}^{\prime}$ form an orthonormal frame for $g^{\prime}$. Note that $\rho^{\prime}\left(e_{j}^{\prime}\right)=\rho\left(e_{j}\right)$.
As shown in [LM89 p. 133-134], we have

$$
\begin{align*}
\nabla_{X}^{\prime} Y & =\nabla_{X} Y+(X f) Y+(Y f) X-g(X, Y)_{\text {grad }}(f), \text { and }  \tag{D.2}\\
\omega_{i j}^{\prime}(V) & =\omega_{i j}(V)+\left(e_{j} f\right) g\left(V, e_{i}\right)-\left(e_{i} f\right) g\left(V, e_{j}\right) .
\end{align*}
$$

Duplicating the definition of $\Omega$, we set

$$
\Omega^{\prime}(V):=\frac{1}{4} \sum_{i, j}\left(\omega_{i j}(V) \rho\left(e_{j}\right) \rho\left(e_{i}\right)+\left(e_{j} f\right) g\left(V, e_{i}\right) \rho\left(e_{j}\right) \rho\left(e_{i}\right)-\left(e_{i} f\right) g\left(V, e_{j}\right) \rho\left(e_{j}\right) \rho\left(e_{i}\right)\right) .
$$

This expression simplifies to

$$
\begin{align*}
\Omega^{\prime}(V) & =\Omega(V)+\frac{1}{4} \rho(\operatorname{grad}(f)) \rho(V)-\frac{1}{4} \rho(V) \rho(\operatorname{grad}(f)) \\
& =\Omega(V)-\frac{1}{2} \rho(V) \rho(\operatorname{grad}(f))-\frac{1}{2} g(V, \operatorname{grad}(f)) . \tag{D.3}
\end{align*}
$$

Let's now check that the connection $\nabla^{\prime}$ induced by $\Omega^{\prime}$ is compatible with the Levi-Civita connection of ( $M, g^{\prime}$ ).
Notice first that

$$
\begin{aligned}
\nabla_{X}\left(\rho^{\prime}(Y) s\right) & =\nabla_{X}\left(e^{f} \rho(Y) s\right) \\
& =(X f) \rho^{\prime}(Y) s+e^{f} \nabla_{X}(\rho(Y) s) .
\end{aligned}
$$

Since $S$ is a Clifford bundle for $(M, g)$, we have

$$
\begin{equation*}
\nabla_{X}\left(\rho^{\prime}(Y) s\right)=\rho^{\prime}((X f) Y) s+\rho^{\prime}\left(\nabla_{X} Y\right) s+\rho^{\prime}(Y) \nabla_{X} s \tag{D.4}
\end{equation*}
$$

Notice now that

$$
\begin{aligned}
\rho(X) \rho(\operatorname{grod}(f)) \rho^{\prime}(Y) s & =-\rho(X) \rho^{\prime}(Y) \rho(\operatorname{grod}(f))-2 e^{f} \rho(X)(Y f) s \\
& =\left(\rho^{\prime}(Y) \rho(X)+2 e^{f} g(X, Y)\right) \rho(\operatorname{grad}(f)) s-2 \rho^{\prime}((Y f) X) s,
\end{aligned}
$$

so that

$$
\begin{equation*}
-\frac{1}{2} \rho(X) \rho(g r o d(f)) \rho^{\prime}(Y) s=-\frac{1}{2} \rho^{\prime}(Y) \rho(X) \rho(\operatorname{grad}(f)) s+\rho^{\prime}((Y f) X-g(X, Y) \operatorname{grad}(f)) s \tag{D.5}
\end{equation*}
$$

Putting all these computations together, we find

$$
\begin{aligned}
& \nabla_{X}^{\prime}\left(\rho^{\prime}(Y) s\right) \stackrel{\text { D. } 3}{=} \nabla_{X}\left(\rho^{\prime}(Y) s\right)-\frac{1}{2} \rho(X) \rho(g r a d \\
& \stackrel{\text { D. } 4}{=} \rho^{\prime}\left(\nabla_{X} Y+(X f) Y\right) s+\frac{1}{2}(X f) \rho^{\prime}(Y) s \\
& \text { D.2D. } 5 \rho^{\prime}(Y) \nabla_{X} s-\frac{1}{2} \rho(X) \rho\left(\nabla_{X}^{\prime} Y\right) s+\rho^{\prime}(y)\left(\nabla_{X} s-\frac{1}{2} \rho(X) \rho(\text { grad }(f)) s-\frac{1}{2}(X f) s\right) \\
&=\frac{1}{2}(X f) \rho^{\prime}(Y) s \\
& \rho^{\prime}\left(\nabla_{X}^{\prime} Y\right) s+\rho^{\prime}(Y) \nabla_{X}^{\prime} s .
\end{aligned}
$$

We also need to check that the connection $\nabla^{\prime}$ is compatible with the hermitian metric on $S$. Notice first that

$$
\begin{aligned}
\left\langle\rho(V) \rho(\operatorname{grad}(f)) s_{1}, s_{2}\right\rangle & +\left\langle s_{1}, \rho(V) \rho(\operatorname{grad}(f)) s_{2}\right\rangle \\
& =\left\langle s_{1}, \rho(\operatorname{grad}(f)) \rho(V) s_{2}\right\rangle+\left\langle s_{1}, \rho(V) \rho(\operatorname{grad}(f)) s_{2}\right\rangle=-2(V f)\left\langle s_{1}, s_{2}\right\rangle .
\end{aligned}
$$

Thus, we have

$$
\begin{aligned}
\left\langle\nabla_{V}^{\prime} s_{1}, s_{2}\right\rangle+\left\langle s_{1}, \nabla_{V}^{\prime} s_{2}\right\rangle= & \left\langle\nabla_{V} s_{1}, s_{2}\right\rangle+\left\langle s_{1}, \nabla_{V} s_{2}\right\rangle \\
& -\frac{1}{2}\left(\left\langle\rho(V) \rho(\text { grad }(f)) s_{1}, s_{2}\right\rangle+\left\langle s_{1}, \rho(V) \rho(\operatorname{grad}(f)) s_{2}\right\rangle\right) \\
& -\frac{1}{2}\left(\left\langle(V f) s_{1}, s_{2}\right\rangle+\left\langle s_{1},(V f) s_{2}\right\rangle\right) \\
= & \left\langle\nabla_{V} s_{1}, s_{2}\right\rangle+\left\langle s_{1}, \nabla_{V} s_{2}\right\rangle+(V f)\left\langle s_{1}, s_{2}\right\rangle-(V f)\left\langle s_{1}, s_{2}\right\rangle \\
= & V\left\langle s_{1}, s_{2}\right\rangle,
\end{aligned}
$$

since $\nabla$ is compatible with $\langle$,$\rangle .$
We are now finished proving that $S$ with the connection $\nabla^{\prime}$ and the Clifford multiplication $\rho^{\prime}$ is really a Clifford bundle for $\left(M, g^{\prime}\right)$.
Let $V:=\operatorname{Hom}_{C l}\left(S^{\prime}, S\right)$. We have that as Clifford bundles for $\left(M, g^{\prime}\right)$, the bundle $S$ and $S^{\prime} \otimes V$ are isomorphic.
By Schur's lemma, the bundle $V$ is a smooth complex line bundle. Since $S$ and $S^{\prime}$ are isomorphic as smooth complex vector bundles, $V$ is trivial. But more than that, the connection on $V$ is trivial as we now show.
Let $\rho^{\prime}$ also denote the Clifford multiplication on $S^{\prime}$ and $\Omega^{\prime}$ denote the spin connection on $S^{\prime}$. Such an abuse of notation is harmless if we are careful. We have again

$$
\Omega^{\prime}(V)=\frac{1}{4} \sum_{i, j} \omega_{i j}^{\prime}(V) \rho^{\prime}\left(e_{j}^{\prime}\right) \rho^{\prime}\left(e_{i}^{\prime}\right) .
$$

Let $f$ be a section of $V$ and $s$ a section of $S^{\prime}$. Then since $f$ commutes with Clifford multiplication, we have

$$
\begin{aligned}
(\nabla f)(s) & =\nabla(f(s))-f(\nabla s) \\
& =d(f(s))+\Omega^{\prime} \cdot f(s)-f(d s)-f\left(\Omega^{\prime} \cdot s\right) \\
& =d f(s),
\end{aligned}
$$

and thus our claim is proved.
So we can really work with the bundle $S$ when studying the relationship between $D$ and $D^{\prime}$. Let's do that.
We have

$$
D^{\prime} \phi=e^{-f}\left(D \phi+\frac{n-1}{2} \rho(\operatorname{grad}(f)) \phi\right) .
$$

Indeed,

$$
\begin{aligned}
D^{\prime} \phi & =\sum \rho^{\prime}\left(e_{i}^{\prime}\right) \nabla_{e_{i}^{\prime}}^{\prime} \phi \\
& =e^{-f} \sum \rho\left(e_{i}\right) \nabla_{e_{i}}^{\prime} \phi \\
& =e^{-f}\left(D \phi-\frac{1}{2} \sum \rho\left(e_{i}\right) \rho\left(e_{i}\right) \rho(\operatorname{grad}(f)) \phi-\frac{1}{2} \sum \rho\left(e_{i}\right) g\left(e_{i}, g r a d(f)\right) \phi\right) \\
& =e^{-f}\left(D \phi+\frac{n}{2} \rho(\operatorname{grad}(f)) \phi-\frac{1}{2} \rho(\operatorname{grad}(f)) \phi\right) \\
& =e^{-f}\left(D \phi+\frac{n-1}{2} \rho(\operatorname{grad}(f)) \phi\right) .
\end{aligned}
$$

But then,

$$
\begin{aligned}
D\left(e^{\frac{n-1}{2} f} \phi\right) & =e^{\frac{n-1}{2} f}\left(D \phi+\frac{n-1}{2} \rho(\operatorname{grad}(f)) \phi\right) \\
& =e^{\frac{n-1}{2} f} e^{f} D^{\prime} \phi
\end{aligned}
$$

whence $D^{\prime} \phi=e^{-\frac{n+1}{2} f} D\left(e^{\frac{n-1}{2} f} \phi\right)$, as wanted.

## A confirmation

The proof just presented should at least convince us that there exists such a formula. To confirm that the factors are correct, suppose that $D$ and $D^{\prime}$ are linked by formula D. 1 and that $D$ is self-adjoint for the $L^{2}$ inner product on $(M, g)$. We want to prove now that $D^{\prime}$ is self-adjoint on $\left(M, g^{\prime}\right)$. Recall that $d v o l^{\prime}=e^{n f} d v o l$. Thus,

$$
\begin{aligned}
\left\langle D^{\prime} \phi, \psi\right\rangle^{\prime} & =\int\left\langle D^{\prime} \phi, \psi\right\rangle d v o l^{\prime} \\
& =\int\left\langle e^{-(n+1) f / 2} D e^{(n-1) f / 2} \phi, \psi\right\rangle e^{n f} d v o l \\
& =\left\langle D e^{(n-1) f / 2} \phi, e^{(n-1) f / 2} \psi\right\rangle \\
& =\left\langle e^{(n-1) f / 2} \phi, D e^{(n-1) f / 2} \psi\right\rangle \\
& =\left\langle\phi, D^{\prime} \psi\right\rangle^{\prime}
\end{aligned}
$$

That is it.
In fact, this computation and the triviality of $V$ show that any formula which has one of the factors must also have the other up to a constant.

## Appendix E

## Weighted Sobolev spaces on $\mathbb{R}^{n}$, Bartnik's presentation

In this appendix, we visit and work through a part of Robert Bartnik's paper "The Mass of an Asymptotically Flat Manifold," [Bar86]. The part we are concerned with deals with weighted Sobolev spaces on $\mathbb{R}^{n}$ for $n \geq 3$, and Fredholmness of certain 2 nd order elliptic partial differential operators. This appendix is an companion to Bartnik's writing, merely adding proofs that were lacking.

## Weighted Sobolev Spaces

Set

$$
\sigma(x):=\left(1+|x|^{2}\right)^{1 / 2}
$$

The space $L_{\delta}^{p}$ is the space of measurable functions in $L_{\text {loc }}^{p}$ which are finite in the $\left\|\left\|\|_{p, \delta}\right.\right.$-norm:

$$
\|u\|_{p, \delta}= \begin{cases}\left(\int_{\mathbb{R}^{n}} \sigma^{-\delta p-n}(x)|u(x)|^{p} d x\right)^{1 / p}, & 1 \geq p<\infty \\ \text { ess } \sup _{\mathbb{R}^{n}} \sigma^{-\delta}|u|, & p=\infty\end{cases}
$$

When $p<\infty$, the usual $L^{p}$ space arise as $L_{-n / p}^{p}$. In fact, we have an even stronger proposition.
Proposition E-1. The following map is an isometry:

$$
\begin{aligned}
L^{p} & \rightarrow L_{\delta}^{p} \\
f & \mapsto \sigma^{\delta+n / p} f
\end{aligned}
$$

The space $W_{\delta}^{k, p}$ is then defined as the space of functions in $L_{\delta}^{p}$ with weak derivatives in the appropriate weighted $L^{p}$ space. The norm is

$$
\|u\|_{k, p, \delta}=\sum_{l \leq k}\left\|D^{l} u\right\|_{p, \delta-l}
$$

These spaces are nicely set-up for the interesting theorems to be in some sense independent of $n$, as will become apparent later with the various inequalities and embeddings.

For one thing, we have that

$$
\left\|\sigma^{a}\right\|_{\infty, \delta}=\text { ess sup } \sigma^{a-\delta}= \begin{cases}1, & \text { if } a \leq \delta \\ \infty, & \text { otherwise }\end{cases}
$$

Similarly, we have, for $1 \leq p<\infty$,

$$
\begin{aligned}
\left\|\sigma^{a}\right\|_{p, \delta}^{p} & =\int_{\mathbb{R}^{n}} \sigma^{(a-\delta) p-n} \\
& =\omega_{n} \int_{0}^{\infty} \sigma(r)^{(a-\delta) p-n} r^{n-1} d r .
\end{aligned}
$$

The integral on the last line is finite if and only if its $[1, \infty)$ part is finite. It is the case if and only if $a<\delta$. In fact, we have the following stronger proposition.

Proposition E-2. We have

$$
\sigma^{a} \in W_{\delta}^{k, \infty} \quad \Longleftrightarrow \quad a \leq \delta
$$

and for $1 \leq p<\infty$, we have

$$
\sigma^{a} \in W_{\delta}^{k, p} \quad \Longleftrightarrow \quad a<\delta
$$

Proof: We already proved the case $k=0$. Notice now that

$$
\frac{\partial}{\partial x_{i}}\left(\sigma^{a}\right)=a \sigma^{a-2} x_{i}
$$

and similarly, these exist homogeneous polynomials $p_{\alpha}$ of degree $|\alpha|$ such that

$$
\frac{\partial}{\partial x^{\alpha}}\left(\sigma^{a}\right)=\sigma^{a-2|\alpha|} p_{\alpha} .
$$

Thus

$$
\left\|\frac{\partial \sigma^{a}}{\partial x^{\alpha}}\right\|_{p, k, \delta-|\alpha|} \leq C\left\|\sigma^{a-|\alpha|}\right\|_{p, \delta-|\alpha|} .
$$

This inequality imply that $\sigma^{a} \in L_{\delta}^{p}$ implies $\sigma^{a} \in W_{\delta}^{k, p}$.
Many different choices for weight function $\sigma$ could have been considered. Apparently some work has been done with exponential weights instead of the "polynomial" weight that we use here.
We define similarly the norm $\left\|\|_{k, p, \delta}^{\prime}\right.$ and the spaces $L_{\delta}^{\prime p}$ and $W_{\delta}^{\prime k, p}$ of functions on $\mathbb{R}^{n} \backslash\{0\}$ by changing the weight function to $r(x):=|x|$.
In these modified Sobolev spaces, scaling becomes homogeneous. Indeed, set

$$
u_{R}(x):=u(R x) ;
$$

then

$$
\left\|u_{R}\right\|_{k, p}^{\prime}=R^{\delta}\|u\|_{k, p, \delta}^{\prime} .
$$

Of course any norm without a weight refers to a usual Sobolev norm.
Set $A_{R}=B_{2 R} \backslash B_{R}$ and $E_{R}=\mathbb{R}^{n} \backslash B_{R}$. We use an obvious notation for restriction over subset of $\mathbb{R}^{n}$. Then the norm $u \mapsto\|u\|_{p, \delta ; A_{R}}$ is equivalent to the norm $u \mapsto R^{-\delta}\left\|u_{R}\right\|_{p ; A_{1}}$, and Bartnik writes

$$
\|u\|_{p, \delta ; A_{R}} \approx R^{-\delta}\left\|u_{R}\right\|_{p ; A_{1}},
$$

with constants not depending on $R$ but depending on $\delta$. This equivalence allows us to rescale and apply local estimates to prove part of the following theorem.

Theorem E-3. We have the following inequalities:

$$
\begin{gather*}
\text { If } p \leq q \text { and } \delta_{2}<\delta_{1} \text { then }\|u\|_{p, \delta_{1}} \leq c\|u\|_{q, \delta_{2}} .  \tag{E.1}\\
\text { If } 1 / p=1 / q+1 / r \text { and } \delta=\delta_{1}+\delta_{2} \text { then }\|u v\|_{p, \delta} \leq\|u\|_{q, \delta_{1}}\|v\|_{r, \delta_{2}} .  \tag{E.2}\\
\text { For any } \epsilon>0 \text {, there exists } C(\epsilon) \text { s.t. }\|u\|_{1, p, \delta} \leq \epsilon\|u\|_{2, p, \delta}+C(\epsilon)\|u\|_{0, p, \delta} . \tag{E.3}
\end{gather*}
$$

Proof: To prove Inequality (E.2), we write, using Proposition E-1,

$$
\|u v\|_{p, \delta}=\left\|\sigma^{-\delta-n / p} u v\right\|_{p}=\left\|\sigma^{-\delta_{1}-n / q} u \sigma^{-\delta_{2}-n / r} v\right\|_{p}
$$

Using the usual Hölder inequality, we obtain the result.
In the conditions under which we wish Inequality (E.1) to be true, there exists $r$ such that $1 / p=$ $1 / q+1 / r$. Thus, we can again use Proposition E-1 and the usual Hölder inequality to get

$$
\begin{aligned}
\|u\|_{p, \delta_{1}} & =\left\|\sigma^{-\delta_{1}-n / p} u\right\|_{p} \\
& =\left\|\left(\sigma^{-n / r} \sigma^{\delta_{2}-\delta_{1}}\right)\left(\sigma^{-\delta_{2}-n / q} u\right)\right\|_{p} \\
& \leq\left\|\sigma^{\delta_{2}-\delta_{1}}\right\|_{r, 0}\|u\|_{q, \delta_{2}}
\end{aligned}
$$

Since $\delta_{2}<\delta_{1}$, we know by Proposition E-2 that $\sigma^{\delta_{2}-\delta_{1}} \in L_{0}^{r}$ hence Inequality (E.1) is true.
We now rescale and apply local estimates to prove Inequality (E.3). We know that there exists constants $C_{1}, C_{2}$ not depending on $R$ such that

$$
\begin{gathered}
\|u\|_{1, p, \delta ; A_{R}} \leq C_{1} R^{-\delta}\left\|u_{R}\right\|_{1, p ; A_{1}} \text {, and } \\
R^{-\delta}\left\|u_{R}\right\|_{2, p ; A_{1}} \leq C_{2}\|u\|_{2, p, \delta ; A_{R}} .
\end{gathered}
$$

Set $\epsilon^{\prime}=\epsilon / C_{1} C_{2}$. We know from the local interpolation inequality that for some $C\left(\epsilon^{\prime}\right)$,

$$
\left\|u_{R}\right\|_{1, p ; A_{1}} \leq \epsilon^{\prime}\left\|u_{R}\right\|_{2, p ; A_{1}}+C\left\|u_{R}\right\|_{0, p ; A_{1}} .
$$

Thus

$$
\begin{aligned}
\|u\|_{1, p, \delta ; A_{R}} & \leq C_{1} R^{-\delta}\left\|u_{R}\right\|_{1, p ; A_{1}} \\
& \leq C_{1} R^{-\delta} \epsilon^{\prime}\left\|u_{R}\right\|_{2, p ; A_{1}}+R^{-\delta} C\left\|u_{R}\right\|_{0, p ; A_{1}} \\
& \leq \epsilon\|u\|_{2, p, \delta ; A_{R}}+C\|u\|_{0, p, \delta ; A_{R}} .
\end{aligned}
$$

Now, we need to patch all these interpolations together. Recall first that for all $p>0$, we have

$$
(a+b)^{p} \leq 2^{p-1}\left(a^{p}+b^{p}\right)
$$

Set $D_{0}:=B_{1}$ and, for $i>0$, set $D_{i}=A_{2^{i-1}}$. The $D_{i}$ are disjoint and cover $\mathbb{R}^{n}$. Set $u_{i}=\left.u\right|_{D_{i}}$.

The fact that

$$
\|u\|_{1, p, \delta} \leq \epsilon\|u\|_{2, p, \delta}+C(\epsilon)\|u\|_{0, p, \delta}
$$

follows from the computation

$$
\begin{aligned}
\|u\|_{1, p, \delta}^{p} & =\sum\left\|u_{i}\right\|_{1, p, \delta ; D_{i}}^{p} \\
& \leq \sum\left(\epsilon\left\|u_{i}\right\|_{2, p, \delta ; D_{i}}+C\left\|u_{i}\right\|_{0, p, \delta ; D_{i}}\right)^{p} \\
& \leq 2^{p-1} \epsilon^{p}\|u\|_{2, p, \delta}^{p}+C\|u\|_{0, p, \delta}^{p} .
\end{aligned}
$$

The proof for $p=\infty$ goes along the same lines.

The inequalities are great and useful, and, apart from the first one, are generalizations of what we have for usual Sobolev spaces.

Before studying the generalization of other classical powerful estimates, the Sobolev embedding theorems, we must generalize yet another type of space, the Hölder spaces. Define first, for $x \in \mathbb{R}^{n}$, the punctured ball $B(x)$ to be the set of all $y$ such that $0<4|x-y|<\sigma(x)$. For $0<\alpha \leq 1$, the weighted Hölder norm is defined by the equation

$$
\|u\|_{C_{\delta}^{0, \alpha}}=\sup _{x \in \mathbb{R}^{n}}\left(\sigma^{-\delta}(x)|u(x)|\right)+\sup _{x \in \mathbb{R}^{n}}\left(\sigma^{-\delta+\alpha}(x) \sup _{y \in B(x)} \frac{|u(x)-u(y)|}{|x-y|^{\alpha}}\right)
$$

Note that if $\|u\|_{C_{\delta}^{0, \alpha}}$ is finite, then $u$ is continuous. Indeed, close to $x_{0}$, we have

$$
\left|u\left(x_{0}\right)-u(y)\right| \leq C \sigma\left(x_{0}\right)^{-\alpha+\delta}\left|x_{0}-y\right|^{\alpha}
$$

We may now proceed.

Theorem E-4. Suppose $u \in W_{\delta}^{k, p}$. We have the following inequalities:

$$
\begin{align*}
& \text { If } n-k p>0 \text { and } p \leq q \leq n p /(n-k p) \text { then }\|u\|_{n p /(n-k p), \delta} \leq C\|u\|_{k, q, \delta}  \tag{E.4}\\
& \qquad \begin{array}{c}
\text { If } n-k p<0 \text { then }\|u\|_{\infty, \delta} \leq C\|u\|_{k, p, \delta} \\
\text { and }|u(x)|=o\left(r^{\delta}\right) \text { as } r \rightarrow \infty \\
\text { If } n-k p<0 \text { and } 0<\alpha \leq \min (1, k-n / p) \text { then }\|u\|_{C_{\delta}^{0, \alpha}} \leq C\|u\|_{k, p, \delta} \\
\text { and }\|u\|_{C_{\delta}^{0, \alpha}\left(A_{R}\right)}=o(1) \text { as } R \rightarrow \infty
\end{array} \tag{E.5}
\end{align*}
$$

Proof: Suppose first that $n-k p>0$. Then set $p^{*}:=n p /(n-k p)<\infty$. We have

$$
\begin{aligned}
\|u\|_{p^{*}, \delta ; A_{R}} & \leq C R^{-\delta}\left\|u_{R}\right\|_{p^{*} ; A_{1}} \\
& \leq C R^{-\delta}\left\|u_{R}\right\|_{k, q ; A_{1}} \text { (by the usual Sobolev inequality) } \\
& \leq C\|u\|_{k, q, \delta ; A_{R}}
\end{aligned}
$$

Writing $u=\sum u_{i}$ as in the proof of Theorem [E-3] we obtain

$$
\begin{aligned}
\|u\|_{p^{*}, \delta} & =\left(\sum\left\|u_{i}\right\|_{p^{*}, \delta}^{p^{*}}\right)^{1 / p^{*}} \\
& \leq C\left(\sum\left\|u_{i}\right\|_{k, q, \delta}^{p^{*}}\right)^{1 / p^{*}} \\
& \leq C\left(\sum\left\|u_{i}\right\|_{k, q, \delta}^{q}\right)^{1 / q} \\
& =C\|u\|_{k, q, \delta} .
\end{aligned}
$$

Inequality ( $\mathbb{E} .4$ is now proved.
Maybe it is worthwhile noting down the proof of the last inequality of this proof. In fact, it is sufficient to prove that $\left(1+x^{p^{*}}\right)^{1 / p^{*}} \leq\left(1+x^{q}\right)^{1 / q}$ for $x \geq 0$. The function

$$
f(x)=\frac{\left(1+x^{p^{*}}\right)^{q}}{\left(1+x^{q}\right)^{p^{*}}}
$$

has derivative

$$
\frac{x^{p^{*}-1} p^{*} q\left(1+x^{p^{*}}\right)^{q-1}\left(1+x^{q}\right)^{p^{*}}-p^{*} q x^{q-1}\left(1+x^{p^{*}}\right)^{q}\left(1+x^{q}\right)^{p^{*}-1}}{\left(1+x^{q}\right)^{2 p^{*}}} .
$$

Once we remove the common factors, which are anyway strictly positive, we have

$$
x^{p^{*}-1}\left(1+x^{q}\right)-x^{q-1}\left(1+x^{p^{*}}\right)=x^{p^{*}-1}-x^{q-1} \geq 0
$$

since $q \leq p^{*}$. Thus $f$ always increases. But obviously, it takes the value 1 at infinity.
The same scaling argument and application of the usual Sobolev inequality apply to prove Inequality (E.5). Of course, $u \in W_{\delta}^{k, p}$ imply then that

$$
\left|r(x)^{-\delta} u(x)\right| \leq C\|u\|_{\infty, \delta ; A_{|x|}} \leq C\|u\|_{k, p, \delta ; A_{|x|}}
$$

which converges to 0 . Thus the asymptotic behavior of Equation (E.6) is now proved.
To prove Inequality (E.7), we would like use the decomposition $u=\sum u_{i}$. That cannot work however. Yet something of the sort works.
As before, we have that

$$
\|u\|_{C_{\delta}^{0, \alpha}\left(A_{R}\right)} \approx R^{-\delta}\left\|u_{R}\right\|_{C_{\delta}^{0, \alpha}\left(A_{1}\right)}
$$

with constants not depending on $R$. But the usual theorem can be used and we have

$$
\begin{aligned}
\|u\|_{C_{\delta}^{0, \alpha}\left(A_{R}\right)} & \leq C R^{-\delta}\left\|u_{R}\right\|_{C_{\delta}^{0, \alpha}\left(A_{1}\right)} \\
& \leq C R^{-\delta}\left\|u_{R}\right\|_{k, p ; A_{1}} \\
& \leq C\|u\|_{k, p, \delta ; A_{R}} .
\end{aligned}
$$

Since $\left\|\|_{k, p ; B_{1}}\right.$ is equivalent to $\| \|_{k, p, \delta ; B_{1}}$, we have

$$
\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)} \leq C\|u\|_{k, p, \delta ; D_{i}},
$$

with the constant $C$ not depending on $i$.
Now the condition $y \in B(x)$ in the inner supremum becomes really essential. Indeed, we can see that if $x \in D_{i}$ then $B(x)$ is contained in $D_{i-1} \cup D_{i} \cup D_{i+1}$. Because of that, we can actually bound the $C_{\delta}^{0, \alpha}$-norm of a function by norms of restrictions.
Let's write $a(x, y)$ for $\sigma(x)^{-\delta+\alpha}|u(x)-u(y)| /|x-y|^{\alpha}$. Thus,

$$
\|u\|_{C_{\delta}^{0, \alpha}}=\|u\|_{\infty, \delta}+\sup _{x} \sup _{y \in B(x)} a(x, y) .
$$

Let $x \in D_{i}$. We split the "ball" $B(x)$ in three parts.
Suppose first that $y \in D_{i-1} \cap B(x)$. There is a point $z \in[x, y] \cap D_{i-1} \cap D_{i}$. As for any point in $D_{i}$, we have the relationship $4^{-1} \sigma(z) \leq \sigma(x) \leq 4 \sigma(z)$. Thus

$$
\begin{aligned}
a(x, y) & \leq a(x, z)+4^{|-\delta+\alpha|} a(z, y) \\
& \leq \sup _{z \in B(x) \cap D_{i}} a(x, z)+4^{|-\delta+\alpha|} \sup _{y \in B(x) \cap D_{i-1}} a(z, y) .
\end{aligned}
$$

But $B(x) \cap D_{i-1} \subset B(z) \cap D_{i-1}$. Thus

$$
a(x, y) \leq\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)}+4^{|-\delta+\alpha|}\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i-1}\right)} .
$$

Similarly, we have for $y \in D_{i+1} \cap B(x)$ that

$$
a(x, y) \leq\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)}+4^{|-\delta+\alpha|}\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i+1}\right)},
$$

and for $y \in D_{i} \cap B(x)$ that $a(x, y) \leq\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)}$.
Hence, for $x \in D_{i}$,

$$
\begin{aligned}
\sup _{y \in B(x)} a(x, y) & \leq \max \left(1,4^{|-\delta+\alpha|}\right)\left(\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i-1}\right)}+\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)}+\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i+1}\right)}\right) \\
& \leq C \sum\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)} .
\end{aligned}
$$

whence $\|u\|_{C_{\delta}^{0, \alpha}} \leq C \sum\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)}$.
But then,

$$
\begin{aligned}
\|u\|_{C_{\delta}^{0, \alpha}} & \leq C \sum\|u\|_{C_{\delta}^{0, \alpha}\left(D_{i}\right)} \\
& \leq C \sum\|u\|_{k, p, \delta ; D_{i}} \\
& \leq C\|u\|_{k, p, \delta},
\end{aligned}
$$

proving Equation (E.7).
Estimate ( $\overline{\text { E.8 }}$ ) is a consequence of Inequality ( $\overline{\text { E.7 }}$ ) for the domain $A_{R}$, and of the finiteness of $p$, which implies that $\|u\|_{k, p, \delta ; A_{R}}$ tends to 0 .

## Fredholm theory for second order operators asymptotic to $\Delta$

Of course, we introduce this big machinery in order to do some Fredholm theory for certain partial differential operators on $\mathbb{R}^{n}$. We consider in Bartnik's paper second order operator which are "asymptotic" to the Laplacian in the following sense.
Definition: The operator $u \rightarrow P u$ defined by

$$
P u=\sum_{i, j} a^{i j}(x) \partial_{i, j}^{2} u+\sum_{i} b^{i}(x) \partial_{i} u+c(x) u
$$

is asymptotic to $\Delta$ at rate $\tau \geq 0$ if there exist $n<q<\infty$ and constants $C_{1}, \lambda$ such that

$$
\begin{gathered}
\lambda|\xi|^{2} \leq \sum_{i, j} a^{i j} \xi_{i} \xi_{j} \leq \lambda^{-1}|\xi|^{2}, \text { and } \\
\left\|a^{i j}-\delta_{i j}\right\|_{1, q,-\tau}+\left\|b^{i}\right\|_{0, q,-1-\tau}+\|c\|_{0, q,-2-\tau} \leq C_{1} .
\end{gathered}
$$

Let's note that Bartnik uses the norm $L_{-2-\tau}^{q / 2}$-norm for $c$ instead of the $L_{-2-\tau}^{q}$-norm like we do. We should first check what are natural domain and codomain for such operators.

Theorem E-5. If $P$ is asymptotic to $\Delta$, then

$$
P: W_{\delta}^{2, p} \rightarrow W_{\delta-2}^{0, p}
$$

is bounded for $1 \leq p \leq q$ and every $\delta \in \mathbb{R}$.
Proof: We show that each piece is bounded. In the first place, we have

$$
\begin{aligned}
\left\|a^{i j} \partial_{i j}^{2} u\right\|_{0, p, \delta-2} & \leq\left\|\left(a^{i j}-\delta_{i j}\right) \partial_{i j}^{2} u\right\|_{0, p, \delta-2}+\left\|\delta_{i j} \partial_{i j}^{2} u\right\|_{0, p, \delta-2} \\
& \leq\left(\left\|a^{i j}-\delta_{i j}\right\|_{\infty, 0}+1\right)\left\|\partial_{i j}^{2} u\right\|_{0, p, \delta-2} \\
& \leq C\left(\left\|a^{i j}-\delta_{i j}\right\|_{\infty,-\tau}+1\right)\|u\|_{2, p, \delta} \\
& \leq C\left(\left\|a^{i j}-\delta_{i j}\right\|_{1, q,-\tau}+1\right)\|u\|_{2, p, \delta} \\
& \leq C\|u\|_{2, p, \delta} .
\end{aligned}
$$

When $p=q$, we have

$$
\begin{aligned}
\left\|b^{i} \partial_{i} u\right\|_{0, p, \delta-2} & \leq\left\|b^{i}\right\|_{p,-1}\left\|\partial_{i} u\right\|_{\infty, \delta-1} \\
& \leq C\left\|b^{i}\right\|_{q,-1-\tau}\left\|\partial_{i} u\right\|_{1, p, \delta-1} \\
& \leq C\|u\|_{2, p, \delta} .
\end{aligned}
$$

Suppose now that $p<q$ and let $1 / p=1 / r+1 / q$. Note that $r-p=r p / q$ thus the inequality $n<q=r p /(r-p)$ implies $n r /(n+r)<p$. Of course, we have $p<r$. Thus the Sobolev embedding theorem says

$$
\begin{gathered}
\left\|\partial_{i} u\right\|_{r, \delta-1+\tau} \leq C\left\|\partial_{i} u\right\|_{r, \delta-1} \leq\left\|\partial_{i} u\right\|_{1, p, \delta-1} \leq C\|u\|_{2, p, \delta} \text {, and } \\
\|u\|_{r, \delta+\tau} \leq C\|u\|_{r, \delta} \leq C\|u\|_{1, p, \delta} \leq C\|u\|_{2, p, \delta}
\end{gathered}
$$

Note that these inequalities are true even for $\tau=0$. We can now complete the proof by seeing that

$$
\begin{aligned}
\left\|b^{i} \partial_{i} u\right\|_{0, p, \delta-2} & \leq\left\|b^{i}\right\|_{q,-1-\tau}\left\|\partial_{i} u\right\|_{r, \delta-1+\tau} \\
& \leq C\|u\|_{2, p, \delta},
\end{aligned}
$$

and $\|c u\|_{0, p, \delta-2} \leq\|c\|_{0, q,-2-\tau}\|u\|_{0, r, \delta+\tau} \leq C\|u\|_{2, p, \delta}$.

Since $P$ is asymptotic to $\Delta$ and is continuous, it is natural to hope and daringly expect that even though we are not on a compact set, some elliptic estimate is anyhow available for $P$.

Proposition E-6 ([]Bar86, Prop 1.6]). Suppose $P$ is asymptotic to $\Delta$, and $1<p \leq q$, and $\delta \in \mathbb{R}$. There is a constant $C=C\left(n, p, q, \delta, C_{1}, \lambda\right)$ such that if $u \in L_{\delta}^{p}$ and $P u \in L_{\delta-2}^{p}$ then $u \in W_{\delta}^{2, p}$ and

$$
\begin{equation*}
\|u\|_{2, p, \delta} \leq C\left(\|P u\|_{0, p, \delta-2}+\|u\|_{0, p, \delta}\right) \tag{E.9}
\end{equation*}
$$

Proof: Let's define

$$
P_{R}=a_{R}^{i j} \partial_{i j}+R b_{R}^{i} \partial_{i}+R^{2} c_{R},
$$

and note that $P_{R} u_{R}=R^{2}(P u)_{R}$.
Here we need to use the usual $L^{p}$ estimates for $A_{1}$. We have, see [GT83] p. 235, Thm 9.11], that for a fattened domain $\tilde{A}_{1}$, we have

$$
\left\|u_{R}\right\|_{2, p ; A_{1}} \leq C\left(\left\|u_{R}\right\|_{p ; \tilde{A}_{1}}+\left\|P_{R} u_{R}\right\|_{p ; \tilde{A}_{1}}\right) .
$$

Going through the proof of [GT83] p. 235, Thm 9.11], we see that the constant in this inequality depends in particular on

$$
\left\|a_{R}^{i j}-\delta_{i j}\right\|_{\infty ; \tilde{A}_{1}} \leq C\left\|a^{i j}-\delta_{i j}\right\|_{\infty, 0} \leq C\left\|a^{i j}-\delta_{i j}\right\|_{C_{-\tau}^{0, \alpha}} \leq C\left(C_{1}\right) .
$$

It turns out that this dependance is independent of $R$.
It also depends on

$$
\left\|R b_{R}^{i}\right\|_{\infty ; \tilde{A}_{1}} \leq C\left\|b^{i}\right\|_{\infty,-1} \leq C\left\|b^{i}\right\|_{0, q,-1-\tau},
$$

which is bounded independently of $R$ as well.
The constant also depends on $c$, but this term is harder to bound. Reading the proof, we reach a point where we want to reduce the number of derivative on $u_{R}$ to use some interpolation estimate. We have

$$
\left\|D^{2} u_{R}\right\|_{p} \leq C\left\|a_{R}^{i j} \partial_{i j}\right\|_{p} \leq C\left(\left\|P_{R} u_{R}\right\|_{p}+\left\|R b_{R}^{i} \partial_{u} u_{R}\right\|_{p}+\left\|R^{2} c_{R} u\right\|_{p}\right),
$$

and the last term must be bounded somehow. We already did the needed work while proving the continuity of $P$ : our proof that multiplication by $c$ is bounded $W^{2, p} \rightarrow L^{p}$ actually works for $W^{1, p} \rightarrow L^{p}$. Thus the constant in the $L^{p}$ estimate depends on

$$
\left\|R^{2} c_{R}\right\|_{q ; \tilde{A}_{1}} \leq C\|c\|_{q,-2} \leq C\|c\|_{q,-2-\tau}
$$

which is bounded independently of $R$.

But then, we have

$$
\begin{aligned}
\|u\|_{2, p, \delta ; A_{R}} & \leq C R^{-\delta}\left\|u_{R}\right\|_{2, p ; \tilde{A}_{1}} \\
& \leq C R^{-\delta}\left(\left\|u_{R}\right\|_{p ; \tilde{A}_{1}}+\left\|P_{R} u_{R}\right\|_{p ; \tilde{A}_{1}}\right) \\
& =C R^{-\delta}\left(\left\|u_{R}\right\|_{p ; \tilde{A}_{1}}+R^{2}\left\|(P u)_{R}\right\|_{p ; \tilde{A}_{1}}\right) \\
& \leq C\left(\|u\|_{p, \delta ; \tilde{A}_{R}}+\|P u\|_{p, \delta-2 ; \tilde{A}_{R}}\right) .
\end{aligned}
$$

The trick we have done so often now with the domains $D_{i}$ completes the proof.
We are interested in the "Fredholmness" of $P$. But the estimate given to us in the previous proposition is not sufficient: we need some compactness of the right-hand-side term.
To understand the Fredholmness of $P$, we first deal with the Laplacian. The orders of growth of harmonic functions in $\mathbb{R}^{n} \backslash B_{1}$ are $\mathbb{Z} \backslash\{-1, \ldots, 3-n\}$ and are called exceptional values.
For nonexceptional weighing parameter $\delta$, we have a very strong Fredholmness result given by the next theorem. Before reaching it, we extract a lemma from a paper of Nirenberg and Walker.

Lemma E-7 ([NW73], lemma 2.1]). Fix $p \in(1, \infty)$, and set $p^{\prime}=p /(p-1)$. Let $a, b \in \mathbb{R}$ be such that $a+b>0$. Set

$$
K^{\prime}(x, y)=|x|^{-a}|x-y|^{a+b-n}|y|^{-b} \text { for } x \neq y
$$

and for $u \in L^{p}$ define

$$
K^{\prime} u(x)=\int K^{\prime}(x, y) u(y) d y
$$

Then there is a constant $c=c(n, p, a, b)$ such that

$$
\left\|K^{\prime} u\right\|_{p} \leq c\|u\|_{p}
$$

if and only if $a<n / p$ and $b<n / p^{\prime}$.
Proof: See [NW73, p. 273] for the proof that the conditions on $a$ and $b$ are necessary.
We can assume that $a$ and $b$ are nonnegative. Indeed, at least one of then is, say $b$. Suppose $a<0$. Then the inequality

$$
\frac{|x|}{|x-y|} \leq 1+\frac{|y|}{|x-y|}
$$

implies that

$$
\begin{aligned}
K^{\prime}(x, y) & \leq|x-y|^{b-n}|y|^{-b}\left(1+\frac{|y|}{|x-y|}\right)^{-a} \\
& \leq C|x-y|^{b-n}|y|^{-b}+C|x-y|^{a+b-n}|y|^{-a-b} .
\end{aligned}
$$

For nonnegative $a$ and $b$ satisfying $a<n / p$ and $b<n / p^{\prime}$, the inequality $|x|^{n} \geq \Pi\left|x_{i}\right|$ yields

$$
K^{\prime}(x, y) \leq \prod_{i=1}^{n}\left|x_{i}\right|^{-a / n}\left|x_{i}-y_{i}\right|^{(a+b) / n-1}\left|y_{i}\right|^{-b / n} .
$$

The problem is thus reduced to one dimensional.
Now for the one-dimensional result, [NW73] cite a lemma whose origin is really unclear.
Suppose that $K(x, y)$ is nonnegative and homogeneous of degree -1 for $x \geq 0$ and $y \geq 0$, and that
the (necessarily identical) quantities

$$
\int_{0}^{\infty} K(x, 1) x^{-1 / p^{\prime}} d x \quad \text { and } \quad \int_{0}^{\infty} K(1, y) x^{-1 / p} d y
$$

are equal to some number $C<\infty$. Then the integral operator

$$
K u(x)=\int_{0}^{\infty} K(x, y) u(y) d y
$$

is bounded on $L_{p}((0, \infty))$ with norm not greater than $C$.
Before proving this result, let's see that $K^{\prime}$ satisfies the hypotheses. It is obviously positive and of the correct homogeneity. Now let $I_{s}(\alpha, \beta):=\int_{0}^{s} r^{\alpha}(1-r)^{\beta} d r$. Set $\alpha=-1 / p^{\prime}-a$ and $\beta=a+b-1$. We then have

$$
\begin{aligned}
\int_{0}^{\infty} K^{\prime}(x, 1) x^{-1 / p^{\prime}} d x & =\int_{0}^{\infty} r^{\alpha}|r-1|^{\beta} d r \\
& =\int_{0}^{1} r^{\alpha}(1-r)^{\beta} d r+\int_{1}^{\infty} r^{\alpha}(r-1)^{\beta} d r \\
& =I_{1}(\alpha, \beta)+I_{1}(\beta, \alpha) \\
& =I_{1 / 2}(\alpha, \beta)+I_{1 / 2}(\beta, \alpha)+I_{1 / 2}(-\alpha-\beta, \beta)+I_{1 / 2}(\beta,-\alpha-\beta)
\end{aligned}
$$

as $\int_{1 / 2}^{1} r^{\alpha}(1-r)^{\beta} d r=\int_{0}^{1 / 2}(1-s)^{\alpha} s^{\beta} d s$.
But on $[0,1 / 2]$, we have $1 / 2<1-r<1$ hence $I_{1 / 2}(\alpha, \beta)$ is comparable to $\int_{0}^{1 / 2} r^{\alpha} d r$ which converges if and only if $\alpha>-1$. Thus $\int_{0}^{\infty} K^{\prime}(x, 1) x^{-1 / p^{\prime}} d x$ converges iff $\alpha, \beta>-1$ and $\alpha+\beta<$ 1. Certainly, $b<1 / p^{\prime}$ and $a<1 / p$ imply all these requirements.

To prove the general result for $K$, we use the homogeneity of $K$ to turn the problem into a convolution problem. Let $\mathcal{M}$ be the multiplicative group $\mathbb{R}_{>0}$. Let $L_{\mathcal{M}}^{P}$ denote the $L^{p}$-space for the Haar measure $d x / x$ on $\mathcal{M}$. Note that multiplication by $x^{1 / p}$ is an isometry $L^{p} \rightarrow L_{\mathcal{M}}^{p}$. Furthermore, since the $L_{\mathcal{M}}^{1}$-norm of $K(x, 1) x^{1 / p}$ is finite, convolution with that function is continuous $L_{\mathcal{M}}^{p} \rightarrow L_{\mathcal{M}}^{p}$. Thus,

$$
\begin{aligned}
\|K u\|_{L^{p}}^{p} & =\int\left(\int K(x, y) u(y) d y\right)^{p} d x \\
& =\int\left(\int K(x / y, 1) u(y) \frac{d y}{y}\right)^{p} d x \\
& =\int\left(x^{-1 / p} \int K(x / y, 1)(x / y)^{1 / p} y^{1 / p} u(y) \frac{d y}{y}\right)^{p} d x \\
& \left.=\| x^{-1 / p}\left(K(\cdot, 1)(\cdot)^{1 / p}\right) *\left(y^{1 / p} u\right)\right) \|_{L^{p}}^{p} \\
& \left.=\| K(\cdot, 1)(\cdot)^{1 / p}\right) *\left(y^{1 / p} u\right) \|_{L_{\mathcal{M}}^{p}}^{p} \\
& \leq C^{p}\left\|y^{1 / p} u\right\|_{L_{\mathcal{M}}^{p}}^{p} \\
& =C^{p}\|u\|_{L^{p}}^{p},
\end{aligned}
$$

and the claim is proved.

We define

$$
k^{-}(\delta)=\max \{k \text { exceptional, } k<\delta\}
$$

and then move on to the theorem.
Theorem E-8 ([|Bar86, Thm 1.7]). Suppose $\delta$ is nonexceptional, $1<p<\infty$ and $s \in \mathbb{N}$. Then

$$
\Delta: W_{\delta}^{\prime s+2, p} \rightarrow W_{\delta-2}^{\prime s, p}
$$

is an isomorphism and there is a constant $C=C(n, p, \delta, s)$ such that

$$
\|u\|_{s+2, p, \delta}^{\prime} \leq C\|\Delta u\|_{s, p, \delta-2}^{\prime} .
$$

Proof: Set $k=k^{-}(\delta)$. Let $\mu=(x \cdot y) /|x||y|$ and $P_{j}^{\lambda}$ denote the ultraspherical function arising in the Taylor expansion of $|x-y|^{2-n}$ with respect to $|y| /|x|$ when $|y|<|x|$ :

$$
\begin{equation*}
|x-y|^{-2 \lambda}=|x|^{-2 \lambda} \sum_{0}^{\infty} P_{j}^{\lambda}(\mu)(|y| /|x|)^{j} . \tag{E.10}
\end{equation*}
$$

Set $\lambda=(n-2) / 2$.
We first show that the inverse of $\Delta: W_{\delta}^{\prime 2, p} \rightarrow W_{\delta-2}^{\prime 0, p}$ has kernel $K(x, y)$ :

$$
c_{n} K(x, y)= \begin{cases}|x-y|^{2-n}, & \text { if } 2-n<\delta<0 ; \\ |x-y|^{2-n}-|y|^{2-n} \sum_{0}^{k} P_{j}^{\lambda}(\mu)(|x| /|y|)^{j}, & \text { if } k \geq 0 ; \\ |x-y|^{2-n}-|x|^{2-n} \sum_{0}^{2-n-k} P_{j}^{\lambda}(\mu)(|y| /|x|)^{j}, & \text { if } k<2-n .\end{cases}
$$

We will refer to these three cases and the three corresponding definitions of $K(x, y)$ as K 1 , K 2 and K3. We now go through a series of step that lead to the proof that $K(x, y)$ defines a bounded operator from $W_{\delta-2}^{\prime 0, p}$ to $W_{\delta}^{\prime 0, p}$.
Note first that in the cases K2 and K3, we have $k<\delta<k+1$, and that in the case K1, we have $k=2-n$.

We have the estimates

$$
|K(x, y)| \leq c(n, k)|x-y|^{2-n} \begin{cases}(|x| /|y|)^{k+1}, & \text { if }|x|<|y| / 2  \tag{E.11}\\ (|x| /|y|)^{n+k-2}, & \text { if }|x| \geq|y| / 2\end{cases}
$$

We need here $n \geq 3$. Then the estimates for K1 are trivial. Indeed, $k+1 \leq 0$ thus $|x|<|y| / 2$ imply $1<2^{3-n}(|x| /|y|)^{k+1}$ and $n+k-2=0$ thus $|x| \geq|y| / 2$ imply $1 \leq 2^{n+k-2}(|x| /|y|)^{n+k-2}$. Let's first prove Estimate (E.11) for K2 in the case $|x| /|y|<1 / 2$. We have in that case that $|x-y| \leq 3|y| / 2$, hence $|x-y|^{2}-n \geq c(n)|y|^{2-n}$. Since we are exactly in the case where the expansion of Equation (E.10) converges (swapping $y$ and $x$ ), we have
as wanted.
We prove the case $|x|>|y| / 2$ for $K 2$ term by term. As $n+k-2 \geq 1$, it must be that
$(|x| /|y|)^{n+k-2} \geq 2^{2-k-n}$ thus $|x-y|^{2-n} \leq c(n, k)|x-y|^{2-n}(|x| /|y|)^{n+k+2}$. For $j \leq k$, we have

$$
\begin{gathered}
\left||y|^{2-n} P_{j}^{\lambda}(\mu)(|x| /|y|)^{j}\right| \leq\left(\max _{i \leq k} \max _{\mu \in B_{1}}\left|P_{i}^{\lambda}(\mu)\right|\right) \frac{|x|^{j}}{|y|^{j+n-2}}, \text { and } \\
\frac{|x|^{j}}{|y|^{j+n-2}} \leq 2 \frac{|x|^{j+1}}{|y|^{j+1+n-2}} \leq 2^{k-j} \frac{|x|^{k}}{|y|^{n+k-2}}
\end{gathered}
$$

These inequalities, along with the fact that $|x-y| \leq 3|x| / 2$ implies

$$
1 \leq c(n)|x-y|^{2-n}|x|^{n-2}
$$

can be used to prove the second estimate for K 2 .
Now let $K_{1}$ be the operator kernel $|x-y|^{2-n}(|x| /|y|)^{\alpha}$. Set $a=\delta+n / p-\alpha$ and $b=-b-n / p+$ $2+\alpha$. Then $a+b=2$ and LemmaE-7 shows that

$$
K_{1}^{\prime}(x, y)=|x|^{-\delta-n / p} K_{1}(x, y)|y|^{\delta-2+n / p}
$$

defines a bounded operator $L^{p} \rightarrow L^{p}$ when $a<n / p$ and $b<n / p^{\prime}$, that is when $\delta<\alpha$ and $\delta>2-n+\alpha$.

Since the composition

$$
L_{\delta-2}^{\prime p} \xrightarrow{|y|^{-\delta+2-n / p}} L^{p} \xrightarrow{K_{1}^{\prime}} L^{p} \xrightarrow{|x|^{\delta+n / p}} L_{\delta}^{\prime p}
$$

is precisely $K_{1}$, then $K_{1}$ is continuous when $\alpha+2-n<\delta<\alpha$.
We use simultaneously $\alpha=k+1$ and $\alpha=n+k-2$ along with Estimate (E.11) to see that $K$ is bounded when $k<\delta<k+1$, which correspond to the cases K2 and K3. The case K1 is dealt with in a slightly different fashion, without the use of Estimate (E.11): we just use $\alpha=0$.
Now that we know that $K$ is bounded $L_{\delta-2}^{\prime p} \rightarrow L_{\delta}^{\prime p}$, we use $K$ to show the surjectivity of $\Delta$. First recall that

$$
\Delta_{x}|x-y|^{2-n}=\Delta_{y}|x-y|^{2-n}=\delta(x-y)
$$

Furthermore, the right-hand-side terms in K2 and K3 are harmonic in $\mathbb{R}^{n} \backslash\{0\}$. Thus

$$
\Delta_{x} K=\Delta_{y} K=\delta(x-y) \text { in } D^{\prime}\left(\mathbb{R}^{n} \backslash\{0\}\right)
$$

Hence $K(\Delta u)=u$ for all $u \in C_{c}^{\infty}\left(\mathbb{R}^{n} \backslash\{0\}\right)$. Since $C_{c}^{\infty}\left(\mathbb{R}^{n} \backslash\{0\}\right)$ is dense in $W_{\delta}^{\prime k, p}$, the continuity of $K$ implies

$$
\|u\|_{p, \delta}^{\prime} \leq C\|\Delta u\|_{p, \delta-2}^{\prime}
$$

Using Estimate E.9, we find

$$
\begin{aligned}
\|u\|_{2, p, \delta}^{\prime} & \leq C\left(\|\Delta u\|_{p, \delta-2}^{\prime}+\|u\|_{p, \delta}^{\prime}\right) \\
& \leq C\|\Delta u\|_{\delta-2}^{\prime}
\end{aligned}
$$

as wanted for the case $s=0$. This $\Delta$ is injective and has closed range. Since it is surjective on $C_{c}^{\infty}\left(\mathbb{R}^{n} \backslash\{0\}\right)$, it must be surjective by density.

With this result in our pocket, we are close to seeing the Fredholmness of $P$. The more powerful
elliptic estimate presented in the next theorem is the tool we need.

Theorem E-9 ([|Bar86] Thm 1.10]). Suppose $P$ is asymptotic to $\Delta$ and $\delta \in \mathbb{R}$ is nonexceptional. For $1<p \leq q$, the map $P: W_{\delta}^{2, p} \rightarrow W_{\delta-2}^{0, p}$ has finite-dimensional kernel and closed range and, for $u \in W_{\delta}^{2, p}$, we have constants $C$ and $R$ depending only on $P, \delta, n, p, q$ such that

$$
\|u\|_{2, p, \delta} \leq C\left(\|P u\|_{0, p, \delta-2}+\|u\|_{L^{p}\left(B_{R}\right)}\right) .
$$

Proof: Let $\left\|\|_{o p}\right.$ denote the operator norm for bounded linear functions $W_{\delta}^{2, p} \rightarrow W_{\delta-2}^{0, p}$ and $\| \|_{o p, R}$ denote the same norm but restricted to functions with support in $E_{R}=\mathbb{R}^{n} \backslash B_{R}$.
Suppose $\operatorname{supp}(u) \subset E_{R}$, then

$$
\|(P-\Delta) u\|_{0, p, \delta-2} \leq\left(\sup _{|x|>R}\left\{\left|a^{i j}(x)-\delta_{i j}\right|\right\}+C\|b\|_{0, q,-1 ; E_{R}}+C\|c\|_{0, q,-2 ; E_{R}}\right)\|u\|_{2, p, \delta} .
$$

Since $P$ is asymptotic to $\Delta$, we thus have

$$
\|P-\Delta\|_{o p, R}=o(1) \text { as } R \rightarrow \infty .
$$

Let $\chi \in C_{c}^{\infty}\left(B_{2}\right)$ be such that $o \leq \chi \leq 1$ with $\chi=1$ in $B_{1}$. Set $\chi_{R}(x)=\chi(x / R)$. Given $u$, write $u_{0}=\chi_{R} u$ and $u_{\infty}=\left(1-\chi_{R}\right) u$. Thus $u=u_{0}+u_{\infty}$.
We have

$$
\begin{aligned}
\left\|u_{\infty}\right\|_{2, p, \delta} & \leq C\left\|\Delta u_{\infty}\right\|_{0, p, \delta-2} \\
& \leq C\left(\left\|P u_{\infty}\right\|_{0, p, \delta-2}+\|P-\Delta\|_{o p, R}\left\|u_{\infty}\right\|_{2, p, \delta}\right)
\end{aligned}
$$

and we estimate

$$
\begin{aligned}
\left\|P u_{\infty}\right\|_{0, p, \delta-2} & \leq\|P u\|_{0, p, \delta-2}+\left\|P u_{0}\right\|_{0, p, \delta-2} \\
& \leq\|P u\|_{0, p, \delta-2}+\left\|\chi_{R} P u\right\|_{0, p, \delta-2}+\left\|\left[P, \chi_{R}\right] u\right\|_{0, p, \delta-2} \\
& \leq C\|P u\|_{0, p, \delta-2}+C\|u\|_{1, p, \delta-1 ; A_{R}} .
\end{aligned}
$$

By throwing in a factor of $R$ in $C$, this last norm can be considered with weight $\delta$. Since $\| P-$ $\Delta \|_{o p, R}=o(1)$, for $R$ sufficiently large we have

$$
\left\|u_{\infty}\right\|_{2, p, \delta} \leq C\left(\|P u\|_{0, p, \delta-2}+\|u\|_{1, p, \delta ; A_{R}}\right) .
$$

We have the exact same estimate for $u_{0}$.
But then

$$
\begin{aligned}
\|u\|_{2, p, \delta} & \leq\left\|u_{0}\right\|_{2, p, \delta}+\left\|u_{\infty}\right\|_{2, p, \delta} \\
& \leq C\left(\|P u\|_{0, p, \delta-2}+\|u\|_{1, p, \delta ; A_{R}}\right) .
\end{aligned}
$$

Using the Interpolation Inequality (E.3), we get the wanted estimate.
Now suppose $\left\{u_{k}\right\} \in \operatorname{ker} P$ satisfy $\left\|u_{k}\right\|_{2, p, \delta}=1$. By Rellich we may assume that $\left\{u_{k}\right\}$ converges
in $L^{p}\left(B_{R}\right)$. Thus

$$
\left\|u_{j}-u_{k}\right\|_{2, p, \delta} \rightarrow 0 \text { as } \min (j, k) \rightarrow \infty
$$

and $\left\{u_{k}\right\}$ is Cauchy hence convergent in $W_{\delta}^{2, p}$. Hence ker $P$ is finite dimensional.
Since $\operatorname{dim}$ ker $P<\infty$, there is a closed subspace $Z$ such that $W_{\delta}^{2, p}=Z+\operatorname{ker} P$ and

$$
\|u\|_{2, p, \delta} \leq C\|P u\|_{0, p, \delta-2} \text { for all } u \in Z
$$

Indeed, should there be no such bound, we could find a sequence $\left\{u_{i}\right\} \in Z$ with $\left\|u_{i}\right\|_{2, p, \delta}=1$ but $P u_{i} \rightarrow 0$. But then using the estimate proved earlier and the Rellich lemma on $B_{R}$, there would be a subsequence of the $u_{i}$ which is Cauchy. By closedness, the $\operatorname{limit} u=\lim u_{i}$ is in $Z$. But then $P u=0$ and $\|u\|_{2, p, \delta}=1$ : contradiction!
The fact that $P$ has closed range follows directly.
We are interested in the dimension of the kernel of $P$.
Theorem E-10. The number $\operatorname{dim} \operatorname{ker}\left(P: W_{\delta}^{2, p} \rightarrow W_{\delta-2}^{0, p}\right)$ is independent of $p$ for $1<p \leq q$.
Proof: We split the range $1<p \leq q$ into three zones:
zone 1: $n<p \leq q$,
zone 2: $n / 2<p \leq n$, and
zone 3: $1<p \leq n / 2$.

Suppose that $P u=0$ and $u \in W_{\delta}^{2, p}$.
Suppose first that $p$ is in zone 1. As $n / p<1$, we have $0<1-n / p \leq 1$. Take any $\alpha$ with $0<\alpha \leq 1-n / p$. Then $u \in W_{\delta}^{2, p}$ implies $\|u\|_{C_{\delta}^{0, \alpha}} \leq C\|u\|_{1, p, \delta}$ hence $u$ is continuous.
Also, since $n-p<0$, we have $|u(x)|=o\left(r^{\delta}\right)$ as $r \rightarrow \infty$. In conjunction with the continuity of $u$, this asymptotic behavior indicates that $u \in L_{\delta}^{s}$ for every $s$. Hence $u \in W_{\delta}^{2, s}$ for every $s$ by Proposition E-6
Suppose now that $p$ is in zone 2 . Then $n-2 p<0$ and $2-n / p \leq 1$. Thus, again,

$$
\|u\|_{C_{\delta}^{0, \alpha}} \leq C\|u\|_{2, p, \delta}
$$

and $u$ is continuous, and

$$
|u(x)|=o\left(r^{\delta}\right) \text { as } r \rightarrow \infty
$$

Again, we have $u \in W_{\delta}^{2, s}$ for every $s$.
Suppose now that $p$ is in zone 3 . Then $n-p \geq n / 2$. Thus $p<2 p \leq n p /(n-p)$ and

$$
\|u\|_{n p /(n-p), \delta} \leq C\|u\|_{1, p, \delta} .
$$

Since $P u=0$, we have by Proposition E-6 that $u \in W_{\delta}^{2, n p /(n-p)}$. Iterating this reasoning a finite number of time, we push $p$ out of zone 3 and once in zone 1 or 2 , we know that $u \in W_{\delta}^{2, s}$ for every $s$.

Note that the role of $q$ is absolutely artificial here. The only reason we need it is to be able to use Proposition E-6.

Because of this last theorem, it is natural to define

$$
N(P, \delta):=\operatorname{dim} \operatorname{ker}\left(P: W_{\delta}^{2, p} \rightarrow W_{\delta-2}^{0, p}\right) .
$$

While there is more in Bartnik's paper that could be done, we end by studying a Theorem quite similar to Theorem 7.2-1 of this thesis.

Theorem E-11 ([Bar86] Thm 1.17]). Suppose $P \sim \Delta$ at rate $\tau>0$. Suppose $\delta$ is nonexceptional and that $u \in W_{\delta}^{2, q}$ satisfies $P u=0$ in $E_{R}$. Then there is an exceptional value $k \leq k^{-}(\delta)$ and $h_{k} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ such that $h_{k}$ is harmonic and homogeneous of degree $k$ in $E_{R}$ and

$$
u-h_{k}=o\left(r^{k-\tau}\right) \text { as } r \rightarrow \infty .
$$

Proof: Set $F:=\Delta u$. Since $P u=0$ in $E_{R}$, we have that $F:=\left(\delta_{i j}-a^{i j}\right) \partial_{i j}^{2} u-b^{i} \partial_{i} u-c u$ in $|x|>R$; thus $F \in W_{\delta-2 \tau}^{0, q}$.
We can take $\epsilon<\tau / 2$ small enough so that $\delta-\tau+\epsilon$ and $\delta-\tau+2 \epsilon$ are nonexceptional. Then $F \in W_{\delta-2-\tau}^{0, q}$ implies that $F \in W_{\delta-2-\tau+\epsilon}^{0, q}$. But $\Delta: W_{\delta-\tau+\epsilon}^{2, q} \rightarrow W_{\delta-2-\tau+\epsilon}^{0, q}$ is Fredholm. So let $\beta_{1}, \ldots, \beta_{n}$ be a basis of $\operatorname{ker}\left(\Delta^{*}\right) \subset\left(W_{\delta-2-\tau+\epsilon}^{0, q}\right)^{*}$. An element $f$ is in $\operatorname{Im}(\Delta)$ if and only if $\beta_{1}(f)=\cdots=\beta_{n}(f)=0$.
Notice that $\left(W_{\delta-2-\tau+\epsilon}^{0, q}\right)^{*}=\left(L_{\delta-2-\tau+\epsilon}^{q}\right)^{*}=L_{-\delta+2+\tau-\epsilon-n}^{q^{\prime}}$ by integration against each other. So the $\beta_{i}$ are functions.
We want to modify $F$ in $B_{R}$ so that it becomes an element of $\operatorname{Im}(\Delta)$. We are thus looking for $f$ with $f=0$ in $E_{R}$ such that

$$
\beta_{i}(f)=\beta_{i}(F), \text { for } i=1, \ldots, n .
$$

Restrict $\beta_{i}$ to $B_{R}$. Since $L_{-\delta+2+\tau-\epsilon-n}^{q^{\prime}}\left(B_{R}\right)=L_{\delta-2-\tau+\epsilon}^{q}\left(B_{R}\right)^{*}$, there are $f_{i} \in L_{\delta-2-\tau+\epsilon}^{q}\left(B_{R}\right)$ with

$$
\beta_{i}\left(f_{j}\right)=\delta_{i j}
$$

Extend $f_{i}$ to $\mathbb{R}^{n}$ by 0 on $E_{R}$. Then $f_{i} \in L_{\delta-2-\tau+\epsilon}^{q}$. The function

$$
F-\beta_{1}(F) f_{1}-\cdots-\beta_{n}(F) f_{n}
$$

is killed by all the $\beta_{i}$ thus it is in the image of $\Delta$.
Thus, there exists a $v$ in $W_{\delta-\tau+\epsilon}^{2, q}$ such that

$$
\Delta(u-v)=0, \text { for }|x|>R .
$$

The classical expansion for harmonic functions now shows that

$$
u-v=h_{k}+O\left(r^{k-1}\right)
$$

for some $k \leq k^{-}(\delta)$ and $h_{k}$ harmonic and of degree $k$ in $E_{R}$. The decay estimate for $v$ is improved by iteration: $u-h_{k} \in W_{\delta-\tau+\epsilon}^{2, q}$ implies $F \in W_{\delta-2-2 \tau+\epsilon}^{0, q}$

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