# Łojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions 

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Paul Feehan dedicates this monograph to the memory of his parents, Martin and Odilla Feehan. Manousos Maridakis dedicates this monograph to his parents, Petros and Pavlina Maridakis.

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## Preface

Our primary goal in this monograph is to prove Eojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions using Sobolev spaces that impose minimal regularity requirements on pairs of connections and sections. Our Lojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions generalize that of the pure Yang-Mills energy function due to the first author [32, Theorems 23.1 and 23.17] for base manifolds of arbitrary dimension and due to Råde [100, Proposition 7.2] for dimensions two and three.

Our Łojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions (for example, Theorems 4 and 6 and the many other examples in Section 1.3) are proved by applying the Lojasiewicz-Simon gradient inequality for an abstract analytic function on a Banach space given by Theorem 1 from our article [38. Now Theorem 1 requires that the Hessian operator for the analytic function be Fredholm with index zero. While the Hessian operator for the coupled Yang-Mills energy function is a linear second-order partial differential operator, it only becomes elliptic when combined with a Coulomb gauge condition [27, 43]. Coupled YangMills energy functions are invariant under the action of gauge transformations (or bundle automorphisms) and so, in principle, one can always find a gauge transformation to produce the required Coulomb gauge condition with the aid of a slice theorem. However, in order to prove the most useful version of the LojasiewiczSimon gradient inequality, one must have a stronger version of the slice theorem for the action of the group of gauge transformations that goes beyond the usual statements found in standard references such as Donaldson and Kronheimer [27] or Freed and Uhlenbeck 43 for connections over four-dimensional manifolds and proved by applying the Implicit Function Theorem. Therefore, a secondary goal of this monograph is to prove a slice theorem using Sobolev norms with borderline (or critical) Sobolev exponents and which is valid in all dimensions.

Since its discovery by Lojasiewicz in the context of analytic functions on Euclidean space 87 and generalization by Simon to a class of analytic functions on certain Hölder spaces [105, Theorem 3], Łojasiewicz-Simon gradient inequalities have played a major role in analyzing questions such as a) global existence, convergence, and analysis of singularities for solutions to nonlinear evolution equations that are realizable as gradient or gradient-like systems for an energy function, b) uniqueness of tangent cones, and c) energy gaps and discreteness of energies.

There are essentially four approaches to establishing a Łojasiewicz-Simon gradient inequality for a particular energy function arising in geometric analysis or mathematical physics: (1) establish the inequality from first principles without relying on the Lojasiewicz inequality for analytic functions on Euclidean space, (2) employ Lyapunov-Schmidt reduction to deduce the gradient inequality for an analytic function on a Banach space from the Lojasiewicz inequality for analytic
functions on Euclidean space, (3) adapt the argument employed by Simon in the proof of his [105. Theorem 3], or (4) apply an abstract version of the ŁojasiewiczSimon gradient inequality for an analytic or Morse-Bott function on a Banach space. Approach (2) is exactly that employed by Simon in $\mathbf{1 0 5}$ and by Råde for the Yang-Mills energy function $\mathbf{1 0 0}$. Occasionally a development from first principles may be necessary, as discussed by Colding and Minicozzi in [25, or may yield the best results, as discussed by the first author in 33. However, in most cases one can derive a Łojasiewicz-Simon gradient inequality for a specific energy function from an abstract version for an analytic or Morse-Bott function on a Banach space. For this strategy to work well, one desires an abstract Łojasiewicz-Simon gradient inequality with the weakest possible hypotheses and proofs of such gradient inequalities were provided by the authors in our article [38, Theorems 1-4].

Appendix $D$ describes the application of Łojasiewicz-Simon gradient inequalities to long-time existence and convergence for solutions to gradient systems for an analytic function while our monograph [32] provides a comprehensive development of these ideas. Given a constant $T>0$, an open subset $\mathscr{U}$ of a Banach space $\mathscr{X}$, a smooth function $\mathscr{E}: \mathscr{U} \rightarrow \mathbb{R}$, and a point $x_{0} \in \mathscr{U}$, a smooth map, $u:[0, T) \rightarrow \mathscr{X}$, is called a gradient flow for $\mathscr{E}$ if it is a solution to the Cauchy problem for the gradient system,

$$
\dot{u}(t)=-\mathscr{E}^{\prime}(u(t)), \quad \text { for all } t \in[0, T), \quad u(0)=x_{0}
$$

as an identity in $\mathscr{X}^{*}$ (the continuous dual space of $\mathscr{X}$ ), where we abbreviate $\dot{u}=$ $d t / d t$ and $\mathscr{E}^{\prime}(x): \mathscr{X} \rightarrow \mathscr{X}^{*}$ is the differential of $\mathscr{E}$ at a point $x \in \mathscr{U}$. The best known examples of gradient flows occurring in geometric analysis include pure and coupled Yang-Mills flows, harmonic map flow, Ricci curvature flow, mean curvature flow, and Yamabe scalar curvature flow.

Simon's approach 105 to studying long-time existence and convergence of a solution to the preceding gradient system relies on his celebrated generalization to infinite dimensions of the Łojasiewicz gradient inequality to a specific class of analytic energy functions on $C^{2, \alpha}$ Hölder spaces of sections of a vector bundle over a closed, smooth Riemannian manifold. Over the intervening years, his ŁojasiewiczSimon gradient inequality has since been generalized by many authors - see our article 38 and references cited therein and our Theorems 11 2 and 3 in this monograph. For example, if $\mathscr{X}$ is continuously embedded in a Hilbert space $\mathscr{H}$ and $\mathscr{E}$ is analytic and $x_{\infty} \in \mathscr{X}$ is a critical point such that the Hessian operator, $\mathscr{E}^{\prime \prime}\left(x_{\infty}\right): \mathscr{X} \rightarrow \mathscr{X}^{*}$, is Fredholm with index zero, then there exist constants $Z \in(0, \infty), \sigma \in(0,1]$, and $\theta \in[1 / 2,1)$ such that gradient map obeys 38, Theorem 1]

$$
\left\|\mathscr{E}^{\prime}(x)\right\|_{\mathscr{X} *} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}, \quad \text { for all } x \in \mathscr{X} \text { such that }\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma .
$$

The preceding inequality is precisely that asserted by Theorem 2 here. We shall apply our more general Theorems 1 and 3 to derive Łojasiewicz-Simon gradient inequalities for all of the energy functions considered in this monograph, namely the a) pure Yang-Mills energy function, b) boson coupled Yang-Mills energy function, c) fermion coupled Yang-Mills energy function, d) Yang-Mills-Higgs energy function, $e$ ) Seiberg-Witten energy function, $f$ ) non-Abelian monopole energy function, and the $g$ ) multiple spinor Seiberg-Witten energy function.

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## CHAPTER 1

## Introduction

Our primary goal in this work is to prove Łojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions. A key feature of our results is that we use systems of Sobolev norms that are as weak as possible. This property is very useful in applications to the analysis of gradient flows, the primary example of an application of Łojasiewicz-Simon gradient inequalities in geometric analysis as illustrated by results of the first author in 32. Our gradient inequalities use $W^{1, p}$ Sobolev norms for coupled Yang-Mills pairs over manifolds of arbitrary dimension $d \geq 2$, including the case $p=d / 2$, where the Sobolev exponent is borderline (or critical) in sense that we explain later in this Introduction.

In the remainder of our Introduction, we outline the history of LojasiewiczSimon gradient inequalities in Section 1.1 and survey their applications in geometric analysis, mathematical physics, and applied mathematics. In Section 1.2, we review our abstract Łojasiewicz-Simon gradient inequality for an analytic function on a Banach space. We state our results on Łojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions in Section 1.3. Unlike the case of the harmonic map energy function considered in 38, one must restrict the Hessian of a coupled Yang-Mills energy function to a suitable slice for the action of the group of gauge transformations in order to obtain an elliptic operator that has the Fredholm property required by our abstract Łojasiewicz-Simon gradient inequality 38, Theorem 2]. In order to obtain the strongest possible version of the resulting Łojasiewicz-Simon gradient inequality for a coupled Yang-Mills energy function, we therefore need to prove existence of a global transformation to Coulomb gauge valid for borderline Sobolev exponents - going beyond standard results described in [27, 43] or previous results due to the first author [34] - and we state the required theorem in Section 1.5 .

### 1.1. A brief history of Łojasiewicz-Simon gradient inequalities and their applications to gradient flows and energy gaps in geometric analysis

Since its discovery by Łojasiewicz in the context of analytic functions on Euclidean space [87, Proposition 1, p. 92] and subsequent generalization by Simon to a class of analytic functions on certain Hölder spaces [105, Theorem 3], the Eojasiewicz-Simon gradient inequality has played a significant role in analyzing questions such as a) global existence, convergence, and analysis of singularities for solutions to nonlinear evolution equations that are realizable as gradient-like systems for an energy function, b) uniqueness of tangent cones, and $c$ ) energy gaps and discreteness of energies. For applications of the Łojasiewicz-Simon gradient inequality to gradient flows arising in geometric analysis, beginning with the harmonic map energy function, we refer to Irwin [70], Kwon 81, Liu and Yang [86],

Simon [106, and Topping [111, 112; for applications to gradient flow for the Chern-Simons function, see Morgan, Mrowka, and Ruberman 92; for applications to gradient flow for the Yamabe function, see Brendle 13, Lemma 6.5 and Equation (100)] and Carlotto, Chodosh, and Rubinstein 16; for applications to Yang-Mills gradient flow, we refer to our monograph [32], Råde [100, and Yang [124; for applications to mean curvature flow, we refer to the survey by Colding and Minicozzi [25; and for applications to Ricci curvature flow, see Ache [3, Haslhofer 58, Haslhofer and Müller [59, and Kröncke [79, 78 .

For applications of the Łojasiewicz-Simon gradient inequality to proofs of global existence, convergence, convergence rate, and stability of nonlinear evolution equations arising in other areas of mathematical physics (including the Cahn-Hilliard, Ginzburg-Landau, Kirchoff-Carrier, porous medium, reaction-diffusion, and semilinear heat and wave equations), we refer to the monograph by Huang 68 for a comprehensive introduction and to the articles by Chill [18, 19], Chill and Fiorenza [20, Chill, Haraux, and Jendoubi 21, Chill and Jendoubi 22, 23, Feireisl and Simondon 40, Feireisl and Takáč 41, Grasselli, Wu, and Zheng [50], Haraux [53], Haraux and Jendoubi [54, 55, 56, Haraux, Jendoubi, and Kavian 57, Huang and Takáč 69, Jendoubi 71, Rybka and Hoffmann 102, 103 , Simon 105, and Takáč [109. For applications to fluid dynamics, see the articles by Feireisl, Laurençot, and Petzeltová [39, Frigeri, Grasselli, and Krejčí 45, Grasselli and Wu 49, and Wu and Xu (122].

For applications of the Łojasiewicz-Simon gradient inequality to proofs of energy gaps and discreteness of energies for Yang-Mills connections, we refer to our article [35. A key feature of our versions of the Łojasiewicz-Simon gradient inequality for the pure Yang-Mills energy function [32, Theorems 23.1 and 23.17] is that they hold for $W^{1,2}$ Sobolev norms for base manifolds of dimensions two, three or four and $W^{2, p}$ Sobolev norms for base manifolds of arbitrary dimension. Those norms are considerably weaker than the $C^{2, \alpha}$ Hölder norms originally employed by Simon in 105, Theorem 3] and this affords considerably greater flexibility in applications. For example, when $(X, g)$ is a closed, four-dimensional, Riemannian manifold, the $W^{1,2}$ Sobolev norm on (bundle-valued) one-forms is (in a suitable sense) quasi-conformally invariant with respect to conformal changes in the Riemannian metric $g$.

## 1.2. Łojasiewicz-Simon gradient inequalities for analytic functions on Banach spaces

There are essentially three approaches to establishing a Łojasiewicz-Simon gradient inequality for a particular energy function arising in geometric analysis or mathematical physics: 1) establish the inequality from first principles, 2) adapt the argument employed by Simon in the proof of his [105, Theorem 3], or 3) apply an abstract version of the Łojasiewicz-Simon gradient inequality for an analytic or Morse-Bott function on a Banach space. Most famously, the first approach is exactly that employed by Simon in 105, although this is also the avenue followed by Kwon [81, Liu and Yang [86] and Topping [111, 112 for the harmonic map energy function and by Råde for the Yang-Mills energy function. Occasionally a development from first principles may be necessary, as discussed by Colding and Minicozzi in 25. However, in almost all of the examples cited in Section 1.1] one can derive a Łojasiewicz-Simon gradient inequality for a specific application from
an abstract version for an analytic or Morse-Bott function on a Banach space. For this strategy to work well, one desires an abstract Łojasiewicz-Simon gradient inequality with the weakest possible hypotheses and a proof of such a gradient inequality (quoted as Theorem here) was the one of the goals of our article [38].

We now recall from 38 a generalization of Simon's infinite-dimensional version [105, Theorem 3] of the Łojasiewicz gradient inequality 87. As we explained in detail in 38, Theorem 1 generalizes Huang's 68, Theorems 2.4.2 (i) and 2.4.5] and other previously published versions of the Łojasiewicz-Simon gradient inequality for analytic functions on Banach spaces.

We begin with the concept of a gradient map [68, Section 2.1B], [7, Section 2.5].

Definition 1.2.1 (Gradient map). (See [68, Definition 2.1.1].) Let $\mathscr{U} \subset \mathscr{X}$ be an open subset of a Banach space, $\mathscr{X}$, and let $\tilde{\mathscr{X}}$ be a Banach space with continuous embedding, $\tilde{\mathscr{X}} \subseteq \mathscr{X}^{*}$. A continuous map, $\mathscr{M}: \mathscr{U} \rightarrow \tilde{\mathscr{X}}$, is called a gradient map if there exists a $C^{1}$ function, $\mathscr{E}: \mathscr{U} \rightarrow \mathbb{R}$, such that

$$
\begin{equation*}
\mathscr{E}^{\prime}(x) v=\langle v, \mathscr{M}(x)\rangle_{\mathscr{X} \times \mathscr{X}^{*}}, \quad \forall x \in \mathscr{U}, \quad v \in \mathscr{X}, \tag{1.2.1}
\end{equation*}
$$

where $\langle\cdot, \cdot\rangle_{\mathscr{X} \times \mathscr{X} *}$ is the canonical bilinear form on $\mathscr{X} \times \mathscr{X}^{*}$. The real-valued function, $\mathscr{E}$, is called a potential for the gradient map, $\mathscr{M}$.

When $\tilde{\mathscr{X}}=\mathscr{X}^{*}$ in Definition 1.2.1, then the differential and gradient maps coincide.

Let $\mathscr{X}$ be a Banach space and let $\mathscr{X}^{*}$ denote its continuous dual space. We call a bilinear form $1, b: \mathscr{X} \times \mathscr{X} \rightarrow \mathbb{R}$, definite if $b(x, x) \neq 0$ for all $x \in \mathscr{X} \backslash\{0\}$. We say that a continuous embedding of a Banach space into its continuous dual space, $\jmath: \mathscr{X} \rightarrow \mathscr{X}^{*}$, is definite if the pullback of the canonical pairing, $\mathscr{X} \times \mathscr{X} \ni$ $(x, y) \mapsto\langle x, \jmath(y)\rangle_{\mathscr{X} \times \mathscr{X}^{*}} \rightarrow \mathbb{R}$, is a definite bilinear form.

Theorem 1 (Łojasiewicz-Simon gradient inequality for analytic functions on Banach spaces). (See [38, Theorem 2].) Let $\mathscr{X}$ and $\tilde{\mathscr{X}}$ be Banach spaces with continuous embeddings, $\mathscr{X} \subset \tilde{\mathscr{X}} \subset \mathscr{X}^{*}$, and such that the embedding, $\mathscr{X} \subset \mathscr{X}^{*}$, is definite. Let $\mathscr{U} \subset \mathscr{X}$ be an open subset, $\mathscr{E}: \mathscr{U} \rightarrow \mathbb{R}$ be a $C^{2}$ function with real analytic gradient map, $\mathscr{M}: \mathscr{U} \rightarrow \tilde{\mathscr{X}}$, and $x_{\infty} \in \mathscr{U}$ be a critical point of $\mathscr{E}$, that is, $\mathscr{M}\left(x_{\infty}\right)=0$. If $\mathscr{M}^{\prime}\left(x_{\infty}\right): \mathscr{X} \rightarrow \tilde{\mathscr{X}}$ is a Fredholm operator with index zero, then there are constants, $Z \in(0, \infty)$, and $\sigma \in(0,1]$, and $\theta \in[1 / 2,1)$, with the following significance. If $x \in \mathscr{U}$ obeys

$$
\begin{equation*}
\left\|x-x_{\infty}\right\| \mathscr{X}<\sigma \tag{1.2.2}
\end{equation*}
$$

then

$$
\begin{equation*}
\|\mathscr{M}(x)\|_{\tilde{X}} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta} . \tag{1.2.3}
\end{equation*}
$$

Remark 1.2.2 (Comments on the embedding hypothesis in Theorem (1). The hypothesis in Theorem 1 on the continuous embedding, $\mathscr{X} \subset \mathscr{X}^{*}$, is easily achieved given a continuous embedding of $\mathscr{X}$ into a Hilbert space $\mathscr{H}$.

REmARK 1.2.3 (On the choice of Banach spaces in applications of Theorem 11). The hypotheses of Theorem 1 are designed to give the most flexibility in applications of a Łojasiewicz-Simon gradient inequality to analytic functions on Banach spaces. An example of a convenient choice of Banach spaces modeled as Sobolev spaces,

[^1]when $\mathscr{M}^{\prime}\left(x_{\infty}\right)$ is realized as an elliptic partial differential operator of order $m$, would be
$$
\mathscr{X}=W^{k, p}(X ; V), \quad \tilde{\mathscr{X}}=W^{k-m, p}(X ; V), \quad \text { and } \quad \mathscr{X}^{*}=W^{-k, p^{\prime}}(X ; V),
$$
where $k \in \mathbb{Z}$ is an integer, $p \in(1, \infty)$ is a constant with dual Hölder exponent $p^{\prime} \in(1, \infty)$ defined by $1 / p+1 / p^{\prime}=1$, while $X$ is a closed Riemannian manifold of dimension $d \geq 2$ and $V$ is a Riemannian vector bundle with a compatible connection, $\nabla: C^{\infty}(X ; V) \rightarrow C^{\infty}\left(X ; T^{*} X \otimes V\right)$, and $W^{k, p}(X ; V)$ denotes a Sobolev space defined in the standard way [6]. When the integer $k$ is chosen large enough, the verification of analyticity of the gradient map, $\mathscr{M}: \mathscr{U} \rightarrow \tilde{\mathscr{X}}$, is straightforward. Normally, that is the case when $k \geq m+1$ and $(k-m) p>d$ or $k-m=d$ and $p=1$, since $W^{k-m, p}(X ; \mathbb{C})$ is then a Banach algebra by [4, Theorem 4.39]. If the Banach spaces are instead modeled as Hölder spaces, as in Simon [105, a convenient choice of Banach spaces would be
$$
\mathscr{X}=C^{k, \alpha}(X ; V) \quad \text { and } \quad \tilde{\mathscr{X}}=C^{k-m, \alpha}(X ; V)
$$
where $\alpha \in(0,1)$ and $k \geq m$, and these Hölder spaces are defined in the standard way [6. Following Remark 1.2.2 the definiteness of the embedding $C^{k, \alpha}(X ; V)=$ $\mathscr{X} \subset \mathscr{X}^{*}$ in this case is the achieved by observing that $C^{k, \alpha}(X ; V) \subset L^{2}(X ; V)$.

We refer the reader to [38, Theorem 4] for a statement and proof of our abstract Łojasiewicz-Simon gradient inequality for Morse-Bott functions on Banach spaces.

Theorem 1 appears to us to be the most widely applicable abstract version of the Łojasiewicz-Simon gradient inequality that we are aware of in the literature. However, for applications where $\mathscr{M}^{\prime}\left(x_{\infty}\right)$ is realized as an elliptic partial differential operator of even order, $m=2 n$, and the nonlinearity of the gradient map is sufficiently mild, it often suffices to choose $\mathscr{X}$ to be the Banach space, $W^{n, 2}(X ; V)$, and choose $\tilde{\mathscr{X}}=\mathscr{X}^{*}$ to be the Banach space, $W^{-n, 2}(X ; V)$. The distinction between the differential, $\mathscr{E}^{\prime}(x) \in \mathscr{X}^{*}$, and the gradient, $\mathscr{M}(x) \in \tilde{\mathscr{X}}$, then disappears. Similarly, the distinction between the Hessian, $\mathscr{E}^{\prime \prime}\left(x_{\infty}\right) \in(\mathscr{X} \times \mathscr{X})^{*}$, and the Hessian operator, $\mathscr{M}^{\prime}\left(x_{\infty}\right) \in \mathscr{L}(\mathscr{X}, \tilde{\mathscr{X}})$, disappears. Finally, if $\mathscr{E}: \mathscr{X} \supset \mathscr{U} \rightarrow \mathbb{R}$ is real analytic, then the simpler Theorem 2 is often adequate for applications.

THEOREM 2 (Łojasiewicz-Simon gradient inequality for analytic functions on Banach spaces). (See 38, Theorem 1].) Let $\mathscr{X} \subset \mathscr{X}^{*}$ be a continuous, definite embedding of a Banach space into its dual space. Let $\mathscr{U} \subset \mathscr{X}$ be an open subset, $\mathscr{E}: \mathscr{U} \rightarrow \mathbb{R}$ be an analytic function, and $x_{\infty} \in \mathscr{U}$ be a critical point of $\mathscr{E}$, that is, $\mathscr{E}^{\prime}\left(x_{\infty}\right)=0$. Assume that $\mathscr{E}^{\prime \prime}\left(x_{\infty}\right): \mathscr{X} \rightarrow \mathscr{X}^{*}$ is a Fredholm operator with index zero. Then there are constants $Z \in(0, \infty)$, and $\sigma \in(0,1]$, and $\theta \in[1 / 2,1)$, with the following significance. If $x \in \mathscr{U}$ obeys

$$
\begin{equation*}
\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma \tag{1.2.4}
\end{equation*}
$$

then

$$
\begin{equation*}
\left\|\mathscr{E}^{\prime}(x)\right\|_{\mathscr{X}^{*}} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta} \tag{1.2.5}
\end{equation*}
$$

While Theorem 1 has important applications to proofs of global existence, convergence, convergence rates, and stability of gradient flows defined by an energy function, $\mathscr{E}: \mathscr{X} \supset \mathscr{U} \rightarrow \mathbb{R}$, with gradient map, $\mathscr{M}: \mathscr{X} \supset \mathscr{U} \rightarrow \mathscr{X}$, (see 32,

Section 2.1] for an introduction and Simon 105 for his pioneering development), the gradient inequality (1.2.3) is most useful when it has the form,

$$
\|\mathscr{M}(x)\|_{\mathscr{H}} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}, \quad \forall x \in \mathscr{U} \text { with }\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma
$$

where $\mathscr{H}$ is a Hilbert space and the Banach space, $\mathscr{X}$, is a dense subspace of $\mathscr{H}$ with continuous embedding, $\mathscr{X} \subset \mathscr{H}$, and so $\mathscr{H}^{*} \subset \mathscr{X}^{*}$ is also a continuous embedding. We refer to Appendix $D$ for applications of this version of the gradient inequality to prove convergence of gradient flows and to Feehan and Maridakis 38, Section 1.2] for further discussion.

As we shall explain further in Section 1.3.2 an $L^{2}$ gradient inequality for a coupled Yang-Mills energy function like (1.3.1) or (1.3.3) does not follow from Theorem 1 when $X$ has dimension $d \geq 5$. However, the desired $L^{2}$ gradient inequalities are implied by the forthcoming Theorem 3.

Theorem 3 (Generalized Lojasiewicz-Simon gradient inequality for analytic functions on Banach spaces). (See [38, Theorem 3].) Let $\mathscr{X}$ and $\tilde{\mathscr{X}}$ be Banach spaces with continuous embeddings, $\mathscr{X} \subset \tilde{\mathscr{X}} \subset \mathscr{X}^{*}$, and such that the embedding, $\mathscr{X} \subset \mathscr{X}^{*}$, is definite. Let $\mathscr{U} \subset \mathscr{X}$ be an open subset, $\mathscr{E}: \mathscr{U} \rightarrow \mathbb{R}$ be an analytic function, and $x_{\infty} \in \mathscr{U}$ be a critical point of $\mathscr{E}$, that is, $\mathscr{E}^{\prime}\left(x_{\infty}\right)=0$. Let

$$
\mathscr{X} \subset \mathscr{G} \subset \tilde{\mathscr{G}} \quad \text { and } \quad \tilde{\mathscr{X}} \subset \tilde{\mathscr{G}} \subset \mathscr{X}^{*}
$$

be continuous embeddings of Banach spaces such that the compositions,

$$
\mathscr{X} \subset \mathscr{G} \subset \tilde{\mathscr{G}} \quad \text { and } \quad \mathscr{X} \subset \tilde{\mathscr{X}} \subset \tilde{\mathscr{G}}
$$

induce the same embedding, $\mathscr{X} \subset \tilde{\mathscr{G}}$. Let $\mathscr{M}: \mathscr{U} \rightarrow \tilde{\mathscr{X}}$ be a gradient map for $\mathscr{E}$ in the sense of Definition 1.2.1. Suppose that for each $x \in \mathscr{U}$, the bounded, linear operator,

$$
\mathscr{M}^{\prime}(x): \mathscr{X} \rightarrow \tilde{\mathscr{X}}
$$

has an extension

$$
\mathscr{M}_{1}(x): \mathscr{G} \rightarrow \tilde{\mathscr{G}}
$$

such that the map

$$
\mathscr{U} \ni x \mapsto \mathscr{M}_{1}(x) \in \mathscr{L}(\mathscr{G}, \tilde{\mathscr{G}}) \quad \text { is continuous. }
$$

If $\mathscr{M}^{\prime}\left(x_{\infty}\right): \mathscr{X} \rightarrow \tilde{\mathscr{X}}$ and $\mathscr{M}_{1}\left(x_{\infty}\right): \mathscr{G} \rightarrow \tilde{\mathscr{G}}$ are Fredholm operators with index zero, then there are constants, $Z \in(0, \infty)$ and $\sigma \in(0,1]$ and $\theta \in[1 / 2,1)$, with the following significance. If $x \in \mathscr{U}$ obeys

$$
\begin{equation*}
\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma \tag{1.2.6}
\end{equation*}
$$

then

$$
\begin{equation*}
\|\mathscr{M}(x)\|_{\mathscr{G}} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta} \tag{1.2.7}
\end{equation*}
$$

Remark 1.2.4 (Generalized Łojasiewicz-Simon gradient inequality for analytic functions on Banach spaces with gradient map valued in a Hilbert space). Suppose now that $\tilde{\mathscr{G}}=\mathscr{H}$, a Hilbert space, so that the embedding $\mathscr{G} \subset \mathscr{H}$ in Theorem 3 factors through $\mathscr{G} \subset \mathscr{H} \simeq \mathscr{H}^{*}$ and therefore

$$
\mathscr{E}^{\prime}(x) v=\langle v, \mathscr{M}(x)\rangle_{\mathscr{X} \times \mathscr{X}^{*}}=(v, \mathscr{M}(x))_{\mathscr{H}}, \quad \forall x \in \mathscr{U} \text { and } v \in \mathscr{X},
$$

using the continuous embeddings, $\tilde{\mathscr{X}} \subset \mathscr{H} \subset \mathscr{X}^{*}$. As we noted in Remark 1.2.2 the hypothesis in Theorem 3 that the embedding, $\mathscr{X} \subset \mathscr{X}^{*}$, is definite is implied
by the assumption that $\mathscr{X} \subset \mathscr{H}$ is a continuous embedding into a Hilbert space. By Theorem 3 if $x \in \mathscr{U}$ obeys

$$
\begin{equation*}
\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma \tag{1.2.8}
\end{equation*}
$$

then

$$
\begin{equation*}
\|\mathscr{M}(x)\|_{\mathscr{H}} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta} \tag{1.2.9}
\end{equation*}
$$

as desired.

## 1.3. Łojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions

In this subsection, we summarize consequences of Theorem 1 for coupled YangMills energy functions.
1.3.1. Łojasiewicz-Simon gradient inequalities for boson and fermion coupled Yang-Mills energy functions. We begin with a definition (due to Parker 97) of two coupled Yang-Mills energy functions.

Definition 1.3.1 (Boson and fermion coupled Yang-Mills energy functions). 97, Section 2] Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $\mathbb{E}$ be a complex finite-dimensional $G$-module equipped with a $G$-invariant Hermitian inner product, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$ be a unitary representation [15, Definitions 2.1.1 and 2.16], and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$, and $m$ and $s$ be smooth real-valued functions on $X$.

We define the boson coupled Yang-Mills energy function by

$$
\begin{equation*}
\mathscr{E}_{g}(A, \Phi):=\frac{1}{2} \int_{X}\left(\left|F_{A}\right|^{2}+\left|\nabla_{A} \Phi\right|^{2}-m|\Phi|^{2}-s|\Phi|^{4}\right) d \operatorname{vol}_{g} \tag{1.3.1}
\end{equation*}
$$

for all smooth connections, $A$ on $P$, and smooth sections, $\Phi$ of $E$, where

$$
\nabla_{A}: C^{\infty}(X ; E) \rightarrow C^{\infty}\left(T^{*} X \otimes E\right)
$$

is the covariant derivative induced on $E$ by the connection $A$ on $P$ and $F_{A} \in$ $\Omega^{2}(X ; \operatorname{ad} P)$ is the curvature of $A$ and $\operatorname{ad} P:=P \times$ ad $\mathfrak{g}$ denotes the real vector bundle associated to $P$ by the adjoint representation of $G$ on its Lie algebra, Ad : $G \ni u \rightarrow \operatorname{Ad}_{u} \in \operatorname{Aut}(\mathfrak{g})$, with fiber metric defined through the Killing form on $\mathfrak{g}$.

Suppose that $X$ admits a $\operatorname{spin}^{c}$ structure comprising a Hermitian vector bundle $W$ over $X$ and a Clifford multiplication map, $c: T^{*} X \rightarrow \operatorname{End}_{\mathbb{C}}(W)$, thus

$$
\begin{equation*}
c(\alpha)^{2}=-g(\alpha, \alpha) \operatorname{id}_{W}, \quad \forall \alpha \in \Omega^{1}(X) \tag{1.3.2}
\end{equation*}
$$

and

$$
D_{A}:=c \circ \nabla_{A}: C^{\infty}(X ; W \otimes E) \rightarrow C^{\infty}(X ; W \otimes E),
$$

is the corresponding Dirac operator [84, Appendix D], 80, Sections 1.1 and 1.2], where $\nabla_{A}$ denotes the covariant derivative induced on $\otimes^{n}\left(T^{*} X\right) \otimes E$ (for $\left.n \geq 0\right)$ and $W \otimes E$ by the connection $A$ on $P$ and Levi-Civita connection for the metric $g$ on $T X$.

We define the fermion coupled Yang-Mills energy function by

$$
\begin{equation*}
\mathscr{F}_{g}(A, \Psi):=\frac{1}{2} \int_{X}\left(\left|F_{A}\right|^{2}+\left\langle\Psi, D_{A} \Psi\right\rangle-m|\Psi|^{2}\right) d \operatorname{vol}_{g} \tag{1.3.3}
\end{equation*}
$$

for all smooth connections, $A$ on $P$, and smooth sections, $\Psi$ of $W \otimes E$.

We recall from [84 Corollary D.4] that a closed orientable smooth manifold $X$ admits a $\operatorname{spin}^{c}$ structure if and only if the second Stiefel-Whitney class $w_{2}(X) \in H^{2}(X ; \mathbb{Z} / 2 \mathbb{Z})$ is the $\bmod 2$ reduction of an integral class. One calls $W$ the fundamental spinor bundle and it carries irreducible representations of $\operatorname{Spin}^{c}(d)$; when $X$ is even-dimensional, there is a splitting $W=W^{+} \oplus W^{-}$and Clifford multiplication restricts to give $\rho: T^{*} X \rightarrow \operatorname{Hom}_{\mathbb{C}}\left(W^{ \pm}, W^{\mp}\right)$ [84, Definition D.9].

Although initially defined for smooth connections and sections, the energy functions $\mathscr{E}_{g}$ and $\mathscr{F}_{g}$ in Definition 1.3.1 extend to the case of Sobolev connections and sections of class $W^{1,2}$.

A short calculation shows that the gradient of the boson coupled Yang-Mills energy function $\mathscr{E}_{g}$ in (1.3.1) with respect to the $L^{2}$ metric on $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$,

$$
\begin{equation*}
\left(\mathscr{M}_{g}(A, \Phi),(a, \phi)\right)_{L^{2}(X, g)}:=\left.\frac{d}{d t} \mathscr{E}_{g}(A+t a, \Phi+t \phi)\right|_{t=0}=\mathscr{E}_{g}^{\prime}(A, \Phi)(a, \phi) \tag{1.3.4}
\end{equation*}
$$

for all $(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, is given by

$$
\begin{align*}
& \left(\mathscr{M}_{g}(A, \Phi),(a, \phi)\right)_{L^{2}(X, g)}  \tag{1.3.5}\\
& =\left(d_{A}^{*} F_{A}, a\right)_{L^{2}(X)}+\operatorname{Re}\left(\nabla_{A}^{*} \nabla_{A} \Phi, \phi\right)_{L^{2}(X)}+\operatorname{Re}\left(\nabla_{A} \Phi, \rho(a) \Phi\right)_{L^{2}(X)} \\
& \quad-\operatorname{Re}(m \Phi, \phi)_{L^{2}(X)}-2 \operatorname{Re} \int_{X} s|\Phi|^{2}\langle\Phi, \phi\rangle d \operatorname{vol}_{g}
\end{align*}
$$

where $d_{A}^{*}=d_{A}^{*, g}: \Omega^{l}(X ; \operatorname{ad} P) \rightarrow \Omega^{l-1}(X ; \operatorname{ad} P)$ is the $L^{2}$ adjoint of the exterior covariant derivative $d_{A}: \Omega^{l}(X ; \operatorname{ad} P) \rightarrow \Omega^{l+1}(X ; \operatorname{ad} P)$, for integers $l \geq 0$. As customary, we let

$$
\Lambda^{l}=\Lambda^{l}\left(T^{*} X\right)
$$

denote the vector bundle over $X$ whose fiber $\Lambda^{l}\left(T_{x}^{*} X\right)$ over each point $x \in X$ is the $l$-th exterior power of the cotangent space, $T_{x}^{*} X$, with $\Lambda^{0}\left(T^{*} X\right):=X \times \mathbb{R}$ and $\Lambda^{1}\left(T^{*} X\right)=T^{*} X$.

We call $(A, \Phi)$ a boson Yang-Mills pair (with respect to the Riemannian metric $g$ on $X)$ if it is a critical point for $\mathscr{E}_{g}$, that is, $\mathscr{M}_{g}(A, \Phi)=0$.

Similarly, one finds that the gradient of the fermion coupled Yang-Mills energy function $\mathscr{F}_{g}$ in (1.3.3) with respect to the $L^{2}$ metric on $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)$,

$$
\begin{align*}
\left(\mathscr{M}_{g}(A, \Psi),(a, \psi)\right)_{L^{2}(X, g)}:=\left.\frac{d}{d t} \mathscr{F}_{g}(A+t a, \Psi+t \psi)\right|_{t=0} &  \tag{1.3.6}\\
& =\mathscr{F}_{g}^{\prime}(A, \Psi)(a, \psi)
\end{align*}
$$

for all $(a, \psi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)$, is given by

$$
\begin{align*}
&\left(\mathscr{M}_{g}(A, \Psi),(a, \psi)\right)_{L^{2}(X, g)}=\left(d_{A}^{*} F_{A}, a\right)_{L^{2}(X)}+\operatorname{Re}\left(D_{A} \Psi-m \Psi, \psi\right)_{L^{2}(X)}  \tag{1.3.7}\\
&+\frac{1}{2}(\Psi, \rho(a) \Psi)_{L^{2}(X)}
\end{align*}
$$

where the action of $a \in \Omega^{1}(X ; \operatorname{ad} P) \equiv C^{\infty}\left(T^{*} X \otimes \operatorname{ad} P\right)$ on $\Psi \in C^{\infty}(X ; W \otimes E)$ is defined by

$$
\begin{aligned}
& \rho(\alpha \otimes \xi)(\phi \otimes \eta):=c(\alpha) \phi \otimes \varrho_{*}(\xi) \eta \\
& \quad \forall \alpha \in \Omega^{1}(X), \quad \xi \in C^{\infty}(X ; \operatorname{ad} P), \quad \phi \in C^{\infty}(X ; W), \quad \eta \in C^{\infty}(X ; E),
\end{aligned}
$$

where $\varrho_{*}: \mathfrak{g} \rightarrow \operatorname{End}_{\mathbb{C}}(\mathbb{E})$ is the representation of the Lie algebra induced by the representation $\varrho: G \rightarrow \operatorname{End}_{\mathbb{C}}(\mathbb{E})$ of the Lie group.

We call $(A, \Psi)$ a fermion Yang-Mills pair (with respect to the Riemannian metric $g$ on $X$ ) if it is a critical point for $\mathscr{F}_{g}$, that is, $\mathscr{M}_{g}(A, \Psi)=0$.

Note that both the boson and fermion coupled Yang-Mills energy functions reduce to the pure Yang-Mills energy function when $\Phi \equiv 0$ or $\Psi \equiv 0$, respectively,

$$
\begin{equation*}
\mathscr{E}_{g}(A):=\frac{1}{2} \int_{X}\left|F_{A}\right|^{2} d \operatorname{vol}_{g} \tag{1.3.8}
\end{equation*}
$$

and $A$ is a Yang-Mills connection (with respect to the Riemannian metric $g$ on $X$ ) if it is a critical point for $\mathscr{E}_{g}$, that is,

$$
\mathscr{M}_{g}(A)=d_{A}^{*, g} F_{A}=0
$$

Given a Hermitian or Riemannian vector bundle, $V$, over $X$ and covariant derivative, $\nabla_{A}$, which is compatible with the fiber metric on $V$, we denote the Banach space of sections of $V$ of Sobolev class $W^{k, p}$, for any $k \in \mathbb{N}$ and $p \in[1, \infty]$, by $W_{A}^{k, p}(X ; V)$, with norm,

$$
\begin{equation*}
\|v\|_{W_{A}^{k, p}(X)}:=\left(\sum_{j=0}^{k} \int_{X}\left|\nabla_{A}^{j} v\right|^{p} d \operatorname{vol}_{g}\right)^{1 / p} \tag{1.3.9}
\end{equation*}
$$

when $1 \leq p<\infty$ and

$$
\begin{equation*}
\|v\|_{W_{A}^{k, \infty}(X)}:=\sum_{j=0}^{k} \underset{X}{\operatorname{ess} \sup }\left|\nabla_{A}^{j} v\right|, \tag{1.3.10}
\end{equation*}
$$

when $p=\infty$, where $v \in W_{A}^{k, p}(X ; V)$. If $k=0$, then we denote $\|v\|_{W^{0, p}(X)}=$ $\|v\|_{L^{p}(X)}$. For $p \in[1, \infty)$ and nonnegative integers $k$, we use 4, Theorem 3.12] (applied to $W_{A}^{k, p}(X ; V)$ and noting that $X$ is a closed manifold) and Banach space duality to define

$$
W_{A}^{-k, p^{\prime}}(X ; V):=\left(W_{A}^{k, p}(X ; V)\right)^{*}
$$

where $p^{\prime} \in(1, \infty]$ is the dual exponent defined by $1 / p+1 / p^{\prime}=1$ and we use the fiber metric on $V$ to replace $V^{*}$ by $V$ on the left-hand side. Elements of the Banach space dual $\left(W_{A}^{k, p}(X ; V)\right)^{*}$ may be characterized via [4, Section 3.10] as distributions in the Schwartz space $\mathscr{D}^{\prime}(X ; V)$ 4, Section 1.57].

As our first application of Theorem 1 we have the following generalization of [32, Theorem 23.17] from the case of the pure Yang-Mills energy function (1.3.8), when $p=2$ and $X$ has dimension $d=2,3$, or 4, and Råde's [100, Proposition 7.2], when $p=2$ and $X$ has dimension $d=2$ or 3 . Because gauge transformations of class $W^{2,2}$ are continuous when $d=2$ or 3 and standard versions of the slice theorem [27, Proposition 2.3.4], 43, Theorem 3.2], 83, Theorem 10.4] for the action of gauge transformations are applicable, the proof of the analogue of Theorem 4 for the pure Yang-Mills energy function due to Råde is simpler for $d=2,3$ and $p=2$.

ThEOREM 4 (Łojasiewicz-Simon gradient inequality for the boson coupled Yang-Mills energy function). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. Let $A_{1}$ be a $C^{\infty}$ reference connection on $P$, and $\left(A_{\infty}, \Phi_{\infty}\right)$ a boson coupled Yang-Mills pair
on $(P, E)$ for $g$ of class $W^{1, q}$, with $q \in[2, \infty)$ obeying $q>d / 2$. If $p \in[2, \infty)$ obeys $d / 2 \leq p \leq q$, then the gradient map,

$$
\mathscr{M}_{g}:\left(A_{1}, 0\right)+W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is real analytic and there are constants $Z \in(0, \infty)$, and $\sigma \in(0,1]$, and $\theta \in[1 / 2,1)$, depending on $A_{1},\left(A_{\infty}, \Phi_{\infty}\right), g, G, p$, and $q$ with the following significance. If $(A, \Phi)$ is a $W^{1, q}$ Sobolev pair on $(P, E)$ obeying the Łojasiewicz-Simon neighborhood condition,

$$
\begin{equation*}
\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A}^{1, p}(X)}<\sigma, \tag{1.3.11}
\end{equation*}
$$

then the boson coupled Yang-Mills energy function (1.3.1) obeys the ŁojasiewiczSimon gradient inequality

$$
\begin{equation*}
\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{W_{A_{1}}^{-1, p}(X)} \geq Z\left|\mathscr{E}_{g}(A, \Phi)-\mathscr{E}_{g}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \tag{1.3.12}
\end{equation*}
$$

The statement of Theorem 4 simplifies with the addition of the rather mild assumption that $A_{1}=A_{\infty}$ and that $\left(A_{\infty}, \Phi_{\infty}\right)$ is $C^{\infty}$ (which can be assumed, modulo a $W^{2, q}$ gauge transformation, provided by the regularity Theorem 2.11.2).

Corollary 5 (Łojasiewicz-Simon gradient inequality for the boson coupled Yang-Mills energy function). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$ bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E}) . \operatorname{Let}\left(A_{\infty}, \Phi_{\infty}\right)$ be a smooth boson coupled Yang-Mills pair for $g$ on $(P, E)$. If $p \in[2, \infty)$ obeys $p \geq d / 2$, then the gradient map,

$$
\mathscr{M}_{g}:\left(A_{\infty}, 0\right)+W_{A_{\infty}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{\infty}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is real analytic and, for $d / 2<q<\infty$ obeying $q \geq p$, there are constants $Z \in(0, \infty)$, and $\sigma \in(0,1]$, and $\theta \in[1 / 2,1)$, depending on $\left(A_{\infty}, \Phi_{\infty}\right), g, G$, $p$, and $q$ with the following significance. If $(A, \Phi)$ is a $W^{1, q}$ Sobolev pair on $(P, E)$ that obeys the Łojasiewicz-Simon neighborhood condition,

$$
\begin{equation*}
\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{\infty}}^{1, p}(X)}<\sigma \tag{1.3.13}
\end{equation*}
$$

then the boson coupled Yang-Mills energy function (1.3.1) obeys the ŁojasiewiczSimon gradient inequality,

$$
\begin{equation*}
\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{W_{A_{\infty}}^{-1, p}(X)} \geq Z\left|\mathscr{E}_{g}(A, \Phi)-\mathscr{E}_{g}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \tag{1.3.14}
\end{equation*}
$$

Similarly, for the fermion coupled Yang-Mills energy function, we have the
Theorem 6 (Łojasiewicz-Simon gradient inequality for the fermion coupled Yang-Mills energy function). Assume the hypotheses of Theorem 4, except that we require that $X$ admit a spinc structure $(\rho, W)$, replace the role of $\mathscr{E}_{g}$ in (1.3.1) by $\mathscr{F}_{g}$ in (1.3.3), and replace the role of the pair $(A, \Phi)$ and critical point $\left(A_{\infty}, \Phi_{\infty}\right)$ of $\mathscr{E}_{g}$ by the pair $(A, \Psi)$ and critical point $\left(A_{\infty}, \Psi_{\infty}\right)$ of $\mathscr{F}_{g}$, where $\Psi$ and $\Psi_{\infty}$ are sections of $W \otimes E$. Then the conclusions of Theorem 4 hold mutatis mutandis for $\mathscr{F}_{g}$.

REMARK 1.3.2 (Łojasiewicz-Simon gradient inequality for coupled Yang-Mills energy functions on quotient spaces). We recall that the space of all smooth connections on $P$ is an affine space, $\mathscr{A}(p)=A_{1}+\Omega^{1}(X ; \operatorname{ad} P)$. While the energy functions
$\mathscr{E}_{g}$ and $\mathscr{F}_{g}$ in Definition 1.3 .1 were initially defined on affine spaces modeled on $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$ or $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)$, the functions are invariant under the action of the group of gauge transformations, $\operatorname{Aut}(P)$, and thus descend to the corresponding quotient spaces. The resulting configuration spaces may be given the structure of smooth Banach manifolds in a standard way [27, Sections 4.2.1], [43, Chapter 3], [51, Section 1] and, with minor modifications of standard proofs, the structure of real analytic Banach manifolds as discussed in Section 1.5.2

Remark 1.3.3 (Lojasiewicz-Simon gradient inequality for the Yang-Mills energy function over a Riemann surface). When $d=2$, it is known in many cases (see [33) that the pure Yang-Mills energy function obeys the Morse-Bott condition in the sense of [38, Definition 1.9] and so by [38, Theorem 4] (our abstract Łojasiewicz-Simon gradient inequality for Morse-Bott functions on Banach spaces), one has the optimal Łojasiewicz-Simon exponent, $\theta=1 / 2$.

We have chosen to derive the Łojasiewicz-Simon gradient inequalities (in Theorems 4 and 6) for two specific coupled Yang-Mills energy functions, motivated by physical considerations, namely the properties of regularity, naturality, and conformal invariance (in dimension four) described by Parker in [97, Section 2].

However, it is clear from the proofs of Theorems 4 and 6 that one can expect the same conclusions for any energy function on pairs of connections and sections with the same nonlinearity structure. Indeed, proofs of such results can be obtained by simple modifications of our proof of the Łojasiewicz-Simon gradient inequality for the boson coupled Yang-Mills energy function, just as we do in this monograph for the case of the fermion coupled Yang-Mills energy function.
1.3.2. Łojasiewicz-Simon $W^{-1,2}$ gradient inequalities for boson and fermion coupled Yang-Mills energy functions. For reasons noted prior to the statement of Theorem 3 it is desirable to replace the Lojasiewicz-Simon gradient inequality (1.3.14) in Theorem 4 with one where the term $\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{W_{A \infty}^{-1, p}(X)}$ is replaced by $\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{L^{2}(X)}$. Such a modification is trivial when we have a continuous embedding, $L^{2}(X) \subset W^{-1, p}(X)$, or equivalently when $W^{1, p^{\prime}}(X) \subset L^{2}(X)$, where $p \in[1, \infty)$ and $p^{\prime}=p /(p-1) \in(1, \infty]$ is the dual Sobolev exponent. By [4, Theorem 4.12], with $d \geq 2$, we have a continuous Sobolev embedding $W^{1, s}(X) \subset L^{s^{*}}(X)$ for $1 \leq s<d$. When $d=2$, 3 , or 4 , then $p=2$ is the smallest value allowed by Theorem 4 and because $p^{\prime}=2$, the continuous Sobolev embedding $W^{1,2}(X) \subset L^{2}(X)$ and, in particular, its dual embedding $L^{2}(X) \subset W^{-1,2}(X)$ guarantee that the desired replacement described above is possible. For $d \geq 5$, the smallest and thus most favorable choice is $p=d / 2$, in which case $p^{\prime}=(d / 2) /((d / 2)-1)=d /(d-2)<d$. We need $\left(p^{\prime}\right)^{*} \geq 2$. Now,

$$
\begin{aligned}
\left(p^{\prime}\right)^{*}=d p^{\prime} /\left(d-p^{\prime}\right)=\left(d^{2} /(d-2)\right) / & (d-d /(d-2)) \\
& =(d /(d-2)) /(1-1 /(d-2))=d /(d-3)
\end{aligned}
$$

and thus $\left(p^{\prime}\right)^{*} \geq 2$ is equivalent to $d \geq 2(d-3)$, that is, $d \leq 6$. Consequently, Theorem 4 only allows replacement of the term $\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{W_{A_{\infty}}^{-1, p}(X)}$ by $\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{L^{2}(X)}$ when $d \leq 6$.

However, the following Łojasiewicz-Simon $W^{-1,2}$ gradient inequalities for boson and fermion coupled Yang-Mills energy functions are corollaries of the proofs of Theorems 4 and 6, respectively, that are valid for any $d \geq 2$. They are proved
in Chapter 4 and $L^{2}$ gradient inequalities follow immediately from these using the continuous Sobolev embedding $L^{2}(X) \subset W^{-1,2}(X)$.

Corollary 7 (Lojasiewicz-Simon $W^{-1,2}$ gradient inequality for the boson coupled Yang-Mills energy function). Assume the hypotheses of Theorem \& If $(A, \Phi)$ is a $W^{1, q}$ Sobolev pair on $(P, E)$ obeying the Łojasiewicz-Simon neighborhood condition (1.3.11), then the boson coupled Yang-Mills energy function (1.3.1) obeys the Lojasiewicz-Simon gradient inequality,

$$
\begin{equation*}
\left\|\mathscr{M}_{g}(A, \Phi)\right\|_{W^{-1,2}(X)} \geq Z\left|\mathscr{E}(A, \Phi)-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \tag{1.3.15}
\end{equation*}
$$

Corollary 8 (Łojasiewicz-Simon $W^{-1,2}$ gradient inequality for the fermion coupled Yang-Mills energy function). Assume the hypotheses of Theorem 6. Then the conclusions of Corollary 7 hold mutatis mutandis for the fermion coupled YangMills energy function $\mathscr{F}_{g}$ in (1.3.3).
1.3.3. Łojasiewicz-Simon gradient inequality for Yang-Mills-Higgs energy functions. A well-known example in complex differential geometry of a coupled Yang-Mills energy function is the Yang-Mills-Higgs function, which we now describe. See Bradlow [10, 11], Bradlow and García-Prada [12], Hitchin [64], Hong [65, Li and Zhang [85, and Simpson 107 for additional details and further references.

We shall follow the description by Bradlow and García-Prada [12, Section 3], but refer the reader to Hong [65, Li and Zhang [85], and the cited references for variants of the Yang-Mills-Higgs function described here. Let $E$ be a complex vector bundle with Hermitian metric $H$ over a compact Kähler manifold $(X, \omega)$. Let $\mathscr{A}_{H}$ denote the affine space of smooth connections on $E$ that are unitary (that is, compatible with the metric $H$ ), and $\Omega^{0}(X ; E)$ denote the vector space of smooth sections of $E$, and $\tau \in \mathbb{R}$.

One defines the Yang-Mills-Higgs energy function on $(A, \Phi) \in \mathscr{A}_{H} \times \Omega^{0}(X ; E)$ by

$$
\begin{equation*}
\mathscr{E}_{H, \omega, \tau}(A, \Phi):=\frac{1}{2} \int_{X}\left(\left|F_{A}\right|^{2}+2\left|\nabla_{A} \Phi\right|^{2}+\left|\Phi \otimes \Phi^{*}-\tau \operatorname{id}_{E}\right|^{2}\right) d \operatorname{vol}_{\omega} \tag{1.3.16}
\end{equation*}
$$

where $\Phi^{*}:=\langle\cdot, \Phi\rangle_{H}$, the dual of $\Phi$ with respect to the metric $H$.
By definition, a Yang-Mills-Higgs pair $(A, \Phi)$ is a critical point of the Yang-Mills-Higgs function $\mathscr{E}_{H, \omega, \tau}$, so $\mathscr{M}_{H, \omega, \tau}(A, \Phi)=0$, or equivalently $(A, \Phi)$ satisfies the second-order Yang-Mills-Higgs equations (the Euler-Lagrange equations defined by the function (1.3.16)). A calculation reveals that a pair is an absolute minimum of $\mathscr{E}_{H, \omega, \tau}$ if and only if it obeys the first-order vortex equations,

$$
\begin{align*}
F_{A}^{0,2} & =0 \\
\bar{\partial} \Phi & =0  \tag{1.3.17}\\
\Lambda F_{A} & =\sqrt{-1}\left(\Phi \otimes \Phi^{*}-\tau \operatorname{id}_{E}\right)
\end{align*}
$$

where $\Lambda F_{A}$ denotes contraction of $F_{A}$ with $\omega$. Let $\mathfrak{u}(E) \subset \operatorname{End}_{\mathbb{C}}(E)$ denote the subbundle of skew-Hermitian endomorphisms of $E$.

The proof of Theorem 4 carries over mutatis mutandis to give
Theorem 9 (Lojasiewicz-Simon gradient inequality for the Yang-Mills-Higgs energy function). Let $X$ be a compact, Kähler manifold of complex dimension $n \geq 1$ and $E$ be a complex vector bundle with Hermitian metric $H$ over $X$. Let $A_{1}$ be a
smooth reference connection on the principal frame bundle for E. Assume that $d=2 n \geq 2$ and $p \in(1, \infty)$ obey one of the conditions in Theorem 4 Then the gradient map,

$$
\mathscr{M}_{H, \omega, \tau}:\left(A_{1}, 0\right)+W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \mathfrak{u}(E) \oplus E\right) \rightarrow W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \mathfrak{u}(E) \oplus E\right)
$$

is real analytic and the remaining conclusions of Theorem 4 and Corollary 7 hold mutatis mutandis for the Yang-Mills-Higgs energy function (1.3.16) near a Yang-Mills-Higgs pair, $\left(A_{\infty}, \Phi_{\infty}\right)$.
1.3.4. Łojasiewicz-Simon gradient inequality for the Seiberg-Witten energy function. For another example of a coupled Yang-Mills energy function whose absolute minima can be readily identified, we consider the Seiberg-Witten equations.

Expositions of the Seiberg-Witten equations are now provided by many authors but, for the sake of consistency, we shall follow our development in 37. Let $(\rho, W)$ denote a $\operatorname{spin}^{c}$ structure on a four-dimensional manifold, $X$, with Riemannian metric, $g$. We recall from [37, Equation (2.55)] that a pair $(B, \Psi)$, comprising a $\operatorname{spin}^{c}$ connection, $B$, on $W=W^{+} \oplus W^{-}$and a section, $\Psi$, of $W^{+}$is a Seiberg-Witten monopole if

$$
\begin{align*}
\operatorname{tr}\left(F_{B}^{+}\right)-\rho^{-1}\left(\Psi \otimes \Psi^{*}\right)_{0} & =0  \tag{1.3.18}\\
D_{B} \Psi & =0
\end{align*}
$$

recalling that $\rho: \Lambda^{+} \cong \mathfrak{s u}\left(W^{+}\right)$is the isomorphism of Riemannian vector bundles induced by Clifford multiplication, $D_{B}: C^{\infty}\left(X ; W^{+}\right) \rightarrow C^{\infty}\left(X ; W^{-}\right)$is the Dirac operator, and $(\cdot)_{0}$ denotes the trace-free part of $\Psi \otimes \Psi^{*} \in \operatorname{End}_{\mathbb{C}}\left(W^{+}\right)$. We have restricted $B$ to $W^{+}$, so $F_{B}^{+} \in C^{\infty}\left(X ; \mathfrak{u}\left(W^{+}\right) \otimes \Lambda^{+}\right)=\Omega^{+}\left(X ; \mathfrak{u}\left(W^{+}\right)\right)$and $\operatorname{tr}\left(F_{B}^{+}\right) \in C^{\infty}\left(X ; i \Lambda^{+}\right)=\Omega^{+}(X ; i \mathbb{R})$, using the fiberwise trace homomorphism, $\operatorname{tr}: \mathfrak{u}\left(W^{+}\right) \rightarrow i \underline{\mathbb{R}}$. The Seiberg-Witten equations (1.3.18) are a system of firstorder partial differential equations in $(B, \Psi)$ and thus cannot be the Euler-Lagrange equations of any action function. However, as we recall from [29, 30, $\mathbf{6 6}, \mathbf{7 2}, \mathbf{9 4}$, Seiberg-Witten monopoles have a variational interpretation by an argument which is the reverse of those provided by Bradlow and García-Prada [12, Section 3] or Hong 65, Section 1] in their derivations of the vortex equations or Li and Zhang 85. Section 1] for the Hermitian-Einstein equations.

Thus, from [66, Equation (1.6)] or [94, Propsition 2.1.4], the Seiberg-Witten energy function is

$$
\begin{align*}
& \mathscr{E}_{g}(B, \Psi)=\int_{X}\left(\left|\nabla_{B} \Psi\right|^{2}+\frac{1}{2}\left|\operatorname{tr}\left(F_{B}\right)\right|^{2}+\frac{R}{4}|\Psi|^{2}+\frac{1}{8}|\Psi|^{2}\right) d \operatorname{vol}_{g}  \tag{1.3.19}\\
& +2 \pi^{2} c_{1}\left(W^{+}\right)^{2},
\end{align*}
$$

where $c_{1}\left(W^{+}\right)^{2}:=\int_{X} c_{1}\left(W^{+}\right)^{2}$. The topological term, $2 \pi^{2} c_{1}\left(W^{+}\right)^{2}$, is independent of the pair $(B, \Psi)$ and does not affect the critical points. In particular,

$$
\mathscr{E}(B, \Psi) \geq 2 \pi^{2} c_{1}\left(W^{+}\right)^{2}
$$

and a pair $(B, \Psi)$ is a Seiberg-Witten monopole if and only if equality is achieved.
Hong and Schabrun derive a version of the Lojasiewicz-Simon gradient inequality [66, Lemma 5.3] based in part on an earlier proof due to Wilkin for the Yang-Mills-Higgs function over a Riemann surface [121, Proposition 3.5]. However, the proof of Theorem 4 carries over mutatis mutandis to give

ThEOREM 10 (Łojasiewicz-Simon gradient inequality for the Seiberg-Witten energy function). Let $(X, g)$ be a closed, four-dimensional, oriented, Riemannian smooth manifold with spin ${ }^{c}$ structure $(\rho, W)$. Let $B_{1}$ be a smooth reference spin $^{c}$ connection on $W$. Assume that $p \in(1, \infty)$ obeys the hypotheses of Theorem 4 with $d=4$. Then the gradient map,

$$
\mathscr{M}_{g}:\left(B_{1}, 0\right)+W_{B_{1}}^{1, p}\left(X ; i \Lambda^{1} \oplus W^{+}\right) \rightarrow W_{B_{1}}^{-1, p}\left(X ; i \Lambda^{1} \oplus W^{+}\right)
$$

is real analytic and the remaining conclusions of Theorem 4 and Corollary 7 hold mutatis mutandis for the Seiberg-Witten energy function (1.3.19) near a SeibergWitten monopole, $\left(B_{\infty}, \Psi_{\infty}\right)$.
1.3.5. Łojasiewicz-Simon gradient inequality for non-Abelian monopole energy functions. For our final example of a coupled Yang-Mills energy function whose absolute minima can be readily identified, we have the non-Abelian monopoles arising in the work of the first author and Leness [37, Okonek and Teleman 95, and Pidstrigatch and Tyurin 98.

Following [37, we consider pairs $(A, \Phi)$ obeying

$$
\begin{align*}
\left(F_{A}^{+}\right)_{0}-\rho^{-1}\left(\Phi \otimes \Phi^{*}\right)_{00} & =0 \\
D_{A} \Phi & =0 \tag{1.3.20}
\end{align*}
$$

where $A$ is a unitary connection on a Hermitian vector bundle, $E$, with curvature $F_{A} \in C^{\infty}\left(X ; \Lambda^{2} \otimes \mathfrak{u}(E)\right)=\Omega^{2}(X ; \mathfrak{u}(E))$ and $\left(F_{A}^{+}\right)_{0} \in C^{\infty}\left(X ; \Lambda^{+} \otimes \mathfrak{s u}(E)\right)=$ $\Omega^{+}(X ; \mathfrak{s u}(E))$, while $\rho: \Lambda^{+} \otimes \mathfrak{s u}(E) \cong \mathfrak{s u}\left(W^{+}\right) \otimes \mathfrak{s u}(E)$ is the isomorphism of Riemannian vector bundles induced by Clifford multiplication, $D_{A}: C^{\infty}\left(X ; W^{+} \otimes\right.$ $E) \rightarrow C^{\infty}\left(X ; W^{-} \otimes E\right)$ is the Dirac operator, and $(\cdot)_{00}$ denotes the trace-free part of $\Phi \otimes \Phi^{*} \in \operatorname{End}_{\mathbb{C}}\left(W^{+} \otimes E\right)$. Let $\mathfrak{s u}(E) \subset \operatorname{End}_{\mathbb{C}}(E)$ denote the subbundle of skew-Hermitian, trace-free endomorphisms of $E$.

By extending the derivations of the Seiberg-Witten energy function in [66] or [94], we find that the non-Abelian monopole energy function is

$$
\begin{align*}
\mathscr{E}_{g}(A, \Phi)= & \int_{X}\left(\left|\nabla_{A} \Phi\right|^{2}+\frac{1}{2}\left|F_{A}\right|^{2}+\frac{R}{4}|\Phi|^{2}+\frac{1}{8}|\Phi|^{2}\right) d \operatorname{vol}_{g}  \tag{1.3.21}\\
& -4 \pi^{2} c^{2}(E)+\frac{1}{2}\left\|F_{A_{w}}^{+}\right\|_{L^{2}(X)}^{2}-2\left\|F_{A_{e}}^{+}\right\|_{L^{2}(X)}^{2}
\end{align*}
$$

The connections $A_{e}$ on $\operatorname{det} E$ and $A_{w}$ on $\operatorname{det} W^{+}$are fixed, with no dynamical role, so the true variables in the $\mathrm{SO}(3)$-monopole equations are the $\mathrm{SO}(3)$ connection $\hat{A}$ induced by $A$ on the bundle $\mathfrak{s u}(E)$ and the (spinor) section $\Phi$ of $W^{+} \otimes E$. The action function, $\mathscr{E}_{g}(A, \Phi)$, again has a universal lower bound and is achieved if and only if $(A, \Phi)$ is a non-Abelian monopole, namely a solution to (1.3.20). Again the proof of Theorem 4 carries over mutatis mutandis to give

Theorem 11 (Łojasiewicz-Simon gradient inequality for the non-Abelian monopole energy function). Let $(X, g)$ be a closed, four-dimensional, oriented, smooth Riemannian manifold with spin ${ }^{c}$ structure $(\rho, W)$. Let $E$ be a Hermitian vector bundle over $X$, and $A_{e}$ be a smooth connection on $\operatorname{det} E$, and $B$ be a smooth spin $^{c}$ connection on $W$, and $A_{1}$ be a smooth reference connection on $E$ inducing $A_{e}$ on $\operatorname{det} E$. Assume that $p \in(1, \infty)$ obeys the hypotheses of Theorem 4 with $d=4$.

Then the gradient map,

$$
\begin{aligned}
\mathscr{M}_{g}:\left(A_{1}, 0\right)+W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \mathfrak{s u}(E) \oplus W^{+}\right. & \otimes E) \\
& \rightarrow W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \mathfrak{s u}(E) \oplus W^{+} \otimes E\right)
\end{aligned}
$$

is real analytic and the remaining conclusions of Theorem 4 and Corollary 7 hold mutatis mutandis for the non-Abelian monopole energy function (1.3.21) near a non-Abelian monopole, $\left(A_{\infty}, \Phi_{\infty}\right)$.
1.3.6. Łojasiewicz-Simon gradient inequality for the multiple spinor Seiberg-Witten energy function. In our next example of a coupled Yang-Mills energy function, the absolute minima can be identified as solutions to the multiple spinor Seiberg-Witten equations arising in the work of Haydys 60, 61, and Haydys and Walpuski 62].

Let $X$ be a closed, oriented Riemannian three-manifold with fixed spin structure $(\rho, W)$, and $L$ be a Hermitian line bundle over $X$, and $E$ be a Hermitian vector bundle of rank $n \geq 1$, trivial determinant $\operatorname{det} E$, and a compatible fixed connection $B$ on $E$. Following [62], we consider pairs $(A, \Psi)$ obeying

$$
\begin{align*}
\rho\left(F_{A}\right)-\mu(\Psi) & =0 \\
D_{A \otimes B} \Psi & =0 \tag{1.3.22}
\end{align*}
$$

where $A$ is a unitary connection on $L$, with curvature $F_{A} \in \Omega^{2}(X ; \mathfrak{u}(L))$, while $\rho: \Lambda^{2} \otimes \mathfrak{u}(L) \cong \mathfrak{s u}(W) \otimes \mathfrak{s u}(L)$ is the isomorphism of Riemannian vector bundles induced by Clifford multiplication,

$$
D_{A \otimes B}: C^{\infty}(X ; \operatorname{Hom}(E ; W \otimes L)) \rightarrow C^{\infty}(X ; \operatorname{Hom}(E ; W \otimes L))
$$

is the Dirac operator, and

$$
\mu: \operatorname{Hom}(E ; W \otimes L) \rightarrow \mathfrak{u}(L) \otimes \mathfrak{s u}(W) \simeq i \mathfrak{s u}(W)
$$

sends $\Psi \in \operatorname{Hom}(E ; W \otimes L)$ to the trace-free part of $\Psi \otimes \Psi^{*} \in \operatorname{End}_{\mathbb{C}}(W \otimes L)$.
By extending the derivations of the Seiberg-Witten energy function in 66 or 94, we find that the Seiberg-Witten multiple spinor energy function is

$$
\begin{align*}
\mathscr{E}_{g}(A, \Psi)=\int_{X}\left(\left|\nabla_{A} \Psi\right|^{2}+\frac{1}{2}\left|F_{A}\right|^{2}+\frac{R}{4}|\Psi|^{2}+\right. & \frac{1}{2}|\mu(\Psi)|^{2}  \tag{1.3.23}\\
& \left.+\left\langle\rho\left(F_{B}\right) \Psi, \Psi\right\rangle\right) d \operatorname{vol}_{g}
\end{align*}
$$

The connection $B$ on $E$ is fixed, with no dynamical role, so the true variables in multiple spinor Seiberg-Witten equations (1.3.22) are the unitary connection $A$ and the (spinor) section $\Psi$ of $\operatorname{End}(E ; W \otimes L) \simeq E^{*} \otimes W \otimes L$. The energy function, $\mathscr{E}_{g}(A, \Psi)$ has universal lower bound zero and that is achieved if and only if $(A, \Psi)$ is solution to (1.3.22). Note that the term $\left\langle\rho\left(F_{B}\right) \Psi, \Psi\right\rangle$ is quadratic in $\Psi$ while $|\mu(\Psi)|^{2}$ is fourth order in $\Psi$. Again the proof of Theorem 4 carries over mutatis mutandis to give

Theorem 12 (Łojasiewicz-Simon gradient inequality for the Seiberg-Witten multiple spinor energy function). Let $X$ be a closed, oriented Riemannian threemanifold with fixed spin structure $(\rho, W)$, and $L$ be a Hermitian line bundle over $X$, and $E$ be a Hermitian vector bundle of rank $n \geq 1$, trivial determinant $\operatorname{det} E$, and a compatible fixed connection $B$ on $E$. Let $A_{1}$ be a fixed $C^{\infty}$ reference connection on
$L$ and endow $W$ with the unique smooth spin connection. Assume that $p \in(1, \infty)$ obeys the hypotheses of Theorem 4 with $d=3$. Then the gradient map,

$$
\begin{aligned}
\mathscr{M}_{g}:\left(A_{1}, 0\right)+W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \mathfrak{u}(L) \oplus\right. & \left.E^{*} \otimes W \otimes L\right) \\
& \rightarrow W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \mathfrak{u}(L) \oplus E^{*} \otimes W \otimes L\right)
\end{aligned}
$$

is real analytic and the remaining conclusions of Theorem 4 and Corollary 7 hold mutatis mutandis for the Seiberg-Witten multiple spinor energy function (1.3.23) near a Seiberg-Witten multi-monopole, $\left(A_{\infty}, \Psi_{\infty}\right)$.
1.3.7. Łojasiewicz-Simon gradient inequality for the pure Yang-Mills energy function near $G_{2}$-instantons. Our final example is concerned with $G_{2^{-}}$ instantons, which are known to be absolute minima of the pure Yang-Mills energy function in dimension seven; for a development of the Łojasiewicz-Simon gradient inequality for the Yang-Mills energy function near anti-self-dual connections over four-dimensional manifolds or flat connections over manifolds of arbitrary dimension, we refer the reader to [32] or [33], respectively. For an introduction to Yang-Mills gauge theory in dimension seven, see Donaldson and Segal [28], Joyce [74] and Kovalev [77]. For recent progress concerning $G_{2}$-instantons, see Walpuski 118 and references therein.

Let $(X, \phi)$ be a closed $G_{2}$-manifold with non-degenerate three-form $\phi \in \Omega^{3}(X)$ (the associative form), $G$ be a compact Lie group, and $P$ be a principal $G$-bundle over $X$. Following [118, a connection $A$ on $P$ is a $G_{2}$-instanton if obeys the anti-self-duality equation,

$$
\begin{equation*}
*\left(F_{A} \wedge \phi\right)=-F_{A} . \tag{1.3.24}
\end{equation*}
$$

It is proven in Walpuski 118 , Proposition 1.97] that when $A$ is a $G_{2}$-instanton, then $A$ is an absolute minimum of the Yang-Mills energy function,

$$
\begin{equation*}
\mathscr{E}_{g}(A, \Psi)=\frac{1}{2} \int_{X}\left|F_{A}\right|^{2} d \operatorname{vol}_{g} \tag{1.3.25}
\end{equation*}
$$

with minimum value $-4 \pi^{2}\left\langle p_{1}(\operatorname{ad} P) \smile[\phi],[X]\right\rangle$. (Note that in our convention, the definition of the Yang-Mills energy function differs from that of $\mathbf{1 1 8}$ by a factor of $1 / 2$.) In this case, a direct application of Theorem 4 and Corollary 7 gives

THEOREM 13 (Łojasiewicz-Simon gradient inequality for $G_{2}$-instantons). Let $(X, \phi)$ be a closed $G_{2}$-manifold, $G$ be a compact Lie group and $P$ be a principal $G$ bundle over $X$ with a fixed $C^{\infty}$ reference connection $A_{1}$. Assume that $p, q \in(1, \infty)$ satisfy $7 / 2 \leq p \leq q$ with $q>7 / 2$ and that $A_{\infty}$ is $W^{1, q}$ connection on $P$ satisfying the $G_{2}$-instanton equation (1.3.24). Then there are constants $Z \in(0, \infty)$, and $\sigma \in(0,1]$, and $\theta \in[1 / 2,1)$, depending on $A_{1}, A_{\infty}, g, G p$, and $q$ with the following significance. For any connection $A$ of class $W^{1, q}$ obeying the Łojasiewicz-Simon neighborhood condition,

$$
\begin{equation*}
\left\|A-A_{\infty}\right\|_{W_{A_{1}}^{1, p}(X)}<\sigma \tag{1.3.26}
\end{equation*}
$$

then the Yang-Mills energy function (1.3.25) obeys the Łojasiewicz-Simon gradient inequality

$$
\begin{equation*}
\left.\left\|d_{A}^{*} F_{A}\right\|_{L^{2}(X)} \geq Z \mid \mathscr{E}_{g}(A)\right)+\left.4 \pi^{2}\left\langle p_{1}(\operatorname{ad} P) \smile[\phi],[X]\right\rangle\right|^{\theta} . \tag{1.3.27}
\end{equation*}
$$

### 1.4. Applications of the Łojasiewicz-Simon gradient inequality for coupled Yang-Mills energy functions

Our interest in Łojasiewicz-Simon gradient inequalities for coupled Yang-Mills and harmonic map energy functions is motivated by the wealth of potential applications. We shall survey some of those applications below.

In [32], we apply the Łojasiewicz-Simon gradient inequality for the pure YangMills energy function [32, Theorems 23.1 and 23.17] to prove global existence, convergence, convergence rate, and stability results for solutions $A(t)$ to the associated gradient flow,

$$
\frac{\partial A}{\partial t}=-\mathscr{M}_{g}(A(t)), \quad A(0)=A_{0}
$$

that is,

$$
\frac{\partial A}{\partial t}=-d_{A(t)}^{*, g} F_{A(t)}, \quad A(0)=A_{0}
$$

Given our Łojasiewicz-Simon gradient inequalities for the boson and fermion coupled Yang-Mills energy functions, Theorems 4 and 6, the main conclusions in 32 for pure Yang-Mills gradient flow should extend easily to the more general case of coupled Yang-Mills gradient flows.

In 35, we applied the Łojasiewicz-Simon gradient inequality to prove an energy gap result for Yang-Mills connections with small $L^{d / 2}$ energy. The proof of that result should extend without difficulty to the case of solutions to the coupled Yang-Mills equations.

### 1.5. Automorphisms and transformation to Coulomb gauge

For some energy functions, the associated Hessian is already an elliptic secondorder partial differential operator on a Sobolev space, but for others the Hessian is only elliptic when combined with a type of Coulomb gauge condition [27, 43] and it is only then that one can apply Theorem 1. For example, in the first category, one has the harmonic map energy and Yamabe functions, while in the second category one has the Yang-Mills and coupled Yang-Mills energy functions.

The Yang-Mills energy function is invariant under the action of gauge transformations (or bundle automorphisms) and so, in principle, one can always find a gauge transformation to produce the required Coulomb gauge condition with the aid of a slice theorem. However, in order to prove the most useful version of the Łojasiewicz-Simon gradient inequality, it is convenient to have a stronger version of the slice theorem for the action of the group of gauge transformations, going beyond the usual statements found in standard references such as Donaldson and Kronheimer [27] or Freed and Uhlenbeck [43] and proved by applying the Implicit Function Theorem. One stronger version of a slice theorem, valid in dimension four, was proved by the first author as [34, Theorem 1.1] but it nevertheless falls short of what we need for our application to the proofs of our Łojasiewicz-Simon gradient inequalities (even when translated to the setting of pairs). Thus, a second purpose of this monograph is to prove a stronger version of [34, Theorem 1.1] for both connections and pairs rather than just connections as in [34, but using standard Sobolev norms with borderline Sobolev exponents rather than the critical-exponent norms employed in 34 and valid in all dimensions.
1.5.1. Transformation to Coulomb gauge. We first state the desired result for connections and then its analogue for pairs.

Theorem 14 (Existence of $W^{2, q}$ Coulomb gauge transformations for $W^{1, q}$ connections that are $W^{1, \frac{d}{2}}$ close to a reference connection). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$. If $A_{1}$ is a $C^{\infty}$ connection on $P$, and $A_{0}$ is a Sobolev connection on $P$ of class $W^{1, q}$ with $d / 2<q<\infty$, and $p \in(1, \infty)$ obeys $d / 2 \leq p \leq q$, then there exists a constant $\zeta=\zeta\left(A_{0}, A_{1}, g, G, p, q\right) \in(0,1]$ with the following significance. If $A$ is a $W^{1, q}$ connection on $P$ that obeys

$$
\begin{equation*}
\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta \tag{1.5.1}
\end{equation*}
$$

then there exists a gauge transformation $u \in \operatorname{Aut}(P)$ of class $W^{2, q}$ such that

$$
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0,
$$

and

$$
\left\|u(A)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

where $N=N\left(A_{0}, A_{1}, g, G, p, q\right) \in[1, \infty)$ is the constant in the forthcoming Proposition 2.5.1.

For a description of the action of the group of gauge transformations in Theorem 14 and the definition of the Coulomb gauge condition for connections, we refer the reader to Section 2.6. and for an explanation of the remainder of the notation in Theorem 14 we refer the reader to Section 1.3.1

The essential point in Theorem 14 is that the result holds for the critical exponent, $p=d / 2$ with $d \geq 3$, when the Sobolev space $W^{2, p}(X)$ fails to embed in $C(X)$ (see [4, Theorem 4.12]) and a proof of Theorem 14 by the Implicit Function Theorem in the case $p>d / 2$ fails when $p=d / 2$. In this situation, a $W^{2, \frac{d}{2}}$ gauge transformation $u$ of $P$ is not continuous, the set $\mathrm{Aut}^{2, \frac{d}{2}}(P)$ of $W^{2, \frac{d}{2}}$ gauge transformations is not a manifold, and Aut ${ }^{2, \frac{d}{2}}(P)$ cannot act smoothly on the affine space $\mathscr{A}^{1, \frac{d}{2}}(P)$ of $W^{1, \frac{d}{2}}$ connections on $P$. When $d=4$ and $p \geq 2$, this phenomenon is discussed by Freed and Uhlenbeck in [43, Appendix A].

While the Quantitative Implicit Function Theorem F. 1 could be applied to give a shorter proof of Theorem [14 when $p>d / 2$, it is unclear whether this method could yield the borderline case $p=d / 2$. See Remark 1.5.1 for an explanation. However, to illustrate use of the Quantitative Implicit Function Theorem in an application that is of interest in its own right, we shall apply Theorem F. 1 to prove

ThEOREM 15 (Existence of $W^{2, q}$ Coulomb gauge transformations for $W^{1, q}$ connections that are $L^{r}$ close to a reference connection). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$. If $A_{1}$ is a $C^{\infty}$ connection on $P$, and $A_{0}$ is a Sobolev connection on $P$ of class $W^{1, q}$ with $d / 2<q<\infty$, and $r$ is a constant that obeys $d<r \leq q^{*}=d q /(d-q)$ if $q<d$ or $d<r<\infty$ if $q \geq d$, then there are constants $\delta=\delta\left(A_{0}, A_{1}, g, G, r\right) \in(0,1]$ and $C=C\left(A_{0}, A_{1}, g, G, r\right) \in(0, \infty)$ with the following significance. If $A$ is a $W^{1, q}$ connection on $P$ that obeys

$$
\begin{equation*}
\left\|A-A_{0}\right\|_{L^{r}(X)}<\delta \tag{1.5.2}
\end{equation*}
$$

then there exists a gauge transformation $u \in \operatorname{Aut}(P)$ of class $W^{2, q}$ such that

$$
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0,
$$

and

$$
\left\|u(A)-A_{0}\right\|_{L^{r}(X)}<C\left\|A-A_{0}\right\|_{L^{r}(X)}
$$

Remark 1.5.1 (Borderline Sobolev exponents). The constant $C$ in Theorem 15 and the $L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ norm depend continuously on $r \geq d$ (see 4, Theorem 2.17] or [47, Problem 7.1]) and so (as in the proof of [114, Corollary 2.2]), we may take the limit as $r \downarrow d$ to give

$$
\left\|u(A)-A_{0}\right\|_{L^{d}(X)} \leq C\left\|A-A_{0}\right\|_{L^{d}(X)}
$$

However, it is unclear whether one can also relax the condition (1.5.2) by allowing $r=d$ because our direct proof in Section 2.7 of Theorem 15 via the Quantitative Implicit Function Theorem F. 1 precludes that choice. It is also unclear that one could achieve such a result by applying the Method of Continuity (by analogy with our proof of Theorem (14) since there is no obvious analogue of Proposition 2.5.1 The analogous comments apply to the forthcoming Theorem 17.

The proof of Theorem 14 adapts mutatis mutandis to establish the following refinement of [36, Proposition 2.8] and [97, Theorem 4.1].

Theorem 16 (Existence of $W^{2, q}$ Coulomb gauge transformations for $W^{1, q}$ pairs that are $W^{1, \frac{d}{2}}$ close to a reference pair). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $A_{1}$ is a $C^{\infty}$ connection on $P$, and $\left(A_{0}, \Phi_{0}\right)$ is a Sobolev pair on $(P, E)$ of class $W^{1, q}$ with $d / 2<q<\infty$, and $p \in(1, \infty)$ obeys $d / 2 \leq p \leq q$, then there exists a constant $\zeta=\zeta\left(A_{1}, A_{0}, \Phi_{0}, g, G, p, q\right) \in(0,1]$ with the following significance. If $(A, \Phi)$ is a $W^{1, q}$ pair on $(P, E)$ that obeys

$$
\begin{equation*}
\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta \tag{1.5.3}
\end{equation*}
$$

then there exists a gauge transformation $u \in \operatorname{Aut}(P)$ of class $W^{2, q}$ such that

$$
d_{A_{0}, \Phi_{0}}^{*}\left(u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0,
$$

and

$$
\left\|u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)}
$$

where $N=N\left(A_{1}, A_{0}, \Phi_{0}, g, G, p, q\right) \in[1, \infty)$ is the constant in the forthcoming Proposition 2.9.5.

For a description of the action of the group of gauge transformations in Theorem 16 and the definition of the Coulomb gauge condition for pairs, we refer the reader to Section 2.9. The proof of Theorem 15 adapts mutatis mutandis to establish the following refinement of Theorem 16,

Theorem 17 (Existence of $W^{2, q}$ Coulomb gauge transformations for $W^{1, q}$ pairs that are $L^{r}$ close to a reference pair). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $A_{1}$ is
a $C^{\infty}$ connection on $P$, and $\left(A_{0}, \Phi_{0}\right)$ is a Sobolev pair on $(P, E)$ of class $W^{1, q}$ with $d / 2<q<\infty$, and $r$ is a constant that obeys $d<r \leq q^{*}=d q /(d-q)$ if $q<d$ or $d<r<\infty$ if $q \geq d$, then there are constants $\delta=\delta\left(A_{0}, \Phi_{0}, A_{1}, g, G, r\right) \in(0,1]$ and $C=C\left(A_{0}, \Phi_{0}, A_{1}, g, G, r\right) \in(0, \infty)$ with the following significance. If $(A, \Phi)$ is a $W^{1, q}$ pair on $(P, E)$ that obeys

$$
\begin{equation*}
\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{L^{r}(X)}<\delta \tag{1.5.4}
\end{equation*}
$$

then there exists a gauge transformation $u \in \operatorname{Aut}(P)$ of class $W^{2, q}$ such that

$$
d_{A_{0}, \Phi_{0}}^{*}\left(u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0,
$$

and

$$
\left\|u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{L^{r}(X)}<C\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{L^{r}(X)} .
$$

1.5.2. Real analytic Banach manifold structures on quotient spaces. In order to establish the analyticity of the pure or coupled Yang-Mills energy functions on affine spaces of $W^{1, q}$ connections $\mathscr{A}^{1, q}(P)$ or pairs $\mathscr{P}^{1, q}(P, E)$, respectively, it is not necessary to know that their quotient spaces with respect to the action of the group $\mathrm{Aut}^{2, q}(P)$ of gauge transformations are analytic Banach manifolds. Nevertheless, because this readily follows from the proofs of Theorem 14 and Theorem 16. respectively, we include the relevant statements here for the case of connections, noting that the analogous statements for pairs are similar.

Theorem 14 provides the essential ingredient one needs to show not only that the quotient space $\mathscr{B}(P):=\mathscr{A}^{1, q}(P) / \operatorname{Aut}^{2, q}(P)$ is a $C^{\infty}$ but also a real analytic Banach manifold away from orbits $[A]=\left\{u(A): u \in \operatorname{Aut}^{2, q}(P)\right\}$ corresponding to $W^{1, q}$ connections $A$ on $P$ whose stabilizers (or isotropy groups), $\operatorname{Stab}(A):=\{u \in$ $\left.\operatorname{Aut}^{2, q}(P): u(A)=A\right\}$, are non-minimal, that is, contain the $\operatorname{Center}(G)$ as a proper subgroup. To show that

$$
\mathscr{B}^{*}(P)=\left\{A \in \mathscr{A}^{1, q}(P): \operatorname{Stab}(A)=\operatorname{Center}(G)\right\} / \operatorname{Aut}^{2, q}(P)
$$

is a $C^{\infty}$ Banach manifold, one only needs the 'easy case' of Theorem 14 where $p=q$, as the condition $q>d / 2$ ensures that the proofs using $H^{k+1}(X)$ Sobolev spaces (with $d=4$ and $k \geq 2$ ) due to Donaldson and Kronheimer [27, Sections 4.2.1 and 4.2.2] or Freed and Uhlenbeck [43, pp. 48-51] apply mutatis mutandis. We have the following analogue of [27, Proposition 4.2.9], 43, Corollary, p. 50], for real analytic Banach manifolds and $X$ of dimension $d \geq 2$ rather than $C^{\infty}$ Hilbert manifolds and $X$ of dimension four.

Corollary 18 (Real analytic Banach manifold structure on the quotient space of $W^{1, q}$ connections). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$, and $q$ obey $d / 2<q<\infty$. If $A_{1}$ is a $C^{\infty}$ reference connection on $P$ and $[A] \in \mathscr{B}(P)$, then there is a constant $\varepsilon=\varepsilon\left(A_{1},[A], g, G, q\right) \in(0,1]$ with the following significance. If

$$
\mathbf{B}_{A}(\varepsilon):=\left\{a \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right): d_{A}^{*} a=0 \text { and }\|a\|_{W_{A_{1}}^{1, q}(X)}<\varepsilon\right\}
$$

then the map,

$$
\pi_{A}: \mathbf{B}_{A}(\varepsilon) / \operatorname{Stab}(A) \ni[a] \mapsto[A+a] \in \mathscr{B}(P)
$$

is a homeomorphism onto an open neighborhood of $[A] \in \mathscr{B}(P)$. For $a \in \mathbf{B}_{A}(\varepsilon)$, the stabilizer of $a$ in $\operatorname{Stab}(A)$ is naturally isomorphic to that of $\pi_{A}(a)$ in $\operatorname{Aut}^{2, q}(P)$.

In particular, the inverse coordinate charts, $\pi_{A}$, determine real analytic transition functions for $\mathscr{B}^{*}(P)$, giving it the structure of a real analytic Banach manifold, and each map $\pi_{A}$ is a real analytic diffeomorphism from the open subset of points $[a] \in \mathbf{B}_{A}(\varepsilon) / \operatorname{Stab}(A)$ where $\pi_{A}(a)$ has stabilizer isomorphic to Center $(G)$.

As in [110, p. 328], one may consider the quotient space of framed connections modulo gauge transformations, $\mathscr{B}^{\prime}(P):=\left(\mathscr{A}^{1, q}(P) \times\left. P\right|_{x_{0}}\right) /$ Aut $^{2, q}(P)$, for some fixed base point $x_{0} \in X$, and now the obvious analogue of Theorem 14 shows that $\mathscr{B}^{\prime}(P)$ is a real analytic Banach manifold.

Corollary 18 may be easily extended to the setting of pairs by applying Theorem 16 in place of Theorem 14. We leave such extensions to the reader, but refer to 97 , Theorem 4.2] and [36, Proposition 2.8] for statements and proofs of $C^{\infty}$ Banach manifold structures for quotient spaces of pairs.

### 1.6. Outline of the monograph

To apply Theorem to pure or coupled Yang-Mills energy functions and obtain the best possible results in those applications, one requires the global Coulomb gauge constructions provided by Theorems 14 or [16, and those results are proved in Chapter 2, In Chapter 3, we derive Łojasiewicz-Simon gradient inequalities for the coupled Yang-Mills energy functions, proving Theorems $4,6,19,10,11$, and Corollary 5. In Chapter 4 we extend the methods of Chapter 3 to derive the stronger Łojasiewicz-Simon $L^{2}$ gradient inequalities for boson and fermion coupled Yang-Mills energy functions, proving Corollaries 7 and 8 . Appendix A establishes the Fredholm properties and computes the index of an elliptic partial differential operator with smooth coefficients acting on Sobolev spaces, Appendix B discusses the equivalence of Sobolev norms defined by Sobolev and smooth connections, and Appendix Cestablishes the Fredholm properties and computes the index of a Hodge Laplacian with smooth and Sobolev coefficients. Appendix D contains a summary of key convergence results for gradient flows under the validity of a ŁojasiewiczSimon gradient inequality. Appendix Econtains a review of Huang's 68, Theorem 2.4 .2 (i)] for the Łojasiewcz-Simon gradient inequality for analytic functions on Banach spaces. Because precise statements in useful generality are difficult to find in the literature, in Appendix F we provide statements and proofs of quantitative implicit and inverse function theorems for maps of Banach spaces.

### 1.7. Notation and conventions

For the notation of function spaces, we follow Adams and Fournier [4], and for functional analysis, Brezis [14] and Rudin 101]. We let $\mathbb{N}:=\{0,1,2,3, \ldots\}$ denote the set of non-negative integers. We use $C=C(*, \ldots, *)$ to denote a constant which depends at most on the quantities appearing on the parentheses. In a given context, a constant denoted by $C$ may have different values depending on the same set of arguments and may increase from one inequality to the next. If $\mathscr{X}, \mathscr{Y}$ is a pair of Banach spaces, then $\mathscr{L}(\mathscr{X}, \mathscr{Y})$ denotes the Banach space of all continuous linear operators from $\mathscr{X}$ to $\mathscr{Y}$. We denote the continuous dual space of $\mathscr{X}$ by $\mathscr{X}^{*}=\mathscr{L}(\mathscr{X}, \mathbb{R})$. We write $\alpha(x)=\langle x, \alpha\rangle_{\mathscr{X}} \times \mathscr{X} *$ for the canonical pairing between $\mathscr{X}$ and its dual space, where $x \in \mathscr{X}$ and $\alpha \in \mathscr{X}^{*}$. If $T \in \mathscr{L}(\mathscr{X}, \mathscr{Y})$, then its adjoint is denoted by $T^{*} \in \mathscr{L}\left(\mathscr{Y}^{*}, \mathscr{X}^{*}\right)$, where $\left(T^{*} \beta\right)(x):=\beta(T x)$ for all $x \in \mathscr{X}$ and $\beta \in \mathscr{Y}^{*}$. If $x_{0} \in \mathscr{X}$ and $r \in(0, \infty)$, we let $B_{r}\left(x_{0}\right)$ denote the open ball of radius $r$ and, if $x_{0}$ is the origin, denote the ball by $B_{r}$.

## CHAPTER 2

## Existence of Coulomb gauge transformations for connections and pairs

In Section 2.6, we prove our refinement, Theorem 14 of the standard construction of a $W^{2, q}$ Coulomb gauge transformation $u$, with $q>d / 2$, for a $W^{1, q}$ connection $A$ on a principal $G$-bundle $P$ over a closed Riemannian smooth manifold of dimension $d \geq 2$. We extend this result in Section 2.9 to the action of gauge transformations on affine spaces of $W^{1, q}$ pairs, obtaining our refinement, Theorem 16 of the standard constructions of Coulomb gauge transformations in that context due to Parker 97 ] and the first author and Leness [36. Finally, in Section 2.11] we extend known regularity results for solutions to the Yang-Mills equations in dimensions greater than or equal to two and coupled Yang-Mills equations in dimension four to the case of solutions to the coupled Yang-Mills equations in dimensions greater than or equal to two.

### 2.1. Action of Sobolev gauge transformations on Sobolev connections

Suppose that $P$ is a smooth principal $G$-bundle over a smooth manifold $X=$ $P / G$ of dimension $d$, where $P \times G \rightarrow P$ is a right action of $G$ on $P$. For $q>d / 2$, let $\mathrm{Aut}^{2, q}(P)$ denote the Banach Lie group of Sobolev $W^{2, q}$ automorphisms (or gauge transformations) of $P$ 27, Section 2.3.1], 43, Appendix A and p. 32 and pp. 45-51], [44, Section 3.1.2]. We recall that there is a smooth left action,

$$
\operatorname{Aut}^{2, q}(P) \times P \rightarrow P,
$$

which commutes with the right action of $G$ on $P$. This induces a smooth right (affine) action on the affine space $\mathscr{A}^{1, q}(P)$ of Sobolev $W^{1, q}$ connections on $P$,

$$
\begin{equation*}
\mathscr{A}^{1, q}(P) \times \operatorname{Aut}^{2, q}(P) \ni(A, u) \rightarrow u(A) \in \mathscr{A}^{1, q}(P) \tag{2.1.1}
\end{equation*}
$$

defined by pull-back,

$$
u(A):=u^{*} A, \quad \forall u \in \operatorname{Aut}^{2, q}(P) \text { and } A \in \mathscr{A}^{1, q}(P)
$$

We recall from Definition 1.3.1 that we (implicitly) assume that $G$ is a matrix Lie group, equipped with a unitary representation, $\varrho: G \rightarrow$ Aut $_{\mathbb{C}}(\mathbb{E})$. The constraint $q>d / 2$ is required to ensure that $W^{2, q}(X) \subset C(X)$ by the Sobolev Embedding [4, Theorem 4.12] and thus $u \in \operatorname{Aut}^{2, q}(P)$ is a continuous gauge transformation of $P$ and that $W^{2, q}(X)$ is a Banach algebra by [4, Theorem 4.39].

Given a $W^{1, q}$ connection $A_{0}$ on $P$ and a $C^{\infty}$ reference connection $A_{1}$ on $P$, the standard construction of a slice for the action of $\mathrm{Aut}^{2, q}(P)$ on $\mathscr{A}^{1, q}(P)$ provides constants $\varepsilon=\varepsilon\left(A_{0}, A_{1}, g, P\right) \in(0,1]$ and $C=C\left(A_{0}, A_{1}, g, P\right) \in[1, \infty)$ such that if $A$ is close to $A_{0}$ in the sense that,

$$
\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}<\varepsilon
$$

then there exists $u \in \operatorname{Aut}^{2, q}(P)$ such that $u(A)$ is in Coulomb gauge relative to $A_{0}$, that is,

$$
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0
$$

and $u(A)$ is close to $A_{0}$,

$$
\left\|u(A)-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}<C \varepsilon
$$

For example, see [27, Proposition 2.3.4], 43, Theorem 3.2], 83, Theorem 10.4] or [34, Theorem 1.1] for statements of the Slice Theorem and their proofs using the Implicit Function Theorem for smooth maps of Banach spaces.

Our Theorem 14 relaxes the condition that $A$ be $W_{A_{1}}^{1, q}$ close to $A_{0}$ for $q>d / 2$ to $W_{A_{1}}^{1, p}$ close for $p$ obeying $d / 2 \leq p \leq q$ when $d \geq 3$ and provides $W^{1, p}$ bounds for $u(A)-A_{0}$ in terms of $A-A_{0}$. This is significant since Aut ${ }^{2, \frac{d}{2}}(P)$ is not a smooth manifold and the action (2.1.1) cannot be smooth when $q=d / 2$, so the Implicit Function Theorem does not apply.

### 2.2. A priori estimates for Laplace operators with Sobolev coefficients and existence and uniqueness of strong solutions

Before proceeding to the proof of Theorem 14 we begin with some preparatory lemmata and remarks that have some interest in their own right. Standard theory for existence and uniqueness of strong solutions to (scalar) second-order elliptic partial differential equations, such as 47, Chapter 9], requires that the secondorder coefficients be continuous and the lower-order coefficients be bounded. Here, we observe that one can relax those requirements on the lower-order coefficients and accommodate the setting we employ in this monograph.

For a smooth connection $A$ on $P$ and integers $l \geq 0$, we let

$$
\begin{equation*}
\Delta_{A}=d_{A}^{*} d_{A}+d_{A} d_{A}^{*} \quad \text { on } \Omega^{l}(X ; \operatorname{ad} P) \tag{2.2.1}
\end{equation*}
$$

denote the Hodge Laplace operator. Our proof of Theorem 14 will require a priori $L^{p}$ estimates, existence and uniqueness results, Fredholm properties, and Hodge decompositions involving the Hodge Laplacian (2.2.1) when $A$ is a $W^{1, q}$ Sobolev connection. When $A$ is a $C^{\infty}$ connection and we restrict our attention to $p=2$, those properties are immediate consequences of more general results (for example, see Gilkey [48]) for elliptic operators on sections of vector bundles over closed manifolds.

Proposition 2.2.1 (A priori $L^{p}$ estimate for a Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$, and $l \geq 0$ be an integer. If $A$ is a $W^{1, q}$ connection on $P$ with $q>d / 2$, and $A_{1}$ is a $C^{\infty}$ connection on $P$, and $p$ obeys $d / 2 \leq p \leq q$, then

$$
\begin{equation*}
\Delta_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{2.2.2}
\end{equation*}
$$

is a bounded operator. If in addition $p \in(1, \infty)$, then there is a constant $C=$ $C\left(A, A_{1}, g, G, l, p, q\right) \in[1, \infty)$ such that

$$
\begin{equation*}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C\left(\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{p}(X)}\right), \quad \forall \xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{2.2.3}
\end{equation*}
$$

REMARK 2.2.2 (Regularity of distributional solutions to elliptic partial differential equations). Suppose as in the hypotheses of Proposition 2.2.1 that $A_{1}$ is a smooth connection on $P$. By analogy with [89, Definition 2.56], we call $\xi \in L^{1}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ a distributional solution to the equation $\Delta_{A_{1}} \xi=0$ if

$$
\left(\xi, \Delta_{A_{1}} \eta\right)_{L^{2}(X)}=0, \quad \forall \eta \in C^{\infty}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

In the case of the scalar Laplace operator on functions, $C^{\infty}$-smoothness of distributional solutions is provided by Weyl's Lemma [127, Theorem 18.G]. More generally, the $C^{\infty}$-smoothness of a solution $u \in L_{\mathrm{loc}}^{1}(\Omega)$ to a scalar (second-order) elliptic equation on an open subset $\Omega \subset \mathbb{R}^{d}$ is a consequence of regularity theory for solutions in $H_{\text {loc }}^{s}(\Omega)$, for $s \in \mathbb{R}$ [42, Theorem 6.33]. Such regularity results extend to the case of elliptic systems (see [32 and references therein) and so we conclude that if $\xi$ is a distributional solution to the equation $\Delta_{A_{1}} \xi=0$, then $\xi \in C^{\infty}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$.

Proposition 2.2.1 implies that the domain of the unbounded operator $\Delta_{A}$ on $L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is

$$
\mathscr{D}_{p}\left(\Delta_{A}\right)=W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

(We omit the subscript $p$ when that is clear from the context.) In order to give criteria for when the term $\|\xi\|_{L^{p}(X)}$ can be eliminated from the right-hand side of the a priori estimate (2.2.3), we need to analyze the spectrum of the Hodge Laplacian with Sobolev coefficients. The forthcoming Proposition 2.2.3 is an analogue of 47, Theorem 8.6], for a scalar, second-order, strictly elliptic equation in divergence form with homogeneous Dirichlet boundary condition over a bounded domain $\Omega \subset \mathbb{R}^{d}$. However, it is not a direct consequence since the first and zeroth-order coefficients of the Laplace operator $\Delta_{A}$ on $\Omega^{l}(X ; \operatorname{ad} P)$ are not necessarily bounded unless $q>d$, which we do not wish to assume, for a $W^{1, q}$ connection $A$ on $P$.

Let $\mathscr{X}$ be a Banach space and $T: \mathscr{X} \rightarrow \mathscr{X}$ a bounded operator. Recall from 101. Definition 4.17 (c)] that the spectrum, $\sigma(T)$, of $T$ is the set of all $\lambda \in \mathbb{C}$ such that $T-\lambda$ is not invertible. Thus $\lambda \in \sigma(T)$ if and only if at least one of the following two statements is true: a) The range of $T-\lambda$ is not all of $\mathscr{X}$, or $b$ ) $T-\lambda$ is not one-to-one. In the latter case, $\lambda$ is an eigenvalue of $T$; the corresponding eigenspace is $\operatorname{Ker}(T-\lambda$ ); each $x \in \operatorname{Ker}(T-\lambda)$ (except $x=0$ ) is an eigenvector of $T$ and satisfies the equation $T x=\lambda x$.

If $T: \mathscr{D}(\mathscr{X}) \subset \mathscr{X} \rightarrow \mathscr{X}$ is a closed operator with dense domain, $\mathscr{D}(\mathscr{X}) \subset \mathscr{X}$, then we say that $\lambda \notin \sigma(T)$ if the operator $T-\lambda: \mathscr{D}(\mathscr{X}) \rightarrow \mathscr{X}$ has a bounded inverse and otherwise that $\lambda \in \sigma(T)$ [75, Section 5.6, p. 357], [76, Section III.6.1, pp. 174-175].

Proposition 2.2.3 (Spectral properties of a Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $l \geq 0$ be an integer. If $A$ is a $W^{1, q}$ connection on $P$ with $d / 2<q<\infty$, and $A_{1}$ is a $C^{\infty}$ reference connection on $P$, and $p \in(1, \infty)$ obeys $d / 2 \leq p \leq q$, then the spectrum, $\sigma\left(\Delta_{A}\right)$, of the unbounded operator,

$$
\Delta_{A}: \mathscr{D}\left(\Delta_{A}\right) \subset L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

is countable without accumulation points, consisting of non-negative, real eigenvalues, $\lambda$, with finite multiplicities equal to $\operatorname{dim} \operatorname{Ker}\left(\Delta_{A}-\lambda\right)$.

Proof. Corollary C. 2 implies that the operator,

$$
\Delta_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right),
$$

is Fredholm and, setting

$$
K:=\operatorname{Ker}\left(\Delta_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right),
$$

that

$$
\Delta_{A}: K^{\perp} \cap W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow K^{\perp} \cap L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right),
$$

is invertible. We denote the Green's operator of $\Delta_{A}$ by

$$
G_{A}: L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right),
$$

so that $G_{A} \Delta_{A}=1-\Pi_{A}$, where

$$
\Pi_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow K
$$

is $L^{2}$-orthogonal projection, and $\Delta_{A} G_{A}=1-\Pi_{A}$, where

$$
\Pi_{A}: L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow K
$$

again denotes $L^{2}$-orthogonal projection.
The Sobolev embedding,

$$
W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \Subset L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right),
$$

is compact by [4, Theorem 6.3] and hence the composition of $G_{A}$ with this embedding,

$$
G_{A}: L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right),
$$

is compact by [14, Proposition 6.3]. But then [101, Theorem 4.25] implies that preceding operator has the spectral properties, aside from reality and non-negativity, described in the conclusion of Proposition 2.2.3.

To relate the spectra of $G_{A}$ and $\Delta_{A}$, observe that for any $\lambda \in \mathbb{C} \backslash\{0\}$ and $\chi \in K^{\perp} \cap L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, the equation,

$$
\left(\Delta_{A}-\lambda\right) \xi=\chi,
$$

for $\xi \in K^{\perp} \cap W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is equivalent to the equation,

$$
\left(G_{A} \Delta_{A}-\lambda G_{A}\right) \xi=G_{A} \chi,
$$

that is,

$$
\left(G_{A}-\lambda^{-1}\right) \xi=-\lambda^{-1} G_{A} \chi .
$$

In other words, $\lambda \neq 0$ is in the spectrum of $\Delta_{A}$ if and only if $\lambda^{-1}$ is in the spectrum of $G_{A}$.

To see that the eigenvalues of $G_{A}$ are real as claimed, note that $\Delta_{A}$ has $L^{2}$ adjoint, $\Delta_{A}^{*}=\Delta_{A}$, and so $G_{A}$ has $L^{2}$-adjoint, $G_{A}^{*}=G_{A}$, and the operator,

$$
G_{A}: L^{2}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{2}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right),
$$

is bounded and self-adjoint. Thus, $\sigma\left(G_{A}\right) \subset \mathbb{R}$ by [99, Theorem VI.8] and hence $\sigma\left(\Delta_{A}\right) \subset \mathbb{R}$.

Finally, to see that the eigenvalues of $\Delta_{A}$ are non-negative, observe that if $\Delta_{A} \xi=\lambda \xi$ for $\lambda \in \sigma\left(\Delta_{A}\right) \backslash\{0\}$ and $\xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \backslash\{0\}$, then

$$
\begin{aligned}
& \lambda\|\xi\|_{L^{2}(X)}^{2}=(\lambda \xi, \xi)_{L^{2}(X)}=\left(\Delta_{A} \xi, \xi\right)_{L^{2}(X)} \\
& \quad=\left(d_{A} \xi, d_{A} \xi\right)_{L^{2}(X)}+\left(d_{A}^{*} \xi, d_{A}^{*} \xi\right)_{L^{2}(X)} \geq 0
\end{aligned}
$$

by (2.2.1), and thus $\lambda \geq 0$, as claimed.

REMARK 2.2.4 (Spectral properties of a Laplace operator with Sobolev coefficients on $L^{p}$ spaces and compact perturbations). We recall from Weyl's Theorem [76, Theorem IV.5.35] that if $T$ is a closed operator on a Banach space $\mathscr{X}$ and $K$ is an operator on $\mathscr{X}$ that is compact relative to $T$, then $T$ and $T+K$ have the same essential spectrum. In particular, under the hypotheses of Corollary C.2 the operator,

$$
\Delta_{A}-\Delta_{A_{1}}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

is compact by the proof of that corollary. Therefore, the essential spectrum of $\Delta_{A}$ as an unbounded operator on $L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is empty and hence the spectrum of $\Delta_{A}$ consists purely of real eigenvalues with finite multiplicity, since the same is true of $\Delta_{A_{1}}$. These observations could be used to give an alternative proof of Proposition 2.2.3, in place of the one that we provide.

Corollary 2.2.5 (A priori $L^{p}$ estimate for a Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$, and $l \geq 0$ be an integer. If $A$ is a $W^{1, q}$ connection on $P$ with $d / 2<q<\infty$, and $A_{1}$ is a $C^{\infty}$ connection on $P$, and $p \in(1, \infty)$ obeys $d / 2 \leq p \leq q$, then the kernel $\operatorname{Ker} \Delta_{A} \cap W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ of the operator (2.2.2) is finite-dimensional and

$$
\begin{equation*}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}, \quad \forall \xi \in\left(\operatorname{Ker} \Delta_{A}\right)^{\perp} \cap W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{2.2.4}
\end{equation*}
$$

where $\perp$ denotes $L^{2}$-orthogonal complement and $C=C\left(A, A_{1}, g, G, l, p, q\right) \in[1, \infty)$.
Before proceeding to the proofs of these results proper, we begin with the
Lemma 2.2.6 ( $A$ priori $L^{p}$ estimate for a Laplace operator with smooth coefficients). Assume the hypotheses on $A_{1}, d, G, l, P$, and $(X, g)$ in Proposition 2.2.1] and let $p \in(1, \infty)$. If $A$ is $C^{\infty}$, then there is a constant $C=C\left(A, A_{1}, g, G, l, p\right) \in$ $[1, \infty)$ such that

$$
\begin{equation*}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C\left(\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{p}(X)}\right), \quad \forall \xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{2.2.5}
\end{equation*}
$$

Proof. Suppose first that $\Delta_{g}$ is the Laplace-Beltrami operator on $C^{\infty}(X)$ defined by the Riemannian metric $g$. The a priori $L^{p}$ estimate for $\Delta_{g}$ analogous to (2.2.5) can be obtained from the a priori interior $L^{p}$ estimate provided by 47, Theorem 9.11] for a scalar, second-order, strictly elliptic operator with $C^{\infty}$ coefficients defined on a bounded domain $\Omega \subset \mathbb{R}^{d}$ with the aid of a $C^{\infty}$ partition of unity subordinate to a finite set of coordinate charts covering the closed manifold, $X$. For the general case, one first chooses in addition a set of local trivializations for $\Lambda^{l} \otimes \operatorname{ad} P$ corresponding to the coordinate neighborhoods, after shrinking those neighborhoods if needed. The Bochner-Weitzenböck formula 43, Equation (C.7)], 83, Equation (II.1)] for $\Delta_{A}$ implies that $\Delta_{A}-\nabla_{A}^{*} \nabla_{A}$ is a first-order differential operator with $C^{\infty}$ coefficients and that $\Delta_{A}$ has principal symbol given by the $C^{\infty}$ Riemannian metric $g$ times the identity endomorphism of $\Lambda^{l} \otimes \operatorname{ad} P$. (In fact, $\Delta_{A}=\nabla_{A}^{*} \nabla_{A}$ when $l=0$.) The (first-order) covariant derivative of $\xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ may be estimated with the following analogue of the interpolation inequality [47, Theorem 7.27], valid for $p \in[1, \infty)$,

$$
\begin{equation*}
\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq \varepsilon\|\xi\|_{W_{A_{1}}^{2, p}(X)}+C \varepsilon^{-1}\|\xi\|_{L^{p}(X)} \tag{2.2.6}
\end{equation*}
$$

where $C=C\left(A_{1}, g, G, l, p\right) \in[1, \infty)$ and $\varepsilon$ is any positive constant. The conclusion now follows by combining the preceding observations and using rearrangement with small $\varepsilon$ to remove the term $\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)}$ from the right-hand side.

We can now proceed to the
Proof of Proposition 2.2.1. We choose a $C^{\infty}$ connection, $A_{s}$, on $P$ that we regard as a smooth approximation to $A$. We write $A=A_{s}+a$, with $a \in$ $W_{A_{1}}^{1, q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ obeying a bound $\|a\|_{W_{A_{1}}^{1, q}(X)} \leq \varepsilon$ with small constant $\varepsilon \in(0,1]$ to be chosen during the proof, and write $A_{s}=A_{1}+a_{1}$, where $a_{1} \in C^{\infty}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ may be 'large'. We expand $\Delta_{A}=\Delta_{A_{s}+a}$ to give

$$
\begin{equation*}
\Delta_{A} \xi=\Delta_{A_{s}} \xi+\nabla_{A_{s}} a \times \xi+a \times \nabla_{A_{s}} \xi+a \times a \times \xi \tag{2.2.7}
\end{equation*}
$$

and thus, for $\xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$,

$$
\begin{equation*}
\Delta_{A} \xi=\Delta_{A_{s}} \xi+\nabla_{A_{1}} a \times \xi+a_{1} \times a \times \xi+a \times \nabla_{A_{1}} \xi+a \times a \times \xi \tag{2.2.8}
\end{equation*}
$$

We define $r \in[p, \infty]$ by $1 / p=1 / q+1 / r$ and recall that by [4, Theorem 4.12] we have i) $W^{2, p}(X) \subset L^{r}(X)$ for any $r \in[1, \infty)$ when $p=d / 2$, and ii) $W^{2, p}(X) \subset L^{\infty}(X)$ when $p>d / 2$. The expansion (2.2.8) and continuous Sobolev multiplication map, $L^{q}(X) \times L^{r}(X) \rightarrow L^{p}(X)$, yield

$$
\begin{align*}
&\left\|\left(\Delta_{A}-\Delta_{A_{s}}\right) \xi\right\|_{L^{p}(X)} \leq z\left\|\nabla_{A_{1}} a\right\|_{L^{q}(X)}\|\xi\|_{L^{r}(X)}+\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)}  \tag{2.2.9}\\
&+z\left\|a_{1}\right\|_{C(X)}\|a\|_{L^{2 p}(X)}\|\xi\|_{L^{2 p}(X)}+z\left\|\left.a\right|^{2}\right\|_{L^{q}(X)}\|\xi\|_{L^{r}(X)}
\end{align*}
$$

where $z=z(g, G, l) \in[1, \infty)$. To ensure a continuous Sobolev embedding, $W^{1, p}(X) \subset$ $L^{d}(X)$, by 4. Theorem 4.12], we need $p^{*}=d p /(d-p) \geq d$, that is, $p \geq d-p$ or $p \geq d / 2$, which we assume in our hypotheses.

To ensure a continuous Sobolev embedding, $W^{1, q}(X) \subset L^{2 q}(X)$, when $q<d$, we need $q^{*}=d q /(d-q) \geq 2 q$, that is, $d \geq 2 d-2 q$ or $2 q \geq d$ or $q \geq d / 2$, which follows from our hypothesis that $q \geq p \geq d / 2$; when $q \geq d$, the fact that $W^{1, q}(X) \subset L^{2 q}(X)$ is a continuous Sobolev embedding is immediate from [4, Theorem 4.12].

Consequently, by the applying these continuous Sobolev embeddings and the Kato Inequality [43, Equation (6.20)] to the preceding inequality, we obtain

$$
\begin{aligned}
\left\|\left(\Delta_{A}-\Delta_{A_{s}}\right) \xi\right\|_{L^{p}(X)} \leq & z\left(\left\|\nabla_{A_{1}} a\right\|_{L^{q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}^{2}\right)\|\xi\|_{W_{A_{1}}^{2, p}(X)} \\
& +z\left\|a_{1}\right\|_{C(X)}\|a\|_{W_{A_{1}}^{1, q}(X)}\|\xi\|_{W_{A_{1}}^{1, p}(X)}+\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)}
\end{aligned}
$$

where $z=z(g, G, l, p, q) \in[1, \infty)$.
When $q<d$, we recall from [4, Theorem 4.12] that there is a continuous embedding $W^{1, q}(X) \subset L^{q^{*}}(X)$, where $q^{*}=d q /(d-q)$. Hence, $1 / q^{*}=1 / q-1 / d$ or $1 / q=1 / q^{*}+1 / d$ and so, using $p \leq q$,

$$
\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq z\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{q}(X)} \leq z\|a\|_{L^{q^{*}}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{d}(X)}
$$

where $z=z(g, p, q) \in[1, \infty)$, and therefore, by the preceding continuous Sobolev embeddings,

$$
\begin{equation*}
\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq C\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{W_{A_{1}}^{1, p}(X)} \tag{2.2.10}
\end{equation*}
$$

where $C=C(g, G, l, p, q) \in[1, \infty)$. When $q=d$ and $d / 2 \leq p<d$, we can define $t \in[d, \infty)$ by $1 / p=1 / t+1 / d$ and apply the continuous Sobolev embedding $W^{1, d}(X) \subset L^{t}(X)$ from [4, Theorem 4.12] to give

$$
\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq z\|a\|_{L^{t}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{d}(X)} \leq C\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{W_{A_{1}}^{1, p}(X)},
$$

and (2.2.10) again holds; when $q=d=p$, we can simply use the embedding $W^{1, d}(X) \subset L^{t}(X)$ for any $t \in[1, \infty)$ and observe that (2.2.10) holds from

$$
\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{d}(X)} \leq z\|a\|_{L^{2 d}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{2 d}(X)} \leq C\|a\|_{W_{A_{1}}^{1, d}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{W_{A_{1}}^{1, d}(X)}
$$

Finally, when $q>d$ we have the continuous Sobolev embedding $W^{1, q}(X) \subset C(X)$ from [4, Theorem 4.12] and so

$$
\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq z\|a\|_{C(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq C\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)},
$$

which also yields (2.2.10).
Combining our previous $L^{p}$ bound for $\left(\Delta_{A}-\Delta_{A_{s}}\right) \xi$ with the inequality (2.2.10) gives

$$
\begin{align*}
\left\|\left(\Delta_{A}-\Delta_{A_{s}}\right) \xi\right\|_{L^{p}(X)} \leq & z\left(\left\|\nabla_{A_{1}} a\right\|_{L^{q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}^{2}\right)\|\xi\|_{W_{A_{1}}^{2, p}(X)}  \tag{2.2.11}\\
& +z\left\|a_{1}\right\|_{C(X)}\|a\|_{W_{A_{1}}^{1, q}(X)}\|\xi\|_{W_{A_{1}}^{1, p}(X)} \\
& +z\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{W_{A_{1}}^{1, p}(X)},
\end{align*}
$$

where $z=z(g, G, l, p, q) \in[1, \infty)$. Combining the preceding bound with the a priori estimate (2.2.5) for $\Delta_{A_{s}}$ provided by Lemma 2.2.6.

$$
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C_{0}\left(\left\|\Delta_{A_{s}} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{p}(X)}\right)
$$

with constant denoted by $C_{0}=C_{0}\left(A_{1}, A_{s}, g, p\right) \in[1, \infty)$ for clarity, yields

$$
\begin{aligned}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq & C_{0}\left(\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{p}(X)}\right) \\
& +z\left(\|a\|_{W_{A_{1}}^{1, q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}^{2}\right)\|\xi\|_{W_{A_{1}}^{2, p}(X)} \\
& +z\left\|a_{1}\right\|_{C(X)}\|a\|_{W_{A_{1}}^{1, q}(X)}\|\xi\|_{W_{A_{1}}^{1, p}(X)} .
\end{aligned}
$$

We can choose $a=A-A_{s}$ so that $\|a\|_{W_{A_{1}}^{1, q}(X)} \leq \varepsilon$ for small constant $\varepsilon \in(0,1]$, but we are not at liberty to choose $a_{1}=A_{s}-A_{1}$ to be $W_{A_{1}}^{1, q}(X)$-small (since $A_{1}$ is a fixed $C^{\infty}$ connection on $P$ and $A$ is an arbitrary $W^{1, q}$ connection on $P$ ). Thus in our forthcoming rearrangement arguments we first apply the interpolation inequality (2.2.6),

$$
\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq \delta\|\xi\|_{W_{A_{1}}^{2, p}(X)}+C_{1} \delta^{-1}\|\xi\|_{L^{p}(X)}
$$

where $C_{1}=C_{1}\left(A_{1}, g\right) \in[1, \infty)$ and $\delta=\delta\left(A_{1},\left\|A_{s}-A_{1}\right\|_{C(X)}, g, G, l, p, q\right) \in(0,1]$ is a constant chosen small enough that

$$
\delta z\left\|a_{1}\right\|_{C(X)} \leq 1 / 2
$$

and thus, for a constant $C_{2}=C_{2}\left(A_{1}, A_{s}, g, G, l, p, q\right) \in[1, \infty)$,

$$
\begin{aligned}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq 2 C_{0} & \left(\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{p}(X)}\right) \\
& +2 z\left(\|a\|_{W_{A_{1}}^{1, q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}^{2}\right)\|\xi\|_{W_{A_{1}}^{2, p}(X)}+C_{2}\|\xi\|_{L^{p}(X)}
\end{aligned}
$$

Provided $\|a\|_{W_{A_{1}^{1, q}(X)}} \leq \varepsilon$ and we choose $\varepsilon \equiv \varepsilon(g, G, l, p, q)=1 /(8 z) \in(0,1]$ in the preceding inequality, rearrangement yields the desired estimate (2.2.3).

Our proof of (2.2.3) also verifies that the operator $\Delta_{A}$ in (2.2.2) is bounded since $\Delta_{A_{s}}$ is bounded with the same domain and range spaces. This completes the proof of Proposition 2.2.1.

## Next, we have the

Proof of Corollary 2.2.5. The fact that $\operatorname{Ker} \Delta_{A} \cap W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is finite-dimensional follows from Corollary C.2. We observe that, by increasing the constant $C$ as needed, the term $\|\xi\|_{L^{p}(X)}$ appearing on the right-hand side of the inequality (2.2.3) can be replaced by $\|\xi\|_{L^{2}(X)}$. This is clear when $p \leq 2$, while if $p>2$, we can choose $s \in(p, \infty)$ and apply the interpolation inequality 47, Equation (7.10)],

$$
\|\xi\|_{L^{p}(X)} \leq \delta\|\xi\|_{L^{s}(X)}+\delta^{-\nu}\|\xi\|_{L^{2}(X)}
$$

for $\nu:=(1 / 2-1 / p) /(1 / p-1 / s)>0$ and arbitrary positive $\delta$. Because $p \geq d / 2$, we have a continuous Sobolev embedding $W^{2, p}(X) \subset L^{s}(X)$ as already observed in the proof of Proposition 2.2.1, so

$$
\|\xi\|_{L^{p}(X)} \leq C_{1} \delta\|\xi\|_{W_{A_{1}}^{2, p}(X)}+\delta^{-\nu}\|\xi\|_{L^{2}(X)}
$$

where $C_{1}=C_{1}\left(A_{1}, g, l, p\right) \in[1, \infty)$. Hence, for $\delta\left(A_{1}, g, l, p\right) \in(0,1]$ given by $\delta=1 /\left(2 C C_{1}\right)$, we can use rearrangement in (2.2.3) to replace $\|\xi\|_{L^{p}(X)}$ by $\|\xi\|_{L^{2}(X)}$. Therefore, the estimate (2.2.3) implies

$$
\begin{equation*}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C\left(\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{2}(X)}\right), \quad \forall \xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{2.2.12}
\end{equation*}
$$

Proposition 2.2.3 implies that the spectrum $\sigma\left(\Delta_{A}\right)$ of $\Delta_{A}$ on $L^{2}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ consists purely of non-negative eigenvalues and is discrete with no accumulation points. Let $\mu[A]$ denote the least positive eigenvalue of $\Delta_{A}$ on $L^{2}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ and recall from [17, Rayleigh's Theorem, p. 16] (or more generally [108, Theorem 6.5.1], applied to the Green's operator $G_{A}$ for $\Delta_{A}$ ) that

$$
\mu[A]=\inf _{\xi \in\left(\operatorname{Ker} \Delta_{A}\right)^{\perp}} \frac{\left(\xi, \Delta_{A} \xi\right)_{L^{2}(X)}}{\|\xi\|_{L^{2}(X)}^{2}}
$$

where $\left(\operatorname{Ker} \Delta_{A}\right)^{\perp}$ is the $L^{2}$-orthogonal complement of $\operatorname{Ker} \Delta_{A} \subset H_{A_{1}}^{1}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, with equality achieved in the infimum if and only if $\xi$ is an eigenvector with eigenvalue $\mu[A]$. Therefore, if $\xi \in\left(\operatorname{Ker} \Delta_{A}\right)^{\perp} \cap W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, then

$$
\|\xi\|_{L^{2}(X)} \leq \mu[A]^{-1}\left\|\Delta_{A} \xi\right\|_{L^{2}(X)}
$$

Hence, the inequality (2.2.4) follows from the preceding eigenvalue bound and (2.2.12) when $p \geq 2$.

For the case $d / 2 \leq p<2$ (which forces $d=2$ or 3 ), we observe that

$$
\|\xi\|_{L^{2}(X)}^{2} \leq\left.\frac{1}{\mu[A]}\left(\Delta_{A} \xi, \xi\right)\right|_{L^{2}(X)} \leq \frac{1}{\mu[A]}\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}\|\xi\|_{L^{p^{\prime}}(X)}
$$

where $p^{\prime}$ is defined by $1 / p+1 / p^{\prime}=1$, and so

$$
\begin{aligned}
&\|\xi\|_{L^{2}(X)} \leq \frac{1}{\sqrt{\mu[A]}}\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}^{1 / 2}\|\xi\|_{L^{p^{\prime}}(X)}^{1 / 2} \\
& \leq \frac{1}{2 \sqrt{\mu[A]}}\left(\delta^{-1}\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\delta\|\xi\|_{L^{p^{\prime}}(X)}\right)
\end{aligned}
$$

for arbitrary positive $\delta$. If $d=2$ and $p \in(1,2)$ (so $p^{\prime} \in(2, \infty)$, then $W^{2, p}(X) \subset$ $C^{0}(X)$ by [4, Theorem 4.12] and thus $W^{2, p}(X) \subset L^{p^{\prime}}(X)$; if $d=3$ and $p \in[3 / 2,2)$ (so $p^{\prime} \in(2,3]$, then $W^{2, p}(X) \subset L^{q}(X)$ for $p=3 / 2$ and any $q \in[1, \infty)$ while $W^{2, p}(X) \subset C^{0}(X)$ for $p>3 / 2$ by [4, Theorem 4.12] and thus $W^{2, p}(X) \subset L^{p^{\prime}}(X)$ in either case. Hence, for $d=2,3$ and $p>1$ obeying $d / 2 \leq p<2$, we have the bound

$$
\|\xi\|_{L^{p^{\prime}}(X)} \leq C_{2}\|\xi\|_{W_{A_{1}}^{2, p}(X)}
$$

where $C_{2}=C_{2}\left(A_{1}, g, l, p\right) \in[1, \infty)$. Combining the preceding inequalities for the case $d / 2 \leq p<2$ gives

$$
\|\xi\|_{L^{2}(X)} \leq \frac{1}{2 \sqrt{\mu[A]}}\left(\delta^{-1}\left\|\Delta_{A} \xi\right\|_{L^{p}(X)}+\delta C_{2}\|\xi\|_{W_{A_{1}}^{2, p}(X)}\right)
$$

Combining the preceding inequality with (2.2.12) and applying rearrangement by choosing $\delta=\sqrt{\mu[A]} /\left(C C_{2}\right)$ yields the desired the inequality (2.2.4) for this case too.

### 2.3. Regularity for distributional solutions to an elliptic equation with Sobolev coefficients

We shall need to address a complication that arises when establishing regularity for distributional solutions to an elliptic equation with Sobolev coefficients. We shall confine our discussion to the Hodge Laplace operator, though one can clearly establish more general results of this kind.

Lemma 2.3.1 (Regularity for distributional solutions to an equation defined by the Hodge Laplace operator for a Sobolev connection). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group and $P$ be a smooth principal $G$-bundle over $X$. Let $A_{1}$ be a $C^{\infty}$ connection on $P$, and $A$ be a $W^{1, q}$ connection on $P$ with $d / 2<q<\infty$, and $l \geq 0$ be an integer. If $q^{\prime} \in(1, \infty)$ is the dual exponent defined by $1 / q+1 / q^{\prime}=1$ and $\eta \in L^{q^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is a distributional solution to

$$
\Delta_{A} \eta=0
$$

then $\eta \in W_{A_{1}}^{2, q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$.
Proof. We recall from equation (C.3) in the proof of Corollary C. 2 (with $p=q$ ) that

$$
\Delta_{A}-\Delta_{A_{1}}: W_{A_{1}}^{1, u}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

[^2]is a bounded operator, where allowable values of $u \in(1, \infty)$ are given by
\[

u= $$
\begin{cases}d+\varepsilon & \text { if } d / 2<q<d  \tag{2.3.1}\\ 2 d & \text { if } q=d \\ q & \text { if } q>d\end{cases}
$$
\]

and $\varepsilon \in(0,1]$ is chosen small enough that $q^{*}=d q /(d-q) \geq d+\varepsilon$, which is possible when $q>d / 2$. Consequently, the dual operator,

$$
\Delta_{A}-\Delta_{A_{1}}: L^{q^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{-1, u^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

is bounded, where the dual exponent $u^{\prime} \in(1, \infty)$ is defined by $1 / u+1 / u^{\prime}=1$.
We write $\Delta_{A}=\Delta_{A_{1}}+\left(\Delta_{A}-\Delta_{A_{1}}\right)$, set

$$
\alpha:=\left(\Delta_{A_{1}}-\Delta_{A}\right) \eta \in W_{A_{1}}^{-1, u^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

and observe that $\eta$ is a distributional solution to

$$
\begin{equation*}
\Delta_{A_{1}} \eta=\alpha \tag{2.3.2}
\end{equation*}
$$

We now appeal to regularity for distributional solutions to an equation (namely, (2.3.2)) defined by an elliptic operator $\Delta_{A_{1}}$ with $C^{\infty}$ coefficients (see Remark 2.2.2 and [32]) to conclude that $\eta \in W_{A_{1}}^{1, u}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$. The range of exponents $u \in$ $(d, \infty)$ given by (2.3.1) ensures that $\eta \in C\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, since $W^{1, u}(X) \subset C(X)$ by [4, Theorem 4.12] when $u>d$. Moreover, the estimate (C.5) for $\left(\Delta_{A_{1}}-\Delta_{A}\right) \eta$ in terms of $a=A-A_{1}$ and $\eta$ ensures that

$$
\left(\Delta_{A_{1}}-\Delta_{A}\right) \eta \in L^{q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

In particular, $\alpha \in L^{q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ and regularity for solutions to an elliptic equation (that is, (2.3.2)) with $C^{\infty}$ coefficients implies that $\eta \in W_{A_{1}}^{2, q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, as desired.

### 2.4. Surjectivity of a perturbed Laplace operator

We now consider surjectivity properties of a perturbation of a Laplace operator, namely

Lemma 2.4.1 (Surjectivity of a perturbed Laplace operator). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group and $P$ be a smooth principal $G$-bundle over $X$. Let $A_{1}$ be a $C^{\infty}$ connection on $P$ and $A$ be a $W^{1, q}$ connection on $P$ with $d / 2<q<\infty$. Then there is a constant $\delta=\delta(A, g) \in(0,1]$ with the following significance. If $a \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ obeys

$$
\begin{equation*}
\|a\|_{L^{d}(X)}<\delta \quad \text { when } d \geq 3 \quad \text { or } \quad\|a\|_{L^{4}(X)}<\delta \quad \text { when } d=2 \tag{2.4.1}
\end{equation*}
$$

then the operator,

$$
\begin{equation*}
d_{A}^{*} d_{A+a}:\left(\operatorname{Ker} \Delta_{A}\right)^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow\left(\operatorname{Ker} \Delta_{A}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P) \tag{2.4.2}
\end{equation*}
$$

is well-defined and surjective.
REmARK 2.4.2 (Comparison with the argument due to Donaldson and Kronheimer). Our proof of Lemma 2.4.1 is based on a similar argument arising in Donaldson and Kronheimer [27, p. 66], as part of their version of the proof of Uhlenbeck's local Coulomb gauge-fixing theorem [27, Theorem 2.3.7]. We make some adjustments to their argument for reasons which we briefly explain here. When $B$ is a Sobolev connection matrix [27, p. 66], one has to take into account the
possibility that their operator $d^{*} d_{B}$, for a $\mathfrak{g}$-valued Sobolev one-form $B$ over $S^{4}$, could have dense range but still fail to be surjective, so one would first have to verify, for example, that $d^{*} d_{B}$ has closed range. Similarly, while elliptic regularity ensures that if $\eta \in L^{2}\left(S^{4} ; \operatorname{ad} P\right) \cap \operatorname{Ker} d^{*} d_{B}$ then $\eta$ is $C^{\infty}\left(S^{4} ; \operatorname{ad} P\right)$ when $B$ is $C^{\infty}$, those regularity issues become more subtle when $B$ is merely a Sobolev one-form.

REMARK 2.4.3 (Dual spaces and direct sums of subspaces of Banach spaces). Our proof of Lemma 2.4.1 is clarified by a few observations concerning the dual space of a finite-dimensional subspace $K$ of a Banach space $\mathscr{X}$ that is continuously embedded in a Hilbert space, $\mathscr{H}$. Since $K$ has finite dimension, it has a closed complement, $\mathscr{X}_{0} \subset \mathscr{X}$, such that $\mathscr{X}=\mathscr{X}_{0} \oplus K$ (vector space direct sum) by 101, Definition 4.20 and Lemma 4.21(a)]. We may also simply define $\mathscr{X}_{0}:=K^{\perp} \cap \mathscr{X}$, where $K^{\perp} \subset \mathscr{H}$ is the orthogonal complement of $K \subset \mathscr{H}$ and a closed subspace of $\mathscr{H}$ [101, Theorem 12.4]. Therefore, $K^{\perp} \cap \mathscr{X} \subset \mathscr{X}$ is also a closed subspace and a closed complement for $K \subset \mathscr{X}$ in the sense of [101, Definition 4.20]. In summary:

$$
\begin{equation*}
\mathscr{X}=\mathscr{X}_{0} \oplus K, \quad \text { with } \mathscr{X}_{0}:=K^{\perp} \cap \mathscr{X}, \tag{2.4.3}
\end{equation*}
$$

is an orthogonal direct sum of a finite-dimensional and a closed subspace.
Recall that if $M \subset \mathscr{X}$ is any subspace, then $M^{\circ}:=\left\{\alpha \in \mathscr{X}^{*}:\langle u, \alpha\rangle_{\mathscr{X} \times \mathscr{X}^{*}}=\right.$ $0, \forall u \in M\}$ denotes the annihilator of $M$ in $\mathscr{X}^{*}$ [101, Section 4.6]. Since $\mathscr{X}=$ $\mathscr{X}_{0} \oplus K$, we have $\mathscr{X}^{*}=\mathscr{X}_{0}^{\circ} \oplus K^{\circ}$, the direct sum of the annihilators in $\mathscr{X}^{*}$ of the subspaces $\mathscr{X}_{0}$ and $K$ of $\mathscr{X}$.

Because $K \subset \mathscr{X}$ and $\mathscr{X}_{0} \subset \mathscr{X}$ are closed subspaces, we have $K^{\circ} \cong(\mathscr{X} / K)^{*}$ and $\mathscr{X}_{0}^{\circ} \cong K^{*}$ by [101, Theorem 4.9]. Since $\mathscr{X}=\mathscr{X}_{0} \oplus K$, then $\mathscr{X} / K \cong \mathscr{X}_{0}$ and $\mathscr{X} / \mathscr{X}_{0} \cong K$, so $K^{\circ} \cong \mathscr{X}_{0}^{*}$ and $\mathscr{X}_{0}^{\circ} \cong K^{*} \cong K$, where the final isomorphism follows from the fact that $K$ is finite-dimensional. In particular, $\mathscr{X}^{*} \cong K^{*} \oplus \mathscr{X}_{0}^{*} \cong$ $K \oplus \mathscr{X}_{0}^{*}$. Since $K^{*}$ is finite-dimensional, then $K^{*} \subset \mathscr{X}^{*}$ is a closed subspace by [14, Proposition 11.1] and so the complement $\mathscr{X}_{0}^{*} \subset \mathscr{X}^{*}$ is a closed subspace by 101, Definition 4.20 and Lemma 4.21]. In summary:

$$
\begin{equation*}
\mathscr{X}^{*}=\mathscr{X}_{0}^{*} \oplus K \tag{2.4.4}
\end{equation*}
$$

is a direct sum of a finite-dimensional and a closed subspace.
Proof of Lemma 2.4.1. Note that when $a=0$, the operator (2.4.2) is invertible by Corollary C.2 To see that the operator (2.4.2) is well-defined, observe that if $\xi=d_{A}^{*} d_{A+a} \chi \in L^{q}(X ; \operatorname{ad} P)$ for some $\chi \in W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$ and $\eta \in$ Ker $\Delta_{A} \cap L^{q}(X ; \operatorname{ad} P)$, then

$$
(\xi, \eta)_{L^{2}(X)}=\left(d_{A}^{*} d_{A+a} \chi, \eta\right)_{L^{2}(X)}=\left(d_{A+a} \chi, d_{A} \eta\right)_{L^{2}(X)}=0
$$

since $\Delta_{A} \eta=d_{A}^{*} d_{A} \eta=0$ and hence

$$
\left(d_{A}^{*} d_{A} \eta, \eta\right)_{L^{2}(X)}=\left\|d_{A} \eta\right\|_{L^{2}(X)}^{2}=0
$$

so that $d_{A} \eta=0$. Thus,

$$
\begin{align*}
\operatorname{Ran}\left(d_{A}^{*} d_{A+a}: W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow L^{q}(X ; \operatorname{ad} P)\right) &  \tag{2.4.5}\\
& \subset\left(\operatorname{Ker} \Delta_{A}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)
\end{align*}
$$

and so the operator (2.4.2) is well-defined.
It is convenient to abbreviate $\mathscr{X}:=W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$, and $K:=\operatorname{Ker} \Delta_{A} \cap$ $W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$, and $\mathscr{Y}:=L^{q}(X ; \operatorname{ad} P)$, and $\mathscr{H}:=L^{2}(X ; \operatorname{ad} P)$, and denote $T=$
$d_{A}^{*} d_{A+a}$ in (2.4.2). The kernel $K \subset \mathscr{X}$ is finite-dimensional by Corollary C.2 and its $L^{2}$-orthogonal complements $K^{\perp} \cap \mathscr{X}$ and $K^{\perp} \cap \mathscr{Y}$ and $K^{\perp}$ provide closed complements of $K$ in $\mathscr{X}$ and $\mathscr{Y}$ and $\mathscr{H}$, respectively. Similarly, we let $K^{\perp}$ and $K^{\perp} \cap \mathscr{Y}^{*}$ and $K^{\perp} \cap \mathscr{X}^{*}$ denote the closed complements of $K^{*} \cong K$ in $\mathscr{H}^{*} \cong \mathscr{H}$ and $\mathscr{Y}^{*}$ and $\mathscr{X}^{*}$, respectively, where $\mathscr{Y}^{*}=L^{q^{\prime}}(X ; \operatorname{ad} P)$, with $q^{\prime} \in(1, \infty)$ defined by $1 / q+1 / q^{\prime}=1$, and $\mathscr{X}^{*}=W_{A_{1}}^{-2, q^{\prime}}(X ; \operatorname{ad} P)$.

If $N \subset \mathscr{Y}^{*}$ is any subspace, we recall that the annihilator [101, Section 4.6] of $N$ in $\mathscr{Y}$ is

$$
{ }^{\circ} N:=\left\{\xi \in \mathscr{Y}:\langle\xi, \alpha\rangle_{\mathscr{Y}} \times \mathscr{Y}^{*}=0, \forall \alpha \in \mathscr{Y}^{*}\right\}
$$

and $\langle\cdot, \cdot\rangle_{\mathscr{Y} \times \mathscr{Y}^{*}}: \mathscr{Y} \times \mathscr{Y}^{*} \rightarrow \mathbb{R}$ is the canonical pairing. The operator, $T: \mathscr{X} \rightarrow \mathscr{Y}$, is Fredholm by Lemma C. 4 and we can identify its range using

$$
\begin{aligned}
\operatorname{Ran}(T: \mathscr{X} \rightarrow \mathscr{Y}) & =\overline{\operatorname{Ran}}(T: \mathscr{X} \rightarrow \mathscr{Y}) \quad \text { (by closed range) } \\
& \left.={ }^{\circ} \operatorname{Ker}\left(T^{*}: \mathscr{Y}^{*} \rightarrow \mathscr{X}^{*}\right) \quad(\text { by }[\mathbf{1 4}, \text { Corollary } 2.18(i v))]\right) .
\end{aligned}
$$

Therefore, we have shown that

$$
\operatorname{Ran}(T: \mathscr{X} \rightarrow \mathscr{Y})={ }^{\circ} \operatorname{Ker}\left(T^{*}: \mathscr{Y}^{*} \rightarrow \mathscr{X}^{*}\right)
$$

and so $\operatorname{Ran}(T: \mathscr{X} \rightarrow \mathscr{Y})=\mathscr{Y}$ if and only if $\operatorname{Ker}\left(T^{*}: \mathscr{Y}^{*} \rightarrow \mathscr{X}^{*}\right)=0$. Similarly, because $T: \mathscr{X} \rightarrow \mathscr{Y}$ has closed range and $\operatorname{Ran}(T: \mathscr{X} \rightarrow \mathscr{Y}) \subset K^{\perp} \cap \mathscr{Y}$ by (2.4.5), the operator,

$$
T: K^{\perp} \cap \mathscr{X} \rightarrow K^{\perp} \cap \mathscr{Y},
$$

also has closed range and we obtain

$$
\operatorname{Ran}\left(T: K^{\perp} \cap \mathscr{X} \rightarrow K^{\perp} \cap \mathscr{Y}\right)={ }^{\circ} \operatorname{Ker}\left(T^{*}: K^{\perp} \cap \mathscr{Y}^{*} \rightarrow K^{\perp} \cap \mathscr{X}^{*}\right) .
$$

Consequently,

$$
\begin{align*}
\operatorname{Ran}\left(T: K^{\perp} \cap \mathscr{X} \rightarrow K^{\perp} \cap \mathscr{Y}\right) & =K^{\perp} \cap \mathscr{Y}  \tag{2.4.6}\\
& \Longleftrightarrow \operatorname{Ker}\left(T^{*}: K^{\perp} \cap \mathscr{Y}^{*} \rightarrow K^{\perp} \cap \mathscr{X}^{*}\right)=0
\end{align*}
$$

If $T^{*}: K^{\perp} \cap \mathscr{Y}^{*} \rightarrow K^{\perp} \cap \mathscr{X}^{*}$ were not injective, there would be a non-zero $\eta \in K^{\perp} \cap \mathscr{Y}^{*}$ such that $T^{*} \eta=0$. In other words, because $T^{*}=d_{A+a}^{*} d_{A}$, there would be a non-zero $\eta \in K^{\perp} \cap L^{q^{\prime}}(X ; \operatorname{ad} P)$ such that

$$
d_{A+a}^{*} d_{A} \eta=0 \in W_{A_{1}}^{-2, q^{\prime}}(X ; \operatorname{ad} P)
$$

that is, $\left(d_{A+a}^{*} d_{A} \eta\right)(\xi)=0$ for all $\xi \in K^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$ or equivalently,

$$
\begin{aligned}
&\left\langle\chi, d_{A+a}^{*} d_{A} \eta\right\rangle_{\mathscr{X} \times \mathscr{X}^{*}}=\left\langle\eta, d_{A}^{*} d_{A+a} \chi\right\rangle_{\mathscr{Y} \times \mathscr{Y}^{*}}=\left(\eta, d_{A}^{*} d_{A+a} \chi\right)_{L^{2}(X)}=0 \\
& \forall \chi \in K^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) .
\end{aligned}
$$

Lemma 2.3.1 implies that $\eta \in W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$. Observe that, by writing $d_{A+a} \chi=$ $d_{A} \chi+[a, \chi]$,

$$
\begin{aligned}
0 & =\left(d_{A}^{*} d_{A+a} \chi, \eta\right)_{L^{2}(X)} \\
& =\left(d_{A+a} \chi, d_{A} \eta\right)_{L^{2}(X)} \\
& =\left(d_{A} \chi, d_{A} \eta\right)_{L^{2}(X)}+\left([a, \chi], d_{A} \eta\right)_{L^{2}(X)}
\end{aligned}
$$

Since $\eta \perp \operatorname{Ker} \Delta_{A} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$ and letting $\mu[A]$ denote the least positive eigenvalue of the Laplace operator $\Delta_{A}$ on $L^{2}(X ; \operatorname{ad} P)$ provided by Proposition 2.2.3, we have

$$
\mu[A] \leq \frac{\left(\eta, d_{A}^{*} d_{A} \eta\right)_{L^{2}(X)}}{\|\eta\|_{L^{2}(X)}^{2}}=\frac{\left\|d_{A} \eta\right\|_{L^{2}(X)}^{2}}{\|\eta\|_{L^{2}(X)}^{2}}
$$

and thus,

$$
\|\eta\|_{L^{2}(X)} \leq \mu[A]^{-1 / 2}\left\|d_{A} \eta\right\|_{L^{2}(X)}
$$

Hence, we obtain

$$
\begin{equation*}
\|\eta\|_{W_{A}^{1,2}(X)} \leq C_{1}\left\|d_{A} \eta\right\|_{L^{2}(X)} \tag{2.4.7}
\end{equation*}
$$

for a constant $C_{1}=C_{1}(A, g)=1+\mu[A]^{-1 / 2} \in[1, \infty)$. For $d \geq 3$ and using $1 / 2=$ $(d-2) / 2 d+1 / d$ and the continuous multiplication $L^{2 d /(d-2)}(X) \times L^{d}(X) \rightarrow L^{2}(X)$, we see that

$$
\begin{aligned}
\left|\left([a, \chi], d_{A} \eta\right)_{L^{2}(X)}\right| & \leq\|[a, \chi]\|_{L^{2}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)} \\
& \leq c\|a\|_{L^{d}(X)}\|\chi\|_{L^{2 d /(d-2)}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)} \\
& \leq C_{2}\|a\|_{L^{d}(X)}\|\chi\|_{W_{A}^{1,2}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)}
\end{aligned}
$$

where $c \in[1, \infty)$ and the constant $C_{2}=C_{2}(g) \in[1, \infty)$ is a multiple by $c$ of the norm of the continuous Sobolev embedding $W^{1,2}(X) \subset L^{2 d /(d-2)}(X)$ provided by [4. Theorem 4.12]. Hence, setting $\chi=\eta$ and applying the a priori estimate (2.4.7), the preceding identity and inequalities yield

$$
\begin{aligned}
\left\|d_{A} \eta\right\|_{L^{2}(X)}^{2} & =\left|\left([a, \chi], d_{A} \eta\right)_{L^{2}(X)}\right| \\
& \leq C_{1} C_{2}\|a\|_{L^{d}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)}^{2}
\end{aligned}
$$

and so, if $\eta \not \equiv 0$, we have

$$
\|a\|_{L^{d}(X)} \geq C_{1} C_{2}
$$

contradicting our hypothesis (2.4.1) that $\|a\|_{L^{d}(X)}<\delta$, with $\delta$ small.
For the case $d=2$, we instead use the continuous multiplication $L^{4}(X) \times$ $L^{4}(X) \rightarrow L^{2}(X)$ and continuous Sobolev embedding $W^{1,2}(X) \subset L^{r}(X)$ for $1 \leq r<$ $\infty$ provided by [4, Theorem 4.12]. In particular, using the embedding $W^{1,2}(X) \subset$ $L^{4}(X)$ with constant $C_{3}=C_{3}(g) \in[1, \infty)$ a multiple by $c \in[1, \infty)$ of the norm of that embedding, we obtain

$$
\begin{aligned}
\left\|d_{A} \eta\right\|_{L^{2}(X)}^{2} & \leq\left|\left([a, \eta], d_{A} \eta\right)_{L^{2}(X)}\right| \\
& \leq c\|[a, \eta]\|_{L^{2}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)} \\
& \leq c\|a\|_{L^{4}(X)}\|\eta\|_{L^{4}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)} \\
& \leq C_{3}\|a\|_{L^{4}(X)}\|\eta\|_{W_{A}^{1,2}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)} \\
& \left.\leq C_{1} C_{3}\|a\|_{L^{4}(X)}\left\|d_{A} \eta\right\|_{L^{2}(X)}^{2} \quad \text { (by (2.4.7) }\right),
\end{aligned}
$$

which again yields a contradiction to our hypothesis (2.4.1), just as in the case $d \geq 3$.

### 2.5. A priori estimates for Coulomb gauge transformations

We now establish a generalization of [34, Lemma 6.6], which is in turn an analogue of [27, Lemma 2.3.10]. We allow $p \geq 2$ and any $d \geq 2$ rather than assume $d=4$, as in [34, Lemma 6.6], but we use standard Sobolev norms rather than the 'critical exponent' Sobolev norms employed in the statement and proof of 34, Lemma 6.6], since we do not seek an explicit optimal dependence of constants on the reference connection, $A_{0}$.

Before stating our generalization of [34, Lemma 6.6], we digress to recall from [44, p. 231], that a gauge transformation, $u \in \operatorname{Aut}(P)$, may be viewed as a section of the fiber bundle $\operatorname{Ad} P:=P \times_{\text {Ad }} G \rightarrow X$, where we denote $\operatorname{Ad}(g): G \ni h \rightarrow$ $g^{-1} h g \in G$, for all $h \in G$. With the aid of a choice of a unitary representation, $\varrho$ : $G \subset \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$, we may therefore consider $\operatorname{Ad} P$ to be a subbundle of the Hermitian vector bundle $P \times_{\varrho} \operatorname{End}_{\mathbb{C}}(\mathbb{E})$. We can alternatively replace $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$ by $\operatorname{Ad}: G \rightarrow \operatorname{Aut}(\mathfrak{g})$ and $\operatorname{End}_{\mathbb{C}}(\mathbb{E})$ by $\operatorname{End}(\mathfrak{g})$. A choice of connection $A$ on $P$ induces covariant derivatives on all associated vector bundles, such as $E=P \times \varrho \mathbb{E}$ and $P \times{ }_{\varrho} \operatorname{End}_{\mathbb{C}}(\mathbb{E})$ or $\operatorname{ad} P=P \times$ ad $\mathfrak{g}$ and $P \times_{\mathrm{Ad}} \operatorname{End}(\mathfrak{g})$. We can thus define Sobolev norms of sections of $\operatorname{Ad} P$, generalizing the construction of Freed and Uhlenbeck in [43, Appendix A]

A similar construction is described by Parker [97, Section 4], but we note that while the center of $\operatorname{Aut}(P)$ - which is given by $P \times_{\mathrm{Ad}} \operatorname{Center}(G)$ - acts trivially on $\mathscr{A}(P)$, it does not act trivially on $C^{\infty}(X ; E)$.

Proposition 2.5.1 (A priori $W^{1, p}$ estimate for $u(A)-A_{0}$ in terms of $A-A_{0}$ ). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$. Let $A_{1}$ be a $C^{\infty}$ connection on $P$, and $A_{0}$ be a $W^{1, q}$ connection on $P$ with $d / 2<q<\infty$, and $p \in(1, \infty)$ obey $d / 2 \leq p \leq q$, and $\delta \in(0,1]$. Then there are constants $N \in[1, \infty)$ and $\varepsilon=\varepsilon\left(A_{0}, A_{1}, g, G, p, q\right) \in(0,1]$ (with dependence on $p$ replaced by dependence on $\delta$ when $p=d$ ) with the following significance. If $A$ is $a W^{1, q}$ connection on $P$ and $u \in \operatorname{Aut}(P)$ is a gauge transformation of class $W^{2, q}$ such that

$$
\begin{equation*}
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0 \tag{2.5.1}
\end{equation*}
$$

then the following hold. If

$$
\begin{equation*}
\left\|A-A_{0}\right\|_{L^{s}(X)} \leq \varepsilon \quad \text { and } \quad\left\|u(A)-A_{0}\right\|_{L^{s}(X)} \leq \varepsilon \tag{2.5.2}
\end{equation*}
$$

where

$$
s(p):= \begin{cases}d & \text { if } p<d  \tag{2.5.3}\\ d+\delta & \text { if } p=d \\ p & \text { if } p>d\end{cases}
$$

then

$$
\begin{equation*}
\left\|u(A)-A_{0}\right\|_{L^{p}(X)} \leq N\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \tag{2.5.4}
\end{equation*}
$$

where $N=N\left(A_{0}, A_{1}, g, G, p, q\right)$. If in addition $A$ obey $\$^{3}$

$$
\begin{equation*}
\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq M \tag{2.5.5}
\end{equation*}
$$

[^3]for some constant $M \in[1, \infty)$, then
\[

$$
\begin{equation*}
\left\|u(A)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq N\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \tag{2.5.6}
\end{equation*}
$$

\]

where $N=N\left(A_{0}, A_{1}, g, G, M, p, q\right)$.
Proof. Following the convention of [114, p. 32] for the action of $u \in \operatorname{Aut}(P)$ on connections $A$ on $P$ and setting $B:=u(A)$ for convenience, we have

$$
\begin{equation*}
B-A_{0}=u^{-1}\left(A-A_{0}\right) u+u^{-1} d_{A_{0}} u \tag{2.5.7}
\end{equation*}
$$

Our task is thus to estimate the term $d_{A_{0}} u$. Rewriting the preceding equality gives a first-order, linear elliptic equation in $u$ with $W^{1, q}$ coefficients,

$$
\begin{equation*}
d_{A_{0}} u=u\left(B-A_{0}\right)-\left(A-A_{0}\right) u \tag{2.5.8}
\end{equation*}
$$

Corollary C. 2 implies that the kernel,

$$
K:=\operatorname{Ker}\left(\Delta_{A_{0}}: W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow L^{q}(X ; \operatorname{ad} P)\right)
$$

is finite-dimensional. Let

$$
\Pi: L^{2}(X ; \operatorname{ad} P) \rightarrow K \subset L^{2}(X ; \operatorname{ad} P)
$$

denote the $L^{2}$-orthogonal projection and denote

$$
\begin{aligned}
\gamma & :=\Pi u \in K \subset W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \\
u_{0} & :=u-\gamma \in K^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)
\end{aligned}
$$

where $\perp$ is $L^{2}$-orthogonal complement. We may assume without loss of generality that we have a unitary representation, $G \subset \mathrm{U}(n)$. Recall that (due to $\mathbf{1 1 9}$ Equation (6.2)]),

$$
\begin{equation*}
d_{A_{0}}^{*}=(-1)^{-d(k+1)+1} * d_{A_{0}} * \quad \text { on } \Omega^{k}(X ; \operatorname{ad} P) \tag{2.5.9}
\end{equation*}
$$

where $*: \Omega^{k}(X) \rightarrow \Omega^{d-k}$ is the Hodge star operator on $k$-forms. Because $d_{A_{0}}^{*}(B-$ $\left.A_{0}\right)=0$ and $d_{A_{0}} u=d_{A_{0}} u_{0}$, an application of $d_{A_{0}}^{*}$ to (2.5.8) yields

$$
\begin{aligned}
d_{A_{0}}^{*} d_{A_{0}} u_{0}= & -*\left(d_{A_{0}} u \wedge *\left(B-A_{0}\right)\right)+u d_{A_{0}}^{*}\left(B-A_{0}\right) \\
& -\left(d_{A_{0}}^{*}\left(A-A_{0}\right)\right) u+*\left(*\left(A-A_{0}\right) \wedge d_{A_{0}} u\right) \\
= & -*\left(d_{A_{0}} u_{0} \wedge *\left(B-A_{0}\right)\right)-\left(d_{A_{0}}^{*}\left(A-A_{0}\right)\right) u+*\left(*\left(A-A_{0}\right) \wedge d_{A_{0}} u_{0}\right)
\end{aligned}
$$

We shall use the bound $\|u\|_{C(X)} \leq 1$ for any $u \in \operatorname{Aut}(P)$ of class $W^{2, q}$, implied by the fact that the representation for $G$ is unitary. Recall that $\Delta_{A_{0}}=d_{A_{0}}^{*} d_{A_{0}}$. We first consider the case $p<d$ and setting $p^{*}=d p /(d-p)$, we have $1 / p=1 / d+1 / p^{*}$ and the continuous multiplication $L^{d}(X) \times L^{p^{*}}(X) \rightarrow L^{p}(X)$, so

$$
\begin{aligned}
\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)} \leq & \left\|d_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)}\left\|B-A_{0}\right\|_{L^{d}(X)}+\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}\|u\|_{C(X)} \\
& +\left\|A-A_{0}\right\|_{L^{d}(X)}\left\|d_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)} \\
\leq & \left(\left\|B-A_{0}\right\|_{L^{d}(X)}+\left\|A-A_{0}\right\|_{L^{d}(X)}\right)\left\|d_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)} \\
& +\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} .
\end{aligned}
$$

Second, for the case $p>d$, we use the continuous multiplication $L^{p}(X) \times L^{\infty}(X) \rightarrow$ $L^{p}(X)$, so

$$
\begin{aligned}
\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)} \leq & \left(\left\|B-A_{0}\right\|_{L^{p}(X)}+\left\|A-A_{0}\right\|_{L^{p}(X)}\right)\left\|d_{A_{0}} u_{0}\right\|_{L^{\infty}(X)} \\
& +\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}
\end{aligned}
$$

Third, for the case $p=d$, we can instead use $r \in(d, \infty)$ defined by $1 / d=1 /(d+$ $\delta)+1 / r$ and the resulting continuous multiplication, $L^{d+\delta}(X) \times L^{r}(X) \rightarrow L^{d}(X)$, to give

$$
\begin{aligned}
\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{d}(X)} \leq & \left(\left\|B-A_{0}\right\|_{L^{d+\delta}(X)}+\left\|A-A_{0}\right\|_{L^{d+\delta}(X)}\right)\left\|d_{A_{0}} u_{0}\right\|_{L^{r}(X)} \\
& +\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{d}(X)}
\end{aligned}
$$

By [4. Theorem 4.12], we have continuous Sobolev embeddings,

$$
W^{1, p}(X) \subset \begin{cases}L^{d p /(d-p)}(X) & \text { if } 1 \leq p<d  \tag{2.5.10}\\ L^{r}(X) & \text { if } p=d \text { and } 1 \leq r<\infty \\ C(X) & \text { if } p>d\end{cases}
$$

Therefore, the Sobolev Embedding Theorem and Kato Inequality 43, Equation (6.20)] give, for $r \in(d, \delta)$ determined by $\delta$ as above,

$$
\begin{align*}
\left\|d_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)} \leq C_{0}\left\|d_{A_{0}} u_{0}\right\|_{W_{A_{1}}^{1, p}(X)}, & p<d \\
\left\|d_{A_{0}} u_{0}\right\|_{L^{r}(X)} \leq C_{0}\left\|d_{A_{0}} u_{0}\right\|_{W_{A_{1}}^{1, d}(X)}, & p=d  \tag{2.5.11}\\
\left\|d_{A_{0}} u_{0}\right\|_{L^{\infty}(X)} \leq C_{0}\left\|d_{A_{0}} u_{0}\right\|_{W_{A_{1}}^{1, p}(X)}, & p>d
\end{align*}
$$

where $C_{0}=C_{0}(g, p)$ or $C_{0}(g, \delta) \in[1, \infty)$ is bounded below by the norm of the Sobolev embedding (2.5.10). Writing $A_{0}=A_{1}+a_{0}$, for $a_{0} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, so $d_{A_{0}} u_{0}=d_{A_{1}} u_{0}+\left[a_{0}, u_{0}\right]$, we see that

$$
\begin{aligned}
\left\|\nabla_{A_{1}} d_{A_{0}} u_{0}\right\|_{L^{p}(X)} & \leq\left\|\nabla_{A_{1}} d_{A_{1}} u_{0}\right\|_{L^{p}(X)}+\left\|\nabla_{A_{1}}\left[a_{0}, u_{0}\right]\right\|_{L^{p}(X)} \\
& \leq\left\|\nabla_{A_{1}}^{2} u_{0}\right\|_{L^{p}(X)}+\left\|\nabla_{A_{1}} a_{0} \times u_{0}+a_{0} \times \nabla_{A_{1}} u_{0}\right\|_{L^{p}(X)}
\end{aligned}
$$

By hypothesis, we have $p=d / 2<q$ or $d / 2<p \leq q$, and thus, defining $r \in[p, \infty]$ by $1 / p=1 / q+1 / r$ and using the continuous Sobolev embeddings $W^{1, p} \subset L^{2 p}(X)$ and $W^{2, p}(X) \subset L^{r}(X)$ (with norm bounded above by $C_{0}$ ) for $p \geq d / 2$,

$$
\begin{aligned}
& \left\|\nabla_{A_{1}} d_{A_{0}} u_{0}\right\|_{L^{p}(X)} \\
& \quad \leq\left\|\nabla_{A_{1}}^{2} u_{0}\right\|_{L^{p}(X)}+z\left\|\nabla_{A_{1}} a_{0}\right\|_{L^{q}(X)}\left\|u_{0}\right\|_{L^{r}(X)}+z\left\|a_{0}\right\|_{L^{2 p}(X)}\left\|\nabla_{A_{1}} u_{0}\right\|_{L^{2 p}(X)} \\
& \quad \leq\left\|\nabla_{A_{1}}^{2} u_{0}\right\|_{L^{p}(X)}+z C_{0}\left\|\nabla_{A_{1}} a_{0}\right\|_{L^{q}(X)}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)} \\
& \quad+z C_{0}^{2}\left\|a_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\left\|\nabla_{A_{1}} u_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad \leq\left\|\nabla_{A_{1}}^{2} u_{0}\right\|_{L^{p}(X)}+z C_{0}\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}+z C_{0}^{2}\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)},
\end{aligned}
$$

where $z=z(g) \in[1, \infty)$ and now $C_{0}=C_{0}\left(A_{1}, g, p, q\right) \in[1, \infty)$. By substituting the preceding bound into (2.5.11), we find that

$$
\begin{aligned}
\left\|d_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)} \leq C_{1}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}, & p<d \\
\left\|d_{A_{0}} u_{0}\right\|_{L^{r}(X)} \leq C_{1}\left\|u_{0}\right\|_{W_{A_{1}}^{2, d}(X)}, & p=d \\
\left\|d_{A_{0}} u_{0}\right\|_{L^{\infty}(X)} \leq C_{1}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}, & p>d
\end{aligned}
$$

for a constant $C_{1}=C_{1}\left(A_{1}, g, p, q,\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}\right) \in[1, \infty)$, with dependence on $p$ replaced by $\delta$ when $p=d$. By combining the preceding three cases $(p<d$, and
$p=d$, and $p>d$ ), we obtain

$$
\begin{align*}
&\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)} \leq C_{1}\left(\left\|B-A_{0}\right\|_{L^{s}(X)}+\left\|A-A_{0}\right\|_{L^{s}(X)}\right)\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}  \tag{2.5.12}\\
&+\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}
\end{align*}
$$

where $s=s(p)$ is as in (2.5.3). From the a priori estimate (2.2.4) in Corollary 2.2.5 - and noting that this lemma also holds for $\operatorname{Ad} P$ in place of $\operatorname{ad} P$ via the definition of Sobolev norms of $u \in \operatorname{Aut}(P)$ described earlier - we have the a priori estimate,

$$
\begin{equation*}
\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)} \leq C_{3}\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)}, \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q \tag{2.5.13}
\end{equation*}
$$

where $C_{3}=C_{3}\left(A_{1}, A_{0}, g, G, p, q\right) \in[1, \infty)$. Substituting the a priori estimate (2.5.13) into our $L^{p}$ bound (2.5.12) for $\Delta_{A_{0}} u_{0}$ gives, for $s$ as in (2.5.3),

$$
\begin{aligned}
\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)} \leq & C_{1} C_{3}\left(\left\|B-A_{0}\right\|_{L^{s}(X)}+\left\|A-A_{0}\right\|_{L^{s}(X)}\right)\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)} \\
& +\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}, \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
\end{aligned}
$$

Provided

$$
\begin{equation*}
\left\|B-A_{0}\right\|_{L^{s}(X)} \leq 1 /\left(4 C_{1} C_{3}\right) \quad \text { and } \quad\left\|A-A_{0}\right\|_{L^{s}(X)} \leq 1 /\left(4 C_{1} C_{3}\right) \tag{2.5.14}
\end{equation*}
$$

as assured by (2.5.2), then rearrangement in the preceding inequality yields

$$
\begin{equation*}
\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)} \leq 2\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \tag{2.5.15}
\end{equation*}
$$

$$
\text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
$$

Therefore, by combining the inequalities (2.5.15) and (2.5.13) we find that

$$
\begin{align*}
& \left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)} \leq 2 C_{3}\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}  \tag{2.5.16}\\
& \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
\end{align*}
$$

Using $d_{A_{0}} u=d_{A_{0}} u_{0}$ and (2.5.7) and the facts that $|u| \leq 1$ and $\left|u^{-1}\right| \leq 1$ on $X$, we obtain

$$
\begin{equation*}
\left\|B-A_{0}\right\|_{L^{p}(X)} \leq\left\|A-A_{0}\right\|_{L^{p}(X)}+\left\|d_{A_{0}} u_{0}\right\|_{L^{p}(X)} \tag{2.5.17}
\end{equation*}
$$

From (2.5.17) and (2.5.16), we see that

$$
\begin{aligned}
\left\|B-A_{0}\right\|_{L^{p}(X)} & \leq\left\|A-A_{0}\right\|_{L^{p}(X)}+\left\|u_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \leq\left\|A-A_{0}\right\|_{L^{p}(X)}+2 C_{3}\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}
\end{aligned}
$$

Using $A_{0}=A_{1}+a_{0}$, we have

$$
d_{A_{0}}^{*}\left(A-A_{0}\right)=d_{A_{1}}^{*}\left(A-A_{0}\right)+a_{0} \times\left(A-A_{0}\right)
$$

and

$$
\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \leq z\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}+z\left\|a_{0}\right\|_{L^{2 p}(X)}\left\|A-A_{0}\right\|_{L^{2 p}(X)}
$$

Applying the continuous Sobolev embedding, $W^{1, p}(X) \subset L^{2 p}(X)$, with norm $C_{0}=$ $C_{0}(g, p) \in[1, \infty)$ and the Kato Inequality [43, Equation (6.20)],

$$
\begin{align*}
&\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \leq z\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}  \tag{2.5.18}\\
&+z C_{0}^{2}\left\|a_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}
\end{align*}
$$

Thus,

$$
\begin{equation*}
\left\|B-A_{0}\right\|_{L^{p}(X)} \leq C_{4}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}, \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q \tag{2.5.19}
\end{equation*}
$$

where $C_{4}=C_{4}\left(A_{0}, A_{1}, g, G, p, q\right) \in[1, \infty)$, giving the desired $L^{p}$ estimate (2.5.4) for $B-A_{0}$.

We now estimate the $L^{p}$ norms of the covariant derivatives of the right-hand side of the identity (2.5.7). Considering the term $u^{-1} d_{A_{0}} u$ in the right-hand side of the identity (2.5.7) and recalling that $\nabla_{A_{0}} u=d_{A_{0}} u=d_{A_{0}} u_{0}$, we have

$$
\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)=-u^{-1}\left(\nabla_{A_{0}} u_{0}\right) u^{-1} \otimes d_{A_{0}} u_{0}+u^{-1} \nabla_{A_{0}} d_{A_{0}} u_{0}
$$

First, if $d / 2 \leq p<d$ and using the continuous multiplication, $L^{p^{*}}(X) \times L^{d}(X) \rightarrow$ $L^{p}(X)$,

$$
\begin{aligned}
\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)} & \leq\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)}\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{d}(X)}+\left\|\nabla_{A_{0}}^{2} u_{0}\right\|_{L^{p}(X)} \\
& \leq C_{0}^{2}\left\|\nabla_{A_{0}} u_{0}\right\|_{W_{A_{0}}^{1, p}(X)}^{2}+\left\|\nabla_{A_{0}}^{2} u_{0}\right\|_{L^{p}(X)} \\
& \leq C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{0}}^{2, p}(X)}^{2}+\left\|u_{0}\right\|_{W_{A_{0}}^{2, p}(X)}
\end{aligned}
$$

Second, if $p=d$ and using the continuous multiplication, $L^{2 d}(X) \times L^{2 d}(X) \rightarrow$ $L^{d}(X)$,

$$
\begin{aligned}
\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{d}(X)} & \leq\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{2 d}(X)}\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{2 d}(X)}+\left\|\nabla_{A_{0}}^{2} u_{0}\right\|_{L^{d}(X)} \\
& \leq C_{0}^{2}\left\|\nabla_{A_{0}} u_{0}\right\|_{W_{A_{0}}^{1, d}(X)}^{2}+\left\|\nabla_{A_{0}}^{2} u_{0}\right\|_{L^{d}(X)} \\
& \leq C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{0}}^{2, d}(X)}^{2}+\left\|u_{0}\right\|_{W_{A_{0}}^{2, d}(X)}
\end{aligned}
$$

Third, if $p>d$,

$$
\begin{aligned}
\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)} & \leq\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{\infty}(X)}\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{p}(X)}+\left\|\nabla_{A_{0}}^{2} u_{0}\right\|_{L^{p}(X)} \\
& \leq C_{0}^{2}\left\|\nabla_{A_{0}} u_{0}\right\|_{W_{A_{0}}^{1, p}(X)}^{2}+\left\|\nabla_{A_{0}}^{2} u_{0}\right\|_{L^{p}(X)} \\
& \leq C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{0}}^{2, p}(X)}^{2}+\left\|u_{0}\right\|_{W_{A_{0}}^{2, p}(X)}
\end{aligned}
$$

Thus, by combining the three preceding cases and applying Lemma B. 2 (5),

$$
\begin{array}{r}
\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)} \leq C_{6}^{2} C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}^{2}+C_{6}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}  \tag{2.5.20}\\
\text { for } d / 2 \leq p \leq q
\end{array}
$$

where $C_{6}=C_{6}\left(A_{0}, A_{1}, g, p, q\right) \in[1, \infty)$ is the constant in Lemma B. 2 (5). By combining the bound (2.5.20) for $\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)}$ with (2.5.16), we find that

$$
\begin{aligned}
\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)} \leq & 2 C_{6}^{2} C_{0}^{2} C_{3}\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}^{2} \\
& +2 C_{6} C_{3}\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \\
& \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
\end{aligned}
$$

But

$$
\begin{equation*}
\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \leq z\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, p}(X)} \tag{2.5.21}
\end{equation*}
$$

for a generic constant $z=z(g) \in[1, \infty)$, and by Lemma B.2 (3),

$$
\begin{equation*}
\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, p}(X)} \leq C_{7}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \tag{2.5.22}
\end{equation*}
$$

where $C_{7}=C_{7}\left(A_{0}, A_{1}, g, p\right) \in[1, \infty)$ is the constant in Lemma (3), and because $A$ is now assumed to obey (2.5.5), that is,

$$
\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq M
$$

we obtain

$$
\begin{array}{r}
\left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)} \leq 2 C_{6} C_{3}\left(z C_{6} C_{7} C_{0}^{2} M+1\right)\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \\
\text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
\end{array}
$$

and thus by (2.5.21) and (2.5.22),

$$
\begin{align*}
& \left\|\nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right)\right\|_{L^{p}(X)}  \tag{2.5.23}\\
& \quad \leq 2 z^{2} C_{6} C_{7} C_{3}\left(z C_{6} C_{7} C_{0}^{2} M+1\right)\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
\end{align*}
$$

Considering the term $u^{-1}\left(A-A_{0}\right) u$ in the right-hand side of (2.5.7), we discover that

$$
\begin{aligned}
\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)= & -u^{-1}\left(\nabla_{A_{0}} u\right) u^{-1} \otimes\left(A-A_{0}\right) u+u^{-1}\left(\nabla_{A_{0}}\left(A-A_{0}\right)\right) u \\
& +u^{-1}\left(A-A_{0}\right) \otimes \nabla_{A_{0}} u
\end{aligned}
$$

Noting that $\nabla_{A_{0}} u=\nabla_{A_{0}} u_{0}$ and $\|u\|_{C(X)} \leq 1$, the preceding identity gives, for $d / 2 \leq p<d$,

$$
\begin{aligned}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{p^{*}}(X)}\left\|A-A_{0}\right\|_{L^{d}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2 C_{0}^{2}\left\|\nabla_{A_{0}} u_{0}\right\|_{W_{A_{0}}^{1, p}(X)}\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, p}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2 C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{0}}^{2, p}(X)}\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, p}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)} .
\end{aligned}
$$

Second, for the case $p=d$,

$$
\begin{aligned}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{d}(X)} \\
& \quad \leq 2\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{2 d}(X)}\left\|A-A_{0}\right\|_{L^{2 d}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{d}(X)} \\
& \quad \leq 2 C_{0}^{2}\left\|\nabla_{A_{0}} u_{0}\right\|_{W_{A_{0}}^{1, d}(X)}\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, d}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{d}(X)} \\
& \quad \leq 2 C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{0}}^{2, p}(X)}\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, d}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{d}(X)}
\end{aligned}
$$

Third, for the case $d<p<\infty$,

$$
\begin{aligned}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2\left\|\nabla_{A_{0}} u_{0}\right\|_{L^{\infty}(X)}\left\|A-A_{0}\right\|_{L^{p}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2 C_{0}\left\|\nabla_{A_{0}} u_{0}\right\|_{W_{A_{0}}^{1, p}(X)}\left\|A-A_{0}\right\|_{L^{p}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2 C_{0}\left\|u_{0}\right\|_{W_{A_{0}^{2}}^{2, p}(X)}\left\|A-A_{0}\right\|_{L^{p}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)}
\end{aligned}
$$

Hence, the combination of the preceding three cases gives

$$
\begin{aligned}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{p}(X)} \\
& \quad \leq 2 C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{0}^{2}(X)}^{2, p}}\left\|A-A_{0}\right\|_{W_{A_{0}}^{1, p}(X)}+\left\|\nabla_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{p}(X)}
\end{aligned}
$$

$$
\text { for } p<\infty \text { and } d / 2 \leq p \leq q
$$

Applying Lemma B.2. Items (3) and (5), yields

$$
\begin{align*}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{p}(X)}  \tag{2.5.24}\\
& \leq 2 C_{6} C_{7} C_{0}^{2}\left\|u_{0}\right\|_{W_{A_{1}}^{2, p}(X)}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad+C_{7}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}, \quad \text { for } p<\infty \text { and } d / 2 \leq p \leq q .
\end{align*}
$$

Therefore, combining the inequalities (2.5.16) and (2.5.24) yields

$$
\begin{align*}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{p}(X)}  \tag{2.5.25}\\
& \quad \leq 4 C_{6} C_{7} C_{0}^{2} C_{3}\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& +C_{7}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q .
\end{align*}
$$

From (2.5.5) and (2.5.21), we see that (2.5.25) simplifies to give

$$
\begin{align*}
& \left\|\nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\right)\right\|_{L^{p}(X)}  \tag{2.5.26}\\
& \leq\left(4 z C_{6} C_{7} C_{0}^{2} C_{3} M+1\right) C_{7}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q
\end{align*}
$$

From the identity (2.5.7) (noting again that $d_{A_{0}} u=d_{A_{0}} u_{0}$ ) we have

$$
\left\|\nabla_{A_{0}}\left(B-A_{0}\right)\right\|_{L^{p}(X)} \leq \| \nabla_{A_{0}}\left(u^{-1}\left(A-A_{0}\right) u\left\|_{L^{p}(X)}+\right\| \nabla_{A_{0}}\left(u^{-1} d_{A_{0}} u_{0}\right) \|_{L^{p}(X)}\right.
$$

Combining the preceding estimate with the inequalities (2.5.23) and (2.5.26) gives

$$
\begin{align*}
&\left\|\nabla_{A_{0}}\left(B-A_{0}\right)\right\|_{L^{p}(X)} \leq C_{5}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}  \tag{2.5.27}\\
& \text { for } 1<p<\infty \text { and } d / 2 \leq p \leq q,
\end{align*}
$$

with a constant $C_{5}=C_{5}\left(A_{0}, A_{1}, g, G, M, p, q\right) \in[1, \infty)$. Finally, from (2.5.19) and (2.5.27) and Lemma B.2 (3) we obtain the desired $W_{A_{1}}^{1, p}$ bound (2.5.6) for $u(A)-A_{0}$ in terms of $A-A_{0}$, recalling that $B=u(A)$, with large enough constant $N=N\left(A_{0}, A_{1}, g, G, M, p, q\right) \in[1, \infty)$ and under the hypothesis (2.5.2) with small enough constant $\varepsilon=\varepsilon\left(A_{0}, A_{1}, g, G, p, q\right) \in(0,1]$.

The proof of Proposition 2.5.1 also yields the following useful
Lemma 2.5.2 (A priori $W^{2, p}$ estimate for a $W^{2, q}$ gauge transformation $u$ intertwining two $W^{1, q}$ connections). Assume the hypotheses of Proposition 2.5.1, excluding those on the connection $A$. Then there is a constant $C=C\left(A_{0}, A_{1}, g, G, p, q\right) \in$ $[1, \infty)$ with the following significance. If $A$ obeys the hypotheses of Proposition 2.5.1 and $u \in \operatorname{Aut}^{2, q}(P)$ is the resulting gauge transformation, depending on $A$ and $A_{0}$, such that

$$
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0
$$

and furthermore

$$
\begin{equation*}
\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon \tag{2.5.28}
\end{equation*}
$$

then

$$
\|u\|_{W_{A_{1}}^{2, p}(X)} \leq C
$$

Proof. Write $u=u_{0}+\gamma$ as in the proof of Proposition 2.5.1 with $u_{0} \in$ $\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp}$ and $\gamma \in \operatorname{Ker} \Delta_{A_{0}}$, and observe that

$$
\begin{aligned}
\|u\|_{W_{A_{1}}^{2, p}(X)} & \leq C\left(\left\|\Delta_{A_{0}} u\right\|_{L^{p}(X)}+\|u\|_{L^{p}(X)}\right) \quad(\text { by Proposition 2.2.1) } \\
& =C\left(\left\|\Delta_{A_{0}} u_{0}\right\|_{L^{p}(X)}+\|u\|_{L^{p}(X)}\right) \\
& \leq C\left(\left\|d_{A_{0}}^{*}\left(A-A_{0}\right)\right\|_{L^{p}(X)}+\|u\|_{L^{p}(X)}\right) \quad(\text { by }(2.5 .15)) \\
& \leq C\left(1+\left\|A_{0}-A_{1}\right\|_{W_{A_{1}}^{1, p}(X)}\right)\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}+C\|u\|_{L^{p}(X)} \quad(\text { by (2.5.18) }) \\
& \leq C\left(1+\left\|A_{0}-A_{1}\right\|_{W_{A_{1}}^{1, p}(X)}\right) \varepsilon+C \operatorname{Vol}_{g}(X)^{1 / p},
\end{aligned}
$$

where the last inequality follows from (2.5.28) and the fact that $|u| \leq 1$ pointwise. This completes the proof.

### 2.6. Existence of Coulomb gauge transformations for connections that are $W^{1, d / 2}$-close to a reference connection

Finally, we can proceed to the
Proof of Theorem [14. We shall apply the method of continuity, modeled on the proofs of [114, Theorem 2.1] due to Uhlenbeck and [27, Proposition 2.3.13] due to Donaldson and Kronheimer. For a related application of the method of continuity, see the proof [34, Theorem 1.1] due to the first author.

We begin by defining a one-parameter family of $W^{1, q}$ connections by setting

$$
\begin{equation*}
A_{t}:=A_{0}+t\left(A-A_{0}\right), \quad \forall t \in[0,1] \tag{2.6.1}
\end{equation*}
$$

and observe that their curvatures are given by

$$
F\left(A_{t}\right)=F_{A_{0}}+t d_{A_{0}}\left(A-A_{0}\right)+\frac{t^{2}}{2}\left[A-A_{0}, A-A_{0}\right]
$$

and they obey the bounds,

$$
\begin{aligned}
\left\|F\left(A_{t}\right)\right\|_{L^{q}(X)} & \leq\left\|F_{A_{0}}\right\|_{L^{q}(X)}+\left\|d_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{q}(X)}+\frac{c_{0}}{2}\left\|A-A_{0}\right\|_{L^{2 q}(X)}^{2} \\
& \leq\left\|F_{A_{0}}\right\|_{L^{q}(X)}+\left\|d_{A_{0}}\left(A-A_{0}\right)\right\|_{L^{q}(X)}+c\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}^{2}
\end{aligned}
$$

with $c_{0} \in[1, \infty)$ and $c=c(g, q) \in[1, \infty)$ and where we use the Sobolev embedding, $W^{1, q}(X) \subset L^{2 q}(X)$ 4. Theorem 4.12] and the Kato Inequality 43, Equation $(6.20)]$ to obtain the last inequality. Note that we have a continuous embedding, $W^{1, q}(X) \subset L^{2 q}(X)$, when $\left.a\right) q<d$ and $2 q \leq q^{*}:=d q /(d-q)$, that is $2 d-2 q \leq d$ or $q \geq d / 2$, as implied by our hypotheses, or b) $q \geq d$. Therefore,

$$
\begin{equation*}
\left\|F\left(A_{t}\right)\right\|_{L^{q}(X)} \leq K, \quad \forall t \in[0,1] \tag{2.6.2}
\end{equation*}
$$

for $K=K\left(A, A_{0}, G, g, q\right) \in[1, \infty)$, noting that $F_{A_{0}} \in L^{q}\left(X ; \Lambda^{2} \otimes \operatorname{ad} P\right)$ since $A_{0}$ is of class $W^{1, q}$.

Let $S$ denote the set of $t \in[0,1]$ such that there exists a $W^{2, q}$ gauge transformation $u_{t} \in \operatorname{Aut}(P)$ with the property that

$$
d_{A_{0}}^{*}\left(u_{t}\left(A_{t}\right)-A_{0}\right)=0 \quad \text { and } \quad\left\|u_{t}\left(A_{t}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|A_{t}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

where $N$ is the constant in Proposition 2.5.1. Clearly, $0 \in S$ since the identity automorphism of $P$ is the required gauge transformation in that case, so $S$ is nonempty. As usual, we need to show that $S$ is an open and closed subset of $[0,1]$.

Step 1 ( $S$ is open). To prove openness, we shall adapt the argument of Donaldson and Kronheimer in [27, Section 2.3.8]. We apply the Implicit Function Theorem to the gauge fixing equation,

$$
d_{A_{0}}^{*}\left(u_{t}\left(A_{t}\right)-A_{0}\right)=d_{A_{0}}^{*}\left(u_{t}^{-1}\left(A_{t}-A_{0}\right) u_{t}+u_{t}^{-1} d_{A_{0}} u_{t}\right)=0
$$

As usual, we denote $B_{t}=u_{t}\left(A_{t}\right)$ for convenience, for $t \in S$. Let $t_{0} \in S$, so we have $d_{A_{0}}^{*}\left(u_{t_{0}}\left(A_{t_{0}}\right)-A_{0}\right)=0 \quad$ and $\quad\left\|u_{t_{0}}\left(A_{t_{0}}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|A_{t_{0}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}$.
Our task is to show that $t_{0}+s \in S$ for $|s|$ sufficiently small, that is, there exists $u_{t_{0}+s} \in \operatorname{Aut}^{2, q}(P)$ such that the preceding two properties hold with $t_{0}$ replaced by $t_{0}+s$. By the hypothesis (1.5.1) of Theorem 14 we have $\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta$ and thus, since

$$
A_{t_{0}}-A_{0}=A_{0}+t_{0}\left(A-A_{0}\right)-A_{0}=t_{0}\left(A-A_{0}\right)
$$

and $t_{0} \in[0,1]$, we see that

$$
\left\|A_{t_{0}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta
$$

and so

$$
\left\|u_{t_{0}}\left(A_{t_{0}}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta
$$

It will be convenient to define $a \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ by

$$
\begin{equation*}
u_{t_{0}}\left(A_{t_{0}}\right)=: A_{0}+a, \tag{2.6.3}
\end{equation*}
$$

The preceding inequality ensures that

$$
\begin{equation*}
\|a\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta . \tag{2.6.4}
\end{equation*}
$$

We shall seek a solution $u_{t_{0}+s} \in \operatorname{Aut}^{2, q}(P)$ to the gauge-fixing equation,

$$
d_{A_{0}}^{*}\left(u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right)=0
$$

In particular, we shall seek a solution in the form

$$
u_{t_{0}+s}=e^{\chi_{s}} u_{t_{0}}, \quad \text { for } \chi_{s} \in W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)
$$

so the gauge-fixing equation becomes

$$
\begin{equation*}
d_{A_{0}}^{*}\left(e^{\chi_{s}} u_{t_{0}}\left(A_{t_{0}+s}\right)-A_{0}\right)=0 \tag{2.6.5}
\end{equation*}
$$

For $s \in \mathbb{R}$, it will be convenient to define $b_{s} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ by

$$
\begin{equation*}
u_{t_{0}}\left(A_{t_{0}+s}\right)=: A_{0}+a+b_{s} \tag{2.6.6}
\end{equation*}
$$

We can determine $b_{s}$ explicitly using $A_{t_{0}+s}=A_{0}+\left(t_{0}+s\right)\left(A-A_{0}\right)$, so that

$$
\begin{aligned}
u_{t_{0}}\left(A_{t_{0}+s}\right)-A_{0} & =u_{t_{0}}^{-1}\left(\left(t_{0}+s\right)\left(A-A_{0}\right)\right) u_{t_{0}}+u_{t_{0}}^{-1} d_{A_{0}} u_{t_{0}} \\
& =u_{t_{0}}^{-1}\left(t_{0}\left(A-A_{0}\right)\right) u_{t_{0}}+u_{t_{0}}^{-1} d_{A_{0}} u_{t_{0}}+u_{t_{0}}^{-1}\left(s\left(A-A_{0}\right)\right) u_{t_{0}} \\
& =u_{t_{0}}^{-1}\left(A_{t_{0}}-A\right) u_{t_{0}}+u_{t_{0}}^{-1} d_{A_{0}} u_{t_{0}}+u_{t_{0}}^{-1}\left(s\left(A-A_{0}\right)\right) u_{t_{0}} \\
& =u_{t_{0}}\left(A_{t_{0}}\right)-A_{0}+s u_{t_{0}}^{-1}\left(A-A_{0}\right) u_{t_{0}} \\
& =a+s u_{t_{0}}^{-1}\left(A-A_{0}\right) u_{t_{0}} \quad(\text { by (2.6.3) }),
\end{aligned}
$$

and thus

$$
\begin{equation*}
b_{s}=s u_{t_{0}}^{-1}\left(A-A_{0}\right) u_{t_{0}}, \quad s \in \mathbb{R} \tag{2.6.7}
\end{equation*}
$$

Note that $u_{t_{0}} \in \operatorname{Aut}^{2, q}(P)$ and so we have the estimate,

$$
\begin{equation*}
\left\|b_{s}\right\|_{W_{A_{1}^{1}(X)}^{1, q}(X)} \leq|s| C_{0}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}, \quad s \in \mathbb{R} \tag{2.6.8}
\end{equation*}
$$

for $C_{0}=C_{0}\left(A_{0}, g, q, t_{0}\right) \in[1, \infty)$. In particular, $b_{s} \rightarrow 0$ in $W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ strongly as $s \rightarrow 0$.

The gauge fixing equation (2.6.5) takes the form

$$
\begin{aligned}
e^{\chi_{s}} u_{t_{0}}\left(A_{t_{0}+s}\right)-A_{0} & =e^{\chi_{s}}\left(A_{0}+a+b_{s}\right)-A_{0} \\
& =e^{-\chi_{s}}\left(a+b_{s}\right) e^{\chi_{s}}+e^{-\chi_{s}} d_{A_{0}}\left(e^{\chi_{s}}\right) .
\end{aligned}
$$

The equation to be solved is then $H\left(\chi_{s}, b_{s}\right)=0$, where

$$
\begin{equation*}
H(\chi, b):=d_{A_{0}}^{*}\left(e^{-\chi}(a+b) e^{\chi}+e^{-\chi} d_{A_{0}}\left(e^{\chi}\right)\right) . \tag{2.6.9}
\end{equation*}
$$

For any $q>d / 2$, the expression (2.6.9) for $H$ defines a smooth map,

$$
\begin{align*}
H:\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \times W_{A_{1}}^{1, q} & \left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)  \tag{2.6.10}\\
& \rightarrow\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)
\end{align*}
$$

Here, we note that if $\xi=d_{A_{0}}^{*} a_{\xi}$ for some $a_{\xi} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then $\xi \perp \operatorname{Ker} \Delta_{A_{0}}$, as implied by the preceding expression for $H$. Indeed, for any $\gamma \in \operatorname{Ker} \Delta_{A_{0}}$,

$$
(\xi, \gamma)_{L^{2}(X)}=\left(d_{A_{0}}^{*} a_{\xi}, \gamma\right)_{L^{2}(X)}=\left(a_{\xi}, d_{A_{0}} \gamma\right)_{L^{2}(X)}=0
$$

Hence, the image of $H$ is contained in $\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)$.
The Implicit Function Theorem asserts that if the partial derivative,

$$
\left(D_{1} H\right)_{(0,0)}:\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)
$$

is surjective, then for small $b \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ there is a small solution $\chi \in$ $W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$ to $H(\chi, b)=0$ and that is $L^{2}$-orthogonal to $\operatorname{Ker} \Delta_{A_{0}}$.

Now the linearization, $\left(D_{1} H\right)_{(0,0)}$, of the map $H$ at the origin $(0,0)$ with respect to variations in $\chi$ is given by

$$
\left(D_{1} H\right)_{(0,0)} \chi=d_{A_{0}}^{*} d_{A_{0}+a} \chi
$$

But $\|a\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta$ by (2.6.4) and because of the continuous Sobolev embeddings provided by [4. Theorem 4.12],

$$
\begin{gathered}
W^{1, p}(X) \subset L^{d}(X), \quad \text { for } d \geq 3 \text { and } p \geq d / 2 \\
W^{1, p}(X) \subset L^{4}(X), \quad \text { for } d=2 \text { and } p \geq 2
\end{gathered}
$$

we obtain,

$$
\|a\|_{L^{d}(X)}<2 C_{1} N \zeta \quad \text { when } d \geq 3 \quad \text { and } \quad\|a\|_{L^{4}(X)}<2 C_{1} N \zeta \quad \text { when } d=2
$$

where $C_{1}=C_{1}(g, p) \in[1, \infty)$ is the norm of the Sobolev embedding employed. By the hypothesis of Theorem [14, we can choose $\zeta \in(0,1]$ as small as desired. Hence, the operator $d_{A_{0}}^{*} d_{A_{0}+a}$ is surjective by Lemma 2.4.1.

To summarize, we have shown that if $|s|$ is small and $t_{0} \in S$, then there exists $u_{t_{0}+s} \in \operatorname{Aut}^{2, q}(P)$ such that

$$
d_{A_{0}}^{*}\left(u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right)=0
$$

It remains to check that the following norm condition holds,

$$
\begin{equation*}
\left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|A_{t_{0}+s}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}, \tag{2.6.11}
\end{equation*}
$$

for small enough $|s|$, to conclude that $t_{0}+s \in S$. To see this, we first note that since $A_{t_{0}+s}=A_{0}+\left(t_{0}+s\right)\left(A-A_{0}\right)$, we have

$$
\left\|A_{t_{0}+s}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}=\left(t_{0}+s\right)\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<\left(t_{0}+s\right) \zeta \quad(\text { by (1.5.1) })
$$

and thus, for $t_{0}+s \leq 1$ and $\zeta \leq \varepsilon / C_{1}$,

$$
\left\|A_{t_{0}+s}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon / C_{1}
$$

and thus,

$$
\begin{aligned}
\left\|A_{t_{0}+s}-A_{0}\right\|_{L^{d}(X)} \leq \varepsilon, & \text { if } p<d \\
\left\|A_{t_{0}+s}-A_{0}\right\|_{L^{d+\delta}(X)} \leq \varepsilon, & \text { if } p=d \\
\left\|A_{t_{0}+s}-A_{0}\right\|_{L^{p}(X)} \leq \varepsilon, & \text { if } p>d
\end{aligned}
$$

where $\varepsilon$ is the constant in Proposition 2.5.1 and $C_{1}=C_{1}(g, p)$ or $C_{1}(\delta, p) \in[1, \infty)$ is the norm (provided by [4, Theorem 4.12]) of the continuous Sobolev embedding, $W^{1, p}(X) \subset L^{d}(X)$ when $d / 2 \leq p<d$ and $W^{1, p}(X) \subset L^{d+\delta}(X)$ when $p=d$. This verifies the hypotheses (2.5.2) and (2.5.5) of Proposition 2.5.1 for $A_{t_{0}+s}-A_{0}$ (in place of $A-A_{0}$ in the statement of that proposition).

On the other hand,

$$
\begin{aligned}
& \left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad=\left\|e^{\chi_{s}} u_{t_{0}}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad \leq\left\|e^{\chi_{s}} u_{t_{0}}\left(A_{t_{0}+s}\right)-u_{t_{0}}\left(A_{t_{0}+s}\right)\right\|_{W_{A_{1}}^{1, p}(X)}+\left\|u_{t_{0}}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \quad=\left\|e^{\chi_{s}} u_{t_{0}}\left(A_{t_{0}+s}\right)-u_{t_{0}}\left(A_{t_{0}+s}\right)\right\|_{W_{A_{1}}^{1, p}(X)}+\left\|a+b_{s}\right\|_{W_{A_{1}}^{1, p}(X)} \quad(\text { by (2.6.6) })
\end{aligned}
$$

The inequalities (1.5.1), (2.6.4), and (2.6.8) (which also holds with $p$ in place of $q$ ) yield the bound

$$
\begin{aligned}
\left\|a+b_{s}\right\|_{W_{A_{1}}^{1, p}(X)} & \leq\|a\|_{W_{A_{1}}^{1, p}(X)}+\left\|b_{s}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& <2 N \zeta+|s| C_{0}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& \leq 2 N \zeta+|s| C_{0} \zeta \\
& \leq \varepsilon /\left(2 C_{1}\right)
\end{aligned}
$$

for small enough $\zeta$. (We could also have used $\left\|b_{s}\right\|_{W_{A_{1}}^{1, q}(X)} \leq C_{0}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}$ and the continuous embedding $W^{1, q}(X) \subset W^{1, p}(X)$ and rely on our freedom to also choose $|s|$ small.) But if $|s|$ is small then so is $\left\|b_{s}\right\|_{W_{A}^{1, q}(X)}$ by (2.6.8) and hence $\left\|\chi_{s}\right\|_{W_{A_{1}}^{2, q}(X)}$ is small by the Implicit Function Theorem and so we may assume that

$$
\left\|e^{\chi_{s}} u_{t_{0}}\left(A_{t_{0}+s}\right)-u_{t_{0}}\left(A_{t_{0}+s}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon /\left(2 C_{1}\right)
$$

for small enough $|s|$. Collecting the preceding inequalities gives

$$
\left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon / C_{1}
$$

[^4]and thus,
\[

$$
\begin{aligned}
\left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{L^{d}(X)} \leq \varepsilon, & \text { if } p<d \\
\left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{L^{d+\delta}(X)} \leq \varepsilon, & \text { if } p=d \\
\left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{L^{p}(X)} \leq \varepsilon, & \text { if } p>d,
\end{aligned}
$$
\]

verifying the hypotheses (2.5.2) and (2.5.5) of Proposition 2.5.1 for $u_{t_{0}+s}\left(A_{t_{0}+s}\right)-$ $A_{0}$ (in place of $u(A)-A_{0}$ in the statement of that proposition).

Hence, Proposition 2.5.1 yields the bound

$$
\begin{aligned}
\left\|u_{t_{0}+s}\left(A_{t_{0}+s}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} & \leq N\left\|A_{t_{0}+s}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& <2 N\left\|A_{t_{0}+s}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}
\end{aligned}
$$

This verifies the norm condition (2.6.11) and we conclude that $t_{0}+s \in S$ and so $S$ is open.

Step 2 ( $S$ is closed). For closedness, we adapt the argument in [27, Section 2.3.7]. Set $B_{t}:=u_{t}\left(A_{t}\right)$ for $t \in S$ and observe that the inequality (2.6.2) yields

$$
\begin{aligned}
& \left\|F\left(B_{t}\right)\right\|_{L^{q}(X)}=\left\|F\left(u_{t}\left(A_{t}\right)\right)\right\|_{L^{q}(X)}=\left\|u_{t}\left(F\left(A_{t}\right)\right)\right\|_{L^{q}(X)} \\
& =\left\|F\left(A_{t}\right)\right\|_{L^{q}(X)} \leq K, \quad \forall t \in S
\end{aligned}
$$

Let $\left\{t_{m}\right\}_{m \in \mathbb{N}} \subset S$ be a sequence and suppose that $t_{m} \rightarrow t_{\infty} \in[0,1]$ as $m \rightarrow \infty$. Since $q>d / 2$ by hypothesis, the Uhlenbeck Weak Compactness [114, Theorem $1.5=3.6$ ] (see also [120, Theorem 7.1] for a recent exposition) implies that there exists a subsequence $\left\{m^{\prime}\right\} \subset\{m\}$ and, after relabeling, a sequence of $W^{2, q}$ gauge transformations, $\left\{u_{t_{m}}\right\}_{m \in \mathbb{N}} \subset \operatorname{Aut}(P)$ and a $W^{1, q}$ connection $B_{\infty}$ on $P$ such that, as $m \rightarrow \infty$, we have

$$
B_{t_{m}} \rightharpoonup B_{\infty} \quad \text { weakly in } W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

Hence, for a) $q<d$ and $r<q *:=d q /(d-q)$, or b) $q \geq d$ and $r<\infty$, the RellichKondrachov Theorem 4. Theorem 6.3] implies that there is a compact embedding of Sobolev spaces, $W^{1, q}(X) \Subset L^{r}(X)$, and hence there exists a subsequence $\left\{m^{\prime \prime}\right\} \subset$ $\{m\}$ such that, after again relabeling, as $m \rightarrow \infty$ we have

$$
B_{t_{m}} \rightarrow B_{\infty} \quad \text { strongly in } L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

The $W^{2, q}$ gauge transformations $u_{t}$ intertwine the $W^{1, q}$ connections $A_{t}$ and $B_{t}$ via the relation (2.5.7),

$$
B_{t}=A_{0}+t u_{t}^{-1}\left(A-A_{0}\right) u_{t}+u_{t}^{-1} d_{A_{0}} u_{t},
$$

and thus

$$
d_{A_{0}} u_{t}=u_{t} B_{t}-u_{t} A_{0}-t\left(A-A_{0}\right) u_{t}, \quad \forall t \in[0,1] .
$$

Because $A_{t_{m}} \rightarrow A_{\infty}:=A_{0}+t_{\infty}\left(A-A_{0}\right)$ strongly in $W^{1, q}$ (see 2.6.1) ) and $B_{t_{m}} \rightharpoonup$ $B_{\infty}$ weakly in $W^{1, q}$ and $B_{t_{m}} \rightarrow B_{\infty}$ strongly in $L^{r}$ as $m \rightarrow \infty$, there exists a $W^{2, q}$ gauge transformation $u_{\infty} \in \operatorname{Aut}(P)$ such that, as $m \rightarrow \infty$,

$$
\begin{aligned}
u_{t_{m}} \rightharpoonup u_{\infty} \quad \text { weakly in } W_{A_{1}}^{2, q}(X ; \operatorname{Ad} P) & \text { and } \\
& u_{t_{m}} \rightarrow u_{\infty} \quad \text { strongly in } W_{A_{1}}^{1, r}(X ; \operatorname{Ad} P) .
\end{aligned}
$$

In particular,

$$
B_{\infty}=u_{\infty}\left(A_{\infty}\right) \quad \text { and } \quad d_{A_{0}}^{*}\left(u_{\infty}\left(A_{\infty}\right)-A_{0}\right)=0
$$

The Coulomb gauge condition follows from the fact that, for any $\xi \in C^{\infty}(X ; \operatorname{ad} P)$, we have

$$
\begin{aligned}
0 & =\lim _{m \rightarrow \infty}\left(d_{A_{0}}^{*}\left(u_{t_{m}}\left(A_{t_{m}}\right)-A_{0}\right), \xi\right)_{L^{2}(X)} \\
& =\lim _{m \rightarrow \infty}\left(u_{t_{m}}\left(A_{t_{m}}\right)-A_{0}, d_{A_{0}} \xi\right)_{L^{2}(X)} \\
& =\left(u_{\infty}\left(A_{\infty}\right)-A_{0}, d_{A_{0}} \xi\right)_{L^{2}(X)} \\
& =\left(d_{A_{0}}^{*}\left(u_{\infty}\left(A_{\infty}\right)-A_{0}\right), \xi\right)_{L^{2}(X)} .
\end{aligned}
$$

Similarly, for any $a \in \Omega^{1}(X ; \operatorname{ad} P)$,

$$
\begin{aligned}
\lim _{m \rightarrow \infty}\left(B_{t_{m}}-A_{0}, a\right)_{L^{2}(X)} & =\lim _{m \rightarrow \infty}\left(u_{t_{m}}\left(A_{t_{m}}\right)-A_{0}, a\right)_{L^{2}(X)} \\
& =\lim _{m \rightarrow \infty}\left(u_{t_{m}}^{-1}\left(A_{t_{m}}-A_{0}\right) u_{t_{m}}+u_{t_{m}}^{-1} d_{A_{0}} u_{t_{m}}, a\right)_{L^{2}(X)} \\
& =\lim _{m \rightarrow \infty}\left(\left(A_{t_{m}}-A_{0}\right) u_{t_{m}}+d_{A_{0}} u_{t_{m}}, u_{t_{m}} a\right)_{L^{2}(X)} \\
& =\left(\left(A_{\infty}-A_{0}\right) u_{\infty}+d_{A_{0}} u_{\infty}, u_{\infty} a\right)_{L^{2}(X)} \\
& =\left(u_{\infty}^{-1}\left(A_{\infty}-A_{0}\right) u_{\infty}+u_{\infty}^{-1} d_{A_{0}} u_{\infty}, a\right)_{L^{2}(X)} \\
& =\left(u_{\infty}\left(A_{\infty}\right)-A_{0}, a\right)_{L^{2}(X)} .
\end{aligned}
$$

We now wish to apply Proposition 2.5.1 to bound $\left\|u_{t_{\infty}}\left(A_{t_{\infty}}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}$ and establish the remaining norm condition required to show that $t_{\infty} \in S$. First, we note that

$$
\begin{aligned}
\left\|A_{t_{m}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} & =\left\|A_{0}+t_{m}\left(A-A_{0}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& =t_{m}\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \\
& <t_{m} \zeta \quad(\text { by hypothesis (1.5.1) }) \\
& \leq \zeta, \quad \forall m \in \mathbb{N} .
\end{aligned}
$$

Because $t_{m} \in S$, we have

$$
\left\|u_{t_{m}}\left(A_{t_{m}}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|A_{t_{m}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}, \quad \forall m \in \mathbb{N}
$$

and combining this inequality with the preceding inequality yields

$$
\left\|u_{t_{m}}\left(A_{t_{m}}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta, \quad \forall m \in \mathbb{N}
$$

Since $u_{\infty}\left(A_{\infty}\right)$ is the weak limit of $u_{t_{m}}\left(A_{t_{m}}\right)$ in $W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, we have (see [31, Appendix D.4])

$$
\left\|u_{\infty}\left(A_{\infty}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq \liminf _{m \rightarrow \infty} 2 N\left\|A_{t_{m}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

But $A_{t_{m}} \rightarrow A_{t_{\infty}}$ strongly in $W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \mathrm{ad} P\right)$ by construction (2.6.1) of the path $A_{t}$ and as $p \leq q$ by hypothesis and $\left\|A_{t_{m}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta$ for all $m \in \mathbb{N}$, then

$$
\left\|A_{t_{\infty}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}=\lim _{m \rightarrow \infty}\left\|A_{t_{m}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq \zeta
$$

Combining the preceding inequalities yields

$$
\left\|u_{\infty}\left(A_{\infty}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq 2 N \zeta
$$

We choose $\zeta \in(0,1]$ small enough that $\zeta \leq \varepsilon / C_{1}$ and $2 N \zeta \leq \varepsilon / C_{1}$, where $C_{1}$ is the norm of the Sobolev embedding $W^{1, p}(X) \subset L^{d}(X)$ when $\bar{p} \neq d$ or $L^{d+\delta}(X)$ when $p=d$ as in Step and then observe that

$$
\left\|A_{t_{\infty}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon / C_{1} \quad \text { and } \quad\left\|u_{\infty}\left(A_{\infty}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<\varepsilon / C_{1}
$$

The first inequality above verifies the hypothesis (2.5.5) of Proposition 2.5.1. Moreover,

$$
\begin{aligned}
\left\|A_{t_{\infty}}-A_{0}\right\|_{L^{d}(X)} \leq \varepsilon & \text { and } \quad
\end{aligned} \quad\left\|u_{\infty}\left(A_{t_{\infty}}\right)-A_{0}\right\|_{L^{d}(X)} \leq \varepsilon, \quad \text { if } p<d, \quad \text { and } \quad\left\|u_{\infty}\left(A_{t_{\infty}}\right)-A_{0}\right\|_{L^{d+\delta}(X)} \leq \varepsilon, \quad \text { if } p=d, ~=A_{t_{\infty}} \|_{L^{d+\delta}(X)} \leq \varepsilon \quad \text { and } \quad\left\|u_{\infty}\left(A_{t_{\infty}}\right)-A_{0}\right\|_{L^{p}(X)} \leq \varepsilon, \quad \text { if } p>d,
$$

which verifies the hypothesis (2.5.2) of Proposition 2.5.1 on norms (for $A_{t_{\infty}}-A_{0}$ in place of $A-A_{0}$ and $u_{\infty}\left(A_{t_{\infty}}\right)-A_{0}$ in place of $u(A)-A_{0}$ in the statement of that proposition). Since $d_{A_{0}}^{*}\left(u_{\infty}\left(A_{\infty}\right)-A_{0}\right)=0$, as required by (2.5.1), Proposition 2.5.1 implies that

$$
\left\|u_{\infty}\left(A_{\infty}\right)-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)} \leq N\left\|A_{t_{\infty}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|A_{t_{\infty}}-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

Thus, $t_{\infty} \in S$ and so $S$ is closed.
Consequently, $S \subset[0,1]$ is non-empty and open and closed by the preceding two steps, so $S=[0,1]$ and this completes the proof of Theorem 14

### 2.7. Existence of Coulomb gauge transformations for connections that are $L^{r}$-close to a reference connection

Rather than modify the proof of Theorem (14, we shall prove Theorem 15 directly for the case $r>d$ using the Quantitaive Implicit Function Theorem F. 1 for smooth maps on Banach spaces. While the resulting proof is much easier than that of Theorem 14 and avoids the Method of Continuity, it does not obviously yield the full borderline case $r=d$ involving $W^{1, d}$ gauge transformations. See Remark 1.5.1 for related explanations. The issues here are well-known, namely, that $W^{1, d}$ gauge transformations need not be continuous and the map $H$ in (2.6.9) need not be smooth. While there versions of the Implicit Function Theorem for Lipschitz functions on Euclidean space, it is unclear whether they hold on Banach spaces: see Clarke [24 for a well-known example, Papi 96 for a more recent example, and references cited therein and references citing them for examples of Implicit Function Theorems for non-smooth functions.

Proof of Theorem 15, For $a:=A-A_{0} \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ and $H$ as in (2.6.9),

$$
H(\chi, a):=d_{A_{0}}^{*}\left(e^{-\chi} a e^{\chi}+e^{-\chi} d_{A_{0}}\left(e^{\chi}\right)\right),
$$

we see that for any $r>d$, we have $W^{1, r}(X) \subset C(X)$ by [4, Theorem 4.12] and so the expression (2.6.9) for $H$ defines a smooth map,

$$
\begin{align*}
H:\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \times L^{r} & \left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)  \tag{2.7.1}\\
& \rightarrow\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
\end{align*}
$$

The proof that the image of $H$ in (2.7.1) is contained in $\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)$ and not just $W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)$ is the same as the proof that the image of $H$ in (2.6.10) is contained in $\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)$ and not just $L^{q}(X ; \operatorname{ad} P)$.

Given $a$, we will solve the equation $H(\chi, a)=0$ for $\chi$ using Theorem F. 1 so that $d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0$ (as asserted by Theorem 15) with $u=e^{\chi}$, using

$$
\begin{aligned}
\mathscr{X} & =\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \\
\mathscr{Y} & =L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right) \\
\mathscr{Z} & =\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P),
\end{aligned}
$$

and $\left(\chi_{0}, a_{0}\right)=(0,0)$ and interchanging the roles of the first and second variables.
From Remark 2.7.1, we compute the partial derivative of $H(\chi, a)$ with respect to $\chi$ at the origin to be $\left(D_{1} H\right)_{(0,0)}=d_{A_{0}}^{*} d_{A_{0}}$ and so we must first verify that the operator

$$
\begin{equation*}
\left(D_{1} H\right)_{(0,0)}:\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \rightarrow\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P) \tag{2.7.2}
\end{equation*}
$$

is an isomorphism of Banach spaces. We temporarily assume that $A_{0}$ is $C^{\infty}$ and observe that

$$
d_{A_{0}}^{*} d_{A_{0}}: W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
$$

is Fredholm with index zero by standard elliptic theory with the usual kernel and range for any $r \in(1, \infty)$ : see the proof of PropositionC. 1 via Theorem A. 1 (While we might allow $A_{0}$ to be a $W^{1, q}$ Sobolev connection in this step with $q \in(d / 2, \infty)$ and $r \in(d, \infty)$, the proof of the analogue of Corollary C.2 the proof of such a refinement would be quite technical.) In particular, Theorem A.1 implies that (2.7.2) is an isomorphism.

In order to verify the remaining hypotheses (F.2) of Theorem F. 1 and prove the estimate in the conclusion of Theorem 15 via (F.3c), we shall need to identify the dependencies of the constants $M$ and $\beta$. Since

$$
\operatorname{Exp}: W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{1, r}(X ; \operatorname{ad} P)
$$

is the smooth map defined by the smooth exponential map, $\exp : \mathfrak{g} \rightarrow G$, for the Lie group $G$ with Lie algebra $\mathfrak{g}=T_{1} G$, and the domain of exp is all of $\mathfrak{g}$ by 63, Proposition 9.2.5], then the domain of Exp is all of $W_{A_{1}}^{1, r}(X ; \operatorname{ad} P)$. The constant $M=M\left(A_{0}, A_{1}, g, G, r\right)$ in (F.2a) is given by

$$
M=\left\|D_{1} H(0,0)\right\|_{\mathscr{L}\left(K^{\perp} \cap W^{1, r}, K^{\perp} \cap W^{-1, r}\right)}=\left\|\left(d_{A_{0}}^{*} d_{A_{0}}\right)^{-1}\right\|_{\mathscr{L}\left(K^{\perp} \cap W^{1, r}, K^{\perp} \cap W^{-1, r}\right)},
$$

where we abbreviate

$$
\begin{aligned}
K & =W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \cap \operatorname{Ker} d_{A_{0}}^{*} d_{A_{0}} \\
W^{1, r} & =W_{A_{1}}^{1, r}(X ; \operatorname{ad} P), \quad \text { and } \\
W^{-1, r} & =W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
\end{aligned}
$$

To satisfy the hypothesis (F.2b), we choose $\zeta=\zeta\left(A_{0}, A_{1}, g, G, r\right) \in(0,1]$ small enough that

$$
\begin{aligned}
\sup _{(\chi, a) \in U \times V} \| & D_{1} H(\chi, a)-D_{1} H(0,0) \|_{\mathscr{L}^{2}\left(K^{\perp} \cap W^{1, r} \times L^{r}, K^{\perp} \cap W^{-1, r}\right)} \\
& \leq \sup _{(\chi, a) \in U \times V}\left\|D_{1} H(\chi, a)-D_{1} H(0,0)\right\|_{\mathscr{L}^{2}\left(W^{1, r} \times L^{r}, W^{-1, r}\right)} \leq \frac{1}{2 M},
\end{aligned}
$$

where we abbreviate

$$
\begin{aligned}
L^{r} & =L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right), \\
U & =\left\{\chi \in W_{A_{1}}^{1, r}(X ; \operatorname{ad} P):\|\chi\|_{W_{A_{1}}^{1, r}(X)}<\zeta\right\}, \\
V & =\left\{a \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right):\|a\|_{L^{r}(X)}<\zeta\right\} .
\end{aligned}
$$

Finally, the constant $\beta=\beta\left(A_{0}, A_{1}, g, G, r\right)$ in (F.2C) is given by

$$
\beta=\sup _{(\chi, a) \in U \times V}\left\|D_{2} H(\chi, a)\right\|_{\mathscr{L}\left(L^{r}, K^{\perp} \cap W^{-1, r}\right)}
$$

Theorem F. 1 now provides constants $\delta=\delta\left(A_{0}, A_{1}, g, G, r\right) \in(0,1 /(2 \beta M)]$ and $C=C\left(A_{0}, A_{1}, g, G, r\right) \in(0,1]$ and a smooth map,

$$
\begin{aligned}
& \chi:\left\{b \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right):\|b\|_{L^{r}(X)}<\delta\right\} \ni a \\
& \mapsto \chi(a) \in\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap\left\{\chi \in W_{A_{1}}^{1, r}(X ; \operatorname{ad} P):\|\chi\|_{W_{A_{1}}^{1, r}(X)}<\zeta\right\}
\end{aligned}
$$

such that

$$
H(\boldsymbol{\chi}(a), a)=0, \quad \forall a \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right) \text { with }\|a\|_{L^{r}(X)}<\delta
$$

that is, the $W^{1, r}$ gauge transformation $u=e^{\chi}$ defined by the $W^{1, r}$ section $\chi=\chi(a)$ obeys

$$
\begin{gathered}
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0 \\
\left\|u(A)-A_{0}\right\|_{L^{r}(X)}<C\left\|A-A_{0}\right\|_{L^{r}(X)}
\end{gathered}
$$

as required by Theorem 15 . We can approximate any $W^{1, q}$ connection $A_{0}$ by a $C^{\infty}$ connection $\tilde{A}_{0}$ (see [4, Theorem 3.17]) and so the preceding conclusions follow by approximation for $W^{1, q}$ connections $A_{0}$ and continuity of the constants $M$ and $\beta$ with respect to variations of $A_{0}$ in the $W^{1, q}$ norm. (The proof of the latter continuity would be very similar to the proof of Corollary C.2, but simpler since we only require continuity and not compactness of Sobolev embeddings and thus is omitted.)

Finally, because $A$ is a $W^{1, q}$ connection, then $a=A-A_{0}$ is in $W^{1, q}$ and to show that $\chi$ and $u$ are in $W^{2, q}$, we can apply the proof of Lemma 2.3.1 to the nonlinear, second-order elliptic equation,

$$
d_{A_{0}}^{*}\left(u^{-1} a u+u^{-1} d_{A_{0}} u\right)=0
$$

to obtain the desired regularity for $u$. (Details of the proof are omitted since they will be very similar to those in the proof of Lemma 2.3.1.) This completes the proof of Theorem 15.

REMARK 2.7.1 (Computation of partial derivatives of $H(\chi, a)$ in (2.6.9)). While unnecessary for the proof of Theorem [15, one can explicitly compute the partial derivatives, $D_{i} H(\chi, a)$, for $i=1,2$. Observe that the expression (2.6.9) yields

$$
\begin{align*}
D_{1} H(\chi, a) \xi= & d_{A_{0}}^{*}\left(-\left(\left(D e^{-\chi}\right) \xi\right) a e^{\chi}+e^{-\chi} a\left(D e^{\chi}\right) \xi\right.  \tag{2.7.3a}\\
& \left.\left.\quad-\left(\left(D e^{-\chi}\right) \xi\right) d_{A_{0}}\left(e^{\chi}\right)+e^{-\chi} d_{A_{0}}\left(D e^{\chi}\right) \xi\right)\right) \\
D_{2} H(\chi, a) b:= & d_{A_{0}}^{*}\left(e^{-\chi} b e^{\chi}\right) \tag{2.7.3b}
\end{align*}
$$

where $\left(D e^{\chi}\right) \xi:=(D \operatorname{Exp})(\chi) \xi \in W_{A_{1}}^{1, r}(X ; \operatorname{ad} P)$. Recall that $(D \exp )(0)=\operatorname{id}_{\mathfrak{g}}$ by [63. Proposition 9.2.5] and so $(D \operatorname{Exp})(0)=\operatorname{id}_{P}$ when $\chi=0$ and we see that (2.7.3a) becomes

$$
D_{1} H(0, a) \xi=d_{A_{0}}^{*}\left(-\xi a+a \xi+d_{A_{0}} \xi\right)=d_{A_{0}}^{*}\left(d_{A_{0}} \xi+[a, \xi]\right)=d_{A_{0}}^{*} d_{A_{0}+a} \xi
$$

using $d_{A_{0}} \mathrm{id}_{P}=0$, and thus $D_{1} H(0,0)=d_{A_{0}}^{*} d_{A_{0}}$ while 2.7.3b becomes $D_{2} H(0, a) b=$ $d_{A_{0}}^{*} b$, for any $a \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$.

### 2.8. Real analytic Banach manifold structure on the quotient space of connections

The statements and proofs of Lemmata 2.8.1 and 2.8 .2 would follow standard lines (see Gilkey 48, Theorem 1.5.2], for example) if the operators

$$
d_{A}: \Omega^{l}(X ; \operatorname{ad} P) \rightarrow \Omega^{l+1}(X ; \operatorname{ad} P), \quad l \geq 0
$$

had $C^{\infty}$ coefficients, rather than Sobolev coefficients as we allow here, and formed an elliptic complex, rather than only satisfying $d_{A} \circ d_{A}=F_{A}$.

LEMMA 2.8.1 (Continuous operators on $L^{p}$ spaces and $L^{2}$-orthogonal decompositions). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$. If $A$ is a connection on $P$ of class $W^{1, q}$ with $q \geq d / 2$, and $A_{1}$ is a $C^{\infty}$ is reference connection on $P$, and $l \geq 1$ is an integer, and $p$ obeys $d / 2 \leq p \leq q$, then the operator

$$
d_{A}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right)
$$

is continuous and, if in addition $q>d / 2$, then the operator

$$
d_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

is also continuous, and there is an $L^{2}$-orthogonal decomposition,

$$
\begin{aligned}
& W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)= \operatorname{Ker} \\
&\left(d_{A}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right)\right) \\
& \oplus \operatorname{Ran}\left(d_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right)
\end{aligned}
$$

Proof. If $\xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right)$ and we write $A=A_{1}+a$, with $a \in$ $W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then $d_{A} \xi=d_{A_{1}} \xi+[a, \xi]$ and using the fact that $W^{1, p}(X) \subset$ $L^{2 p}(X)$ for any $p \geq d / 2$ by [4, Theorem 4.12] and applying the Kato Inequality 43. Equation (6.20],

$$
\begin{aligned}
\left\|d_{A} \xi\right\|_{L^{p}(X)} & \leq z\left(\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)}+\|a\|_{L^{2 p}(X)}\|\xi\|_{L^{2 p}(X)}\right) \\
& \leq z\left(\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\|\xi\|_{W_{A_{1}}^{1, p}(X)}\right) \\
& \leq z\left(1+\|a\|_{W_{A_{1}}^{1, q}(X)}\right)\|\xi\|_{W_{A_{1}}^{1, p}(X)}
\end{aligned}
$$

where $z=z(g, G, p, q) \in[1, \infty)$ and we use the fact that $q \geq p$. Similarly,

$$
\left\|d_{A}^{*} \xi\right\|_{L^{p}(X)} \leq z\left(1+\|a\|_{W_{A_{1}}^{1, q}(X)}\right)\|\xi\|_{W_{A_{1}}^{1, p}(X)}
$$

and so the operator $d_{A}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right)$ is continuous.

Moreover, defining $r \in[p, \infty]$ by $1 / p=1 / q+1 / r$, we recall that by 4, Theorem 4.12] we have $i$ ) $W^{2, p}(X) \subset L^{r}(X)$ for any $r \in[1, \infty)$ when $p=d / 2$, and ii) $W^{2, p}(X) \subset L^{\infty}(X)$ when $p>d / 2$. Thus, using

$$
\nabla_{A_{1}} d_{A} \xi=\nabla_{A_{1}} d_{A_{1}} \xi+\nabla_{A_{1}} a \times \xi+a \times \nabla_{A_{1}} \xi
$$

we see that

$$
\begin{aligned}
& \left\|\nabla_{A_{1}} d_{A} \xi\right\|_{L^{p}(X)} \\
& \leq z\left(\left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{q}(X)}\|\xi\|_{L^{r}(X)}+\|a\|_{L^{2 p}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{2 p}(X)}\right) \\
& \leq z\left(\left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{q}(X)}\|\xi\|_{W_{A_{1}}^{2, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{W_{A_{1}}^{1, p}(X)}\right) \\
& \leq z\left(1+\|a\|_{W_{A_{1}}^{1, q}(X)}\right)\|\xi\|_{W_{A_{1}}^{2, p}(X)} .
\end{aligned}
$$

We conclude that the operator

$$
d_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

is also continuous.
Note that $W^{1, p}(X) \subset L^{2}(X)$ when $p \geq 2$ or, when $1 \leq p<2$, if $p^{*}:=$ $d p /(d-p) \geq 2$, that is, $d p \geq 2 d-2 p$ or $p \geq 2 d /(d+2)$. But $p \geq d / 2$ by hypothesis and $d / 2 \geq 2 d /(d+2)$ for all $d \geq 2$, so we have $W^{1, p}(X) \subset L^{2}(X)$ for all $p \geq d / 2$ and $d \geq 2$. Using $\perp$ to denote $L^{2}$-orthogonal complement and $\oplus$ to denote $L^{2}$ orthogonal decomposition, we have

$$
\begin{aligned}
& W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \\
& =\left(\operatorname{Ran}\left(d_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right)\right)^{\perp} \\
& \quad \oplus \operatorname{Ran}\left(d_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right) .
\end{aligned}
$$

For all $\eta \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ and $\xi \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right)$ we have

$$
\left(\eta, d_{A} \xi\right)_{L^{2}(X)}=\left(d_{A}^{*} \eta, \xi\right)_{L^{2}(X)}
$$

and so

$$
\eta \perp \operatorname{Ran}\left(d_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right)
$$

if and only if

$$
\eta \in \operatorname{Ker}\left(d_{A}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l-1} \otimes \operatorname{ad} P\right)\right)
$$

This concludes the proof of the lemma.
Although not required by the proofs of Lemma 2.8.1 or Corollary 18 it is useful to note that the operator $d_{A}$ in that statement has closed range.

Lemma 2.8.2 (Closed range operators on $L^{p}$ spaces). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$. If $A$ is a connection on $P$ of class $W^{1, q}$ with $d / 2<q<\infty$, and $A_{1}$ is a $C^{\infty}$ reference connection on $P$, and $p$ obeys $d / 2 \leq p \leq q$, then the operator

$$
d_{A}: W_{A_{1}}^{2, p}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

has closed range.

Proof. Note that the operator in the statement of the lemma is bounded by Lemma 2.8.1. Let $\left\{\chi_{n}\right\}_{n \in \mathbb{N}} \subset W_{A_{1}}^{2, p}(X ; \operatorname{ad} P)$ and suppose that $d_{A} \chi_{n} \rightarrow \xi \in$ $W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ as $n \rightarrow \infty$. Thus $d_{A}^{*} d_{A} \chi_{n}=\Delta_{A} \chi_{n} \rightarrow d_{A}^{*} \xi \in L^{p}(X ; \operatorname{ad} P)$ as $n \rightarrow \infty$. We may assume without loss of generality that $\left\{\chi_{n}\right\}_{n \in \mathbb{N}} \subset\left(\operatorname{Ker} \Delta_{A}\right)^{\perp}$, where $\perp$ denotes $L^{2}$-orthogonal complement, and so the a priori estimate (2.2.4) in Corollary 2.2.5 then implies that

$$
\left\|\chi_{n}-\chi_{m}\right\|_{W_{A_{1}}^{2, p}(X)} \leq C\left\|\Delta_{A}\left(\chi_{n}-\chi_{m}\right)\right\|_{L^{p}(X)}, \quad \forall n, m \in \mathbb{N}
$$

Hence, the sequence $\left\{\chi_{n}\right\}_{n \in \mathbb{N}}$ is Cauchy in $W_{A_{1}}^{2, p}(X ; \operatorname{ad} P)$ and thus $\chi_{n} \rightarrow \chi \in$ $W_{A_{1}}^{2, p}(X ; \operatorname{ad} P)$ and $d_{A} \chi_{n} \rightarrow d_{A} \chi \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ as $n \rightarrow \infty$. Therefore, $d_{A}$ on $W_{A_{1}}^{2, p}(X ; \operatorname{ad} P)$ has closed range.

We are now ready to complete the
Proof of Corollary 18. Every compact Lie group has a compatible structure of a real analytic manifold [15, Section III.4, Exercise 1] and this structure is unique by [116, Theorem 2.11.3]. In particular, the exponential map is a real analytic diffeomorphism from an open neighborhood of the origin in the Lie algebra $\mathfrak{g}$ onto an open neighborhood of the identity in $G$. We recall from [43, Proposition A.2] that Aut $^{k+1,2}(P)$ may be given the structure of a Hilbert Lie group when $k \geq 2$ and, because $W^{2, q}(X)$ (with $q>d / 2$ ) and $H^{k+1}(X)=W^{k+1,2}(X)$ are Banach algebras and contained in $C(X)$ (the Banach algebra of continuous functions on $X$ ), the same arguments show that $\operatorname{Aut}^{2, q}(P)$ may be given the structure of a $C^{\infty}$ Banach Lie group and that both $\operatorname{Aut}^{2, q}(P)$ and $\mathrm{Aut}^{k, 2}(P)$ may be given the structure of real analytic manifolds.

According to [43, Proposition A.3], the (right) action of $\mathrm{Aut}^{k+1,2}(P)$ on $\mathscr{A}^{k, 2}(P)$ is $C^{\infty}$ when $k \geq 2$ and the same proof applies mutatis mutandis to show that this action is real analytic and that the action (2.1.1) of $\operatorname{Aut}^{2, q}(P)$ on $\mathscr{A}^{1, q}(P)$ is not only $C^{\infty}$ but also real analytic.

The only additional ingredient one needs to show that $\mathscr{B}^{*}(P)$ is real analytic is the observation that the map $H$ defined in (2.6.9) and (2.6.10) is real analytic and thus, rather than apply the customary $C^{\infty}$ Inverse Function Theorem one can instead apply its real analytic counterpart [38, Section 2.1.1] to show that for each $A_{0} \in \mathscr{A}^{1, q}(P)$, the map defined in the statement of [27, Theorem 3.2],

$$
\begin{align*}
& \Xi_{A_{0}}: \mathscr{A}^{1, q}(P) \supset \mathscr{O}_{A_{0}} \ni A \mapsto\left(u(A)-A_{0}, u\right)  \tag{2.8.1}\\
& \quad \in \operatorname{Ker}\left\{d_{A_{0}}^{*}: W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right) \rightarrow L^{q}(X ; \operatorname{ad} P)\right\} \times \operatorname{Aut}^{2, q}(P),
\end{align*}
$$

is a real analytic diffeomorphism onto an open neighborhood of $\left(0, \mathrm{id}_{P}\right)$, for a small enough open neighborhood $\mathscr{O}_{A_{0}}$ of a $W^{1, q}$ connection $A_{0}$ on $P$ and the gauge transformation $u$ is produced by Theorem 14, so $u(A)$ is in Coulomb gauge with respect to $A_{0}$. The open neighborhood $\mathscr{O}_{A_{0}}$ may be chosen to be $\operatorname{Aut}^{2, q}(P)$-invariant and the map $\Xi_{A_{0}}$ is $\operatorname{Aut}^{2, q}(P)$-equivariant. The proof that the quotient $\mathscr{B}(P)$ is a Hausdorff topological space follows mutatis mutandis either by adapting the proof of 43, Corollary, p. 50] or by adapting the proof of [27, Lemma 4.2.4], using the observation the $L^{2}$ distance function,

$$
\begin{equation*}
\operatorname{dist}_{L^{2}}([A],[B]):=\inf _{u \in \operatorname{Aut}^{2}, q(P)}\|u(A)-B\|_{L^{2}(X)} \tag{2.8.2}
\end{equation*}
$$

is a metric on $\mathscr{B}(P)$ and, in particular, that the quotient topology is metrizable. This completes the proof of Corollary 18

### 2.9. Existence of Coulomb gauge transformations for pairs

We now adapt the construction of Section 2.6 to the case of pairs. In 36, p. 280], we employed a left action of $\operatorname{Aut}(P)$ on the affine space of pairs, $\mathscr{A}(P) \times$ $C^{\infty}(X ; E)$, so $\operatorname{Aut}(P)$ acts on $\mathscr{A}(P)$ by pushforward (consistent with Donaldson and Kronheimer [27]) and on $C^{\infty}(X ; E)$ in the usual way, which is a left action. Here, to be consistent with Section 2.6 we shall use the opposite convention and continue to let $\operatorname{Aut}(P)$ act on $\mathscr{A}(P)$ by pullback (consistent with Freed and Uhlenbeck 43 and Uhlenbeck [114) and use inversion to define a right action on $C^{\infty}(X ; E)$, so that

$$
\begin{align*}
& u(A, \Phi):=\left(u^{*} A, u^{-1} \Phi\right)  \tag{2.9.1}\\
& \forall A \in \mathscr{A}(P), \Phi \in C^{\infty}(X ; E), \text { and } u \in \operatorname{Aut}(P)
\end{align*}
$$

giving a smooth (affine) right action,

$$
\mathscr{A}(P) \times C^{\infty}(X ; E) \times \operatorname{Aut}(P) \rightarrow \mathscr{A}(P) \times C^{\infty}(X ; E)
$$

Passing to Banach space completions, but temporarily suppressing the $W^{1, q}$ reference connection $A_{0}$ (for $q>d / 2$ ) from our notation, the differential of the smooth map,

$$
\operatorname{Aut}^{2, q}(P) \ni u \mapsto u(A, \Phi) \in \mathscr{A}^{1, q}(P) \times W^{1, q}(X ; E)
$$

at $\operatorname{id}_{P} \in \operatorname{Aut}^{2, q}(P)$ is given by
(2.9.2) $\quad W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \ni \xi$

$$
\mapsto d_{A, \Phi} \xi:=\left(d_{A} \xi,-\xi \Phi\right) \in W_{A_{1}}^{1, q}(X ; \operatorname{ad} P) \oplus W_{A_{1}}^{1, q}(X ; E)
$$

 pair $(A, \Phi)$ is in Coulomb gauge relative to $\left(A_{0}, \Phi_{0}\right)$ if

$$
\begin{equation*}
d_{A_{0}, \Phi_{0}}^{*}\left((A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0 \tag{2.9.3}
\end{equation*}
$$

As in the case of Theorem [14, the proof of Theorem 16 is facilitated by preparatory lemmata and a proposition, which we now state. For convenience, we define

$$
\Delta_{A_{0}, \Phi_{0}}:=d_{A_{0}, \Phi_{0}}^{*} d_{A_{0}, \Phi_{0}} \quad \text { on } W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)
$$

The proofs of Propositions 2.2.1 and 2.2.3 and Corollary 2.2.5 adapt mutatis mutandis to establish the following analogues for pairs, specialized to the case $l=0$.

Proposition 2.9.1 (A priori $L^{p}$ estimate for a Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finitedimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $(A, \Phi)$ is a $W^{1, q}$ pair on $(P, E)$ with $q>d / 2$, and $A_{1}$ is a $C^{\infty}$ connection on $P$, and $p$ obeys $d / 2 \leq p \leq q$, then

$$
\begin{equation*}
\Delta_{A, \Phi}: W_{A_{1}}^{2, p}(X ; \operatorname{ad} P) \rightarrow L^{p}(X ; \operatorname{ad} P) \tag{2.9.4}
\end{equation*}
$$

is a bounded operator. If in addition $p \in(1, \infty)$, then there is a constant $C=$ $C\left(A, \Phi, A_{1}, g, G, p, q\right) \in[1, \infty)$ such that

$$
\begin{equation*}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C\left(\left\|\Delta_{A, \Phi} \xi\right\|_{L^{p}(X)}+\|\xi\|_{L^{p}(X)}\right), \quad \forall \xi \in W_{A_{1}}^{2, p}(X ; \operatorname{ad} P) \tag{2.9.5}
\end{equation*}
$$

Proposition 2.9.2 (Spectral properties of a Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $(A, \Phi)$ is a $W^{1, q}$ pair on $(P, E)$ with $d / 2<q<\infty$, and $A_{1}$ is a $C^{\infty}$ reference connection on $P$, and $p \in(1, \infty)$ obeys $d / 2 \leq p \leq q$, then the spectrum, $\sigma\left(\Delta_{A, \Phi}\right)$, of the unbounded operator,

$$
\Delta_{A, \Phi}: \mathscr{D}\left(\Delta_{A, \Phi}\right) \subset L^{p}(X ; \operatorname{ad} P) \rightarrow L^{p}(X ; \operatorname{ad} P)
$$

is countable without accumulation points, consisting of non-negative, real eigenvalues, $\lambda$, with finite multiplicities, $\operatorname{dim} \operatorname{Ker}\left(\Delta_{A, \Phi}-\lambda\right)$.

Corollary 2.9.3 (A priori $L^{p}$ estimate for a Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group and $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, and $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $(A, \Phi)$ is a $W^{1, q}$ pair on $(P, E)$ with $d / 2<q<\infty$, and $A_{1}$ is a $C^{\infty}$ connection on $P$, and $p \in(1, \infty)$ obeys $d / 2 \leq p \leq q$, then the kernel $\operatorname{Ker} \Delta_{A, \Phi} \cap W_{A_{1}}^{2, p}(X ; \operatorname{ad} P)$ of the operator (2.9.4) is finite-dimensional. Moreover,

$$
\begin{equation*}
\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq C\left\|\Delta_{A, \Phi} \xi\right\|_{L^{p}(X)}, \quad \forall \xi \in\left(\operatorname{Ker} \Delta_{A, \Phi}\right)^{\perp} \cap W_{A_{1}}^{2, p}(X ; \operatorname{ad} P) \tag{2.9.6}
\end{equation*}
$$

where $\perp$ denotes $L^{2}$-orthogonal complement and $C=C\left(A, \Phi, A_{1}, g, G, p, q\right) \in[1, \infty)$.
The proof of Lemma 2.4.1 adapts mutatis mutandis to establish the following analogue for pairs.

Lemma 2.9.4 (Surjectivity of a perturbed Laplace operator). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. Let $A_{1}$ be a $C^{\infty}$ connection on $P$ and $(A, \Phi)$ be a $W^{1, q}$ pair on $(P, E)$ with $d / 2<q<\infty$. Then there is a constant $\delta=\delta(A, \Phi, g) \in(0,1]$ with the following significance. If $(a, \phi) \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$ obeys

$$
\begin{equation*}
\|(a, \phi)\|_{L^{d}(X)}<\delta \quad \text { when } d \geq 3 \quad \text { or } \quad\|(a, \phi)\|_{L^{4}(X)}<\delta \quad \text { when } d=2 \tag{2.9.7}
\end{equation*}
$$

then the operator,

$$
d_{A, \Phi}^{*} d_{A+a, \Phi+\phi}:\left(\operatorname{Ker} \Delta_{A, \Phi}\right)^{\perp} \cap W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow\left(\operatorname{Ker} \Delta_{A, \Phi}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)
$$

is surjective.
Finally, the proof of Proposition 2.5.1 adapts mutatis mutandis to establish the following (simplified) analogue for pairs.

Proposition 2.9.5 (A priori $W^{1, p}$ estimate for $u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)$ in terms of $\left.(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)$. Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over
$X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finitedimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. Let $A_{1}$ be a $C^{\infty}$ connection on $P$, and $\left(A_{0}, \Phi_{0}\right)$ be a $W^{1, q}$ pair on $(P, E)$ with $d / 2<q<\infty$ and $p \in(1, \infty)$ obey $d / 2 \leq p \leq q$. Then there are constants $N=N\left(A_{1}, A_{0}, \Phi_{0}, g, G, p, q\right) \in[1, \infty)$ and $\varepsilon=\varepsilon\left(A_{1}, A_{0}, \Phi_{0}, g, G, p, q\right) \in(0,1]$ with the following significance. If $(A, \Phi)$ is a $W^{1, q}$ pair on $(P, E)$ and $u \in \operatorname{Aut}(P)$ is a gauge transformation of class $W^{2, q}$ such that

$$
\begin{equation*}
d_{A_{0}, \Phi_{0}}^{*}\left(u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0 \tag{2.9.8}
\end{equation*}
$$

then the following hold. If $(A, \Phi)$ and $u(A, \Phi)$ obey

$$
\begin{equation*}
\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon \quad \text { and } \quad\left\|u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \leq \varepsilon \tag{2.9.9}
\end{equation*}
$$

then

$$
\begin{equation*}
\left\|u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \leq N\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \tag{2.9.10}
\end{equation*}
$$

Proof of Theorem 16, Given these preliminaries, Corollary 2.9.3, Lemma 2.9.4 and Proposition 2.9.5, the proof of Theorem 16 follows mutatis mutandis from that of Theorem 14

The proof of Proposition 2.9.5yields the following analogue of Lemma 2.5.2 for pairs.

Lemma 2.9.6 (A priori $W^{2, p}$ estimate for a $W^{2, q}$ gauge transformation $u$ intertwining two $W^{1, q}$ pairs). Assume the hypotheses of Proposition 2.9.5, excluding those on the pair $(A, \Phi)$. Then there is a constant $C=C\left(A_{0}, \Phi_{0}, A_{1}, g, G, p, q\right) \in$ $[1, \infty)$ with the following significance. If $(A, \Phi)$ obeys the hypotheses of Proposition 2.9.5 and $u \in \mathrm{Aut}^{2, q}(P)$ is the resulting gauge transformation, depending on $(A, \Phi)$ and $\left(A_{0}, \Phi_{0}\right)$, such that

$$
d_{A_{0}, \Phi_{0}}^{*}\left(u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0
$$

then

$$
\|u\|_{W_{A_{1}}^{2, p}(X)} \leq C
$$

While not required for the proof of Theorem 16 this is a convenient point at which to note that the proof of Lemma 2.8.1 (specialized to the case $l=1$ ) adapts mutatis mutandis to give the following analogue for pairs.

Lemma 2.9.7 (Continuous operators on $L^{p}$ spaces and $L^{2}$-orthogonal decompositions). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times \varrho \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finitedimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $(A, \Phi)$ is a Sobolev pair on $(P, E)$ of class $W^{1, q}$ with $q \geq d / 2$, and $A_{1}$ is a $C^{\infty}$ reference connection on $P$, and $p$ obeys $d / 2 \leq p \leq q$, then the operator

$$
d_{A, \Phi}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{p}(X ; \operatorname{ad} P \oplus E)
$$

is continuous and, if in addition $q>d / 2$, then the operator

$$
d_{A, \Phi}: W_{A_{1}}^{2, p}(X ; \operatorname{ad} P \oplus E) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is also continuous and there is an $L^{2}$-orthogonal decomposition,

$$
\begin{aligned}
& W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
& =\operatorname{Ker}\left(d_{A, \Phi}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{p}(X ; \operatorname{ad} P \oplus E)\right) \\
& \oplus \operatorname{Ran}\left(d_{A, \Phi}: W_{A_{1}^{2}}^{2, p}(X ; \operatorname{ad} P \oplus E) \rightarrow W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)
\end{aligned}
$$

### 2.10. Existence of Coulomb gauge transformations for pairs that are $L^{r}$-close to a reference pair

We now adapt the construction of Section 2.7 to the case of pairs, following the conventions, definitions, and notation introduced in Section 2.9 and complete the

Proof of Theorem 17, By analogy with the definition of $H$ in (2.6.9), now with $(a, \phi):=\left(A-A_{0}, \Phi-\Phi_{0}\right) \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, our definition (2.9.1) of the action of $u \in \operatorname{Aut}(P)$ on pairs $(A, \Phi)$, our definition of $d_{A_{0}, \Phi_{0}}$ in (2.9.2), and our definition (2.9.3) of Coulomb gauge for pairs, we set

$$
\begin{equation*}
H(\chi,(a, \phi)):=d_{A_{0}, \Phi_{0}}^{*}\left(e^{-\chi} a e^{\chi}+e^{-\chi} d_{A_{0}}\left(e^{\chi}\right), e^{-\chi} \phi\right) \tag{2.10.1}
\end{equation*}
$$

and the expression (2.10.1) for $H$ again defines a smooth map,

$$
\begin{align*}
H:\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) & \times L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)  \tag{2.10.2}\\
& \rightarrow\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
\end{align*}
$$

The proof that the image of $H$ in (2.10.2) is contained in $\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)$ and not just $W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)$ is the same as the proof that the image of $H$ in (2.6.10) is contained in $\left(\operatorname{Ker} \Delta_{A_{0}}\right)^{\perp} \cap L^{q}(X ; \operatorname{ad} P)$ and not just $L^{q}(X ; \operatorname{ad} P)$.

Given $(a, \phi)$, we solve the equation $H(\chi,(a, \phi))=0$ for $\chi$ using Theorem F. 1 . so that

$$
d_{A_{0}, \Phi_{0}}^{*}\left(u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0,
$$

as asserted by Theorem 15. with $u=e^{\chi}$, using

$$
\begin{aligned}
\mathscr{X} & =\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \\
\mathscr{Y} & =L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
\mathscr{Z} & =\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P),
\end{aligned}
$$

and $\left(\chi_{0},\left(a_{0}, \phi_{0}\right)\right)=(0,(0,0))$, and interchanging the roles of the first and second variables.

The partial derivative of $H(\chi,(a, \phi))$ with respect to $\chi$ at the origin is

$$
\left(D_{1} H\right)_{(0,(0,0))}=d_{A_{0}, \Phi_{0}}^{*} d_{A_{0}, \Phi_{0}}
$$

and so we must first verify that the operator

$$
\begin{align*}
\left(D_{1} H\right)_{(0,(0,0))}:\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap W_{A_{1}}^{1, r} & (X ; \operatorname{ad} P)  \tag{2.10.3}\\
& \rightarrow\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
\end{align*}
$$

is an isomorphism of Banach spaces. We temporarily assume that $\left(A_{0}, \Phi_{0}\right)$ is $C^{\infty}$ and observe that

$$
d_{A_{0}, \Phi_{0}}^{*} d_{A_{0}, \Phi_{0}}: W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
$$

is Fredholm with index zero by standard elliptic theory with the usual kernel and range for any $r \in(1, \infty)$ : see the proof of Proposition C. 1 via Theorem A.1. In particular, Theorem A.1 implies that (2.10.3) is an isomorphism.

In order to verify the remaining hypotheses (F.2) of Theorem F. 1 and prove the estimate in the conclusion of Theorem 15 via (F.3C), we shall need to identify the dependencies of the constants $M$ and $\beta$. The constant $M=M\left(A_{0}, \Phi_{0}, A_{1}, g, G, r\right)$ in (F.2a) is given by

$$
\begin{aligned}
&\left.M=\left\|D_{1} H(0,(0,0))\right\|_{\mathscr{L}\left(K^{\perp} \cap W^{1, r},\right.} K^{\perp} \cap W^{-1, r}\right) \\
&=\left\|\left(d_{A_{0}, \Phi_{0}}^{*} d_{A_{0}, \Phi_{0}}\right)^{-1}\right\|_{\mathscr{L}\left(K^{\perp} \cap W^{1, r}, K^{\perp} \cap W^{-1, r}\right)},
\end{aligned}
$$

where we abbreviate

$$
\begin{aligned}
K & =W_{A_{1}}^{1, r}(X ; \operatorname{ad} P) \cap \operatorname{Ker} d_{A_{0}, \Phi_{0}}^{*} d_{A_{0}, \Phi_{0}} \\
W^{1, r} & =W_{A_{1}}^{1, r}(X ; \operatorname{ad} P), \quad \text { and } \\
W^{-1, r} & =W_{A_{1}}^{-1, r}(X ; \operatorname{ad} P)
\end{aligned}
$$

To satisfy the hypothesis (F.2b), we choose $\zeta=\zeta\left(A_{0}, A_{1}, g, G, r\right) \in(0,1]$ small enough that

$$
\begin{aligned}
& \sup _{(\chi,(a, \phi)) \in U \times V}\left\|D_{1} H(\chi,(a, \phi))-D_{1} H(0,(0,0))\right\|_{\mathscr{L}^{2}\left(K^{\perp} \cap W^{1, r} \times L^{r}, K^{\perp} \cap W^{-1, r}\right)} \\
& \quad \leq \sup _{(\chi,(a, \phi)) \in U \times V}\left\|D_{1} H(\chi,(a, \phi))-D_{1} H(0,(0,0))\right\|_{\mathscr{L}^{2}\left(W^{1, r} \times L^{r}, W^{-1, r}\right)} \leq \frac{1}{2 M},
\end{aligned}
$$

where we abbreviate

$$
\begin{aligned}
L^{r} & =L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right), \\
U & =\left\{\chi \in W_{A_{1}}^{1, r}(X ; \operatorname{ad} P):\|\chi\|_{W_{A_{1}}^{1, r}(X)}<\zeta\right\}, \quad \text { and } \\
V & =\left\{a \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right):\|(a, \phi)\|_{L^{r}(X)}<\zeta\right\}
\end{aligned}
$$

Finally, the constant $\beta=\beta\left(A_{0}, \Phi_{0}, A_{1}, g, G, r\right)$ in (F.2C) is given by

$$
\beta=\sup _{(\chi,(a, \phi)) \in U \times V}\left\|D_{2} H(\chi,(a, \phi))\right\|_{\mathscr{L}\left(L^{r}, K^{\perp} \cap W^{-1, r}\right)} .
$$

Theorem F. 1 provides constants $\delta=\delta\left(A_{0}, \Phi_{0}, A_{1}, g, G, r\right) \in(0,1 /(2 \beta M)]$ and $C=$ $C\left(A_{0}, \Phi_{0}, A_{1}, g, G, r\right) \in(0,1]$ and a smooth map,

$$
\begin{aligned}
\chi:\{(b, \psi) & \left.\in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right):\|(b, \psi)\|_{L^{r}(X)}<\delta\right\} \ni a \\
& \mapsto \chi(a) \in\left(\operatorname{Ker} \Delta_{A_{0}, \Phi_{0}}\right)^{\perp} \cap\left\{\chi \in W_{A_{1}}^{1, r}(X ; \operatorname{ad} P):\|\chi\|_{W_{A_{1}}^{1, r}(X)}<\zeta\right\}
\end{aligned}
$$

such that

$$
\begin{aligned}
& H(\boldsymbol{\chi}(a, \phi),(a, \phi))=0 \\
& \qquad \forall(a, \phi) \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \text { with }\|(a, \phi)\|_{L^{r}(X)}<\delta
\end{aligned}
$$

that is, the $W^{1, r}$ gauge transformation $u=e^{\chi}$ defined by the $W^{1, r}$ section $\chi=$ $\chi(a, \phi)$ obeys

$$
\begin{gathered}
d_{A_{0}, \Phi_{0}}^{*}\left(u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right)=0 \\
\left\|u(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{L^{r}(X)}<C\left\|(A, \Phi)-\left(A_{0}, \Phi_{0}\right)\right\|_{L^{r}(X)}
\end{gathered}
$$

as required by Theorem 17. We can approximate any $W^{1, q}$ pair $\left(A_{0}, \Phi_{0}\right)$ by a $C^{\infty}$ pair $\left(\tilde{A}_{0}, \tilde{\Phi}_{0}\right)$ (see [4, Theorem 3.17]) and so the preceding conclusions follow by approximation for $W^{1, q}$ pairs $\left(A_{0}, \Phi_{0}\right)$ and continuity of the constants $M$ and $\beta$ with respect to variations of $\left(A_{0}, \Phi_{0}\right)$ in the $W^{1, q}$ norm. (The proof of the latter continuity would be very similar to the proof of Corollary C.2, but simpler since we only require continuity and not compactness of Sobolev embeddings and thus is omitted.)

Finally, because $(A, \Phi)$ is a $W^{1, q}$ pair, then $(a, \phi)=\left(A-A_{0}, \Phi-\Phi_{0}\right)$ is in $W^{1, q}$ and to show that $\chi$ and $u$ are in $W^{2, q}$, we can apply the proof of Lemma 2.3.1 to the nonlinear, second-order elliptic equation,

$$
d_{A_{0}, \Phi_{0}}^{*}\left(u^{-1} a u+u^{-1} d_{A_{0}} u, u^{-1} \phi\right)=0
$$

to obtain the desired regularity for $u$. (Details of the proof are omitted since they will be very similar to those in the proof of Lemma 2.3.1.) This completes the proof of Theorem 17.

### 2.11. Regularity for solutions to the pure and coupled Yang-Mills equations

It is well-known that techniques due to Uhlenbeck [115, 114] can be used to show that, given a weak solution to the Yang-Mills equation, there exists a gauge transformation such that the gauge-transformed solution is smooth. We give a proof of a similar fact here that generalizes easily to the case of coupled Yang-Mills equations.

We have the following generalization of (97, Theorem 5.3], due to Parker, and [36, Proposition 3.7], due to the author and Leness, from the case of $d=4$ to arbitrary $d \geq 2$.

Theorem 2.11.1 (Regularity for solutions to the Yang-Mills equation). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be $a$ compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$. If $q \in(4 / 3, \infty)$ obeys $q>d / 2$ and $A$ is a $W^{1, q}$ connection on $P$ that is a weak solution to the YangMills equation with respect to the Riemannian metric $g$, then there exists a $W^{2, q}$ gauge transformation $u \in \operatorname{Aut}(P)$ such that $u(A)$ is a $C^{\infty}$ Yang-Mills connection on $P$.

Proof. We proceed as in the proof of [36, Proposition 3.7] and note that the affine space $\mathscr{A}(P)$ of $C^{\infty}$ connections on $P$ is dense in the affine space $\mathscr{A}^{1, q}(P)$ of $W^{1, q}$ connections on $P$ and so there exists a $C^{\infty}$ connection $A_{0}$ on $P$ such that

$$
\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}<\zeta
$$

where $\zeta=\zeta\left(A_{0}, A_{1}, g, G, q\right) \in(0,1]$ is the constant in Theorem 14 and $A_{1}$ is any fixed $C^{\infty}$ reference connection on $P$. Hence, there is a $W^{2, q}$ gauge transformation $u \in \operatorname{Aut}(P)$ such that $u(A)$ obeys

$$
d_{A_{0}}^{*}\left(u(A)-A_{0}\right)=0
$$

and

$$
\left\|u(A)-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}<2 N\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, q}(X)}<2 N \zeta,
$$

where $N=N\left(A_{0}, A_{1}, g, G, q\right) \in[1, \infty)$ is the constant in Proposition 2.5.1. Hence, we may assume without loss of generality that $A$ is in Coulomb gauge with respect to $A_{0}$ and that

$$
a:=A-A_{0} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

is a weak solution to

$$
d_{A_{0}+a}^{*} F_{A_{0}+a}=0 \quad \text { and } \quad d_{A_{0}}^{*} a=0,
$$

and thus a weak solution to the quasi-linear, second-order elliptic system,

$$
\begin{equation*}
\left(\Delta_{A_{0}}+\lambda_{0}\right) a+a \times \nabla_{A_{0}} a+a \times a \times a=\lambda_{0} a-F_{A_{0}} \tag{2.11.1}
\end{equation*}
$$

where $\Delta_{A_{0}}=d_{A_{0}}^{*} d_{A_{0}}+d_{A_{0}} d_{A_{0}}^{*}$ is the usual Hodge Laplacian on $\Omega^{1}(X ; \operatorname{ad} P)$ and $\lambda_{0}>0$ is any positive constant. We recall from [32] that, for any $q \in(1, \infty)$ and integer $k \geq 0$, the operator

$$
\begin{equation*}
\Delta_{A_{0}}+\lambda_{0}: W_{A_{1}}^{k+2, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{k, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right) \tag{2.11.2}
\end{equation*}
$$

is invertible. The right-hand side of Equation (2.11.1) belongs to $W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ by hypothesis on $A$ and the fact that $A_{0}$ is $C^{\infty}$. If $q>d$, then $W^{1, q}(X) \subset C(X)$ by [4, Theorem 4.12] and so

$$
a \times \nabla_{A_{0}} a+a \times a \times a \in L^{q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

But then existence and uniqueness of solutions in $W_{A_{1}}^{2, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ to (2.11.1), given a source term in $L^{q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, implies that $a \in W_{A_{1}}^{2, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$. Using the fact that $W^{k, q}(X)$ is a Banach algebra when $k q>d$ and invertibility of (2.11.2), we can iterate in the usual way to show that $a \in W_{A_{1}}^{k+2, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ for all integers $k \geq 0$ and hence that $a \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$.

Therefore, it suffices to consider the case $d / 2<q<d$. Recall that $q^{*}=$ $d q /(d-q)$ and so $q^{*} \in(d, \infty)$. Suppose that we can choosf ${ }^{5} r \in(1, q)$ obeying $1 / q^{*}+1 / q \leq 1 / r$. Then there is a continuous multiplication map $L^{q^{*}}(X) \times L^{q}(X) \rightarrow$ $L^{r}(X)$ and a continuous Sobolev embedding $W^{1, q}(X) \subset L^{q^{*}}(X)$ by 4. Theorem 4.12]. Thus, $a \times \nabla_{A_{0}} a \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ if

$$
1 / r \geq 1 / q^{*}+1 / q=(d-q) /(d q)+1 / q=(2 d-q) /(d q)
$$

that is, $r \leq d q /(2 d-q)$. Now $(2 d-q) /(d q)<1$ (which we need to permit a choice of $r>1)$ and thus $d q /(2 d-q)>1 \Longleftrightarrow d q>2 d-q \Longleftrightarrow q>2 d /(d+1)$. But $q>d / 2$ by hypothesis (when $d \geq 3$ ) and $d / 2 \geq 2 d /(d+1) \Longleftrightarrow d \geq 3$; for the case $d=2$, we need $q>4 / 3$ to ensure that we can choose $r>1$ and this explains the restriction in our hypotheses that $q>4 / 3$ when $d=2$.

Similarly, we have a continuous multiplication map $L^{3 s}(X) \times L^{3 s}(X) \times L^{3 s}(X) \rightarrow$ $L^{s}(X)$ for any $s \in[1, \infty]$ and a continuous Sobolev embedding $W^{1, q}(X) \subset L^{3 s}(X)$ provided $3 s \leq q^{*}$, that is, $s \leq d q /(3(d-q))$ and for this choice of $s$, we have $a \times a \times a \in L^{s}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$.

We now observe that $r \leq s$ when $q \geq d / 2$, as we assume by hypothesis, if we choose $r=d q /(2 d-q)$ and $s=d q /(3(d-q))$. Indeed, we then have

$$
r \leq s \Longleftrightarrow d q /(2 d-q) \leq d q /(3(d-q)) \Longleftrightarrow 3 d-3 q \leq 2 d-q \Longleftrightarrow d \leq 2 q .
$$

[^5]Hence, for $r=d q /(2 d-q)$, we have

$$
a \times \nabla_{A_{0}} a+a \times a \times a \in L^{r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

and thus elliptic regularity theory for (2.11.1) implies that $a \in W_{A_{1}}^{2, r}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, similar to the case $q>d$. By [4, Theorem 4.12], when $r<d$, we have a continuous Sobolev embedding, $W^{2, r}(X) \subset W^{1, r^{*}}(X)$, where $r^{*}=d r /(d-r)$, and thus we obtain

$$
a \in W_{A_{1}}^{1, r^{*}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)
$$

In the limiting case $q=d / 2$ with $d \geq 3$ we would have $r=d q /(2 d-q)=$ $\left(d^{2} / 2\right) /(2 d-d / 2)=d / 3$ and in the limiting case $q=d$ we would have $r=d$, so $r \in(d / 3, d)$ and thus $r^{*}=d r /(d-r) \in(d / 2, \infty)$. In particular,

$$
r^{*}=\frac{d r}{d-r}=\frac{d^{2} q /(2 d-q)}{d-d q /(2 d-q)}=\frac{d q /(2 d-q)}{1-q /(2 d-q)}=\frac{d q}{2(d-q)}
$$

We may write $r^{*}=q+\delta$, where $\delta=\delta(d, q)$ is defined by

$$
\delta:=r^{*}-q=\frac{d q}{2(d-q)}-q=\frac{d q-2(d-q) q}{2(d-q)}=\frac{2 q^{2}-d q}{2(d-q)}=\frac{q(2 q-d)}{2(d-q)}
$$

and thus $\delta(d, q)>0$ for $d / 2<q<d$ and $d \geq 3$. Consequently, we see a regularity improvement and because $\delta(d, q)$ is an increasing function of $q$, only finitely many iterations of this regularity improvement are required to give $r^{*}>d$, at which point we can apply the regularity argument for the case $q>d$ to again conclude that $a \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$.

In the limiting case $q=4 / 3$ with $d=2$, we would have $r=d q /(2 d-q)=$ $(8 / 3) /(4-4 / 3)=(8 / 3) /(8 / 3)=1$, and in the limiting case $q=2$ with $d=2$ we would have $r=4 /(4-2)=2$, so $r \in(1,2)$ and thus $r^{*}=d r /(d-r) \in(2, \infty)$. In particular, $r^{*}>4 / 3+2 / 3$. Consequently, we see a regularity improvement just as we did in the case $d \geq 3$ and the remainder of the argument follows mutatis mutandis as in the case $d \geq 3$.

The proof of Theorem 2.11.1 adapts mutatis mutandis to give the following generalization of [97, Theorem 5.3], due to Parker, and [36, Proposition 3.7], due to the author and Leness, from the case of $d=4$ to arbitrary $d \geq 2$.

Theorem 2.11.2 (Regularity for solutions to the boson coupled Yang-Mills equations). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq$ 2 , and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finitedimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$. If $d / 2<q<\infty$ and $\left(A_{\infty}, \Phi_{\infty}\right)$ is a $W^{1, q}$ pair on $(P, E)$ that is a critical point of the boson coupled Yang-Mills energy function (1.3.1), then there exists a $W^{2, q}$ gauge transformation $u \in \operatorname{Aut}(P)$ such that $u\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair on $(P, E)$.

Proof. We proceed as in the proofs of [97, Theorem 5.3] and [36, Proposition 3.7] and note that the affine space $\mathscr{A}(P) \times C^{\infty}(E)$ of $C^{\infty}$ pairs is dense in the affine space $\mathscr{A}^{1, q}(P) \times W^{1, q}(X ; E)$ of $W^{1, q}$ pairs and so there exists a $C^{\infty}$ pair $\left(A_{0}, \Phi_{0}\right)$ on $(P, E)$ such that

$$
\left\|\left(A_{\infty}, \Phi_{\infty}\right)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, q}(X)}<\zeta
$$

where $\zeta=\zeta\left(A_{1}, A_{0}, \Phi_{0}, g, G, q\right) \in(0,1]$ is the constant in Theorem 16 and $A_{1}$ is any fixed $C^{\infty}$ reference connection on $P$. Hence, there is a gauge transformation $u \in \operatorname{Aut}(P)$ of class $W^{2, q}$ such that $u\left(A_{\infty}, \Phi_{\infty}\right)$ obeys

$$
d_{A_{0}, \Phi_{0}}^{*}\left(u\left(A_{\infty}, \Phi_{\infty}\right)-\left(A_{0}, \Phi_{0}\right)\right)=0
$$

and

$$
\left\|u\left(A_{\infty}, \Phi_{\infty}\right)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<2 N\left\|\left(A_{\infty}, \Phi_{\infty}\right)-\left(A_{0}, \Phi_{0}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta,
$$

where $N=N\left(A_{1}, A_{0}, \Phi_{0}, g, G, q\right) \in[1, \infty)$ is the constant in Proposition 2.9.5 Hence, we may assume without loss of generality that $\left(A_{\infty}, \Phi_{\infty}\right)$ is in Coulomb gauge with respect to $\left(A_{0}, \Phi_{0}\right)$ and that

$$
(a, \phi):=\left(A_{\infty}, \Phi_{\infty}\right)-\left(A_{0}, \Phi_{0}\right) \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is a weak solution to the quasi-linear, second-order elliptic system,

$$
\begin{aligned}
& d_{A_{0}}^{*} d_{A_{0}} a+\nabla_{A_{0}}^{*} \nabla_{A_{0}} \phi+a \times \nabla_{A_{0}} a+a \times \nabla_{A_{0}} \phi+\phi \times \nabla_{A_{0}} \phi \\
& \quad+a \times a \times a+a \times a \times \phi+a \times \phi \times \phi \\
& \quad+m \phi+s \Phi_{0} \times \Phi_{0} \times \phi+s \Phi_{0} \times \phi \times \phi+s \phi \times \phi \times \phi=f\left(m, s, A_{0}, \Phi_{0}\right), \\
& d_{A_{0}, \Phi_{0}} d_{A_{0}, \Phi_{0}}^{*}(a, \phi)=0,
\end{aligned}
$$

where we employ the expression (1.3.5) for the gradient $\mathscr{M}\left(A_{\infty}, \Phi_{\infty}\right)$ of $\mathscr{E}$ at $\left(A_{\infty}, \Phi_{\infty}\right)$ and $f\left(m, s, A_{0}, \Phi_{0}\right) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$ is the tautologically defined right-hand source term. Ellipticity of the preceding system follows by expanding the expression (2.9.2) for $d_{A_{0}, \Phi_{0}}$ on $\Omega^{0}(X ; \operatorname{ad} P)$ to extract the second-order term $d_{A_{0}} d_{A_{0}}^{*} a$ and recalling that the operator $d_{A_{0}}^{*} d_{A_{0}}+d_{A_{0}} d_{A_{0}}^{*}$ is clearly elliptic, being the usual Hodge Laplacian on $\Omega^{1}(X ; \operatorname{ad} P)$. The remainder of the proof now follows the pattern of the proof of Theorem 2.11.1.

## CHAPTER 3

## Lojasiewicz-Simon gradient inequalities for coupled Yang-Mills energy functions

In this section, we apply our Coulomb-gauge transformation result, Theorem 16. and abstract Lojasiewicz-Simon gradient inequality, Theorem [1, to prove the corresponding Łojasiewicz-Simon gradient inequalities for the boson and fermion coupled Yang-Mills energy functions, Theorem 4 in Section 3.1, and Theorem 6 in Section 3.2.

## 3.1. Łojasiewicz-Simon gradient inequality for the boson coupled Yang-Mills energy function

For any $C^{\infty}$ reference connection $A_{1}$ on $P$, let

$$
\begin{equation*}
\mathscr{P}^{k, p}(P, E):=\mathscr{A}^{k, p}(P) \times W_{A_{1}}^{k, p}(X ; E) \tag{3.1.1}
\end{equation*}
$$

denote the affine space of $W^{k, p}$ pairs on $(P, E)$, where $k \in \mathbb{Z}$ is a positive integer and $p \in(1, \infty)$; we write $\mathscr{P}(P, E):=\mathscr{A}(P) \times C^{\infty}(X ; E)$ for the affine space of $C^{\infty}$ pairs on $(P, E)$.
3.1.1. Analyticity of the boson coupled Yang-Mills energy function and its gradient map on the Sobolev space of pairs. We begin by establishing the following result on analyticity of the boson coupled Yang-Mills energy function, $\mathscr{E}: x_{\infty}+\mathscr{X} \rightarrow \mathbb{R}$ (where $x_{\infty}=\left(A_{\infty}, \Phi_{\infty}\right)$ is a critical point); this also serves as a stepping stone towards the proof that its gradient map, $\mathscr{M}: x_{\infty}+\mathscr{X} \rightarrow \tilde{\mathscr{X}}$ is real analytic for suitable choices of the Banach spaces, $\mathscr{X}$ and $\tilde{X}$, in Theorem 1 .

Proposition 3.1.1 (Analyticity of the boson coupled Yang-Mills energy function on the affine space of $W^{1, p}$ pairs). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$, and $A_{1}$ be a $C^{\infty}$ connection on $P$, and $m, s \in C^{\infty}(X)$. If $4 d /(d+4) \leq p<\infty$ wher $d \geq 2$, then the function,

$$
\mathscr{E}: \mathscr{P}^{1, p}(P, E) \rightarrow \mathbb{R}
$$

is real analytic, where $\mathscr{E}$ is as in (1.3.1).
Proof. We fix a pair $(A, \Phi) \in \mathscr{P}^{1, p}(P, E)$ and write $(A, \Phi)=\left(A_{1}, \Phi_{1}\right)+$ $\left(a_{1}, \phi_{1}\right)$, where $\left(a_{1}, \phi_{1}\right) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$. We will show that $\mathscr{E}$ is analytic

[^6]at $(A, \Phi)$. For $(a, \phi) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, we write $A+a=A_{1}+a_{1}+a$ and expand
$$
F_{A+a}=F_{A_{1}+a_{1}+a}=F_{A_{1}}+d_{A_{1}}\left(a_{1}+a\right)+\left(a_{1}+a\right) \times\left(a_{1}+a\right)
$$
and
$$
\nabla_{A+a}(\Phi+\phi)=\nabla_{A+a_{1}+a}(\Phi+\phi)=\nabla_{A_{1}}(\Phi+\phi)+\varrho\left(a_{1}+a\right)(\Phi+\phi)
$$

Using the definition (1.3.1) of $\mathscr{E}$, we compute

$$
2 \mathscr{E}(A+a, \Phi+\phi)=T_{1}+T_{2}+T_{3},
$$

where the terms $T_{i}:=T_{i}(a, \phi)$, for $i=1,2,3$, are given by

$$
\begin{aligned}
T_{1}:= & \left\|F_{A_{1}}\right\|_{L^{2}(X)}^{2}+\left\|d_{A_{1}}\left(a_{1}+a\right)\right\|_{L^{2}(X)}^{2} \\
& +\left\|\left(a_{1}+a\right) \times\left(a_{1}+a\right)\right\|_{L^{2}(X)}^{2}+2\left(F_{A_{1}}, d_{A_{1}}\left(a_{1}+a\right)\right)_{L^{2}(X)} \\
& +2\left(F_{A_{1}},\left(a_{1}+a\right) \times\left(a_{1}+a\right)\right)_{L^{2}(X)} \\
& +2\left(d_{A_{1}}\left(a_{1}+a\right),\left(a_{1}+a\right) \times\left(a_{1}+a\right)\right)_{L^{2}(X)}
\end{aligned}
$$

and

$$
\begin{aligned}
T_{2}:= & \left\|\nabla_{A_{1}}(\Phi+\phi)\right\|_{L^{2}(X)}^{2}+\left\|\varrho\left(a_{1}+a\right)(\Phi+\phi)\right\|_{L^{2}(X)}^{2} \\
& +\left(\nabla_{A_{1}}(\Phi+\phi), \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)} \\
& +\left(\varrho\left(a_{1}+a\right)(\Phi+\phi), \nabla_{A_{1}}(\Phi+\phi)\right)_{L^{2}(X)}, \\
= & \left\|\nabla_{A_{1}}(\Phi+\phi)\right\|_{L^{2}(X)}^{2}+\left\|\varrho\left(a_{1}+a\right)(\Phi+\phi)\right\|_{L^{2}(X)}^{2} \\
& +2 \operatorname{Re}\left(\nabla_{A_{1}}(\Phi+\phi), \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)},
\end{aligned}
$$

and

$$
T_{3}:=-\int_{X}\left(m|\Phi+\phi|^{2}+s|\Phi+\phi|^{4}\right) d \operatorname{vol}_{g}
$$

Hence, we can write the difference as

$$
2 \mathscr{E}(A+a, \Phi+\phi)-2 \mathscr{E}(A, \Phi)=T_{1}^{\prime}+T_{2}^{\prime}+T_{3}^{\prime}
$$

where the difference terms $T_{i}^{\prime}:=T_{i}(a, \phi)-T_{i}(0,0)$, for $i=1,2,3$, are given by

$$
\begin{aligned}
T_{1}^{\prime}= & \left\|d_{A_{1}} a\right\|_{L^{2}(X)}^{2}+2\left(d_{A_{1}} a_{1}, d_{A_{1}} a\right)_{L^{2}(X)} \\
& +\left(a \times\left(a_{1}+a\right),\left(a_{1}+a\right) \times\left(a_{1}+a\right)\right)_{L^{2}(X)} \\
& +2\left(F_{A_{1}}, d_{A_{1}} a\right)_{L^{2}(X)}+\left(F_{A_{1}}, a \times\left(a_{1}+a\right)\right)_{L^{2}(X)} \\
& +\left(d_{A_{1}} a,\left(a_{1}+a\right) \times\left(a_{1}+a\right)\right)_{L^{2}(X)}+\left(d_{A_{1}} a_{1}, a \times\left(a_{1}+a\right)\right)_{L^{2}(X)},
\end{aligned}
$$

and

$$
\begin{align*}
T_{2}^{\prime}= & \left\|\nabla_{A_{1}} \phi\right\|_{L^{2}(X)}^{2}+2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \nabla_{A_{1}} \phi\right)_{L^{2}(X)} \\
& +\|\varrho(a) \Phi\|_{L^{2}(X)}^{2}+\left\|\varrho\left(a_{1}+a\right) \phi\right\|_{L^{2}(X)}^{2} \\
& +2 \operatorname{Re}\left(\varrho\left(a_{1}\right) \Phi, \varrho(a) \Phi\right)_{L^{2}(X)}+2 \operatorname{Re}\left(\varrho\left(a_{1}+a\right) \Phi, \varrho\left(a_{1}+a\right) \phi\right)_{L^{2}(X)}  \tag{3.1.2}\\
& +2 \operatorname{Re}\left(\nabla_{A_{1}} \phi, \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)} \\
& +2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \varrho(a) \Phi+\varrho\left(a_{1}+a\right) \phi\right)_{L^{2}(X)},
\end{align*}
$$

and

$$
\begin{aligned}
T_{3}^{\prime}= & \int_{X}\left(m|\phi|^{2}+2 m \operatorname{Re}\langle\Phi, \phi\rangle+s|\phi|^{4}+4 s(\operatorname{Re}\langle\Phi, \phi\rangle)^{2}\right) d \operatorname{vol}_{g} \\
& +\int_{X}\left(4 s\left(|\Phi|^{2}+|\phi|^{2}\right) \operatorname{Re}\langle\Phi, \phi\rangle+2 s|\Phi|^{2}|\phi|^{2}\right) d \operatorname{vol}_{g}
\end{aligned}
$$

To see the origin of the expression (3.1.2) for $T_{2}^{\prime}$, we observe that

$$
\begin{aligned}
T_{2}^{\prime}:= & \left\|\nabla_{A_{1}+a_{1}+a}(\Phi+\phi)\right\|_{L^{2}(X)}^{2}-\left\|\nabla_{A_{1}+a_{1}} \Phi\right\|_{L^{2}(X)}^{2} \\
= & \left.\left\|\nabla_{A_{1}}(\Phi+\phi)+\varrho\left(a_{1}+a\right)(\Phi+\phi)\right\|_{L^{2}(X)}^{2}-\| \nabla_{A_{1}} \Phi+\varrho\left(a_{1}\right) \Phi\right) \|_{L^{2}(X)}^{2} \\
= & \left\|\nabla_{A_{1}}(\Phi+\phi)\right\|_{L^{2}(X)}^{2}-\left\|\nabla_{A_{1}} \Phi\right\|_{L^{2}(X)}^{2} \\
& +\left\|\varrho\left(a_{1}+a\right)(\Phi+\phi)\right\|_{L^{2}(X)}^{2}-\left\|\varrho\left(a_{1}\right) \Phi\right\|_{L^{2}(X)}^{2} \\
& +2 \operatorname{Re}\left(\nabla_{A_{1}}(\Phi+\phi), \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)}-2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \varrho\left(a_{1}\right) \Phi\right)_{L^{2}(X)} \\
= & T_{21}^{\prime}+T_{22}^{\prime}+T_{23}^{\prime} .
\end{aligned}
$$

For the first term, we have

$$
\begin{aligned}
T_{21}^{\prime}:=\left\|\nabla_{A_{1}}(\Phi+\phi)\right\|_{L^{2}(X)}^{2}-\| \nabla_{A_{1}} \Phi & \|_{L^{2}(X)}^{2} \\
& =\left\|\nabla_{A_{1}} \phi\right\|_{L^{2}(X)}^{2}+2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \nabla_{A_{1}} \phi\right)_{L^{2}(X)}
\end{aligned}
$$

For the second term, we see that

$$
\begin{aligned}
T_{22}^{\prime}:= & \left\|\varrho\left(a_{1}+a\right)(\Phi+\phi)\right\|_{L^{2}(X)}^{2}-\left\|\varrho\left(a_{1}\right) \Phi\right\|_{L^{2}(X)}^{2} \\
= & \left\|\varrho\left(a_{1}\right) \Phi+\varrho(a) \Phi+\varrho\left(a_{1}+a\right) \phi\right\|_{L^{2}(X)}^{2}-\left\|\varrho\left(a_{1}\right) \Phi\right\|_{L^{2}(X)}^{2} \\
= & \|\varrho(a) \Phi\|_{L^{2}(X)}^{2}+\left\|\varrho\left(a_{1}+a\right) \phi\right\|_{L^{2}(X)}^{2}+2 \operatorname{Re}\left(\varrho\left(a_{1}\right) \Phi, \varrho(a) \Phi\right)_{L^{2}(X)} \\
& +2 \operatorname{Re}\left(\varrho\left(a_{1}\right) \Phi, \varrho\left(a_{1}+a\right) \phi\right\rangle+2 \operatorname{Re}\left(\varrho(a) \Phi, \varrho\left(a_{1}+a\right) \phi\right)_{L^{2}(X)} \\
= & \|\varrho(a) \Phi\|_{L^{2}(X)}^{2}+\left\|\varrho\left(a_{1}+a\right) \phi\right\|_{L^{2}(X)}^{2}+2 \operatorname{Re}\left(\varrho\left(a_{1}\right) \Phi, \varrho(a) \Phi\right)_{L^{2}(X)} \\
& +2 \operatorname{Re}\left(\varrho\left(a_{1}+a\right) \Phi, \varrho\left(a_{1}+a\right) \phi\right)_{L^{2}(X)} .
\end{aligned}
$$

For the third term, we have

$$
\begin{aligned}
T_{23}^{\prime}:= & 2 \operatorname{Re}\left(\nabla_{A_{1}}(\Phi+\phi), \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)}-2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \varrho\left(a_{1}\right) \Phi\right) \\
= & 2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \varrho\left(a_{1}\right) \Phi+\varrho(a) \Phi+\varrho\left(a_{1}+a\right) \phi\right)_{L^{2}(X)} \\
& +2 \operatorname{Re}\left(\nabla_{A_{1}} \phi, \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)}-2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \varrho\left(a_{1}\right) \Phi\right)_{L^{2}(X)} \\
= & 2 \operatorname{Re}\left(\nabla_{A_{1}} \Phi, \varrho(a) \Phi+\varrho\left(a_{1}+a\right) \phi\right)_{L^{2}(X)} \\
& +2 \operatorname{Re}\left(\nabla_{A_{1}} \phi, \varrho\left(a_{1}+a\right)(\Phi+\phi)\right)_{L^{2}(X)} .
\end{aligned}
$$

By adding the preceding terms, we obtain the expression (3.1.2) for $T_{2}^{\prime}$.
To complete the proof of analyticity of $\mathscr{E}$ at $(A, \Phi)$, we observe that by 4, Theorem 4.12] there is a continuous Sobolev embedding, $W^{1, p}(X) \subset L^{4}(X)$, if $p$ obeys

$$
\left\{\begin{array}{l}
p^{*}=d p /(d-p) \geq 4 \quad \text { when } 1 \leq p<d, \text { or } \\
p \geq d
\end{array}\right.
$$

When $p<d$, then $p^{*}=d p /(d-p) \geq 4$ if and only if $d p \geq 4(d-p)$, that is, $p(d+4) \geq 4 d$ or equivalently $p \geq 4 d /(d+4)$. Moreover, $4 d /(d+4) \leq d$ for all $d \geq 2$ and thus it suffices to choose $p \geq 4 d /(d+4)$, as in our hypothesis, in order to ensure a continuous Sobolev embedding, $W^{1, p}(X) \subset L^{4}(X)$.

Hence, we obtain a continuous multilinear map, $\otimes_{i=1}^{4} W^{1, p}(X) \rightarrow L^{1}(X)$, by combining the Sobolev embedding $W^{1, p}(X) \rightarrow L^{4}(X)$ with the continuous multiplication map, $\otimes_{i=1}^{4} L^{4}(X) \rightarrow L^{1}(X)$. Combining these observations with the Kato Inequality [43, Equation (6.20)], we obtain an estimate of the form,

$$
|\mathscr{E}(A+a, \Phi+\phi)-\mathscr{E}(A, \Phi)| \leq\left|T_{1}^{\prime}\right|+\left|T_{2}^{\prime}\right|+\left|T_{3}^{\prime}\right|
$$

where, for a constant $C=C(g, G) \in[1, \infty)$,

$$
\begin{aligned}
C^{-1}\left|T_{1}^{\prime}\right| \leq & \|a\|_{W_{A_{1}}^{1, p}(X)}^{2}+\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}\|a\|_{W_{A_{1}}^{1, p}(X)} \\
& +\|a\|_{W_{A_{1}}^{1, p}(X)}\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right)^{3} \\
& +\left\|F_{A_{1}}\right\|_{L^{2}(X)}\|a\|_{W_{A_{1}}^{1, p}(X)} \\
& +\left\|F_{A_{1}}\right\|_{L^{2}(X)}\|a\|_{W_{A_{1}}^{1, p}(X)}\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right) \\
& +\|a\|_{W_{A_{1}}^{1, p}(X)}\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right)^{2} \\
& +\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}\|a\|_{W_{A_{1}}^{1, p}(X)}\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right)
\end{aligned}
$$

noting that $W^{1, p}(X) \subset L^{2}(X)$ if $d p /(d-p) \geq 2$, that is, $d p \geq 2(d-p)$ or $p(d+2) \geq$ $2 d$ or $p \geq 2 d /(d+2)$, and

$$
\begin{aligned}
C^{-1}\left|T_{2}^{\prime}\right| \leq & \|\phi\|_{W_{A_{1}}^{1, p}(X)}^{2}+\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\phi\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}^{2}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2} \\
& +\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right)^{2}\|\phi\|_{W_{A_{1}}^{1, p}(X)}^{2} \\
& +\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}\|a\|_{W_{A_{1}}^{1, p}(X)}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2} \\
& +\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right)_{W_{A_{1}}^{1, p}(X)}^{2}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\phi\|_{W_{A_{1}}^{1, p}(X)} \\
& +\|\phi\|_{W_{A_{1}}^{1, p}(X)}\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right)\left(\|\Phi\|_{W_{A_{1}}^{1, p}(X)}+\|\phi\|_{W_{A_{1}}^{1, p}(X)}\right) \\
& +\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}\|a\|_{W_{A_{1}}^{1, p}(X)} \\
& +\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\phi\|_{W_{A_{1}}^{1, p}(X)}\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}+\|a\|_{W_{A_{1}}^{1, p}(X)}\right),
\end{aligned}
$$

and

$$
\begin{aligned}
C^{-1}\left|T_{3}^{\prime}\right| \leq & \|\phi\|_{W_{A_{1}}^{1, p}(X)}^{2}+\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\phi\|_{W_{A_{1}}^{1, p}(X)} \\
& +\|\phi\|_{W_{A_{1}}^{1, p}(X)}^{4}+\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}\|\phi\|_{W_{A_{1}}^{1, p}(X)}^{2} \\
& +\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{3}\|\phi\|_{W_{A_{1}}^{1, p}(X)}+\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\phi\|_{W_{A_{1}}^{1, p}(X)}^{3} .
\end{aligned}
$$

Note that $4 d /(d+4) \geq 2 d /(d+2)$, so the condition $p \geq 2 d /(d+2)$ is assured by the stronger $p \geq 4 d /(d+4)$. Thus, $\mathscr{E}(A+a, \Phi+\phi)$ is a polynomial of degree four in the variable $(a, \phi) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$. This completes the proof of Proposition 3.1.1.

We now verify the formula (1.3.5) for the differential $\mathscr{E}^{\prime}(A, \Phi)$ and gradient $\mathscr{M}(A, \Phi)$.

Lemma 3.1.2 (Differential and gradient of the boson coupled Yang-Mills energy function). Assume the hypotheses of Proposition 3.1.1 with the dual Hölder exponent $p^{\prime}$ in the range $1<p^{\prime} \leq 4 d /(3 d-4)$ determined by $4 d /(d+4) \leq p<\infty$ and $1 / p+1 / p^{\prime}=1$. Then the expression for $\mathscr{E}^{\prime}(A, \Phi) \in\left(W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)^{*}$ and $\mathscr{M}(A, \Phi) \in W_{A_{1}}^{-1, p^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$ is given by (1.3.5), namely

$$
\begin{aligned}
\mathscr{E}^{\prime}(A, \Phi)(a, \phi)= & (\mathscr{M}(A, \Phi),(a, \phi))_{L^{2}(X)} \\
= & \left(d_{A}^{*} F_{A}, a\right)_{L^{2}(X)}+\operatorname{Re}\left(\nabla_{A}^{*} \nabla_{A} \Phi, \phi\right)_{L^{2}(X)}+\operatorname{Re}\left(\nabla_{A} \Phi, \varrho(a) \Phi\right)_{L^{2}(X)} \\
& -\operatorname{Re}(m \Phi, \phi)_{L^{2}(X)}-2 \operatorname{Re} \int_{X} s|\Phi|^{2}\langle\Phi, \phi\rangle d \operatorname{vol}_{g}, \\
& \forall(a, \phi) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) .
\end{aligned}
$$

Proof. We establish (1.3.5) by extracting the terms that are linear in $(a, \phi)$ from the expressions for $T_{i}^{\prime}$, for $i=1,2,3$, arising in the proof of Proposition 3.1.1 We compute $\mathscr{E}(A+t a, \Phi+t \phi)$ using the identities,

$$
\begin{aligned}
F_{A+t a} & =F_{A}+t d_{A} a+\frac{t^{2}}{2}[a, a] \\
\nabla_{A+t a}(\Phi+t \phi) & =\nabla_{A}(\Phi+t \phi)+\varrho(t a)(\Phi+t \phi),
\end{aligned}
$$

to obtain

$$
\begin{aligned}
\mathscr{E}(A+t a, \Phi+t \phi)= & \frac{1}{2} \int_{X}\left|F_{A}+t d_{A} a+\frac{t^{2}}{2}[a, a]\right|^{2} d \operatorname{vol}_{g} \\
& \left.\left.+\frac{1}{2} \int_{X} \right\rvert\, \nabla_{A} \Phi+t\left(\nabla_{A} \phi+\varrho(a) \Phi\right)+t^{2} \varrho(a) \phi\right)\left.\right|^{2} d \operatorname{vol}_{g} \\
& -\frac{1}{2} \int_{X}\left(m|\Phi+t \phi|^{2}+s|\Phi+t \phi|^{4}\right) d \operatorname{vol}_{g}
\end{aligned}
$$

that is,

$$
\begin{aligned}
\mathscr{E}(A+t a, \Phi+t \phi)= & \frac{1}{2} \int_{X}\left(\left|F_{A}\right|^{2}+2 t\left\langle F_{A}, d_{A} a\right\rangle\right) d \operatorname{vol}_{g} \\
& +\frac{1}{2} \int_{X}\left(\left|\nabla_{A} \Phi\right|^{2}+2 t \operatorname{Re}\left\langle\nabla_{A} \Phi, \nabla_{A} \phi+\varrho(a) \Phi\right\rangle_{L^{2}}\right) d \operatorname{vol}_{g} \\
& -\frac{1}{2} \int_{X} m\left(|\Phi|^{2}+2 t \operatorname{Re}\langle\Phi, \phi\rangle\right) d \operatorname{vol}_{g} \\
& -\frac{1}{2} \int_{X} s\left(|\Phi|^{4}+4 t|\Phi|^{2} \operatorname{Re}\langle\Phi, \phi\rangle\right) d \operatorname{vol}_{g}+\text { higher powers of } t .
\end{aligned}
$$

This yields the expression (1.3.5) for $\mathscr{E}^{\prime}(A, \Phi)(a, \phi)$ and completes the proof of Lemma 3.1.2

It is convenient to define, for $p \in(1, \infty)$ and dual exponent $p^{\prime} \in(1, \infty)$ determined by $1 / p+1 / p^{\prime}=1$, the following Banach spaces,

$$
\begin{gather*}
\mathfrak{X}:=W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \quad \text { and } \quad \mathfrak{X}^{*}=W_{A_{1}}^{-1, p^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right),  \tag{3.1.3}\\
\tilde{\mathfrak{X}}:=W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) .
\end{gather*}
$$

Note that for $p \geq 2$, the inclusion, $\tilde{\mathfrak{X}} \subset \mathfrak{X}^{*}$, is a continuous embedding of Banach spaces.

The expression (1.3.5) defines the gradient as a map, $\mathscr{M}:\left(A_{1}, 0\right)+\mathfrak{X} \rightarrow \tilde{\mathfrak{X}}$, where $\mathscr{M}(A, \Phi) \in \tilde{\mathfrak{X}}$ acts on $(b, \varphi) \in \mathfrak{X}$ by $L^{2}$ inner product via the inclusion $\tilde{\mathfrak{X}} \subset \mathfrak{X}^{*}$. We now have the

Proposition 3.1.3 (Analyticity of the gradient map for the boson coupled Yang-Mills energy function). Assume the hypotheses of Proposition 3.1.1 and, in addition, that $p \geq d / 2$ when $d \geq 3$ and $p>4 / 3$ when $d=2$. Then

$$
\mathscr{M}:\left(A_{1}, 0\right)+\mathfrak{X} \rightarrow \tilde{\mathfrak{X}}
$$

is a real analytic map, where $\mathscr{M}$ is as in (1.3.5) and $\mathfrak{X}$ and $\tilde{\mathfrak{X}}$ are as in (3.1.3).
Proof. As usual, given $p \in(1, \infty)$, we let $p^{\prime} \in(1, \infty)$ denote the dual Hölder exponent defined by $1 / p+1 / p^{\prime}=1$. By [4, Theorem 4.12], we have a continuous Sobolev embedding, $W^{1, p^{\prime}}(X ; \mathbb{C}) \subset L^{r}(X ; \mathbb{C})$, when
(1) $1<p^{\prime}<d$ and $1 \leq r \leq\left(p^{\prime}\right)^{*}:=d p^{\prime} /\left(d-p^{\prime}\right)$, or
(2) $p^{\prime}=d$ and $1 \leq r<\infty$, or
(3) $d<p^{\prime}<\infty$ and $1 \leq r \leq \infty$.

In each case, we obtain a continuous Sobolev embedding, $L^{r^{\prime}}(X ; \mathbb{C}) \subset W^{-1, p}(X ; \mathbb{C})$, by duality and density of the continuous embedding, $W^{1, p^{\prime}}(X ; \mathbb{C}) \subset L^{r}(X ; \mathbb{C})$, where $1 / r+1 / r^{\prime}=1$ and $1 / p+1 / p^{\prime}=1$. We shall require $p^{\prime}, r<\infty$ in order to appeal to the dualities, $\left(W^{1, p^{\prime}}(X ; \mathbb{C})\right)^{*}=W^{-1, p}(X ; \mathbb{C})$ and $\left(L^{r}(X ; \mathbb{C})\right)^{*}=L^{r^{\prime}}(X ; \mathbb{C})$; these dualities fail when $p^{\prime}=\infty$ (and thus $p=1$ ) or $r=\infty$ (and thus $r^{\prime}=1$ ).

If we write $A=A_{1}+a_{1}$, for $a_{1} \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then from the formula (1.3.5) for $\mathscr{M}(A, \Phi)$ we have the formal expression,

$$
\begin{align*}
\mathscr{M}\left(A_{1}+a_{1}, \Phi\right)= & d_{A_{1}}^{*} d_{A_{1}} a_{1}+\nabla_{A_{1}}^{*} \nabla_{A_{1}} \Phi+d_{A_{1}}^{*} F_{A_{1}}+F_{A_{1}} \times a_{1} \\
& +\nabla_{A_{1}} \Phi \times a_{1}+\nabla_{A_{1}} a_{1} \times \Phi+\nabla_{A_{1}} a_{1} \times a_{1} \\
& +\nabla_{A_{1}} \Phi \times \Phi-\left(m+2 s|\Phi|^{2}\right) \Phi  \tag{3.1.4}\\
& +a_{1} \times a_{1} \times \Phi+a_{1} \times \Phi \times \Phi+a_{1} \times a_{1} \times a_{1} .
\end{align*}
$$

Observe that

$$
\begin{aligned}
W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) & \ni\left(a_{1}, \Phi\right) \\
& \rightarrow d_{A_{1}}^{*} d_{A_{1}} a_{1}+\nabla_{A_{1}}^{*} \nabla_{A_{1}} \Phi \in W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

is analytic and the term $d_{A_{1}}^{*} F_{A_{1}}$ is constant with respect to $\left(a_{1}, \Phi\right)$. Therefore, to prove that $\mathscr{M}$ is analytic, it suffices to prove the

Claim 3.1.4. Continue the preceding notation. Then

$$
\begin{align*}
& W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \ni\left(a_{1}, \Phi\right)  \tag{3.1.5}\\
& \mapsto \mathscr{M}\left(A_{1}+a_{1}, \Phi\right)-d_{A_{1}}^{*} d_{A_{1}} a_{1}+\nabla_{A_{1}}^{*} \nabla_{A_{1}} \Phi-d_{A_{1}}^{*} F_{A_{1}} \\
& \quad \in L^{r^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{align*}
$$

is a cubic polynomial in $\left(a_{1}, \Phi\right)$ and its first-order covariant derivatives with respect to $\nabla_{A_{1}}$, with universal coefficients (depending at most on $g$ and $G$ ).

Proof of Claim 3.1.4. We compute an $L^{r^{\prime}}$ bound for each term in equation (3.1.4) for $\mathscr{M}(A, \Phi)$; we consider the cases $p^{\prime}<d, p^{\prime}=d$, and $p^{\prime}>d$ separately, recalling that $A=A_{1}+a_{1}$.

STEP 1 ( $L^{r^{\prime}}$ estimates for $F_{A_{1}} \times a_{1}$ and $\left.\nabla_{A_{1}} \Phi \times a_{1}\right)$. We claim that

$$
\begin{align*}
\left\|F_{A_{1}} \times a_{1}\right\|_{L^{r^{\prime}}(X)} & \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}  \tag{3.1.6}\\
\left\|\nabla_{A_{1}} \Phi \times a_{1}\right\|_{L^{r^{\prime}}(X)} & \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)} \tag{3.1.7}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$.
Case $1\left(p^{\prime}<d\right)$. We choose $r=\left(p^{\prime}\right)^{*}=d p^{\prime} /\left(d-p^{\prime}\right)$, where $p^{\prime}=p /(p-1)$, so $r=(d p /(p-1)) /(d-p /(p-1))=d p /(d(p-1)-p)$. Since $p^{\prime}=p /(p-1)<d$ for Case 1 then $p<d p-d$ or $p(d-1)>d$ or $p>d /(d-1)$. By 4. Theorem 4.12] we have continuous Sobolev embeddings, $a) W^{1, p}(X) \subset L^{p^{*}}(X)$, for $p^{*}=d p /(d-p)$ if $p<d ; b) W^{1, p}(X) \subset L^{q}(X)$ for any $q \in[1, \infty)$ if $p=d$; and c) $W^{1, p}(X) \subset L^{\infty}(X)$ if $p>d$. We consider each of these three subcases in turn.

Consider Subcase (11a), so $d /(d-1)<p<d$, which forces $d \geq 3$. In order to have a continuous multiplication map, $L^{p}(X) \times L^{p^{*}}(X) \rightarrow L^{r^{\prime}}(X)$, we require that the inequality $1 / p+1 / p^{*} \leq 1 / r^{\prime}$ holds, that is

$$
1 / p+1 /(d p /(d-p)) \leq 1 / r^{\prime}=1-1 / r=1-(d p-d-p) / d p
$$

or equivalently,

$$
1 / p+(d-p) /(d p) \leq 1-(1-1 / p-1 / d)
$$

namely,

$$
1 / p+1 / p-1 / d \leq 1 / p+1 / d
$$

Therefore, $1 / p \leq 2 / d$ or $p \geq d / 2$, as we assumed in our hypotheses. Thus,

$$
\left\|F_{A_{1}} \times a_{1}\right\|_{L^{r^{\prime}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{L^{p^{*}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

which yields the first inequality in (3.1.6) for all $d \geq 3$ in this subcase.
Consider subcase (11D), so $p>d /(d-1)$ and $p=d$. Because $p>d /(d-1)$ in this subcase, then $d$ must obey $d>d /(d-1)$ or $d-1>1$ or $d \geq 3$ for this subcase. (If $d=2$, then $p=d=2$ forces $p^{\prime}=2$, so the subcase $p^{\prime}<d$ and $p=d$ cannot occur.) We now only need a continuous multiplication map, $L^{d}(X) \times L^{q}(X) \rightarrow L^{r^{\prime}}(X)$, for large enough $q \in[1, \infty)$ and this requires only that $r^{\prime}<d$, that is, $1 / r^{\prime}=1-1 / r>1 / d$ or

$$
1-(d p-d-p) / d p=1 / p+1 / d=2 / d>1 / p=1 / d
$$

which holds for all positive $d$. Thus,

$$
\left\|F_{A_{1}} \times a_{1}\right\|_{L^{r^{\prime}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{L^{q}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

which yields (3.1.6) for all $d \geq 3$ in this subcase.
Consider Subcase (11四), so $p>d /(d-1)$ and $p>d$. But $p>d /(d-1)$ holds for any $d \geq 2$ for this subcase since $d \geq d /(d-1)$ for any $d \geq 2$. We now only need a continuous multiplication map, $L^{p}(X) \times L^{\infty}(X) \rightarrow L^{r^{\prime}}(X)$, and this requires only that $r^{\prime} \leq p$, that is, $1 / r^{\prime}=1-1 / r \geq 1 / p$ or

$$
1-(d p-d-p) / d p=1 / p+1 / d \geq 1 / p
$$

which holds for all positive $d$. Thus,

$$
\left\|F_{A_{1}} \times a_{1}\right\|_{L^{r^{\prime}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{L^{\infty}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

which yields (3.1.6) for all $d \geq 2$ in this subcase.

Case $2\left(p^{\prime}=d\right)$. Since $p^{\prime}=p /(p-1)=d$ for this case, then $p=d p-d$ or $p(d-1)=d$ and so $p=d /(d-1)$. But our hypotheses in Proposition 3.1.3 requires $p \geq 4 d /(d+4)$ when $d \geq 3$. Observe that

$$
4 d /(d+4) \geq d /(d-1) \Longleftrightarrow 4(d-1) \geq d+4 \Longleftrightarrow 3 d \geq 8 \Longleftrightarrow d \geq 8 / 3
$$

that is, $p \geq 4 d /(d+4)>d /(d-1)$ when $d \geq 3$. Hence, the case $p^{\prime}=d$ and $p=d /(d-1)$ is excluded by our hypotheses unless $d=2$ and thus $p^{\prime}=2=p$.

When $d=2$, we have a continuous Sobolev embedding, $W^{1,2}(X) \subset L^{q}(X)$ for any $q \in[1, \infty)$ by [4, Theorem 4.12]. Hence, we obtain a continuous multiplication map, $L^{2}(X) \times L^{q}(X) \rightarrow L^{r^{\prime}}(X)$, for any $r^{\prime}>2$ and large enough $q \in[1, \infty)$. Thus,

$$
\left\|F_{A_{1}} \times a_{1}\right\|_{L^{r^{\prime}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{2}(X)}\left\|a_{1}\right\|_{L^{q}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{2}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1,2}(X)}
$$

which yields (3.1.6) for the case $p^{\prime}=d=2=p$.
Case $3\left(p^{\prime}>d\right)$. We can again choose $r \in[1, \infty)$ arbitrarily large in the continuous Sobolev embedding, $W^{1, p^{\prime}}(X) \subset L^{r}(X)$, and so we may choose $r^{\prime} \in(1, \infty]$ arbitrarily small or equivalently, $1 / r^{\prime} \in[0,1$ ) arbitrarily close to 1 . (We refrain from choosing $r=\infty$ because continuity of the Sobolev embedding, $W^{1, p^{\prime}}(X) \subset L^{\infty}(X)$, does not imply continuity of $L^{1}(X) \subset W^{-1, p}(X)$.)

Since $p^{\prime}=p /(p-1)>d$ for this case, then $p>d p-d$ or $p(d-1)<d$ and so $p<d /(d-1) \leq d$. We therefore have $p^{*}=d p /(d-p)$ and a continuous Sobolev embedding, $W^{1, p}(X) \subset L^{p^{*}}(X)$. In order to have a continuous multiplication map, $L^{p}(X) \times L^{p^{*}}(X) \rightarrow L^{r^{\prime}}(X)$, we require that the inequality $1 / p+1 / p^{*} \leq 1 / r^{\prime}$ holds for some $r \in[1, \infty)$ and $r^{\prime}=r /(r-1) \in(1, \infty]$. Such a choice of $r^{\prime}$ will be possible if and only if the strict inequality below holds,

$$
1 / p+1 / p^{*}=1 / p+(d-p) /(d p)=2 / p-1 / d<1
$$

that is, if and only if $2 / p<1+1 / d=(d+1) / d$ or equivalently, $p>2 d /(d+1)$. But our hypotheses in Proposition 3.1.3 requires $p \geq 4 d /(d+4)$ when $d \geq 3$ and $p>4 / 3$ when $d=2$. Now $4 d /(d+4)>2 d /(d+1)$ when $d \geq 3$ and $4 / 3=2 d /(d+1)$ when $d=2$, so our hypotheses in Proposition 3.1.3 ensure that $p>2 d /(d+1)$ all $d \geq 2$. Hence, we may choose a large enough $r \in[1, \infty)$ and thus small enough $r^{\prime} \in(1, \infty]$ to give $1 / p+1 / p^{*} \leq 1 / r^{\prime}$.

Hence, the map, $L^{p}(X) \times L^{p^{*}}(X) \rightarrow L^{r^{\prime}}(X)$, is continuous for small enough $r^{\prime}>1$ and

$$
\left\|F_{A_{1}} \times a_{1}\right\|_{L^{r^{\prime}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{L^{p^{*}}(X)} \leq z\left\|F_{A_{1}}\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

which again yields (3.1.6) for this case.
An argument identical to that for (3.1.6) gives

$$
\left\|\nabla_{A_{1}} \Phi \times a_{1}\right\|_{L^{r^{\prime}}(X)} \leq z\left\|\nabla_{A_{1}} \Phi\right\|_{L^{p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}
$$

proving (3.1.7). This completes Step 1 .

STEP $2\left(L^{r^{\prime}}\right.$ estimates for $\nabla_{A_{1}} \Phi \times \Phi$ and $\nabla_{A_{1}} a_{1} \times \Phi$ and $\left.\nabla_{A_{1}} a_{1} \times a_{1}\right)$. We claim that

$$
\begin{align*}
\left\|\nabla_{A_{1}} \Phi \times \Phi\right\|_{L^{r^{\prime}}(X)} & \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}  \tag{3.1.8}\\
\left\|\nabla_{A_{1}} a_{1} \times \Phi\right\|_{L^{r^{\prime}}(X)} & \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}  \tag{3.1.9}\\
\left\|\nabla_{A_{1}} a_{1} \times a_{1}\right\|_{L^{r^{\prime}}(X)} & \leq z\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}^{2} \tag{3.1.10}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$.
From the proof of (3.1.6) in Step 1, we have

$$
\left\|\Phi \times \nabla_{A_{1}} \Phi\right\|_{L^{r^{\prime}}(X)} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\left\|\nabla_{A_{1}} \Phi\right\|_{L^{p}(X)}
$$

and this gives (3.1.8) ; identical arguments give (3.1.9) and (3.1.10). This completes Step 2.

Note that $\|m\|_{C(X)}$ and $\|s\|_{C(X)}$ are finite by hypothesis.
Step 3 ( $L^{r^{\prime}}$ estimate for $\left.m \Phi\right)$. We have

$$
\|m \Phi\|_{L^{r^{\prime}}(X)} \leq\|m\|_{C(X)}\|\Phi\|_{L^{r^{\prime}}(X)},
$$

and because $r^{\prime} \leq p$ by inspection of each of the three cases, $p^{\prime}<d$ and $p^{\prime}=d$ and $p^{\prime}>d$, and subcases ( $p<d$ and $p=d$ and $p>d$ where applicable) we obtain

$$
\begin{equation*}
\|m \Phi\|_{L^{r^{\prime}}(X)} \leq z\|m\|_{C(X)}\|\Phi\|_{L^{p}(X)} \tag{3.1.11}
\end{equation*}
$$

as desired.
Step $4\left(L^{r^{\prime}}\right.$ estimates for $\left.s|\Phi|^{2} \Phi\right)$. We claim that

$$
\begin{equation*}
\left\|s|\Phi|^{2} \Phi\right\|_{L^{r^{\prime}(X)}} \leq z\|s\|_{C(X)}\|\Phi\|_{W_{A_{1}(X)}^{1, p}}^{3} \tag{3.1.12}
\end{equation*}
$$

where $z=z(g, p) \in[1, \infty)$.
From the proof of (3.1.6) in Step 1, we have

$$
\left\|s|\Phi|^{2} \Phi\right\|_{L^{r^{\prime}}(X)} \leq z\|s\|_{C(X)}\left\||\Phi|^{2}\right\|_{L^{p}(X)}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}
$$

Moreover, we have

$$
\left\||\Phi|^{2}\right\|_{L^{p}(X)}=\|\Phi\|_{L^{2 p}(X)}^{2} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}
$$

Combining the preceding inequalities yields (3.1.12).
Step $5\left(L^{r^{\prime}}\right.$ estimates for $a_{1} \times a_{1} \times a_{1}$ and $a_{1} \times a_{1} \times \Phi$ and $\left.a_{1} \times \Phi \times \Phi\right)$. We claim that

$$
\begin{align*}
\left\|a_{1} \times \Phi \times \Phi\right\|_{L^{r^{\prime}}(X)} & \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}  \tag{3.1.13}\\
\left\|a_{1} \times a_{1} \times \Phi\right\|_{L^{r^{\prime}}(X)} & \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}^{2},  \tag{3.1.14}\\
\left\|a_{1} \times a_{1} \times a_{1}\right\|_{L^{r^{\prime}}(X)} & \leq z\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}^{3} \tag{3.1.15}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$.
From the proof of (3.1.6) in Step 1, we have

$$
\left\|a_{1} \times \Phi \times \Phi\right\|_{L^{r^{\prime}}(X)} \leq z\left\|a_{1}\right\|_{W_{A_{1}}^{1, p}(X)}\|\Phi \times \Phi\|_{L^{p}(X)}
$$

Continuity of the multiplication, $L^{2 p}(X) \times L^{2 p}(X) \rightarrow L^{p}(X)$, and continuity of the Sobolev embedding, $W^{1, p}(X) \subset L^{2 p}(X)$, valid for any $p \geq d / 2$ (which we assume by hypothesis), gives

$$
\|\Phi \times \Phi\|_{L^{p}(X)} \leq z\|\Phi\|_{L^{2 p}(X)} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}
$$

Combining the preceding inequalities yields (3.1.13). The proofs of the estimates (3.1.14) and (3.1.15) are the same as that of (3.1.13). This completes Step 5 ,

The estimates obtained in each of the preceding steps show that the map (3.1.5) has the properties asserted in the statement of Claim 3.1.4 and this completes its proof.

In particular, Claim 3.1.4 implies that

$$
\mathscr{M}\left(A_{1}+a_{1}, \Phi\right) \in L^{r^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is a continuous cubic polynomial in $\left(a_{1}, \Phi\right) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, with universal coefficients depending at most on $g$ and $G$. Because

$$
L^{r^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \subset W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is a continuous Sobolev embedding by our choice of $r^{\prime}$, we see that

$$
\mathscr{M}\left(A_{1}+a_{1}, \Phi\right) \in W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is a cubic polynomial in $\left(a_{1}, \Phi\right) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, again with universal coefficients depending at most on $g$ and $G$. This completes the proof of Proposition 3.1.3.

### 3.1.2. Fredholm and index properties of the Hessian operator for the

 boson coupled Yang-Mills energy function on the Sobolev space of pairs. Consider the Hessian map, $\mathscr{M}^{\prime}(A, \Phi): \mathfrak{X} \rightarrow \mathfrak{X}^{*}$.Lemma 3.1.5 (Hessian and Hessian operator for the boson coupled Yang-Mills energy function). Assume the hypotheses of Proposition 3.1.1. Then we have the schematic formula for $\mathscr{M}^{\prime}(A, \Phi) \in \mathscr{L}\left(\mathfrak{X}, \mathfrak{X}^{*}\right)$ and $\mathscr{E}^{\prime \prime}(A, \Phi) \in(\mathfrak{X} \times \mathfrak{X})^{*}$ given by

$$
\begin{align*}
\mathscr{M}^{\prime}(A, \Phi)(a, \phi)= & d_{A}^{*} d_{A} a+\nabla_{A}^{*} \nabla_{A} \phi+F_{A} \times a+\nabla_{A}^{*}(\varrho(a) \Phi)+\Phi \times \nabla_{A} \phi \\
& -\rho(a)^{*} \nabla_{A} \Phi+\nabla_{A} \Phi \times \phi+\varrho(a) \Phi \times \Phi  \tag{3.1.16}\\
& -\left(m+2 s|\Phi|^{2}\right) \phi-4 s\langle\Phi, \phi\rangle \Phi, \quad \forall(a, \phi) \in \mathfrak{X},
\end{align*}
$$

where

$$
\begin{align*}
\mathscr{E}^{\prime \prime}(A, \Phi)(a, \phi)(b, \varphi)=\left((b, \varphi), \mathscr{M}^{\prime}(A, \Phi)(a, \phi)\right)_{L^{2}(X)} &  \tag{3.1.17}\\
& \forall(a, \phi),(b, \varphi) \in \mathfrak{X}
\end{align*}
$$

Proof. Let $\left(a_{i}, \phi_{i}\right) \in \mathfrak{X}$, for $i=1,2$. We compute the terms in

$$
\left(\mathscr{M}\left(A+t a_{2}, \Phi+t \phi_{2}\right),\left(a_{1}, \phi_{1}\right)\right)_{L^{2}(X)}
$$

that are linear in $t$ using the expression (1.3.5) for the gradient. First,

$$
\begin{aligned}
\left(F_{A+t a_{2}}, d_{A+t a_{2}} a_{1}\right)_{L^{2}(X)}= & \left(F_{A}+t d_{A} a_{2}+\frac{t^{2}}{2}\left[a_{2}, a_{2}\right], d_{A} a_{1}+t\left[a_{2}, a_{1}\right]\right)_{L^{2}(X)} \\
= & \left(F_{A}, d_{A} a_{1}\right)_{L^{2}(X)}+t\left(F_{A},\left[a_{2}, a_{1}\right]\right)_{L^{2}(X)} \\
& +t\left(d_{A} a_{2}, d_{A} a_{1}\right)_{L^{2}(X)}+O\left(t^{2}\right)
\end{aligned}
$$

Second,

$$
\begin{aligned}
& \left(\nabla_{A+t a_{2}}\left(\Phi+t \phi_{2}\right), \nabla_{A+t a_{2}} \phi_{1}\right)_{L^{2}(X)} \\
& \quad=\left(\left(\nabla_{A}+t \varrho\left(a_{2}\right)\right)\left(\Phi+t \phi_{2}\right), \nabla_{A} \phi_{1}+t \varrho\left(a_{2}\right) \phi_{1}\right)_{L^{2}(X)} \\
& =\left(\nabla_{A} \Phi, \nabla_{A} \phi_{1}\right)_{L^{2}(X)}+t\left(\nabla_{A} \phi_{2}, \nabla_{A} \phi_{1}\right)_{L^{2}(X)}+t\left(\varrho\left(a_{2}\right) \Phi, \nabla_{A} \phi_{1}\right)_{L^{2}(X)} \\
& \quad+t\left(\nabla_{A} \Phi, \varrho\left(a_{2}\right) \phi_{1}\right)_{L^{2}(X)}+O\left(t^{2}\right)
\end{aligned}
$$

Third,

$$
\begin{aligned}
& \left(\nabla_{A+t a_{2}}\left(\Phi+t \phi_{2}\right), \varrho\left(a_{1}\right)\left(\Phi+t \phi_{2}\right)\right)_{L^{2}(X)} \\
& =\left(\left(\nabla_{A}+t \varrho\left(a_{2}\right)\right)\left(\Phi+t \phi_{2}\right), \varrho\left(a_{1}\right)\left(\Phi+t \phi_{2}\right)\right)_{L^{2}(X)} \\
& =\left(\nabla_{A} \Phi, \varrho\left(a_{1}\right) \Phi\right)_{L^{2}(X)}+t\left(\nabla_{A} \phi_{2}, \varrho\left(a_{1}\right) \Phi\right)_{L^{2}(X)}+t\left(\varrho\left(a_{2}\right) \Phi, \varrho\left(a_{1}\right) \Phi\right)_{L^{2}(X)} \\
& \quad+t\left(\nabla_{A} \Phi, \varrho\left(a_{1}\right) \phi_{2}\right)_{L^{2}(X)}+O\left(t^{2}\right)
\end{aligned}
$$

Fourth,

$$
\left(m\left(\Phi+t \phi_{2}\right), \phi_{1}\right)_{L^{2}(X)}=\left(m \Phi, \phi_{1}\right)_{L^{2}(X)}+t\left(m \phi_{2}, \phi_{1}\right)_{L^{2}(X)}
$$

Fifth,

$$
\begin{array}{rl}
\int_{X} & s\left|\Phi+t \phi_{2}\right|^{2}\left\langle\Phi+t \phi_{2}, \phi_{1}\right\rangle d \operatorname{vol}_{g} \\
& =\int_{X} s\left(|\Phi|^{2}+2 t \operatorname{Re}\left\langle\Phi, \phi_{2}\right\rangle+t^{2}\left|\phi_{2}\right|^{2}\right)\left\langle\Phi+t \phi_{2}, \phi_{1}\right\rangle d \operatorname{vol}_{g} \\
= & \int_{X} s|\Phi|^{2}\left\langle\Phi, \phi_{1}\right\rangle d \operatorname{vol}_{g} \\
& \quad+t \int_{X}\left(s|\Phi|^{2}\left\langle\phi_{2}, \phi_{1}\right\rangle+2 s \operatorname{Re}\left\langle\Phi, \phi_{2}\right\rangle\left\langle\Phi, \phi_{1}\right\rangle\right) d \operatorname{vol}_{g}+O\left(t^{2}\right) .
\end{array}
$$

By subtracting $\left(\mathscr{M}^{\prime}(A, \Phi)\left(a_{2}, \phi_{2}\right),\left(a_{1}, \phi_{1}\right)\right)_{L^{2}(X)}$, collecting all the first-order terms in $t$, and reversing the roles of $\left(a_{1}, \phi_{1}\right)$ and $\left(a_{2}, \phi_{2}\right)$, we see that

$$
\begin{aligned}
& \left(\mathscr{M}^{\prime}(A, \Phi)\left(a_{1}, \phi_{1}\right),\left(a_{2}, \phi_{2}\right)\right)_{L^{2}(X)} \\
& \quad=\left(d_{A} a_{1}, d_{A} a_{2}\right)_{L^{2}(X)}+2\left(F_{A},\left[a_{1}, a_{2}\right]\right)_{L^{2}(X)} \\
& \quad+\operatorname{Re}\left(\nabla_{A} \phi_{1}, \nabla_{A} \phi_{2}\right)_{L^{2}(X)} \\
& \quad+\operatorname{Re}\left(\left(\rho\left(a_{1}\right) \Phi, \nabla_{A} \phi_{2}\right)_{L^{2}(X)}+\left(\rho\left(a_{2}\right) \Phi, \nabla_{A} \phi_{1}\right)_{L^{2}(X)}\right) \\
& \quad+\operatorname{Re}\left(\nabla_{A} \Phi, \rho\left(a_{1}\right) \phi_{2}+\rho\left(a_{2}\right) \phi_{1}\right)_{L^{2}(X)} \\
& \quad+\operatorname{Re}\left(\rho\left(a_{1}\right) \Phi, \rho\left(a_{2}\right) \Phi\right)_{L^{2}(X)} \\
& \quad-\operatorname{Re} \int_{X}\left(\left(m+2 s|\Phi|^{2}\right)\left\langle\phi_{1}, \phi_{2}\right\rangle+4 s\left\langle\Phi, \phi_{1}\right\rangle\left\langle\Phi, \phi_{2}\right\rangle\right) d \operatorname{vol}_{g}
\end{aligned}
$$

By now viewing $\mathscr{M}^{\prime}(A, \Phi)\left(a_{1}, \phi_{1}\right)$ as an element of $\mathfrak{X}^{*}$, we obtain the expression (3.1.16).

When $(A, \Phi)$ is a $C^{\infty}$ pair, we shall need to compare $\mathscr{M}^{\prime}(A, \Phi)$ with the $L^{2}$ -self-adjoint, second-order partial differential operator,
(3.1.18) $\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}: C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, in order to prove that $\mathscr{M}^{\prime}(A, \Phi)$ is Fredholm with index zero upon restriction to

$$
\begin{equation*}
\mathscr{X}:=\operatorname{Ker}\left(d_{A, \Phi}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{p}(X ; \operatorname{ad} P)\right) \tag{3.1.19}
\end{equation*}
$$

We recall from (2.9.2) that

$$
d_{A, \Phi} \xi=\left(d_{A} \xi,-\xi \Phi\right), \quad \forall \xi \in C^{\infty}(X ; \operatorname{ad} P),
$$

with $L^{2}$-adjoint,

$$
\begin{equation*}
d_{A, \Phi}^{*}(a, \phi)=d_{A}^{*} a-\langle\phi, \cdot \Phi\rangle^{*}, \quad \forall(a, \phi) \in C^{\infty}\left(\Lambda^{1} \otimes \operatorname{ad} P \oplus E\right), \tag{3.1.20}
\end{equation*}
$$

for every $(a, \phi) \in C^{\infty}\left(\Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, where the section $\langle\phi, \cdot \Phi\rangle^{*}$ of ad $P$ is defined by

$$
\left(\langle\phi, \cdot \Phi\rangle^{*}, \xi\right)_{L^{2}(X)}=(\phi, \xi \Phi)_{L^{2}(X)}, \quad \forall \xi \in C^{\infty}(X ; \operatorname{ad} P) .
$$

According to Lemma 2.9.7, the operator,

$$
d_{A, \Phi}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{p}(X ; \operatorname{ad} P),
$$

is bounded when $(A, \Phi)$ is a $W^{1, q}$ pair with $q \geq d / 2$ and $p$ obeys $d / 2 \leq p \leq q$; therefore $\mathscr{X}$ in (3.1.19) is a Banach space since it is a closed subspace of the Banach space $\mathfrak{X}=W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$.

When $(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \mathrm{ad} P \oplus E\right)$, the expression (3.1.16) yields, after formally expanding $\nabla_{A}^{*}(\varrho(a) \Phi)=\nabla_{A} a \times \Phi+a \times \nabla_{A} \Phi+a \times \Phi$,

$$
\begin{align*}
& \mathscr{M}^{\prime}(A, \Phi)(a, \phi) \\
& \quad=d_{A}^{*} d_{A} a+\nabla_{A}^{*} \nabla_{A} \phi+F_{A} \times a+\nabla_{A} a \times \Phi+a \times \nabla_{A} \Phi+\Phi \times \nabla_{A} \phi  \tag{3.1.21}\\
& \quad+\varrho(a) \Phi+\nabla_{A} \Phi \times \phi+\varrho(a) \Phi \times \Phi-\left(m+2 s|\Phi|^{2}\right) \phi-4 s\langle\Phi, \phi\rangle \Phi .
\end{align*}
$$

To determine the Fredholm property and index of $\mathscr{M}^{\prime}(A, \Phi)$ upon restriction to a Coulomb-gauge slice, we shall need the following consequence of Theorem A. 1

Proposition 3.1.6 (Fredholm and index zero property of the augmented Hessian operator). Let ( $X, g$ ) be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finitedimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$, and $A_{1}$ be a $C^{\infty}$ reference connection on $P$. If $(A, \Phi)$ is a $C^{\infty}$ pair on $(P, E)$ and $k \in \mathbb{Z}$ is an integer and $1<p<\infty$, then the following operator is Fredholm with index zero,

$$
\begin{align*}
\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}: W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad}\right. & P \oplus E)  \tag{3.1.22}\\
& \rightarrow W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) .
\end{align*}
$$

Proof. We can compare the principal symbol of the connection Laplacian, $\nabla_{A}^{*} \nabla_{A}$, and Hodge Laplacian, $\Delta_{A}=d_{A}^{*} d_{A}+d_{A} d_{A}^{*}$ (2.2.1), on $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)=$ $\Omega^{1}(X ; \operatorname{ad} P)$ using the Bochner-Weitzenböck formula [8, 9, [43, Appendix C], 83, Appendix II] and [123]. From [83, Corollary II.2], one has

$$
\begin{equation*}
\Delta_{A} a=\nabla_{A}^{*} \nabla_{A} a+\operatorname{Ric}_{g} \times a+F_{A} \times a, \quad \forall a \in \Omega^{1}(X ; \operatorname{ad} P), \tag{3.1.23}
\end{equation*}
$$

where $\operatorname{Ric}_{g}$ denotes the Ricci curvature tensor of the Riemannian metric $g$ on the manifold $X$ of dimension $d \geq 2$ and here we employ ' $\times$ ' to denote any universal bilinear expression with constant coefficients depending at most on the Lie group, $G$, or metric $g$. In particular, $\Delta_{A}$ is a second-order, elliptic partial differential operator on $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ with $C^{\infty}$ coefficients and scalar principal symbol given by the Riemannian metric $g$ on $T^{*} M$.

From the expressions (3.1.21) for $\mathscr{M}^{\prime}(A, \Phi)$ and (2.9.2) for $d_{A, \Phi}$ and (3.1.20) for $d_{A, \Phi}^{*}$, we see that

$$
\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}=\Delta_{A} \oplus \nabla_{A}^{*} \nabla_{A}+\text { Lower-order terms. }
$$

Thus, $\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}$ is an elliptic, second-order partial differential operator on $C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$ with $C^{\infty}$ coefficients and principal symbol given by the Riemannian metric $g$ on $T^{*} M$.

Equation (3.1.16) for $\mathscr{M}^{\prime}(A, \Phi)$ implies, for $(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, that

$$
\begin{aligned}
& \left(\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}-\left(\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}\right)^{*}\right)(a, \phi) \\
& \quad=F_{A} \times a+\nabla_{A}^{*}(\varrho(a) \Phi)+\Phi \times \nabla_{A} \phi-\rho(a)^{*} \nabla_{A} \Phi+\nabla_{A} \Phi \times \phi+\varrho(a) \Phi \times \Phi \\
& \quad-\left(m+2 s|\Phi|^{2}\right) \phi-4 s\langle\Phi, \phi\rangle \Phi .
\end{aligned}
$$

Consequently, the following expression defines a first-order partial differential operator,

$$
\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}-\left(\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}\right)^{*}
$$

and TheoremA.1implies that the operator (3.1.22) is Fredholm with index zero.
3.1.3. Gradient map and Hessian operator for the boson coupled Yang-Mills energy on a Coulomb-gauge slice. Suppose that $\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair on $(P, E)$, recall that $\mathfrak{X}$ is as in (3.1.3) and $\mathscr{X}$ is as in (3.1.19), and define

$$
\begin{align*}
& \tilde{\mathfrak{X}}:=W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)  \tag{3.1.24}\\
& \tilde{\mathscr{X}}:=\operatorname{Ker}\left(d_{A_{\infty}, \Phi_{\infty}}^{*}: W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right.  \tag{3.1.25}\\
&\left.\rightarrow W_{A_{1}}^{-2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)
\end{align*}
$$

Because $\left(W^{1, p}(X ; \mathbb{C})\right)^{*}=W^{-1, p^{\prime}}(X ; \mathbb{C})$ and $p \geq p^{\prime}$ for all $p \geq 2$, we see that

$$
\tilde{\mathscr{X}} \subset \mathscr{X}^{*}, \quad \forall p \geq 2
$$

For a $C^{\infty}$ pair $(A, \Phi)$ on $(P, E)$, the Hessian operator, $\mathscr{M}^{\prime}(A, \Phi): \mathfrak{X} \rightarrow \mathfrak{X}^{*}$, is defined by the schematic expression (3.1.16) and related to the Hessian, $\mathscr{E}^{\prime \prime}(A, \Phi)$ : $\mathfrak{X} \times \mathfrak{X} \rightarrow \mathbb{R}$, by

$$
\begin{aligned}
\mathscr{E}^{\prime \prime}(A, \Phi)((a, \phi),(b, \varphi)) & =\left\langle(b, \varphi), \mathscr{M}^{\prime}(A, \Phi)(a, \phi)\right\rangle_{\mathfrak{X} \times \mathfrak{X}^{*}} \\
& =\left((b, \varphi), \mathscr{M}^{\prime}(A, \Phi)(a, \phi)\right)_{\mathscr{H}}, \quad \forall(a, \phi),(b, \varphi) \in \mathfrak{X},
\end{aligned}
$$

where we define

$$
\begin{equation*}
\mathscr{H}:=L^{2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) . \tag{3.1.26}
\end{equation*}
$$

According to Theorem A. 1 the elliptic, linear, second-order partial differential operator,

$$
d_{A, \Phi}^{*} d_{A, \Phi}: W_{A_{1}}^{k+2, p}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{k, p}(X ; \operatorname{ad} P)
$$

is Fredholm for any $k \in \mathbb{Z}$ and $p \in(1, \infty)$ with kernel,

$$
K:=\operatorname{Ker}\left(d_{A, \Phi}^{*} d_{A, \Phi}: C^{\infty}(X ; \operatorname{ad} P) \rightarrow C^{\infty}(X ; \operatorname{ad} P)\right),
$$

and range,

$$
\operatorname{Ran}\left(d_{A, \Phi}^{*} d_{A, \Phi}: W_{A_{1}}^{k+2, p}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{k, p}(X ; \operatorname{ad} P)\right)=K^{\perp} \cap W_{A_{1}}^{k, p}(X ; \operatorname{ad} P)
$$

where $\perp$ denotes $L^{2}$-orthogonal complement. Hence, the operator,

$$
d_{A, \Phi}^{*} d_{A, \Phi}: K^{\perp} \cap W_{A_{1}}^{k+2, p}(X ; \operatorname{ad} P) \rightarrow K^{\perp} \cap W_{A_{1}}^{k, p}(X ; \operatorname{ad} P)
$$

is invertible, with inverse

$$
\left(d_{A, \Phi}^{*} d_{A, \Phi}\right)^{-1}: K^{\perp} \cap W_{A_{1}}^{k, p}(X ; \operatorname{ad} P) \rightarrow K^{\perp} \cap W_{A_{1}}^{k+2, p}(X ; \operatorname{ad} P)
$$

We define the Green's operator,

$$
G_{A, \Phi}: W_{A_{1}}^{k, p}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{k+2, p}(X ; \operatorname{ad} P)
$$

for the Laplacian, $d_{A, \Phi}^{*} d_{A, \Phi}$, by setting

$$
G_{A, \Phi} \xi:= \begin{cases}\left(d_{A, \Phi}^{*} d_{A, \Phi}\right)^{-1} \xi, & \forall \xi \in K^{\perp} \cap W_{A_{1}}^{k, p}(X ; \operatorname{ad} P) \\ 0, & \forall \xi \in K\end{cases}
$$

For any $k \in \mathbb{Z}$ and $p \in(1, \infty)$, we now let

$$
\left.\begin{array}{rl}
\Pi_{A_{\infty}, \Phi_{\infty}}: W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad}\right. & P \tag{3.1.27}
\end{array}\right)
$$

denote the $L^{2}$-orthogonal projection onto the slice through $\left(A_{\infty}, \Phi_{\infty}\right)$. Because

$$
\begin{aligned}
& W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
& \quad=\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
& \quad \oplus \operatorname{Ran}\left(d_{A_{\infty}, \Phi_{\infty}}: W_{A_{1}}^{k+1, p}(X ; \operatorname{ad} P \oplus E) \rightarrow W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)
\end{aligned}
$$

is an $L^{2}$-orthogonal direct sum, we see that

$$
\begin{equation*}
\Pi_{A_{\infty}, \Phi_{\infty}}=\mathrm{id}-d_{A_{\infty}, \Phi_{\infty}} G_{A_{\infty}, \Phi_{\infty}} d_{A_{\infty}, \Phi_{\infty}}^{*} \tag{3.1.28}
\end{equation*}
$$

Because of the invariance of the energy function, $\mathscr{E}: \mathscr{P}(P, E) \rightarrow \mathbb{R}$, in (1.3.1) with respect to the action of $\operatorname{Aut}(P)$ on $C^{\infty}$ pairs, $\mathscr{P}(P, E)=\mathscr{A}(P) \times C^{\infty}(X ; E)$, we have the identity

$$
\begin{equation*}
\mathscr{E}(u(A, \Phi))=\mathscr{E}(A, \Phi), \quad \forall u \in \operatorname{Aut}(P) \text { and }(A, \Phi) \in \mathscr{P}(P, E) \tag{3.1.29}
\end{equation*}
$$

Note that if $u(t) \in \operatorname{Aut}(P)$ is a family of gauge transformations depending smoothly on $t \in \mathbb{R}$ such that $u(0)=\operatorname{id}_{P}$, then the identity (3.1.29) implies that

$$
\mathscr{E}^{\prime}(A, \Phi)\left(d_{A, \Phi} \xi\right)=\left.\frac{d}{d t} \mathscr{E}(u(t)(A, \Phi))\right|_{t=0}=0, \quad \forall(A, \Phi) \in \mathscr{P}(P, E)
$$

where

$$
d_{A, \Phi} \xi=\left.\frac{d}{d t} u(t)(A, \Phi)\right|_{t=0} \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

with

$$
\xi=\dot{u}(0) \in C^{\infty}(X ; \operatorname{ad} P)=T_{\operatorname{id}_{P}} \operatorname{Aut}(P)=T_{\operatorname{id}_{P}} C^{\infty}(X ; \operatorname{Ad} P)
$$

Before considering higher-order derivatives of the identity (3.1.29) with respect to $u \in \operatorname{Aut}(P)$, we digress to discuss the Chain Rule for maps of Banach spaces.

If $F: \mathcal{Y} \rightarrow \mathcal{Z}$ and $G: \mathcal{X} \rightarrow \mathcal{Y}$ are $C^{\infty}$ maps of Banach spaces, then the Chain Rule gives, for all $x \in \mathcal{X}$,
(3.1.30) $(F \circ G)^{\prime}(x)=F^{\prime}(G(x)) \circ G^{\prime}(x) \in \mathscr{L}(\mathcal{X}, \mathcal{Z})$,

$$
\begin{equation*}
(F \circ G)^{\prime \prime}(x)=F^{\prime \prime}(G(x)) \circ G^{\prime}(x)^{2}+F^{\prime}(G(x)) \circ G^{\prime \prime}(x) \in \mathscr{L}(\mathcal{X} \times \mathcal{X}, \mathcal{Z}) \tag{3.1.31}
\end{equation*}
$$

More explicitly, if $u, v \in \mathcal{X}$, then

$$
\begin{aligned}
(F \circ G)^{\prime}(x)(u) & =F^{\prime}(G(x))\left(G^{\prime}(x)(u)\right) \\
(F \circ G)^{\prime \prime}(x)(u, v) & =F^{\prime \prime}(G(x))\left(G^{\prime}(x)(u), G^{\prime}(x)(v)\right)+F^{\prime}(G(x))\left(G^{\prime \prime}(x)(u, v)\right)
\end{aligned}
$$

The expression for the Hessian of the composition simplifies when $F^{\prime}(y)=0$ at $y=G(x)$ to give

$$
\begin{equation*}
(F \circ G)^{\prime \prime}(x)=F^{\prime \prime}(G(x)) \circ G^{\prime}(x)^{2} \in \mathscr{L}(\mathcal{X} \times \mathcal{X}, \mathcal{Z}) \tag{3.1.32}
\end{equation*}
$$

This ends our digression on the Chain Rule for $C^{\infty}$ maps of Banach spaces.
By computing the first-order differential with respect to $u$ of the expression (3.1.29) at $\operatorname{id}_{P} \in \operatorname{Aut}(P)$ in directions $\xi \in T_{\operatorname{id}_{P}} \operatorname{Aut}(P)=C^{\infty}(X ; \operatorname{ad} P)$ and recalling the definition (2.9.3) of $d_{A, \Phi} \xi$, we see that the first-order differential of $\mathscr{E}$ and its gradient map obey

$$
\begin{equation*}
\mathscr{E}^{\prime}(A, \Phi)\left(d_{A, \Phi} \xi\right)=0=\left(d_{A, \Phi} \xi, \mathscr{M}(A, \Phi)\right)_{L^{2}(X)}, \quad \forall \xi \in C^{\infty}(X ; \operatorname{ad} P) \tag{3.1.33}
\end{equation*}
$$

and thus

$$
d_{A, \Phi}^{*} \mathscr{M}(A, \Phi)=0, \quad \forall(A, \Phi) \in \mathscr{P}(P, E)
$$

Next, we observe that (3.1.29) implies that the differential of the energy function satisfies

$$
\begin{align*}
& \mathscr{E}^{\prime}(u(A, \Phi))(u(a, \phi))=\mathscr{E}^{\prime}(A, \Phi)(a, \phi)  \tag{3.1.34}\\
& \forall u \in \operatorname{Aut}(P) \text { and }(A, \Phi) \in \mathscr{P}(P, E) \\
& \quad \text { and } \forall(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) .
\end{align*}
$$

By computing the first-order differentials with respect to $u$ of the expression (3.1.34) at $\operatorname{id}_{P} \in \operatorname{Aut}(P)$ in directions $\xi \in T_{\mathrm{id}_{P}} \operatorname{Aut}(P)=C^{\infty}(X ; \operatorname{ad} P)$ and recalling the definition (2.9.3) of $d_{A, \Phi} \xi$, we see that the gradient map and second-order differential obey

$$
\mathscr{E}^{\prime \prime}(A, \Phi)\left(d_{A, \Phi} \xi,(a, \phi)\right)+\mathscr{E}^{\prime}(A, \Phi)\left(d_{A, \Phi} \xi\right)=0
$$

and thus, by (3.1.33) and the fact that the Hessian $\mathscr{E}^{\prime \prime}(A, \Phi)$ is a symmetric operator,

$$
\begin{aligned}
\mathscr{E}^{\prime \prime}(A, \Phi)\left(d_{A, \Phi} \xi,(a, \phi)\right)=0=\mathscr{E}^{\prime \prime}(A, \Phi)((a, \phi) & \left., d_{A, \Phi} \xi\right) \\
& \forall(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

Hence, by the relation (3.1.17) between $\mathscr{E}^{\prime \prime}(A, \Phi)$ and $\mathscr{M}^{\prime}(A, \Phi)$,

$$
\begin{aligned}
\left(d_{A, \Phi} \xi, \mathscr{M}^{\prime}(A, \Phi)(a, \phi)\right)_{L^{2}(X)}=0=((a, \phi) & \left., \mathscr{M}^{\prime}(A, \Phi) d_{A, \Phi} \xi\right)_{L^{2}(X)} \\
& \forall(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

Consequently, the Hessian operator satisfies

$$
\begin{align*}
& d_{A, \Phi}^{*} \mathscr{M}^{\prime}(A, \Phi)(a, \phi)=0, \forall(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) .  \tag{3.1.35a}\\
& \mathscr{M}^{\prime}(A, \Phi) d_{A, \Phi} \xi=0, \quad \forall \xi \in C^{\infty}(X ; \operatorname{ad} P) . \tag{3.1.35b}
\end{align*}
$$

Thus, for any $k \in \mathbb{Z}$ and $p \in(1, \infty)$ when $\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair, the Hessian operator,

$$
\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right): W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

naturally restricts to the Coulomb-gauge slice domain and range defined by $\left(A_{\infty}, \Phi_{\infty}\right)$,

$$
\begin{align*}
\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right): \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} & \cap W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)  \tag{3.1.36}\\
& \rightarrow \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right),
\end{align*}
$$

when the Hessian operator is also computed at the pair $\left(A_{\infty}, \Phi_{\infty}\right)$. The gradient map,

$$
\mathscr{M}:\left(A_{\infty}, \Phi_{\infty}\right)+W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

restricts to the slice domain and range with the $L^{2}$-orthogonal projection, $\Pi_{A_{\infty}, \Phi_{\infty}}$,

$$
\begin{align*}
& \hat{\mathscr{M} \equiv} \begin{aligned}
& \hat{M} \\
&\left(A_{\infty}, \Phi_{\infty}\right. \\
&\left., \Phi_{\infty}\right)+\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
& \rightarrow \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned} \tag{3.1.37}
\end{align*}
$$

The definition (3.1.27) of the $L^{2}$-orthogonal projection, $\Pi_{A, \Phi}$, and the relation (3.1.33) yield

$$
\Pi_{A, \Phi} \mathscr{M}(A, \Phi)=\mathscr{M}(A, \Phi), \quad \forall(A, \Phi) \in \mathscr{P}(P, E) .
$$

In the definition (3.1.37), we suppress dependence on the pair $\left(A_{\infty}, \Phi_{\infty}\right)$ in the choice of Coulomb-gauge slice from the notation for $\hat{\mathscr{M}}$, which we may regard as a lift to a coordinate chart on an open neighborhood of the point $[A, \Phi] \in \mathscr{C}(P, E):=$ $\mathscr{P}(P, E) / \operatorname{Aut}(P)$ of the gradient map on the quotient.

We now apply Proposition 3.1 .6 to prove the Fredholm and index zero properties of the Hessian operator.

Proposition 3.1.7 (Fredholm and index zero properties of the Hessian operator for the boson coupled Yang-Mills energy function on a Coulomb-gauge slice). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$, and $A_{1}$ be a $C^{\infty}$ reference connection on $P$, and $m, s \in C^{\infty}(X)$. If $(A, \Phi)$ is a $C^{\infty}$ pair on $(P, E)$ and $k \in \mathbb{Z}$ is an integer and $1<p<\infty$, then the following operator is Fredholm with index zero,

$$
\begin{aligned}
\mathscr{M}^{\prime}(A, \Phi): \operatorname{Ker} d_{A, \Phi}^{*} \cap W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes\right. & \operatorname{ad} P \oplus E) \\
& \rightarrow \operatorname{Ker} d_{A, \Phi}^{*} \cap W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

Proof. To reduce notational clutter, we abbreviate the domain and range Sobolev spaces by

$$
W^{k+2, p}=W_{A_{1}}^{k+2, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \quad \text { and } \quad W^{k, p}=W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

We shall adapt the argument of Råde [100, p. 148]. Observe that $\mathscr{M}^{\prime}(A, \Phi)=0$ on $\operatorname{Ran} d_{A, \Phi} \cap W^{k+2, p}$ and $d_{A, \Phi}^{*} \mathscr{M}^{\prime}(A, \Phi)=0$ on $W^{k+2, p}$ by (3.1.35) and so we have a well-defined restriction,

$$
\mathscr{M}^{\prime}(A, \Phi): \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k+2, p} \rightarrow \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k, p}
$$

Clearly, $d_{A, \Phi} d_{A, \Phi}^{*}=0$ on $\operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k+2, p}$ and

$$
\operatorname{Ran}\left(d_{A, \Phi} d_{A, \Phi}^{*}: W^{k+2, p} \rightarrow W^{k, p}\right) \subset \operatorname{Ran} d_{A, \Phi} \cap W^{k, p}
$$

Hence, we obtain a decomposition of the operator,

$$
\begin{equation*}
\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}: W^{k+2, p} \rightarrow W^{k, p} \tag{3.1.38}
\end{equation*}
$$

as a direct sum,

$$
\begin{align*}
& \mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}=\left(\begin{array}{cc}
\mathscr{M}^{\prime}(A, \Phi) & 0 \\
0 & d_{A, \Phi} d_{A, \Phi}^{*}
\end{array}\right):  \tag{3.1.39}\\
& \quad\left(\operatorname{Ker} d_{A, \Phi}^{*} \oplus \operatorname{Ran} d_{A, \Phi}\right) \cap W^{k+2, p} \rightarrow\left(\operatorname{Ker} d_{A, \Phi}^{*} \oplus \operatorname{Ran} d_{A, \Phi}\right) \cap W^{k, p} .
\end{align*}
$$

We now claim that the following operator is invertible,

$$
\begin{equation*}
d_{A, \Phi} d_{A, \Phi}^{*}: \operatorname{Ran} d_{A, \Phi} \cap W^{k+2, p} \rightarrow \operatorname{Ran} d_{A, \Phi} \cap W^{k, p} \tag{3.1.40}
\end{equation*}
$$

Indeed, the operator (3.1.40) is clearly bounded. For injectivity, if $\eta=d_{A, \Phi} \xi \in$ $\operatorname{Ran} d_{A, \Phi} \cap W^{k+2, p}$ for $\xi \in W^{k+3, p}$ and $d_{A, \Phi} d_{A, \Phi}^{*} d_{A, \Phi} \xi=0$, then $\left(d_{A, \Phi}^{*} d_{A, \Phi}\right)^{2} \xi=0$. We may assume without loss of generality that $\xi \perp \operatorname{Ker} d_{A, \Phi}^{*} d_{A, \Phi}$ and because $d_{A, \Phi}^{*} d_{A, \Phi}: W^{k+2, p} \rightarrow W^{k, p}$ is Fredholm with index zero by Theorem A.1, the following operator is invertible,

$$
d_{A, \Phi}^{*} d_{A, \Phi}:\left(\operatorname{Ker} d_{A, \Phi}^{*} d_{A, \Phi}\right)^{\perp} \cap W^{k+2, p} \rightarrow\left(\operatorname{Ker} d_{A, \Phi}^{*} d_{A, \Phi}\right)^{\perp} \cap W^{k, p} .
$$

Thus, $\xi=0$ and the operator (3.1.40) is injective. For surjectivity, suppose $\chi=$ $d_{A, \Phi} \zeta \in \operatorname{Ran} d_{A, \Phi} \cap W^{k, p}$ for $\zeta \in W^{k+1, p}$. We may again assume without loss of generality that $\zeta \perp \operatorname{Ker} d_{A, \Phi}^{*} d_{A, \Phi}$. If $\chi \perp d_{A, \Phi} d_{A, \Phi}^{*} d_{A, \Phi} \xi$ for all $\xi \in W^{k+3, p}$, then $\left(d_{A, \Phi}^{*} d_{A, \Phi}\right)^{2} \zeta=0$, and we again find that $\zeta=0$ and so $\chi=0$. Since $d_{A, \Phi} d_{A, \Phi}^{*}$ : $W^{k+2, p} \rightarrow W^{k, p}$ has closed range (because it is a Fredholm operator), this implies that the operator (3.1.40) is surjective and thus invertible by the Open Mapping Theorem.

According to Proposition 3.1.6, the operator 3.1.38 is Fredholm. Consequently,

$$
\mathscr{M}^{\prime}(A, \Phi): \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k+2, p} \rightarrow \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k, p}
$$

is Fredholm by virtue of the direct sum decomposition (3.1.39) and invertibility of the operator (3.1.40). We compute indices,

$$
\begin{aligned}
& \text { Index }\left\{\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}: W^{k+2, p} \rightarrow W^{k, p}\right\} \\
&= \operatorname{Index}\left\{\left(\begin{array}{cc}
\mathscr{M}^{\prime}(A, \Phi) & 0 \\
0 & d_{A, \Phi} d_{A, \Phi}^{*}
\end{array}\right):\left(\operatorname{Ker} d_{A, \Phi}^{*} \oplus \operatorname{Ran} d_{A, \Phi}\right) \cap W^{k+2, p}\right. \\
&\left.\rightarrow\left(\operatorname{Ker} d_{A, \Phi}^{*} \oplus \operatorname{Ran} d_{A, \Phi}\right) \cap W^{k, p}\right\} \\
&= \operatorname{Index}\left\{\mathscr{M}^{\prime}(A, \Phi): \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k+2, p} \rightarrow \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k, p}\right\} \\
&+\operatorname{Index}\left\{d_{A, \Phi} d_{A, \Phi}^{*}: \operatorname{Ran} d_{A, \Phi} \cap W^{k+2, p} \rightarrow \operatorname{Ran} d_{A, \Phi} \cap W^{k, p}\right\}
\end{aligned}
$$

Therefore, because

$$
\operatorname{Index}\left\{d_{A, \Phi} d_{A, \Phi}^{*}: \operatorname{Ran} d_{A, \Phi} \cap W^{k+2, p} \rightarrow \operatorname{Ran} d_{A, \Phi} \cap W^{k, p}\right\}=0
$$

we have

$$
\begin{aligned}
\text { Index }\left\{\mathscr{M}^{\prime}(A, \Phi): \operatorname{Ker} d_{A, \Phi}^{*}\right. & \left.\cap W^{k+2, p} \rightarrow \operatorname{Ker} d_{A, \Phi}^{*} \cap W^{k, p}\right\} \\
& =\operatorname{Index}\left\{\mathscr{M}^{\prime}(A, \Phi)+d_{A, \Phi} d_{A, \Phi}^{*}: W^{k+2, p} \rightarrow W^{k, p}\right\} .
\end{aligned}
$$

But the latter index is zero by Proposition 3.1 .6 and this completes the proof of Proposition 3.1.7.
3.1.4. Analyticity of the gradient map for the boson coupled YangMills energy function on a Coulomb-gauge slice. Suppose $\left(A_{\infty}, \Phi_{\infty}\right)$ and $(A, \Phi)$ are $C^{\infty}$ pairs on $(P, E)$ and recall from (1.3.4) that the first-order differential and gradient map of the boson coupled Yang-Mills energy function (1.3.1) are related by

$$
\mathscr{E}^{\prime}(A, \Phi)(a, \phi)=((a, \phi), \mathscr{M}(A, \Phi))_{L^{2}(X)}, \quad \forall(a, \phi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

If we now restrict $(a, \phi)$ to be a pair in $\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$, then the preceding relation yields

$$
\begin{aligned}
\mathscr{E}^{\prime}(A, \Phi)(a, \phi)=\left((a, \phi), \Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}(A, \Phi)\right)_{L^{2}(X)}=((a, \phi), \hat{\mathscr{M}}(A, \Phi))_{L^{2}(X)} \\
\forall(a, \phi) \in \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right),
\end{aligned}
$$

where $\Pi_{A_{\infty}, \Phi_{\infty}}$ is the $L^{2}$-orthogonal projection (3.1.27) onto the Coulomb-gauge slice through $\left(A_{\infty}, \Phi_{\infty}\right)$ and we appeal to the definition (3.1.37) of $\hat{\mathscr{M}}$. Consequently,

$$
\begin{array}{r}
\hat{\mathscr{M}}=\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}:\left(A_{\infty}, \Phi_{\infty}\right)+\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)  \tag{3.1.41}\\
\rightarrow \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{array}
$$

is the gradient map for the restriction,

$$
\begin{equation*}
\mathscr{E}:\left(A_{\infty}, \Phi_{\infty}\right)+\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow \mathbb{R} \tag{3.1.42}
\end{equation*}
$$

of the boson coupled Yang-Mills energy function (1.3.1) to the Coulomb-gauge slice through $\left(A_{\infty}, \Phi_{\infty}\right)$. Note that the relations (3.1.35) imply that the Hessian operator (namely, the derivative of the gradient map, $\hat{\mathscr{M}})$ at $\left(A_{\infty}, \Phi_{\infty}\right)$ simplifies to

$$
\begin{aligned}
& \hat{\mathscr{M}}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)=\left(\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}\right)^{\prime}\left(A_{\infty}, \Phi_{\infty}\right) \\
&=\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right): \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
& \rightarrow \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

Proposition 3.1.3 yields the
Corollary 3.1.8 (Analyticity of the gradient map for the boson coupled Yang-Mills energy function on a Coulomb-gauge slice). Assume the hypotheses of Proposition 3.1.3 and let $\left(A_{\infty}, \Phi_{\infty}\right)$ be a $C^{\infty}$ pair on $(P, E)$. Then the following map is analytic,

$$
\begin{aligned}
& \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{1, p} \ni\left(a_{1}, \phi_{1}\right) \\
& \mapsto \hat{\mathscr{M}}\left(A_{\infty}+a_{1}, \Phi_{\infty}+\phi_{1}\right)=\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}\left(A_{\infty}\right.\left.+a_{1}, \Phi_{\infty}+\phi_{1}\right) \\
& \in \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{-1, p}
\end{aligned}
$$

where we abbreviate

$$
W^{ \pm 1, p}=W_{A_{1}}^{ \pm 1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

Proof. The conclusion follows from Proposition 3.1.3 and the fact that the operators, $\Pi_{A_{\infty}, \Phi_{\infty}}$, define continuous projections.
3.1.5. Estimates for gauge transformations intertwining two pairs. We shall require the

Lemma 3.1.9 (Estimate for the action of a $W^{2, q}$ gauge transformation intertwining two $W^{1, q}$ pairs). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times{ }_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow$ Aut $_{\mathbb{C}}(\mathbb{E})$, and $A_{1}$ be a $C^{\infty}$ reference connection on $P, q>d / 2$ and $p$ obeys $d / 2 \leq p \leq q$, then there is a constant $C=C(g, G, p) \in[1, \infty)$ with the following significance. If $(A, \Phi)$ and $\left(A^{\prime}, \Phi^{\prime}\right)$ are $W^{1, q}$ pairs on $(P, E)$ and $u \in \operatorname{Aut}^{2, q}(P)$, then

$$
\begin{aligned}
\left\|u(A, \Phi)-u\left(A^{\prime}, \Phi^{\prime}\right)\right\|_{W_{A_{1}}^{1, p}(X)} & \leq C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)\left\|(A, \Phi)-\left(A^{\prime}, \Phi^{\prime}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \\
\left\|(A, \Phi)-\left(A^{\prime}, \Phi^{\prime}\right)\right\|_{W_{A_{1}}^{1, p}(X)} & \leq C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)\left\|u(A, \Phi)-u\left(A^{\prime}, \Phi^{\prime}\right)\right\|_{W_{A_{1}}^{1, p}(X)}
\end{aligned}
$$

If in addition $p>3 / 2$ if $d=3$ and $p \geq 4 / 3$ if $d=2$ and $p^{\prime} \in[1, \infty)$ is the dual Hölder exponent defined by $1 / p+1 / p^{\prime}=1$, and $(a, \phi) \in W_{A_{1}}^{1, p^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then

$$
\begin{aligned}
\|u(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} & \leq C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)\|(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \\
\|(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} & \leq C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)\|u(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)}
\end{aligned}
$$

Proof. Recall that $u(A)-A_{1}=u^{-1}\left(A-A_{1}\right) u+u^{-1} d_{A_{1}} u$ by (2.5.7) and similarly for $A^{\prime}$, so

$$
u(A)-u\left(A^{\prime}\right)=u^{-1}\left(A-A_{1}\right) u-u^{-1}\left(A^{\prime}-A_{1}\right) u=u^{-1}\left(A-A^{\prime}\right) u
$$

and thus,

$$
u(A, \Phi)-u\left(A^{\prime}, \Phi^{\prime}\right)=\left(u^{-1}\left(A-A^{\prime}\right) u, u\left(\Phi-\Phi^{\prime}\right)\right)
$$

Therefore, writing $a:=A-A^{\prime} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ for convenience,

$$
\nabla_{A_{1}}\left(u(A)-u\left(A^{\prime}\right)\right)=-u^{-1}\left(\nabla_{A_{1}} u\right) u^{-1} a u+u^{-1}\left(\nabla_{A_{1}} a\right) u+u^{-1} a\left(\nabla_{A_{1}} u\right)
$$

By taking $L^{p}$ norms and using the pointwise bound $|u| \leq 1$, the Sobolev embedding $W^{1, p}(X) \subset L^{2 p}(X)$ (valid for $p \geq d / 2$ ), and the Kato Inequality 43, Equation (6.20)], we obtain

$$
\begin{aligned}
\left\|\nabla_{A_{1}}\left(u\left(A-A^{\prime}\right)\right)\right\|_{L^{p}(X)} & \leq 2\left\|\nabla_{A_{1}} u\right\|_{L^{2 p}(X)}\|a\|_{L^{2 p}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{p}(X)} \\
& \leq C\left\|\nabla_{A_{1}} u\right\|_{W_{A_{1}}^{1, p}(X)}\|a\|_{W_{A_{1}}^{1, p}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{p}(X)}
\end{aligned}
$$

Similarly, $\nabla_{A_{1}}\left(u\left(\Phi-\Phi^{\prime}\right)\right)=\left(\nabla_{A_{1}} u\right)\left(\Phi-\Phi^{\prime}\right)+u\left(\nabla_{A_{1}}\left(\Phi-\Phi^{\prime}\right)\right)$ and

$$
\begin{aligned}
& \left\|\nabla_{A_{1}}\left(u\left(\Phi-\Phi^{\prime}\right)\right)\right\|_{L^{p}(X)} \\
& \quad \leq\left\|\nabla_{A_{1}} u\right\|_{L^{2 p}(X)}\left\|\Phi-\Phi^{\prime}\right\|_{L^{2 p}(X)}+\left\|\nabla_{A_{1}}\left(\Phi-\Phi^{\prime}\right)\right\|_{L^{p}(X)} \\
& \quad \leq\left\|\nabla_{A_{1}} u\right\|_{W_{A_{1}}^{1, p}(X)}\left\|\Phi-\Phi^{\prime}\right\|_{W_{A_{1}}^{1, p}(X)}+\left\|\nabla_{A_{1}}\left(\Phi-\Phi^{\prime}\right)\right\|_{L^{p}(X)}
\end{aligned}
$$

By combining the preceding estimates, we obtain the first inequality; the second inequality is proved by a symmetric argument.

For the third inequality, we note that $u(a, \phi)=(u(a), u(\phi))=\left(u^{-1} a u, u^{-1} \phi\right)$ by our convention (2.9.1) and use the identity,

$$
\nabla_{A_{1}}(u(a))=-u^{-1}\left(\nabla_{A_{1}} u\right) u^{-1} a u+u^{-1}\left(\nabla_{A_{1}} a\right) u+u^{-1} a\left(\nabla_{A_{1}} u\right)
$$

We make the

Claim 3.1.10. Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$ and $p \in[d / 2, \infty]$, with $p>3 / 2$ if $d=3$ and $p \geq 4 / 3$ if $d=2$. If $p^{\prime} \in[1, \infty)$ is the dual Hölder exponent defined by $1 / p+1 / p^{\prime}=1$, then there is $a$ continuous Sobolev multiplication map,

$$
\begin{equation*}
W^{1, p^{\prime}}(X) \times W^{1, p}(X) \rightarrow L^{p^{\prime}}(X) \tag{3.1.43}
\end{equation*}
$$

Given Claim 3.1.10, we have

$$
\left\|\nabla_{A_{1}}(u(a))\right\|_{L^{p^{\prime}}(X)} \leq z\left\|\nabla_{A_{1}} u\right\|_{W_{A_{1}}^{1, p}(X)}\|a\|_{W_{A_{1}}^{1, p^{\prime}}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{p^{\prime}}(X)}
$$

for $z=z(g, p) \in[1, \infty)$, while

$$
\|u(a)\|_{L^{p^{\prime}}(X)}=\|a\|_{L^{p^{\prime}}(X)} .
$$

Thus,

$$
\|u(a)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq z\left(1+\left\|\nabla_{A_{1}} u\right\|_{W_{A_{1}}^{1, p}(X)}\right)\|a\|_{W_{A_{1}}^{1, p^{\prime}}(X)},
$$

with the analogous estimate for $\|u(\phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)}$. This yields the third inequality and symmetry yields the fourth inequality.

Proof of Claim 3.1.10. A multiplication, $L^{r}(X) \times L^{s}(X) \rightarrow L^{p^{\prime}}(X)$, is continuous provided $r, s \in\left[p^{\prime}, \infty\right]$ obey $1 / r+1 / s \leq 1 / p^{\prime}$. If these exponents $r, s$ also yield continuous Sobolev embeddings,

$$
W^{1, p^{\prime}}(X) \subset L^{r}(X) \quad \text { and } \quad W^{1, p}(X) \subset L^{s}(X)
$$

then we obtain the desired continuous Sobolev multiplication map (3.1.43). To confirm the existence of suitable exponents, $r, s$, we shall consider the cases $p<d$, $p=d$, and $p>d$ separately.

Case $1(p<d)$. Choose $s=p^{*}=d p /(d-p)$ to give the continuous Sobolev embedding, $W^{1, p}(X) \subset L^{p^{*}}(X)$, provided by 4. Theorem 4.12]. We required that $s \in\left[p^{\prime}, \infty\right]$, so $p$ must obey $p^{*} \geq p^{\prime}$, that is

$$
d p /(d-p) \geq p /(p-1)
$$

or $p>1$ and $d(p-1) \geq d-p$ or $p(d+1) \geq 2 d$ and so we require that $p \geq$ $2 d /(d+1)$. Note that $d / 2 \geq 2 d /(d+1) \Longleftrightarrow d+1 \geq 4$, that is, $d \geq 3$, and $d / 2>2 d /(d+1) \Longleftrightarrow d \geq 4$. (This is why for $d=2,3$ we augment the hypothesis $p \geq d / 2$ in Claim 3.1.10 with additional conditions in certain subcases as noted below.) Because $p<d$, we also have $p^{\prime}=p /(p-1)>d /(d-1)$. We consider three subcases, depending on whether $a) p^{\prime}<d$, b) $p^{\prime}=d$, or $\left.c\right) p^{\prime}>d$.

Consider Subcase (1直), so $p<d$ and $p^{\prime}>d$. Then we have a continuous Sobolev embedding, $W^{1, p^{\prime}}(X) \subset L^{\infty}(X)$, by [4, Theorem 4.12] and because $p^{*} \geq p^{\prime}$ when $p \geq 2 d /(d+1)$ for $d \geq 3$ and $p \geq 4 / 3$ for $d=2$, we may choose any $r \in(1, \infty]$ large enough that $1 / r+1 / p^{*} \leq 1 / p^{\prime}$.

Consider Subcase (11b), so $p<d$ and $p^{\prime}=d$. Then we have a continuous Sobolev embedding, $W^{1, p^{\prime}}(X) \subset L^{r}(X)$, by [4, Theorem 4.12] for any $r \in[1, \infty)$ and because $p^{*}>p^{\prime}$ when $p>2 d /(d+1)$ for $d \geq 4$ and $p>3 / 2$ for $d=3$, we may choose $r \in(1, \infty)$ large enough that $1 / r+1 / p^{*} \leq 1 / p^{\prime}$. (The condition $p^{\prime}>d /(d-1)$ and $p^{\prime}=d$ forces $d \geq 3$, so the subcase $p<d$ and $p^{\prime}=d$ cannot occur for $d=2$.)

Consider Subcase (11直), so $p<d$ and $p^{\prime}<d$. Then the condition $p^{\prime}>d /(d-1)$ implies that $d /(d-1)<d$, that is, $d-1>1$ or $d \geq 3$. (The subcase $p<d$
and $p^{\prime}<d$ cannot occur for $d=2$.) We choose $r=\left(p^{\prime}\right)^{*}=d p^{\prime} /\left(d-p^{\prime}\right)=$ $d(p /(p-1)) /(d-p /(p-1))=d p /(d p-d-p)$ and use the continuous Sobolev embedding, $W^{1, p^{\prime}}(X) \subset L^{\left(p^{\prime}\right)^{*}}(X)$, provided by 4, Theorem 4.12]. To have a continuous multiplication, $L^{\left(p^{\prime}\right)^{*}}(X) \times L^{p^{*}}(X) \rightarrow L^{p^{\prime}}(X)$, we see that $p>1$ must obey

$$
1 /\left(p^{\prime}\right)^{*}+1 / p^{*} \leq 1 / p^{\prime}
$$

that is,

$$
(d p-d-p) / d p+(d-p) / d p \leq(p-1) / p
$$

or

$$
d(p-1)-p+d-p \leq d(p-1)
$$

or $d \leq 2 p$, that is, $p \geq d / 2$. Combining the conclusions of the three subcases verifies the case $p<d$.

CASE $2(p=d)$. For any $s \in[1, \infty)$, we have a continuous Sobolev embedding, $W^{1, p}(X) \subset L^{s}(X)$. We also required that $s \in\left[p^{\prime}, \infty\right]$, so $p$ and $s$ must obey $p^{\prime} \leq s<\infty$. Because $p=d$, we have $p^{\prime}=p /(p-1)=d /(d-1)$, so $\left.a\right) p^{\prime}=d$ for $d=2$, or $b) 1<p^{\prime}<d$ for $d \geq 3$.

Consider Subcase (27a), so $p^{\prime}=d$ and $d=2$. Then $W^{1, p^{\prime}}(X) \subset L^{r}(X)$ is a continuous Sobolev embedding for any $r \in[1, \infty)$. We also required that $r \in\left[p^{\prime}, \infty\right]$ and further restrict to $r \in\left(p^{\prime}, \infty\right)$ since $s$ is finite. In particular, we may choose $r, s \in\left(p^{\prime}, \infty\right)$ such that $1 / r+1 / s \leq 1 / p^{\prime}$.

Consider Subcase (2D), so $p^{\prime}<d$ and $d \geq 3$. Then $W^{1, p^{\prime}}(X) \subset L^{r}(X)$ is a continuous Sobolev embedding for $r=\left(p^{\prime}\right)^{*}=d p /(d p-d-p)$. We also required that $r \in\left[p^{\prime}, \infty\right]$ and further restrict to $r \in\left(p^{\prime}, \infty\right)$ since $s$ is finite, so $p$ must obey $p^{\prime}<\left(p^{\prime}\right)^{*}$, that is

$$
p /(p-1)<d p /(d p-d-p)
$$

or $p>1$ and $d p-d-p<d(p-1)$, or simply $p>1$ (automatic since $p=d \geq 3$ ). In particular, we may choose $r=\left(p^{\prime}\right)^{*} \in\left(p^{\prime}, \infty\right)$ and then $s \in\left(p^{\prime}, \infty\right)$ large enough that $1 / r+1 / s \leq 1 / p^{\prime}$.

Combining the conclusions of each of the two subcases verifies the case $p=d$.
Case $3(p>d)$. For any $s \in[1, \infty]$, we have a continuous Sobolev embedding, $W^{1, p}(X) \subset L^{s}(X)$. We also required that $s \in\left[p^{\prime}, \infty\right]$, so $p$ and $s$ must obey $p^{\prime} \leq s \leq \infty$. Because $p>d$, we have $p^{\prime}=p /(p-1)<d /(d-1)$, so $1<p^{\prime}<d$ for $d \geq 2$. Thus, $W^{1, p^{\prime}}(X) \subset L^{r}(X)$ is a continuous Sobolev embedding for $r=\left(p^{\prime}\right)^{*}=$ $d p /(d p-d-p)$. We also required that $r \in\left[p^{\prime}, \infty\right]$, so $p$ must obey $p^{\prime} \leq\left(p^{\prime}\right)^{*}$, which holds for any $p>1$ from our analysis of the case $p=d$. In particular, we may choose $r=\left(p^{\prime}\right)^{*} \in\left[p^{\prime}, \infty\right)$ and then $s \in\left(p^{\prime}, \infty\right]$ large enough that $1 / r+1 / s \leq 1 / p^{\prime}$. This verifies the case $p>d$.

Combining these three cases completes the proof of Claim 3.1.10.
The third and fourth inequalities follow from Claim 3.1.10 as described earlier, so this completes the proof of Lemma 3.1.9
3.1.6. Completion of the proof of Theorem 4. We can now proceed to the

Proof of Theorem 4. We first consider the simpler case where the pair $(A, \Phi)$ is in Coulomb gauge relative to the critical point $\left(A_{\infty}, \Phi_{\infty}\right)$ and then consider the general case

Case $1\left((A, \Phi)\right.$ in Coulomb gauge relative to $\left.\left(A_{\infty}, \Phi_{\infty}\right)\right)$. By hypothesis of the theorem, $\left(A_{\infty}, \Phi_{\infty}\right)$ is a $W^{1, q}$ pair that is a critical point for the function $\mathscr{E}$ in (1.3.1). By the regularity Theorem 2.11.2, there exists a $W^{2, q}$ gauge transformation $u_{\infty}$ such that $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair. In particular, $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$ is a $W^{2, q}$ pair and $u_{\infty}(A, \Phi)$ is in Coulomb gauge relative to $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$. Following (3.1.19) and (3.1.25), we choose the Banach spaces,

$$
\begin{aligned}
\mathscr{X} & =\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{p}(X ; \operatorname{ad} P)\right) \\
\tilde{X} & =\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{-2, p}(X ; \operatorname{ad} P)\right) .
\end{aligned}
$$

Hence, $\mathscr{X} \subset \tilde{\mathscr{X}}$ is a continuous embedding of Banach spaces and

$$
\mathscr{X}^{*}=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: L^{p^{\prime}}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{-1, p^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)
$$

We observe that $\tilde{\mathscr{X}} \subset \mathscr{X}^{*}$ is a continuous embedding of Banach spaces when $W^{-1, p}(X ; \mathbb{C}) \subset W^{-1, p^{\prime}}(X ; \mathbb{C})$ is a continuous Sobolev embedding and thus when $p$ obeys $p \geq 2$.

Proposition 3.1.7 implies that the Hessian operator with $x_{\infty}=u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$,

$$
\hat{\mathscr{M}}^{\prime}\left(x_{\infty}\right): \mathscr{X} \rightarrow \tilde{\mathscr{X}}
$$

is Fredholm with index zero while Corollary 3.1 .8 implies that the gradient map,

$$
\hat{\mathscr{M}}: x_{\infty}+\mathscr{X} \rightarrow \tilde{\mathscr{X}}
$$

is analytic, where we recall from (3.1.37) that

$$
\hat{\mathscr{M}}=\Pi_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)} \mathscr{M}
$$

Hence, Theorem 1 implies that there exist constants $Z^{\prime} \in(0, \infty)$ and $\sigma^{\prime} \in(0,1]$ and $\theta \in[1 / 2,1)$ (depending on $\left(A_{1}, \Phi_{1}\right)$, and $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$, and $\left.g, G, p, P\right)$ such that if

$$
\begin{equation*}
\left\|u_{\infty}(A, \Phi)-u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\sigma^{\prime} \tag{3.1.44}
\end{equation*}
$$

then

$$
\begin{equation*}
\left|\mathscr{E}\left(u_{\infty}(A, \Phi)\right)-\mathscr{E}\left(u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right)\right|^{\theta} \leq Z^{\prime}\left\|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\right\|_{W_{A_{1}}^{-1, p}(X)} \tag{3.1.45}
\end{equation*}
$$

By Lemma 3.1.9, there exists $C_{1}=C_{1}\left(g, G, p, u_{\infty}\right)=C_{1}\left(A_{\infty}, \Phi_{\infty}, g, p\right) \in[1, \infty)$ so that

$$
\left\|u_{\infty}(A, \Phi)-u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \leq C_{1}\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}
$$

More explicitly, Lemma3.1.9 gives $C_{1}=C\left(1+\left\|u_{\infty}\right\|_{W_{A_{1}}^{2, p}(X)}\right)$, where $C=C(g, G, p) \in$ $[1, \infty)$. Therefore, setting $\sigma:=C_{1}^{-1} \sigma^{\prime}$, we see that if $(A, \Phi)$ obeys the LojasiewiczSimon neighborhood condition (1.3.11), namely

$$
\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\sigma
$$

then (3.1.44) holds and thus also (3.1.45). Moreover,

$$
\begin{aligned}
& \left\|\mathscr{\mathscr { M }}\left(u_{\infty}(A, \Phi)\right)\right\|_{W_{A_{1}}^{-1, p}(X)}=\left\|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\right\|_{\left(W_{A_{1}}^{1, p^{\prime}}(X)\right)^{*}} \\
& \quad=\sup \left\{\left|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\left(u_{\infty}(a, \phi)\right)\right|:\left\|u_{\infty}(a, \phi)\right\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq 1\right\} \\
& \quad=\sup \left\{|\hat{\mathscr{M}}(A, \Phi)(a, \phi)|:\left\|u_{\infty}(a, \phi)\right\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq 1\right\} \quad \text { (by gauge invariance) }
\end{aligned}
$$

where the supremum is over all pairs, $(a, \phi) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, obeying the inequality. But

$$
\|(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq C_{1}\left\|u_{\infty}(a, \phi)\right\|_{W_{A_{1}}^{1, p^{\prime}}(X)}
$$

by Lemma 3.1.9 and therefore,

$$
\begin{aligned}
& \left\{(a, \phi) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right):\left\|u_{\infty}(a, \phi)\right\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq 1\right\} \\
& \qquad\left\{(a, \phi) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right): C_{1}^{-1}\|(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq 1\right\} .
\end{aligned}
$$

Combining the preceding equality and inequality yields,

$$
\begin{aligned}
\left\|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\right\|_{W_{A_{1}}^{-1, p}(X)} & \leq \sup \left\{|\hat{\mathscr{M}}(A, \Phi)(a, \phi)|: C_{1}^{-1}\|(a, \phi)\|_{W_{A_{1}}^{1, p^{\prime}}(X)} \leq 1\right\} \\
& =C_{1}\|\hat{\mathscr{M}}(A, \Phi)\|_{W_{A_{1}}^{-1, p}(X)} .
\end{aligned}
$$

Substituting the preceding inequality into (3.1.45) yields

$$
\begin{aligned}
\left|\mathscr{E}(A, \Phi)-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta}= & \left|\mathscr{E}\left(u_{\infty}(A, \Phi)\right)-\mathscr{E}\left(u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right)\right|^{\theta} \\
& \quad(\text { by gauge invariance }) \\
\leq & Z^{\prime}\left\|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\right\|_{W_{A_{1}}^{-1, p}(X)} \quad(\text { by (3.1.45) }) \\
\leq & Z^{\prime} C_{1}\|\hat{\mathscr{M}}(A, \Phi)\|_{W_{A_{1}}^{-1, p}(X)}
\end{aligned}
$$

But $\hat{\mathscr{M}}(A, \Phi)=\Pi_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)} \mathscr{M}$ and because the projection,

$$
\Pi_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}: W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{k, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

is bounded with norm one (for any $k \in \mathbb{Z}$ and $1<p<\infty$ ), then

$$
\|\hat{\mathscr{M}}(A, \Phi)\|_{W_{A_{1}}^{-1, p}(X)} \leq\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1, p}(X)} .
$$

Hence, the Łojasiewicz-Simon gradient inequality (1.3.12) holds for the pairs $(A, \Phi)$ and $\left(A_{\infty}, \Phi_{\infty}\right)$ with constants $(Z, \theta, \sigma)$, where $Z:=C_{1} Z^{\prime}$.

CASE $2\left((A, \Phi)\right.$ not in Coulomb gauge relative to $\left.\left(A_{\infty}, \Phi_{\infty}\right)\right)$. Let

$$
\begin{aligned}
\zeta & =\zeta\left(A_{1}, A_{\infty}, \Phi_{\infty}, g, G, p, q\right) \in(0,1] \quad \text { and } \\
N & =N\left(A_{1}, A_{\infty}, \Phi_{\infty}, g, G, p, q\right) \in[1, \infty)
\end{aligned}
$$

denote the constants in Theorem 16 and choose $\zeta_{1} \in(0, \zeta]$ small enough that $2 N \zeta_{1}<\sigma_{1}$, where we now use $\sigma_{1}$ to denote the Lojasiewicz-Simon constant from Case 1. If $(A, \Phi)$ obeys

$$
\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta_{1},
$$

then Theorem 16 provides $u \in \operatorname{Aut}^{2, q}(P)$, depending on the pair $(A, \Phi)$, such that

$$
\begin{gathered}
d_{A_{\infty}, \Phi_{\infty}}^{*}\left(u(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right)=0 \\
\left\|u(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta_{1}<\sigma
\end{gathered}
$$

By applying Case 1 to the pairs $u(A, \Phi)$ and $\left(A_{\infty}, \Phi_{\infty}\right)$, we obtain

$$
\left|\mathscr{E}(u(A, \Phi))-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \leq C_{1} Z^{\prime}\|\mathscr{M}(u(A, \Phi))\|_{W_{A_{1}}^{-1, p}(X)}
$$

Estimating as in Case 1 with $u$ replacing $u_{\infty}$, we see that

$$
\|\mathscr{M}(u(A, \Phi))\|_{W_{A_{1}}^{-1, p}(X)} \leq C_{2}\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1, p}(X)}
$$

where $C_{2}=C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)$ and $C=C(g, p) \in[1, \infty)$. According to Lemma 2.9.6, we have

$$
\|u\|_{W_{A_{1}}^{2, p}(X)} \leq C_{3}
$$

where $C_{3}=C_{3}\left(A_{\infty}, \Phi_{\infty}, A_{1}, g, G, p, q\right) \in[1, \infty)$. By combining the preceding inequalities, we obtain

$$
\begin{aligned}
\left|\mathscr{E}(A, \Phi)-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta}= & \left|\mathscr{E}(u(A, \Phi))-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \\
& \quad(\text { by gauge invariance }) \\
\leq & C_{1} Z^{\prime}\|\mathscr{M}(u(A, \Phi))\|_{W_{A_{1}}^{-1, p}(X)} \\
\leq & C_{1} C\left(1+C_{3}\right) Z^{\prime}\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1, p}(X)} .
\end{aligned}
$$

Hence, we obtain the Łojasiewicz-Simon gradient inequality (1.3.12) with constants $(Z, \theta, \sigma)$, where we now choose $Z=C_{1} C\left(1+C_{3}\right) Z^{\prime}$ and $\sigma=\zeta_{1}$.

This completes the proof of Theorem 4.
REmARK 3.1.11 (On the proof of Theorem 4 for $p=2$ and $d=2,3,4$ ). As we discussed prior to the statement of Theorem 2, that version of the abstract Łojasiewicz-Simon gradient inequality, while considerably more restrictive, has the advantage that, when applicable, its hypotheses are much easier to verify than those of Theorem 1 Thus, for $2 \leq d \leq 4$, we may choose $p=2$ and observe that the condition $d / 2 \leq p \leq q$ is met for $q>2$. With the notation used in the proof of Theorem 1 we can then choose

$$
\mathscr{X}=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right),
$$

with dual space,

$$
\begin{aligned}
\mathscr{X}^{*}=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{-1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right.\right. & \oplus E) \\
& \left.\rightarrow W_{A_{1}}^{-2,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)
\end{aligned}
$$

Proposition 3.1.1 implies that the restriction of the energy function to the Coulombgauge slice,

$$
\mathscr{E}: x_{\infty}+\mathscr{X} \rightarrow \mathbb{R}
$$

is real analytic, where $x_{\infty}:=u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$. Proposition 3.1.7 implies that the Hessian operator,

$$
\mathscr{E}^{\prime \prime}\left(x_{\infty}\right) \in \mathscr{L}\left(\mathscr{X}, \mathscr{X}^{*}\right)
$$

is Fredholm with index zero. Case 1 of the proof of Theorem 4 now follows easily, where the gauge transformation, $u_{\infty}$, is again chosen via Theorem 2.11.2 so that $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair. The difficult third and fourth inequalities in the statement of Lemma 3.1.9 are not required. Case 2 of the proof of Theorem 4 is unchanged.

## 3.2. Łojasiewicz-Simon gradient inequality for the fermion coupled Yang-Mills energy function

We assume the notation and conventions of Section 3.1 By analogy with Proposition 3.1.1 we establish the forthcoming Proposition 3.2.1 for the analyticity of the fermion coupled Yang-Mills energy function, $\mathscr{F}$. Again, this serves as a stepping stone towards the proof that its gradient map, $\hat{\mathscr{M}}: x_{\infty}+\mathscr{X} \rightarrow \tilde{\mathscr{X}}$, is real analytic for suitable choices of Banach spaces as in Theorem 1

Proposition 3.2.1 (Analyticity of the fermion coupled Yang-Mills energy function). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $(\rho, W)$ be a spin ${ }^{c}$ structure on $X$, and $G$ be a compact Lie group, $P$ be a smooth principal $G$-bundle over $X$, and $E=P \times_{\varrho} \mathbb{E}$ be a smooth Hermitian vector bundle over $X$ defined by a finite-dimensional unitary representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{C}}(\mathbb{E})$, and $A_{1}$ be a smooth reference connection on $P$, and $m \in C^{\infty}(X)$. If $4 d /(d+4) \leq$ $p<\infty$, then the function (1.3.3),

$$
\mathscr{F}: \mathscr{A}^{1, p}(P) \times W_{A_{1}}^{1, p}(X ; W \otimes E) \rightarrow \mathbb{R},
$$

is real analytic.
Proof. We prove analyticity at a point $(A, \Psi)$ and write $A=A_{1}+a_{1}$, where $a \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$. For any $(a, \psi) \in \mathfrak{X}$, we have

$$
\begin{aligned}
F_{A+a} & =F_{A_{1}}+d_{A_{1}}\left(a_{1}+a\right)+\frac{1}{2}\left[a_{1}+a, a_{1}+a\right], \\
D_{A+a}(\Psi+\psi) & =D_{A_{1}}(\Psi+\psi)+\rho\left(a_{1}+a\right)(\Psi+\psi) .
\end{aligned}
$$

By definition (1.3.3) of the fermion coupled Yang-Mills energy function, we obtain

$$
2 \mathscr{F}(A+a, \Psi+\psi)=T_{1}+T_{2}+T_{3},
$$

where the curvature term,

$$
T_{1}:=\left(F_{A+a}, F_{A+a}\right)_{L^{2}(X)},
$$

has the same expansion as the corresponding term $T_{1}$ in Proposition 3.1.1 for the boson coupled Yang-Mills energy function, while

$$
\begin{aligned}
T_{2}:= & \left(\Psi, D_{A_{1}} \Psi\right)_{L^{2}(X)}+\left(\Psi, D_{A_{1}} \psi\right)_{L^{2}(X)}+\left(\psi, D_{A_{1}} \Psi\right)_{L^{2}(X)}+\left(\psi, D_{A_{1}} \psi\right)_{L^{2}(X)} \\
& +\left(\Psi, \rho\left(a_{1}+a\right) \Psi\right)_{L^{2}(X)}+\left(\Psi, \rho\left(a_{1}+a\right) \psi\right)_{L^{2}(X)}+\left(\psi, \rho\left(a_{1}+a\right) \Psi\right)_{L^{2}(X)} \\
& +\left(\psi, \rho\left(a_{1}+a\right) \psi\right)_{L^{2}(X)},
\end{aligned}
$$

and thus,

$$
\begin{aligned}
T_{2}= & \left(\Psi, D_{A_{1}} \Psi\right)_{L^{2}(X)}+2 \operatorname{Re}\left(\Psi, D_{A_{1}} \psi\right)_{L^{2}(X)}+\left(\psi, D_{A_{1}} \psi\right)_{L^{2}(X)} \\
& +\left(\Psi, \rho\left(a_{1}+a\right) \Psi\right)_{L^{2}(X)}+2 \operatorname{Re}\left(\Psi, \rho\left(a_{1}+a\right) \psi\right)_{L^{2}(X)}+\left(\psi, \rho\left(a_{1}+a\right) \psi\right)_{L^{2}(X)}, \\
T_{3}:= & \int_{X} m\left(|\Psi|^{2}+\langle\Psi, \psi\rangle+\langle\psi, \Psi\rangle+|\psi|^{2}\right) d \operatorname{vol}_{g}
\end{aligned}
$$

The terms in the expression for the difference,

$$
2 \mathscr{F}(A+a, \Psi+\psi)-2 \mathscr{F}(A, \Psi)=T_{1}^{\prime}+T_{2}^{\prime}+T_{3}^{\prime},
$$

are organized in such a way that

$$
T_{1}^{\prime}:=\left(F_{A+a}, F_{A+a}\right)_{L^{2}(X)}-\left(F_{A}, F_{A}\right)_{L^{2}(X)}
$$

has the same expansion as the corresponding term $T_{1}^{\prime}$ for the boson coupled YangMills energy function in the proof of Proposition 3.2.1 and the remaining terms are given by

$$
\begin{aligned}
T_{2}^{\prime}:= & 2 \operatorname{Re}\left(\Psi, D_{A_{1}} \psi\right)_{L^{2}(X)}+\left(\psi, D_{A_{1}} \psi\right)_{L^{2}(X)} \\
& +2 \operatorname{Re}\left(\Psi, \rho\left(a_{1}+a\right) \psi\right)_{L^{2}(X)}+\left(\psi, \rho\left(a_{1}+a\right) \psi\right)_{L^{2}(X)} \\
T_{3}^{\prime}:= & \int_{X} m\left(2 \operatorname{Re}\langle\Psi, \psi\rangle+|\psi|^{2}\right) d \operatorname{vol}_{g} .
\end{aligned}
$$

The proof of analyticity of $\mathscr{F}$ at $(A, \Psi)$ now follows by adapting mutatis mutandis the arguments used to prove Proposition 3.1.1.

We now verify the formula (1.3.7) for the differential $\mathscr{E}^{\prime}(A, \Phi)$ and gradient $\mathscr{M}(A, \Phi)$.

Lemma 3.2.2 (Differential and gradient of the fermion coupled Yang-Mills energy function). Assume the hypotheses of Proposition 3.2.1 with the dual Hölder exponent $p^{\prime}$ in the range $1<p^{\prime} \leq 4 d /(3 d-4)$ determined by $4 d /(d+4) \leq p<\infty$ and $1 / p+1 / p^{\prime}=1$. Then the expression for $\mathscr{E}^{\prime}(A, \Psi) \in\left(W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)\right)^{*}$ and $\mathscr{M}(A, \Psi) \in W_{A_{1}}^{-1, p^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)$ is given by (1.3.7), namely

$$
\begin{aligned}
\mathscr{E}^{\prime}(A, \Psi)(a, \psi)= & ((a, \psi), \mathscr{M}(A, \Psi))_{L^{2}(X)} \\
= & \left(d_{A}^{*} F_{A}, a\right)_{L^{2}(X)}+\operatorname{Re}\left(D_{A} \Psi-m \Psi, \psi\right)_{L^{2}(X)}+\frac{1}{2}(\Psi, \rho(a) \Psi)_{L^{2}(X)}, \\
& \forall(a, \psi) \in) \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right) .
\end{aligned}
$$

Proof. It suffices to extract the terms that are linear in $(a, \psi)$ from the expressions for

$$
T_{1}^{\prime}=2\left(F_{A}, d_{A} a\right)_{L^{2}(X)}+\left(d_{A} a, d_{A} a\right)_{L^{2}(X)}
$$

and $T_{2}^{\prime}$ and $T_{3}^{\prime}$ arising in the proof of Proposition 3.2.1
Remark 3.2.3 (Pointwise self-adjointness and reality). The fact that the term $\langle\Psi, \rho(a) \Psi\rangle$ appearing in (1.3.7) is real could be inferred indirectly by noting the origin of this term and the fact that the Dirac operator, $D_{A}$, is self-adjoint. To see directly that $\langle\Psi, \rho(a) \Psi\rangle$ is real, recall that Clifford multiplication is skew-Hermitian, so $c(\alpha)^{*}=-c(\alpha) \in \operatorname{End}_{\mathbb{C}}(W)$ for all $\alpha \in \Omega^{1}(X)$ (for example, see [52, p. 49]) while if $\xi \in \mathfrak{g}$, then $\varrho_{*}(\xi)^{*}=-\varrho_{*}(\xi)$ since we assume that Lie structure group, $G$, of $P$ acts on the complex, finite-dimensional vector space $\mathbb{E}$ via a unitary representation, $\varrho: G \rightarrow \operatorname{End}_{\mathbb{C}}(\mathbb{E})$, and $\varrho_{*}: \mathfrak{g} \rightarrow \operatorname{End}_{\mathbb{C}}(\mathbb{E})$ is the induced representation of the Lie algebra, $\mathfrak{g}$. Hence, given $\alpha \otimes \xi \in C^{\infty}\left(T^{*} X \otimes \operatorname{ad} P\right)=\Omega^{1}(X ; \operatorname{ad} P)$ and recalling that $E=P \times_{\varrho} \mathbb{E}$, then $\rho(\alpha \otimes \xi)=c(\alpha) \otimes \varrho_{*}(\xi) \in \operatorname{End}_{\mathbb{C}}(W \otimes E)$ obeys

$$
\rho(\alpha \otimes \xi)^{*}=c(\alpha)^{*} \otimes \varrho_{*}(\xi)^{*}=c(\alpha) \otimes \varrho_{*}(\xi)=\rho(\alpha \otimes \xi)
$$

In particular, $\rho(a) \in \operatorname{End}_{\mathbb{C}}(W \otimes E)$ satisfies $\rho(a)^{*}=\rho(a)$ for all $a \in \Omega^{1}(X ; \operatorname{ad} P)$.

We now compute the Hessian operator, $\mathscr{M}^{\prime}(A, \Psi)$, at a $C^{\infty}$ pair $(A, \Psi) \in$ $\mathscr{A}(P) \times C^{\infty}(X ; W \otimes E)$. The gradient, $\mathscr{M}(A, \Psi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)$ in (1.3.7), may be written as

$$
\begin{equation*}
\mathscr{M}(A, \Psi)=d_{A}^{*} F_{A}+\frac{1}{2}\left(\left(D_{A}-m\right) \Psi \cdot+\cdot\left(D_{A}-m\right) \Psi\right)+\frac{1}{2} \rho^{-1}\left(\Psi \otimes \Psi^{*}\right) \tag{3.2.1}
\end{equation*}
$$

where the terms involving $\Psi$ in this expression for $\mathscr{M}(A, \Psi)$ are defined by the $L^{2}$-pairings,

$$
\begin{aligned}
\operatorname{Re}\left(D_{A} \Psi-m \Psi, \psi\right)_{L^{2}(X)}= & \frac{1}{2}\left(\left(\left(D_{A}-m\right) \Psi, \psi\right)_{L^{2}(X)}+\left(\psi,\left(D_{A}-m\right) \Psi\right)_{L^{2}(X)}\right) \\
\frac{1}{2}(\Psi, \rho(a) \Psi)_{L^{2}(X)}= & \frac{1}{2}\left(\Psi \otimes \Psi^{*}, \rho(a)\right)_{L^{2}(X)}=\frac{1}{2}\left(\rho^{-1}\left(\Psi \otimes \Psi^{*}\right), a\right)_{L^{2}(X)} \\
& \forall(a, \psi) \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)
\end{aligned}
$$

Taking the derivative of the gradient $\mathscr{M}(A, \Psi)$ in (3.2.1) with respect to $(A, \Psi)$ in the direction $(a, \psi)$ yields

$$
\begin{align*}
& \mathscr{M}^{\prime}(A, \Psi)(a, \psi)=d_{A}^{*} d_{A} a+\frac{1}{2}\left(D_{A} \psi \cdot+\cdot D_{A} \psi\right)-\frac{1}{2}(m \psi \cdot+\cdot m \psi)  \tag{3.2.2}\\
& \quad+(a \wedge \cdot)^{*} F_{A}+\frac{1}{2}(\rho(a) \Psi \cdot+\cdot \rho(a) \Psi)+\frac{1}{2} \rho^{-1}\left(\Psi \otimes \psi^{*}+\psi \otimes \Psi^{*}\right)
\end{align*}
$$

By virtue of (3.2.2) we may view the Hessian operator for $\mathscr{F}$ at a $C^{\infty}$ pair $(A, \Psi)$ as a linear, second-order partial differential operator,

$$
\mathscr{M}^{\prime}(A, \Psi): C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right) \rightarrow C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right) .
$$

The differential,

$$
d_{A, \Psi}: C^{\infty}(X ; \operatorname{ad} P) \rightarrow C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus W \otimes E\right)
$$

is defined just as in (2.9.2) except that we replace $E$ by $W \otimes E$. As in (3.1.38), the operator $\mathscr{M}^{\prime}(A, \Psi)+d_{A, \Psi} d_{A, \Psi}^{*}$ is elliptic and thus Fredholm on the analogous Sobolev spaces. We can now simply make the following observation to conclude the proof of Theorem 6

Proof of Theorem 6. The argument applies mutatis mutandis the corresponding steps used to prove Theorem 4 and all required intermediate results established in Section 3.1

## CHAPTER 4

## Łojasiewicz-Simon $W^{-1,2}$ gradient inequalities for coupled Yang-Mills energy functions

In this chapter, we shall prove Corollaries 7 and 8, using a method analogous to that of our proof of [38, Corollary 6] for the harmonic map energy function. We shall accomplish this by applying the more general Theorem 3 instead of Theorem 11. which we used to prove Theorem 4.

### 4.1. Technical preparation for the proof of the $W^{-1,2}$ gradient inequalities for coupled Yang-Mills energy functions

Before proceeding to the proof of Corollaries 7 and 8 we first establish several technical lemmas and corollaries needed to verify the hypotheses of Theorem 3. In the forthcoming Lemma 4.1.1, it would suffice to have continuity of the extension $\mathscr{M}_{1}(A, \Phi)$ as a function of $(A, \Phi)$ in the affine Banach space $\mathscr{P}^{1, q}(P, E)$ of $W^{1, q}$ pairs on $(P, E)$ in order to satisfy a key hypothesis of Theorem 3, but the proof yields analyticity as a function of $(A, \Phi)$ in the affine Banach space $\mathscr{P}^{1, p}(P, E)$ of $W^{1, p}$ pairs on $(P, E)$ (compare Proposition 3.1.1).

Lemma 4.1.1 (Analyticity of extended boson coupled Yang-Mills augmented Hessian operator map). Assume the hypotheses of Theorem 4. If $(A, \Phi)$ is a $W^{1, p}$ pair on $(P, E)$ then the operator (3.1.22),

$$
\mathscr{M}^{\prime}(A, \Phi): W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

has a bounded extension

$$
\begin{equation*}
\mathscr{M}_{1}(A, \Phi): W_{A_{1}}^{1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{-1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \tag{4.1.1}
\end{equation*}
$$

that is an analytic function of $(A, \Phi)$ in the affine Banach space $\mathscr{P}^{1, p}(P, E)$ of $W^{1, p}$ pairs on $(P, E)$.

Proof. The fact that the augmented Hessian operator (3.1.22) has a welldefined bounded extension (4.1.1) follows from Proposition 3.1.6 when $(A, \Phi)$ is a $C^{\infty}$ pair by taking $k=-1$ and $p=2$.

For a general pair $(A, \Phi)$ of class $W^{1, p}$, we write $(A, \Phi)=\left(A_{1}+a, \Phi\right)$, where $A_{1}$ is $C^{\infty}$ and $a \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$. Then from the expression (3.1.16) (with $(a, \phi)$ replaced by $(b, \varphi))$, we have the formal expansion

$$
\begin{align*}
\mathscr{M}^{\prime}(A, \Phi)(b, \varphi)= & d_{A}^{*} d_{A} b+\nabla_{A}^{*} \nabla_{A} \varphi+F_{A} \times b+\nabla_{A}^{*}(\varrho(b) \Phi)+\Phi \times \nabla_{A} \varphi \\
& -\rho(b)^{*} \nabla_{A} \Phi+\nabla_{A} \Phi \times \varphi+\varrho(b) \Phi \times \Phi  \tag{4.1.2}\\
& -\left(m+2 s|\Phi|^{2}\right) \varphi-4 s\langle\Phi, \varphi\rangle \Phi, \quad \forall(b, \varphi) \in \mathfrak{X},
\end{align*}
$$

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so that, using $\nabla_{A_{1}}^{*}(\varrho(b) \Phi)=b \times \nabla_{A_{1}} \Phi+\nabla_{A_{1}} b \times \Phi$,

$$
\begin{align*}
& \mathscr{M}^{\prime}\left(A_{1}+a, \Phi\right)(b, \varphi) \\
&= d_{A_{1}}^{*} d_{A_{1}} b+\nabla_{A_{1}}^{*} \nabla_{A_{1}} \varphi+F_{A_{1}} \times b \\
& \quad+a \times d_{A_{1}} b+d_{A_{1}} a \times b+\nabla_{A_{1}} a \times \varphi+a \times \nabla_{A_{1}} \varphi \\
&+a \times a \times b+a \times a \times \varphi . \\
&+\nabla_{A_{1}} b \times \Phi+\Phi \times \nabla_{A_{1}} \varphi  \tag{4.1.3}\\
&+b \times \nabla_{A_{1}} \Phi+\nabla_{A_{1}} \Phi \times \varphi \\
&+a \times \Phi \times b+a \times \Phi \times \varphi \\
&+\varrho(b) \Phi \times \Phi-\left(m+2 s|\Phi|^{2}\right) \varphi-4 s\langle\Phi, \varphi\rangle \Phi
\end{align*}
$$

for every $(b, \varphi) \in W_{A_{1}}^{1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$. We make the
Claim 4.1.2. Continue the preceding notation but abbreviate

$$
W^{ \pm 1,2}=W_{A_{1}}^{ \pm 1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
$$

If $\mathscr{M}_{1}\left(A_{1}+a, \Phi\right)$ denotes the bounded extension of $\mathscr{M}^{\prime}\left(A_{1}+a, \Phi\right)$, then

$$
\begin{align*}
W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \ni(a, \Phi) &  \tag{4.1.4}\\
& \mapsto \mathscr{M}_{1}\left(A_{1}+a, \Phi\right) \in \mathscr{L}\left(W^{1,2}, W^{-1,2}\right)
\end{align*}
$$

is a cubic polynomial in $(a, \Phi)$ and its first-order covariant derivatives with respect to $\nabla_{A_{1}}$, with universal coefficients (depending at most on $g$ and $G$ ).

Given Claim 4.1.2, it follows that $\mathscr{M}_{1}(A, \Phi)$ is a well-defined analytic function of $(A, \Phi)$ in an open $W^{1, p}$ neighborhood of $\left(A_{1}, \Phi_{1}\right)$, thus completing the proof of Lemma 4.1.1

We now turn to the
Proof of Claim 4.1.2. We compute a $W^{-1,2}$ bound for each term in equation 4.1.3).

STEP 1 ( $W^{-1,2}$ estimates for $a \times d_{A_{1}} b$ and $a \times \nabla_{A_{1}} \varphi$ and $\Phi \times \nabla_{A_{1}} \varphi$ and $\left.\nabla_{A_{1}} b \times \Phi\right)$. We claim that

$$
\begin{align*}
&\left\|a \times d_{A_{1}} b\right\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.5a}\\
&\left\|a \times \nabla_{A_{1}} \varphi\right\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|\varphi\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.5b}\\
&\left\|\Phi \times \nabla_{A_{1}} \varphi\right\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\varphi\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.5c}\\
&\left\|\nabla_{A_{1}} b \times \Phi\right\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)} \tag{4.1.5d}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$.
We first prove estimate 4.1.5a). Using duality, together with the Hölder inequality with exponents $r, r^{\prime}$ satisfying $1 / r+1 / r^{\prime}=1$, we estimate

$$
\begin{aligned}
\left\|a \times d_{A_{1}} b\right\|_{W_{A_{1}}^{-1,2}(X)} & =\sup _{0<\|\beta\|_{W_{A_{1}(X)}^{1,2}} \leq 1}\left(a \times d_{A_{1}} b, \beta\right)_{L^{2}(X)} \\
& \leq\left\|a \times d_{A_{1}} b\right\|_{L^{r^{\prime}}(X)} \sup _{0<\|\beta\|_{W_{A_{1}}^{1,2}(X)} \leq 1}\|\beta\|_{L^{r}(X)}
\end{aligned}
$$

Estimate 4.1.5a will follow by showing that we can always choose $r \in(1, \infty)$ so that

$$
\begin{align*}
\left\|a \times d_{A_{1}} b\right\|_{L^{r^{\prime}}(X)} & \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.6a}\\
\|\beta\|_{L^{r}(X)} & \leq z\|\beta\|_{W_{A_{1}}^{1,2}(X)} \tag{4.1.6b}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$. We separately consider the cases $d \geq 3$ with $p \geq d$ or $d / 2 \leq p<d$ and $d=2$ with $p \geq 2$.

CASE $1(d \geq 3$ and $d / 2 \leq p<d)$. We choose $r=2^{*}=2 d /(d-2) \in(2,6]$, so $r^{\prime}=r /(r-1)=2 d /(d+2) \in(1,2)$, and $s=p^{*}=d p /(d-p) \in[d, \infty)$. By 4. Theorem 4.12], we have continuous Sobolev embeddings $W^{1,2}(X) \subset L^{r}(X)$ and $W^{1, p}(X) \subset L^{s}(X)$. Also $s \in(2, \infty)$ and obeys $1 / s+1 / 2 \leq 1 / r^{\prime}$, that is

$$
1 / s+1 / 2=(d-p) /(d p)+1 / 2=1 / p-1 / d+1 / 2 \leq 1 / d+1 / 2=1 / r^{\prime}
$$

since $1 / p \leq 2 / d$ from $p \geq d / 2$, and which we assume by hypothesis on $p$. Therefore,

$$
\left\|a \times d_{A_{1}} b\right\|_{L^{r^{\prime}}(X)} \leq z\|a\|_{L^{s}(X)}\left\|d_{A_{1}} b\right\|_{L^{2}(X)} \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}
$$

and $\|\beta\|_{L^{r}(X)} \leq z\|\beta\|_{W_{A_{1}}^{1,2}(X)}$. This yields the estimates (4.1.6) for this case.
CASE $2(d \geq 3$ and $p \geq d)$. We again choose $r=2^{*}=2 d /(d-2) \in(2,6]$, so $r^{\prime}=$ $r /(r-1)=2 d /(d+2) \in(1,2)$ and $s \in(2, \infty)$ is any sufficiently large number that $1 / s+1 / 2 \leq 1 / r^{\prime}$. By [4, Theorem 4.12], we have a continuous Sobolev embedding $W^{1,2}(X) \subset L^{r}(X)$ and continuous Sobolev embeddings $W^{1, p}(X) \subset L^{s}(X)$ when $p=d$ and $W^{1, p}(X) \subset C(X) \subset L^{s}(X)$ when $p>d$. The estimates (4.1.6) now follow just as in Case 1 .

CASE $3(d=2$ and $p \geq 2)$. We choose $r \in(2, \infty)$ large enough that $r^{\prime}=$ $r /(r-1) \in(1,2)$ is sufficiently close to one and choose $s \in(1, \infty)$ large enough that $1 / s+1 / 2 \leq 1 / r^{\prime}$. By [4, Theorem 4.12], we have a continuous Sobolev embedding $W^{1,2}(X) \subset L^{r}(X)$ and continuous Sobolev embeddings $W^{1, p}(X) \subset L^{s}(X)$ when $p=2$ and $W^{1, p}(X) \subset C(X) \subset L^{s}(X)$ when $p>2$. The estimates 4.1.6) now follow just as in Case 1 .

This completes the proof of estimate 4.1.5a). The proofs of estimates 4.1.5b), (4.1.5c), and 4.1.5d) are analogous to that of estimate 4.1.5a). This completes Step 1 .

STEP $2\left(W^{-1,2}\right.$ estimates for $d_{A_{1}} a \times b$ and $\nabla_{A_{1}} a \times \varphi$ and $b \times \nabla_{A_{1}} \Phi$ and $\nabla_{A_{1}} \Phi \times \varphi$ ). We claim that

$$
\begin{align*}
\left\|d_{A_{1}} a \times b\right\|_{W_{A_{1}}^{-1,2}(X)} & \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.7a}\\
\left\|\nabla_{A_{1}} a \times \varphi\right\|_{W_{A_{1}}^{-1,2}(X)} & \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|\varphi\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.7b}\\
\left\|b \times \nabla_{A_{1}} \Phi\right\|_{W_{A_{1}}^{-1,2}(X)} & \leq z\|\Phi\|_{W_{A_{1}(X)}^{1, p}}\|b\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.7c}\\
\left\|\nabla_{A_{1}} \Phi \times \varphi\right\|_{W_{A_{1}}^{-1,2}(X)} & \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\varphi\|_{W_{A_{1}}^{1,2}(X)} \tag{4.1.7d}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$.
We first prove estimate 4.1.7a). Again, we use duality together with the Hölder inequality, first with exponents $r, r^{\prime}$ satisfying $1 / r+1 / r^{\prime}=1$ and second with
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exponents $s, r$ satisfying $1 / p+1 / r \leq 1 / r^{\prime}$. We estimate

$$
\begin{aligned}
\left\|d_{A_{1}} a \times b\right\|_{W_{A_{1}}^{-1,2}(X)} & =\sup _{0<\|\beta\|_{W_{A_{1}}^{1,2}(X)} \leq 1}\left(d_{A_{1}} a \times b, \beta\right)_{L^{2}(X)} \\
& \leq\left\|d_{A_{1}} a \times b\right\|_{L^{r^{\prime}(X)}} \sup _{0<\|\beta\|_{W_{A_{1}}^{1,2}(X)} \leq 1}\|\beta\|_{L^{r}(X)}
\end{aligned}
$$

We separately consider the cases $d \geq 3$ and $d=2$.
Case $1(d \geq 3)$. As in Step 1 for the two cases with $d \geq 3$, we choose $r=$ $2^{*}=2 d /(d-2) \in(2,6]$, so $r^{\prime}=r /(r-1)=2 d /(d+2) \in(1,2)$ and again we have a continuous Sobolev embedding $W^{1,2}(X) \subset L^{r}(X)$. Moreover,

$$
\begin{aligned}
& 1 / p+1 / r=1 / p+(d-2) /(2 d) \leq 1 /(d / 2)+(d-2) /(2 d) \\
& =2 / d+1 / 2-1 / d=1 / d+1 / 2=(2+d) /(2 d)=1 / r^{\prime}
\end{aligned}
$$

Therefore, using $1 / p+1 / r \leq 1 / r^{\prime}$,

$$
\left\|d_{A_{1}} a \times b\right\|_{L^{r^{\prime}}(X)} \leq z\left\|d_{A_{1}} a\right\|_{L^{p}(X)}\|b\|_{L^{r}(X)} \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}
$$

This proves 4.1.7a in this case.
CASE $2(d=2$ and $p \geq 2)$. Noting that $p<\infty$, we may choose $r \in(2, \infty)$ large enough and thus $r^{\prime}=r /(r-1) \in(1,2)$ close enough to one so $1 / p+1 / r \leq 1 / r^{\prime}$. By 4. Theorem 4.12], we have a continuous Sobolev embedding $W^{1,2}(X) \subset L^{r}(X)$ and thus estimate 4.1.7a) follows just as in Case 1

The proofs of estimates 4.1.7b, 4.1.7c , and 4.1.7d are analogous to that of estimate 4.1.7a). This completes Step 2

STEP 3 ( $W^{-1,2}$ estimates for $a \times \Phi \times b$ and $a \times \Phi \times \varphi$ and $a \times a \times b$ and $a \times a \times \varphi$ and $\Phi \times \Phi \times b$ and $\Phi \times \Phi \times \varphi$ ). We claim that

$$
\begin{align*}
&\|a \times \Phi \times b\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.8a}\\
&\|a \times \Phi \times \varphi\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|\varphi\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.8b}\\
&\|a \times a \times b\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|a\|_{W_{A_{1}(X)}^{1, p}}^{2}\|b\|_{W_{A_{1}(X)}^{1,2}}  \tag{4.1.8c}\\
&\|a \times a \times \varphi\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|a\|_{W_{A_{1}(X)}^{1, p}}^{2}\|\varphi\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.8d}\\
&\|\Phi \times \Phi \times b\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}\|b\|_{W_{A_{1}}^{1,2}(X)}  \tag{4.1.8e}\\
&\|\Phi \times \Phi \times \varphi\|_{W_{A_{1}}^{-1,2}(X)} \leq z\|\Phi\|_{W_{A_{1}}^{1, p}(X)}^{2}\|\varphi\|_{W_{A_{1}}^{1,2}(X)} \tag{4.1.8f}
\end{align*}
$$

where $z=z(g, G, p) \in[1, \infty)$.
We first consider the estimate 4.1.8a). We claim that we can always find $r \in(2, \infty)$, with $r^{\prime}=r /(r-1) \in(1,2)$, and $s \in(1, \infty)$ such that $1 / s+1 / r \leq 1 / r^{\prime}$, and such that we have continuous Sobolev embeddings, $W^{1, p}(X) \subset L^{2 s}(X)$ and
$W^{1,2}(X) \subset L^{r}(X)$. Assuming that claim, we have

$$
\begin{aligned}
\|a \times \Phi \times b\|_{W_{A_{1}}^{-1,2}(X)} & =\sup _{0<\|\beta\|_{W_{A_{1}}^{1,2}(X)} \leq 1}(a \times \Phi \times b, \beta)_{L^{2}(X)} \\
& \leq\|a \times \Phi \times b\|_{L^{r^{\prime}}(X)} \sup _{0<\|\beta\|_{W_{A_{1}}^{1,2}(X)} \leq 1}\|\beta\|_{L^{r}(X)} \\
& \leq z\|a \times \Phi\|_{L^{s}(X)}\|b\|_{L^{r}(X)} \sup _{0<\|\beta\|_{W_{A_{1}(X)}^{1,2}} \leq 1}\|\beta\|_{W_{A_{1}}^{1,2}(X)} \\
& \leq z\|a\|_{L^{2 s}(X)}\|\Phi\|_{L^{2 s}(X)}\|b\|_{L^{r}(X)} \\
& \leq z\|a\|_{W_{A_{1}}^{1, p}(X)}\|\Phi\|_{W_{A_{1}}^{1, p}(X)}\|b\|_{W_{A_{1}}^{1,2}(X)}
\end{aligned}
$$

and thus yielding (4.1.8a). Our task now is to verify our claim that such $r, s$ always exist. We separately consider the cases $d \geq 3$ with $p \geq d$ or $d / 2 \leq p<d$ and $d=2$ with $p \geq 2$.

CASE $1(d \geq 3$ and $d / 2 \leq p<d)$. We choose $r=2^{*}=2 d /(d-2) \in(2, \infty)$, so $r^{\prime}=r /(r-1)=2 d /(d+2) \in(1,2)$ and choose $s \in(1, \infty)$ by setting $2 s=p^{*}=$ $d p /(d-p) \in(2, \infty)$. By [4, Theorem 4.12], we have continuous Sobolev embeddings, $W^{1,2}(X) \subset L^{r}(X)$ and $W^{1, p}(X) \subset L^{2 s}(X)$. Moreover,

$$
\begin{aligned}
& 1 / s+1 / r=2(d-p) /(d p)+(d-2) /(2 d)=2 / p-2 / d+1 / 2-1 / d \\
& =2 / p+1 / 2-3 / d \leq 4 / d+1 / 2-3 / d=1 / d+1 / 2=1 / r^{\prime}
\end{aligned}
$$

using $p \geq d / 2$, which we assume by hypothesis, and thus $1 / s+1 / r \leq 1 / r^{\prime}$. This justifies our choice of $r, s$ in this case.

CASE $2(d \geq 3$ and $p \geq d)$. We again choose $r=2^{*}=2 d /(d-2) \in(2, \infty)$, so $r^{\prime}=r /(r-1)=2 d /(d+2) \in(1,2)$, and choose $s \in(1, \infty)$ large enough that $1 / s+1 / r \leq 1 / r^{\prime}$. By [4, Theorem 4.12], we have again have continuous Sobolev embeddings, $W^{1,2}(X) \subset L^{r}(X)$ and $W^{1, p}(X) \subset L^{2 s}(X)$. This justifies our choice of $r, s$ in this case.

CASE $3(d=2$ and $p \geq 2)$. We choose $r \in(2, \infty)$ large enough that $r^{\prime}=$ $r /(r-1) \in(1,2)$ is close enough to one and choose $s \in(1, \infty)$ large enough that $1 / s+1 / r \leq 1 / r^{\prime}$. By [4, Theorem 4.12], we still have continuous Sobolev embeddings, $W^{1,2}(X) \subset L^{r}(X)$ and $W^{1, p}(X) \subset L^{2 s}(X)$. This justifies our choice of $r, s$ in this case.

The proofs of estimates (4.1.8b), 4.1.8c), 4.1.8d), (4.1.8e), and (4.1.8f) are analogous to that of estimate 4.1.8a). This completes Step 3 ,

Step $4\left(W^{-1,2}\right.$ estimates for $F_{A_{1}} \times b$ and $\left.m \varphi\right)$. We claim that

$$
\begin{gather*}
\left\|F_{A_{1}} \times b\right\|_{W_{A_{1}}^{-1,2}(X)} \leq z\left\|F_{A_{1}}\right\|_{C(X)}\|b\|_{W_{A_{1}}^{1,2}(X)},  \tag{4.1.9a}\\
\|m \varphi\|_{W_{A_{1}}^{-1,2}(X)} \leq\|m\|_{C(X)}\|\varphi\|_{W_{A_{1}}^{1,2}(X)}, \tag{4.1.9b}
\end{gather*}
$$

where $z=z(g, G, p) \in[1, \infty)$.

We consider first the estimate 4.1.9a). Using duality, together with the CauchySchwarz inequality, we estimate

$$
\begin{aligned}
\left\|F_{A_{1}} \times b\right\|_{W_{A_{1}}^{-1,2}(X)} & =\sup _{0<\|\beta\|_{W_{A_{1}(X)}^{1,2}} \leq 1}\left(F_{A_{1}} \times b, \beta\right)_{L^{2}(X)} \\
& \leq z\left\|F_{A_{1}}\right\|_{C(X)}\|b\|_{L^{2}(X)} \sup _{0<\|\beta\|_{W_{A_{1}}^{1,2}(X)} \leq 1}\|\beta\|_{L^{2}(X)} \\
& \leq z\left\|F_{A_{1}}\right\|_{C(X)}\|b\|_{W_{A_{1}}^{1,2}(X)},
\end{aligned}
$$

as required. Estimate 4.1.9b) is proved in the same way. This completes Step 4 .
The estimates obtained in each of the preceding steps verify Claim 4.1.2.
Next, we apply Proposition 3.1.7 with $k=-1$ and $p=2$ to give
Corollary 4.1.3 (Fredholm and index zero properties of the extended Hessian operator for the boson coupled Yang-Mills energy function on a Coulomb-gauge slice). Assume the hypotheses of Proposition 3.1.7. If $\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair on $(P, E)$, then the following operator is Fredholm with index zero,

$$
\begin{aligned}
\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right): \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{1,2} & \left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \\
& \rightarrow \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{-1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

where is the Hessian operator $\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)$ is defined in (3.1.17), with explicit schematic expression in (3.1.16).

Recall from (3.1.41) that

$$
\begin{aligned}
\hat{\mathscr{M}} \equiv \Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M} & : \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \ni(A, \Phi) \\
& \mapsto \Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}(A, \Phi) \in \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)
\end{aligned}
$$

is the gradient at a $W^{1, p}$ pair $(A, \Phi)$ of the restriction,

$$
\mathscr{E}:\left(A_{\infty}, \Phi_{\infty}\right)+\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow \mathbb{R}
$$

of the boson coupled Yang-Mills energy function (1.3.1) to the Coulomb-gauge slice through a $C^{\infty}$ pair $\left(A_{\infty}, \Phi_{\infty}\right)$ on $(P, E)$, where we recall that $\Pi_{A_{\infty}, \Phi_{\infty}}$ is the $L^{2}$ orthogonal projection (3.1.27) onto the Coulomb-gauge slice. The Hessian operators are related by

$$
\hat{\mathscr{M}}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)=\left(\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}\right)^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)=\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right) .
$$

Because the operator $\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)$ preserves the Coulomb-gauge condition, then

$$
\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)=\mathscr{M}^{\prime}\left(A_{\infty}, \Phi_{\infty}\right)
$$

although for arbitrary $W^{1, p}$ pairs $(A, \Phi)$, we may have $\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}^{\prime}(A, \Phi) \neq \mathscr{M}^{\prime}(A, \Phi)$.
Recall from Lemma 4.1.1 that for any $W^{1, p}$ pair on $(P, E)$, the operator

$$
\mathscr{M}_{1}(A, \Phi) \in \mathscr{L}\left(W^{1,2}, W^{-1,2}\right)
$$

was constructed to be the extension of the operator

$$
\mathscr{M}^{\prime}(A, \Phi) \in \mathscr{L}\left(W^{1, p}, W^{-1, p}\right)
$$

where we as usual abbreviate $W^{ \pm 1, p}=W_{A_{1}}^{ \pm 1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)$ and similarly for $W^{ \pm 1,2}$. Therefore,

$$
\begin{align*}
& \hat{\mathscr{M}}_{1}(A, \Phi):=\Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}_{1}(A, \Phi)  \tag{4.1.10}\\
& \in \mathscr{L}\left(\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{1,2}, \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{-1,2}\right)
\end{align*}
$$

is the corresponding extension of

$$
\hat{\mathscr{M}}^{\prime}(A, \Phi) \in \mathscr{L}\left(\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{1, p}, \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{-1, p}\right) .
$$

Note that

$$
\begin{aligned}
& \Pi_{A_{\infty}, \Phi_{\infty}} \mathscr{M}_{1}(A, \Phi)=\Pi_{A_{\infty}, \Phi_{\infty}}\left(\mathscr{M}_{1}(A, \Phi)+d_{A_{\infty}, \Phi_{\infty}} d_{A_{\infty}, \Phi_{\infty}}^{*}\right) \\
& \quad \text { on Ker } d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{-1,2} .
\end{aligned}
$$

As a consequence of Lemma 4.1.1, we have
Corollary 4.1.4 (Analyticity of the extended Hessian operator map for the boson coupled Yang-Mills energy function on a Coulomb-gauge slice). Assume the hypotheses of Theorem 4 and let $\left(A_{\infty}, \Phi_{\infty}\right)$ be a $C^{\infty}$ pair on $(P, E)$. Then the following map is well-defined and analytic,

$$
\begin{aligned}
& \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{1, p} \ni(a, \phi) \\
& \quad \mapsto \hat{\mathscr{M}}_{1}\left(A_{\infty}+a, \Phi_{\infty}+\phi\right) \in \mathscr{L}\left(\operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{1,2}, \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{-1,2}\right)
\end{aligned}
$$

Proof. The conclusion follows from Lemma 4.1.1 and the fact that the $L^{2}{ }_{-}$ orthogonal projection operators $\Pi_{A_{\infty}, \Phi_{\infty}} \in \mathscr{L}\left(W^{-1,2}, \operatorname{Ker} d_{A_{\infty}, \Phi_{\infty}}^{*} \cap W^{-1,2}\right)$ are bounded.

As we noted prior to the statement of Lemma 4.1.1 it would suffice for our application of Theorem 3 to prove Corollary 7 for the map in Corollary 4.1.4 to be continuous on a $W^{1, q}$ open neighborhood of the origin.

Finally, we shall need the following an extension of Lemma 3.1.9 that provides $W^{1,2}$ estimates for $u(a, \phi)$, thus stronger than the existing $W^{1, p^{\prime}}$ estimates when $p=d / 2$ and $p^{\prime}=d /(d-2)<2$ for $d \geq 5$.

Lemma 4.1.5 (Estimate for the action of a $W^{2, q}$ gauge transformation). Assume the hypotheses of Theorem 4. Then there is a constant $C=C(g, G, p) \in[1, \infty)$ with the following significance. If $u \in \operatorname{Aut}^{2, q}(P)$ and $(a, \phi) \in W_{A_{1}}^{1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then

$$
\begin{aligned}
\|u(a, \phi)\|_{W_{A_{1}}^{1,2}(X)} & \leq C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)\|(a, \phi)\|_{W_{A_{1}}^{1,2}(X)} \\
\|(a, \phi)\|_{W_{A_{1}}^{1,2}(X)} & \leq C\left(1+\|u\|_{W_{A_{1}}^{2, p}(X)}\right)\|u(a, \phi)\|_{W_{A_{1}}^{1,2}(X)}
\end{aligned}
$$

Proof. We recall that $u(a, \phi)=(u(a), u(\phi))=\left(u^{-1} a u, u^{-1} \phi\right)$ from the proof of Lemma 3.1.9 and so

$$
\nabla_{A_{1}}(u(a))=-u^{-1}\left(\nabla_{A_{1}} u\right) u^{-1} a u+u^{-1}\left(\nabla_{A_{1}} a\right) u+u^{-1} a\left(\nabla_{A_{1}} u\right)
$$

We make the
Claim 4.1.6. Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$. If $p \in[d / 2, \infty)$, with $p \geq 2$ when $d=2$, then there is a continuous Sobolev multiplication map,

$$
\begin{equation*}
W^{1, p}(X) \times W^{1,2}(X) \rightarrow L^{2}(X) \tag{4.1.11}
\end{equation*}
$$

Given Claim 4.1.6 and the preceding expression for $\nabla_{A_{1}}(u(a))$, we have

$$
\begin{aligned}
& \left\|\nabla_{A_{1}}(u(a))\right\|_{L^{2}(X)} \\
& \quad \leq\left\|u^{-1}\left(\nabla_{A_{1}} u\right) u^{-1} a u\right\|_{L^{2}(X)}+\left\|u^{-1}\left(\nabla_{A_{1}} a\right) u\right\|_{L^{2}(X)}+\left\|u^{-1} a\left(\nabla_{A_{1}} u\right)\right\|_{L^{2}(X)} \\
& \quad \leq\left\|\nabla_{A_{1}} u\right\|_{L^{s}(X)}\|a\|_{L^{r}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{2}(X)}+\|a\|_{L^{r}(X)}\left\|\nabla_{A_{1}} u\right\|_{L^{s}(X)} \\
& \quad \leq z\left\|\nabla_{A_{1}} u\right\|_{W_{A_{1}}^{1, p}(X)}\|a\|_{W_{A_{1}}^{1,2}(X)}+\left\|\nabla_{A_{1}} a\right\|_{L^{2}(X)},
\end{aligned}
$$

for $r, s \in(2, \infty)$ and continuous Sobolev embeddings as in the proof of Claim4.1.6 so that

$$
\|u(a)\|_{W_{A_{1}}^{1,2}(X)} \leq z\left(1+\left\|\nabla_{A_{1}} u\right\|_{W_{A_{1}}^{1, p}(X)}\right)\|a\|_{W_{A_{1}}^{1,2}(X)}
$$

for $z=z(g, p) \in[1, \infty)$. We have the analogous estimate for $\|u(\phi)\|_{W_{A_{1}}^{1,2}(X)}$ and the combination yields the required estimate for $\|u(a, \phi)\|_{W_{A_{1}}^{1,2}(X)}$.

Proof of Claim 4.1.6. The multiplication, $L^{r}(X) \times L^{s}(X) \rightarrow L^{2}(X)$, is continuous for $r, s \in(2, \infty)$ provided $1 / r+1 / s \leq 1 / 2$. If these exponents $r, s$ also yield continuous Sobolev embeddings,

$$
W^{1, p}(X) \subset L^{s}(X) \quad \text { and } \quad W^{1,2}(X) \subset L^{r}(X)
$$

then we obtain the desired continuous Sobolev multiplication map 4.1.11). To confirm the existence of suitable exponents, $r, s$, we shall separately consider the cases $d \geq 3$ with $p \geq d$, and $d \geq 3$ with $d / 2 \leq p<d$, and $d=2$ with $p \geq 2$.

CASE $1(d \geq 3$ and $p \geq d)$. We choose $r=2^{*}=2 d /(d-2) \in(2, \infty)$ and $s=d \in(2, \infty)$. By [4, Theorem 4.12], the preceding Sobolev embeddings hold and additionally

$$
1 / r+1 / s=(d-2) /(2 d)+1 / d=1 / 2
$$

so the proof of (4.1.11) is completed in this case.
CASE $2(d \geq 3$ and $d / 2 \leq p<d)$. We choose $r=2^{*}=2 d /(d-2) \in(2, \infty)$ and $s=p^{*}=(d p) /(d-p) \in(2, \infty)$. By 4. Theorem 4.12], the preceding Sobolev embeddings hold and additionally

$$
1 / r+1 / s=(d-2) /(2 d)+(d-p) /(d p)=1 / 2-1 / d+1 / p-1 / d \leq 1 / 2
$$

using $p \geq d / 2$, which we assume by hypothesis in Corollary 7 Therefore, the proof of (4.1.11) is complete in this case.

CASE $3(d=2$ and $p \geq 2)$. We choose any $r, s \in(2, \infty)$ sufficiently large so that $1 / s+1 / r \leq 1 / 2$. By [4, Theorem 4.12], the preceding Sobolev embeddings hold and the proof of (4.1.11) is complete in this case.

Combining these three cases completes the proof of Claim 4.1.6.

Symmetry yields the corresponding estimate for $\|(a, \phi)\|_{W_{A_{1}}^{1,2}(X)}$ and this completes the proof of Lemma 4.1.5.

### 4.2. Completion of the proofs of the the $W^{-1,2}$ gradient inequalities

We can now proceed to the
Proof of Corollary 7. We proceed by verifying the hypotheses of Theorem 3. As in the proof of Theorem 4 (which relied on Theorem (1), the proof of Corollary 7 is divided into two cases. In the first case, we consider the simpler situation where the pair $(A, \Phi)$ is in Coulomb gauge relative to the critical point $\left(A_{\infty}, \Phi_{\infty}\right)$ and in the second case, we consider the general situation.

Case $1\left((A, \Phi)\right.$ in Coulomb gauge relative to $\left.\left(A_{\infty}, \Phi_{\infty}\right)\right)$. By hypothesis, the $W^{1, q}$ pair $\left(A_{\infty}, \Phi_{\infty}\right)$ is a critical point for the function $\mathscr{E}$ in (1.3.1). By the regularity Theorem 2.11.2 there exists a $W^{2, q}$ gauge transformation $u_{\infty}$ such that $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$ is a $C^{\infty}$ pair; we take $u_{\infty}$ to be the identity if $\left(A_{\infty}, \Phi_{\infty}\right)$ is already $C^{\infty}$. In particular, $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$ is a $W^{2, q}$ pair and $u_{\infty}(A, \Phi)$ is in Coulomb gauge relative to $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$. As in the proof of Theorem 4 in Section 3.1.6, we have the Banach spaces,

$$
\begin{aligned}
& \mathscr{X}=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{p}(X ; \operatorname{ad} P)\right) \\
& \tilde{X}=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{-1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{-2, p}(X ; \operatorname{ad} P)\right),
\end{aligned}
$$

and in addition, in order to apply Theorem 3 we also choose Hilbert spaces,

$$
\begin{aligned}
\mathscr{G} & :=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow L^{2}(X ; \operatorname{ad} P)\right) \\
\tilde{\mathscr{G}} & :=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: W_{A_{1}}^{-1,2}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right) \rightarrow W_{A_{1}}^{-2,2}(X ; \operatorname{ad} P)\right) .
\end{aligned}
$$

Hence, $\mathscr{X} \subset \tilde{X}$ and $\mathscr{G} \subset \tilde{\mathscr{G}}$ are continuous embeddings of Banach spaces. Moreover,

$$
\mathscr{X}^{*}=\operatorname{Ker}\left(d_{u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)}^{*}: L^{p^{\prime}}(X ; \operatorname{ad} P) \rightarrow W_{A_{1}}^{-1, p^{\prime}}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P \oplus E\right)\right)
$$

and $\tilde{X} \subset \mathscr{X}^{*}$ is a continuous embedding of Banach spaces, just as in Section 3.1.6. Since $p \geq 2$ by hypothesis of Corollary 7 then $\tilde{\mathscr{X}} \subset \tilde{\mathscr{G}}$ is also a continuous embedding of Banach spaces and the compositions

$$
\mathscr{X} \subset \mathscr{G} \subset \tilde{\mathscr{G}} \quad \text { and } \quad \mathscr{X} \subset \tilde{\mathscr{X}} \subset \tilde{\mathscr{G}},
$$

induce the same embedding. By Proposition 3.1.1 the function $\mathscr{E}: x_{\infty}+\mathscr{X} \rightarrow$ $\mathbb{R}$ is analytic, where as in Section 3.1.6 we abbreviate $x_{\infty}=u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$. By Corollary 4.1.4 the Hessian operator map,

$$
x_{\infty}+\mathscr{X} \ni x \mapsto \hat{\mathscr{M}}^{\prime}(x) \in \mathscr{L}(\mathscr{X}, \tilde{\mathscr{X}}),
$$

extends to a continuous map,

$$
x_{\infty}+\mathscr{X} \ni x \mapsto \hat{\mathscr{M}}_{1}(x) \in \mathscr{L}(\mathscr{G}, \tilde{\mathscr{G}}) .
$$

By Proposition 3.1.7, the operator $\hat{\mathscr{M}}^{\prime}\left(x_{\infty}\right)$ is Fredholm with index zero and by Corollary 4.1.3, the extension $\hat{\mathscr{M}}_{1}\left(x_{\infty}\right)$ is also Fredholm with index zero. Hence, all the hypotheses of Theorem 3 are fulfilled (see also Remark (1.2.4) and so there are constants $Z^{\prime} \in(0, \infty)$ and $\sigma^{\prime} \in(0,1]$ and $\theta \in[1 / 2,1)$ (depending on $\left(A_{1}, \Phi_{1}\right)$, and $u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)$, and $\left.g, G, p, P\right)$ such that if

$$
\begin{equation*}
\left\|u_{\infty}(A, \Phi)-u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\sigma^{\prime} \tag{4.2.1}
\end{equation*}
$$

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then

$$
\begin{equation*}
\left|\mathscr{E}\left(u_{\infty}(A, \Phi)\right)-\mathscr{E}\left(u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right)\right|^{\theta} \leq Z^{\prime}\left\|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\right\|_{W^{-1,2}(X)} \tag{4.2.2}
\end{equation*}
$$

As in the proof of Theorem4, Lemma 3.1.9 gives

$$
\left\|u_{\infty}(A, \Phi)-u_{\infty}\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)} \leq C_{1}\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}
$$

for $C_{1}=C_{1}\left(g, G, p, u_{\infty}\right)=C_{1}\left(A_{\infty}, \Phi_{\infty}, g, G, p\right) \in(0,1]$. Therefore, setting $\sigma:=$ $C_{1}^{-1} \sigma^{\prime}$, we see that if $(A, \Phi)$ obeys the Łojasiewicz-Simon neighborhood condition (1.3.11), namely

$$
\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\sigma
$$

then (4.2.1) holds and thus also (4.2.2). Moreover,

$$
\begin{aligned}
\left\|\hat{\mathscr{M}}\left(u_{\infty}(A, \Phi)\right)\right\|_{W_{A_{1}}^{-1,2}(X)} & \leq C_{2}\|\hat{\mathscr{M}}(A, \Phi)\|_{W_{A_{1}}^{-1,2}(X)}, \\
\|\hat{\mathscr{M}}(A, \Phi)\|_{W_{A_{1}}^{-1,2}(X)} & \leq\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1,2}(X)}
\end{aligned}
$$

just as in the proof of the corresponding $W^{-1, p}$ estimate in the proof of Theorem 4 but now using Lemma4.1.5 with constant $C_{2}=C_{2}\left(g, G, u_{\infty}\right)=C_{2}\left(A_{\infty}, \Phi_{\infty}, g, G\right) \in$ ( 0,1 ].

Increasing $C_{2}\left(A_{\infty}, \Phi_{\infty}, g, G, p\right) \in(0,1]$ if necessary to give $C_{2} \geq C_{1}$ and substituting the preceding inequality into (4.2.2) and following the same steps as in the proof of Theorem 4 yields

$$
\left|\mathscr{E}(A, \Phi)-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \leq Z^{\prime} C_{2}\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1,2}(X)}
$$

Hence, the Łojasiewicz-Simon gradient inequality (1.3.15) holds for the pairs $(A, \Phi)$ and $\left(A_{\infty}, \Phi_{\infty}\right)$ with constants $(Z, \theta, \sigma)$, where $Z:=C_{2} Z^{\prime}$.

Case $2\left((A, \Phi)\right.$ not in Coulomb gauge relative to $\left.\left(A_{\infty}, \Phi_{\infty}\right)\right)$. Let

$$
\begin{aligned}
\zeta & =\zeta\left(A_{1}, A_{\infty}, \Phi_{\infty}, g, G, p, q\right) \in(0,1] \quad \text { and } \\
N & =N\left(A_{1}, A_{\infty}, \Phi_{\infty}, g, G, p, q\right) \in[1, \infty)
\end{aligned}
$$

denote the constants in Theorem 16 and choose $\zeta_{1} \in(0, \zeta]$ small enough that $2 N \zeta_{1}<\sigma_{1}$, where we now use $\sigma_{1}$ to denote the Lojasiewicz-Simon constant from Case 1. If $(A, \Phi)$ obeys

$$
\left\|(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<\zeta_{1}
$$

then Theorem 16 provides $u \in \operatorname{Aut}^{2, q}(P)$, depending on the pair $(A, \Phi)$, such that

$$
\begin{gathered}
d_{A_{\infty}, \Phi_{\infty}}^{*}\left(u(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right)=0 \\
\left\|u(A, \Phi)-\left(A_{\infty}, \Phi_{\infty}\right)\right\|_{W_{A_{1}}^{1, p}(X)}<2 N \zeta_{1}<\sigma .
\end{gathered}
$$

By applying Case 1 to the pairs $u(A, \Phi)$ and $\left(A_{\infty}, \Phi_{\infty}\right)$, we obtain

$$
\left|\mathscr{E}(u(A, \Phi))-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \leq C_{2} Z^{\prime}\|\mathscr{M}(u(A, \Phi))\|_{W_{A_{1}}^{-1,2}(X)}
$$

Estimating as in Case 1 with $u$ replacing $u_{\infty}$, we see that

$$
\|\mathscr{M}(u(A, \Phi))\|_{W_{A_{1}}^{-1,2}(X)} \leq C_{3}\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1,2}(X)},
$$

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just as in the proof of Theorem 4 where $C_{3}=C_{3}\left(A_{1}, A_{\infty}, \Phi_{\infty}, g, G, p, q\right) \in(0,1]$. By combining the preceding inequalities, we obtain

$$
\begin{aligned}
\left|\mathscr{E}(A, \Phi)-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} & =\left|\mathscr{E}(u(A, \Phi))-\mathscr{E}\left(A_{\infty}, \Phi_{\infty}\right)\right|^{\theta} \\
& \leq C_{2} C_{3} Z^{\prime}\|\mathscr{M}(A, \Phi)\|_{W_{A_{1}}^{-1,2}(X)}
\end{aligned}
$$

Hence, we obtain the Łojasiewicz-Simon gradient inequality (1.3.12) with constants $(Z, \theta, \sigma)$, where we now choose $Z=C_{2} C_{3} Z^{\prime}$ and $\sigma=\zeta_{1}$. This completes the proof of this case.

This completes the proof of Corollary 7
Proof of Corollary 8, The argument follows mutatis mutandis that for Corollary 7 .

## APPENDIX A

## Fredholm and index properties of elliptic operators on Sobolev spaces

We note the following corollary of the standard Fredholm property [48, Lemma 1.4.5], 67, Theorem 19.2.1] of elliptic pseudo-differential operators on Sobolev spaces $W^{k+m, p}(M ; V)$ when ${ }^{1} k \in \mathbb{Z}$ and $p=2$ and which may be compared with 117. Theorem 5.7.2]. We include a detailed proof of Theorem A. 1 because we are unaware of a published reference. The statement of Theorem A. 1 both generalizes 67, Theorem 19.2.1] by allowing $1<p<\infty$ rather than $p=2$ and specializes [67, Theorem 19.2.1] by restricting to partial differential operators (thus of integer order $m$ ) and restricting $k$ to be an integer. In our application, we shall only need the case $m=2$, but we provide the more general version for the benefit of other applications. We leave it to the reader to consider extensions to the general case of pseudo-differential operators with real order $m$ on $W^{k, p}$ spaces with real $k$.

Theorem A. 1 (Fredholm property for elliptic partial differential operators on Sobolev spaces). (See [32, Lemma 41.1].) Let $V$ and $W$ be finite-rank, smooth vector bundles over a closed, smooth manifold, $M$, and $k \in \mathbb{Z}$ an integer, and $p \in(1, \infty)$. If $P: C^{\infty}(M ; V) \rightarrow C^{\infty}(M ; W)$ is an elliptic linear partial differential operator of integer order $m \geq 1$ with $C^{\infty}$ coefficients, then $P: W^{k+m, p}(M ; V) \rightarrow$ $W^{k, p}(M ; W)$ is a Fredholm operator with index,

$$
\begin{aligned}
\text { Index } P= & \operatorname{dim} \operatorname{Ker}\left(P: C^{\infty}(M ; V) \rightarrow C^{\infty}(M ; W)\right) \\
& -\operatorname{dim} \operatorname{Ker}\left(P^{*}: C^{\infty}\left(M ; W^{*}\right) \rightarrow C^{\infty}\left(M ; V^{*}\right)\right),
\end{aligned}
$$

where $P^{*}$ is the formal $\left(L^{2}\right)$ adjoint of $P$, and has range,

$$
\operatorname{Ran}\left(P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)\right)=\left(K^{*}\right)^{\perp} \cap W^{k, p}(M ; W)
$$

where $\perp$ denotes $L^{2}$-orthogonal complement and

$$
K=\operatorname{Ker}\left(P^{*}: C^{\infty}\left(M ; W^{*}\right) \rightarrow C^{\infty}\left(M ; V^{*}\right)\right) \cong K^{*} \subset C^{\infty}(M ; W)
$$

If $V=W^{*}$ and $P-P^{*}: C^{\infty}(M ; V) \rightarrow C^{\infty}(M ; W)$ is a differential operator of order $m-1$, then Index $P=0$.

Remark A. 2 (On the index of elliptic operators with scalar principal symbol). We recall from [113, Proposition 2.4] that if $E$ is a smooth complex vector bundle over a closed, smooth manifold, $M$, and $P: C^{\infty}(M ; E) \rightarrow C^{\infty}(M ; E)$ is an elliptic differential operator whose principal symbol is a scalar multiple of the identity (as often occurs in Geometric Analysis), then Index $P=0$. However, if $P: C^{\infty}(M ; E) \rightarrow C^{\infty}(M ; E)$ is more generally an elliptic pseudo-differential operator, then $P$ need not have index zero [113, Example 2.2].

[^7]We refer to Gilkey 48] and Hörmander 67] for the definitions (over domains $\Omega \subseteq \mathbb{R}^{d}$ ) of a differential operator $P$ of integer order $m \geq 1$ [48, Section 1.1], a pseudo-differential operator $P$ of order $m \in \mathbb{R}$ [48, Section 1.2.1], the vector spaces $\Psi^{m}$ of pseudo-differential operators of order $m \in \mathbb{R}$ and $\Psi^{-\infty}=\cap_{m \in \mathbb{R}} \Psi^{m}$ of infinitely smoothing pseudo-differential operators [48, Section 1.2.1], the $L^{2}$ adjoint $P^{*}$ and composition $P Q$ of pseudo-differential operators [48, Section 1.2.2], the extension of the definition of pseudo-differential operators on complex-valued functions over domains in $\mathbb{R}^{d}$ to sections of finite-rank vector bundles over domains in $\mathbb{R}^{d}$ 48, Section 1.2.7], elliptic pseudo-differential operators on sections of finiterank vector bundles over domains in $\mathbb{R}^{d}$ [48, Section 1.3.1], and elliptic pseudodifferential operators on sections of finite-rank vector bundles over closed manifolds 48, Sections 1.3.3 and 1.3.5]. Similarly, for $s \in \mathbb{R}$, we refer to Gilkey 48 and Hörmander 67 for the definitions of Sobolev spaces $W^{s, 2}\left(\mathbb{R}^{d} ; \mathbb{C}\right)$ 48, Section 1.1.3], $W^{s, 2}(\Omega ; \mathbb{C})$ 48, Section 1.2.6], and $W^{s, 2}(M ; \mathbb{C})$ and $W^{s, 2}(M ; V)$ 48, Sections 1.3.4 and 1.3.5].

When $k \in \mathbb{N}$ is a non-negative integer and $1 \leq p \leq \infty$, we use the standard definitions (without appealing to symbols or the symbol calculus) of Sobolev spaces, $W^{k, p}(\Omega ; \mathbb{R})$ and $W_{0}^{k, p}(\Omega ; \mathbb{R})$, in Adams and Fournier [4, Section 3.2] for domains $\Omega \subseteq \mathbb{R}^{d}$ and Aubin [6, Section 2.3] and Gilkey [48, Sections 1.3.4 and 1.3.5] for their extensions to functions on closed manifolds, $W^{k, p}(M ; \mathbb{R})$, and sections of finiterank vector bundles over closed manifolds, $W^{k, p}(M ; V)$. For integers $k<0$ and $1<p<\infty$, we follow [4, Sections 3.9 to 3.14] and by analogy define

$$
W^{k, p}(M ; V):=\left(W^{-k, p^{\prime}}\left(M ; V^{*}\right)\right)^{*},
$$

where the dual Hölder exponent, $1<p^{\prime}<\infty$, is defined by $1 / p+1 / p^{\prime}=1$ and $\left(W^{-k, p^{\prime}}\left(M ; V^{*}\right)\right)^{*}$ is the continuous dual of the Banach space, $W^{-k, p^{\prime}}\left(M ; V^{*}\right)$.

For a differential operator, $P$, of integer order $m \geq 1$, any integer $k \in \mathbb{Z}$ and $v \in W^{k+m, p}(M ; V)$, one may define $P v \in\left(C^{\infty}(M ; V)\right)^{*}$ in the sense of distributions [4. Section 1.62] by

$$
(P v)(w):=\left(v, P^{*} w\right)_{L^{2}(M)}, \quad \forall w \in C^{\infty}\left(M ; W^{*}\right)
$$

where the $L^{2}$-adjoint, $P^{*}: C^{\infty}\left(M ; W^{*}\right) \rightarrow C^{\infty}\left(M ; V^{*}\right)$, is also a differential operator of order $m \geq 1$. For integer $k \geq 0$, then $P v$ has its usual meaning and we have $P v \in W^{k, p}(M ; V)$, with

$$
P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)
$$

defining a bounded operator by definition of the Sobolev spaces, $W^{k+m, p}(M ; V)$. Similarly, for integers $k \leq-m$, we may define $P v \in W^{k, p}(M ; V)$ by

$$
(P v)(w):=\left(v, P^{*} w\right)_{L^{2}(M)}, \quad \forall v \in W^{k+m, p}(M ; V) \text { and } w \in W^{-k, p^{\prime}}\left(M ; W^{*}\right)
$$

noting that $P^{*}: W^{-k, p^{\prime}}\left(M ; W^{*}\right) \rightarrow W^{-k-m, p^{\prime}}\left(M ; V^{*}\right)$ is a bounded operator and so $P^{*} w \in W^{-k-m, p^{\prime}}\left(M ; V^{*}\right)$ with its usual meaning.

Finally, for $m \geq 2$ and integers $k$ in the range $-m+1 \leq k \leq-1$, we choose a Riemannian metric, $g$, on $M$, a connection, $\nabla: C^{\infty}(M ; V) \rightarrow C^{\infty}\left(M ; T^{*} M \otimes V\right)$, and form the augmented connection Laplace operator,

$$
\nabla^{*} \nabla+1: C^{\infty}(M ; V) \rightarrow C^{\infty}(M ; V)
$$

The resulting operator, $\nabla^{*} \nabla+1: W^{l+2, p}(M ; V) \rightarrow W^{l, p}(M ; V)$, is invertible for integers $l \geq 0$ and $1<p<\infty$. By duality (as above) it extends to an invertible
operator, $\nabla^{*} \nabla+1: W^{l+2, p}(M ; V) \rightarrow W^{l, p}(M ; V)$, for integers $l \leq-2$ and $1<$ $p<\infty$ and because it has divergence form, it also defines an invertible operator, $\nabla^{*} \nabla+1: W^{1, p}(M ; V) \rightarrow W^{-1, p}(M ; V)$, for $1<p<\infty$. (Observe that the bilinear map,

$$
W^{1, p}(M ; V) \times W^{1, p^{\prime}}(M ; V) \ni(u, v) \mapsto\left(\left(\nabla^{*} \nabla+1\right) u, v\right)_{L^{2}(X)}
$$

is continuous since

$$
\left.\mid\left(\nabla^{*} \nabla+1\right) u, v\right)_{L^{2}(X)} \mid \leq\|\nabla u\|_{L^{p}(X)}\|\nabla v\|_{L^{p^{\prime}}(X)}+\|u\|_{L^{p}(X)}\|v\|_{L^{p^{\prime}}(X)}
$$

We therefore obtain a continuous, linear map,

$$
\nabla^{*} \nabla+1: W^{1, p}(M ; V) \rightarrow\left(W^{1, p^{\prime}}\left(M ; V^{*}\right)\right)^{*}=W^{-1, p}(M ; V)
$$

which is clearly injective. The map is surjective by [4, Sections 3.5-3.14] because every $\alpha \in\left(W^{1, p^{\prime}}\left(M ; V^{*}\right)\right)^{*}$ is represented by

$$
\alpha(v)=\left(\nabla v, u_{1}\right)_{L^{2}(X)}+\left(v, u_{2}\right)_{L^{2}(X)}
$$

for $u_{i} \in L^{p}(M ; V)$ for $i=1,2$. If $p=2$, then we find that $u_{2}=u$ and $u_{1}=\nabla u$ for $u \in W^{1,2}(M ; V)$ (compare [47, Section 8.2]) and thus also for $2 \leq p<\infty$; duality yields the case $1<p \leq 2$. Hence, the map is an isomorphism by the Open Mapping Theorem.)

For $m=2 n$ and $n \in \mathbb{N}$, we may thus define $P v \in W^{k, p}(M ; V)$ by

$$
\begin{aligned}
(P v)(w):=\left(\left(\nabla^{*} \nabla+1\right)^{m / 2} v\right. & \left., P^{*} w\right)_{L^{2}\left(M ; W^{*}\right)} \\
& \forall v \in W^{k+m, p}(M ; V) \text { and } w \in W^{-k+m, p^{\prime}}(M ; V)
\end{aligned}
$$

noting that $\left(\nabla^{*} \nabla+1\right)^{m / 2} v \in W^{k, p}(M ; V)$ and $P^{*} w \in W^{-k, p^{\prime}}(M ; V)$. For $m=2 n+$ 1 and $n \in \mathbb{N}$, one must first define the square root22, $\left(\nabla^{*} \nabla+1\right)^{1 / 2}: W^{l+1, p}(M ; V) \rightarrow$ $W^{l, p}(M ; V)$, for integers $l \in \mathbb{Z}$ and $1<p<\infty$; see Feehan 32 and references cited therein for a survey of approaches to the definition of the fractional powers, $\left(\nabla^{*} \nabla+1\right)^{s}$, for $s \in \mathbb{R}$. Once the square root of the augmented connection Laplace operator is defined, then the definition of $P v \in W^{k, p}(M ; V)$ is identical to the case $m=2 n$.

Proof of Theorem A.1. There is a pseudo-differential operator

$$
S: C^{\infty}(M ; W) \rightarrow C^{\infty}(M ; V)
$$

of order $-m$, according to [48, Lemma 1.3.6], so that

$$
S P-I \in \Psi^{-\infty}(M ; V) \quad \text { and } \quad P S-I \in \Psi^{-\infty}(M ; W)
$$

where $\Psi^{-\infty}(M ; W)$ and $\Psi^{-\infty}(M ; V)$ are the vector spaces of infinitely smoothing pseudo-differential operators [48, Sections 1.2 and 1.3].

The operator $P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)$ is continuous since $P$ is an elliptic partial differential operator of order $m \geq 1$ and by definition of the Sobolev space $W^{k+m, p}(M ; V)$. Combining the preceding observation with the $a$ priori elliptic estimate (see [32, Theorem 14.60], for example),

$$
\|v\|_{W^{k+m, p}(M)} \leq C\left(\|v\|_{W^{k, p}(M)}+\|P v\|_{W^{k, p}(M)}\right)
$$

[^8]implies that the expression $\|v\|_{W^{k, p}(M)}+\|P v\|_{W^{k, p}(M)}$ defines a norm for $v \in$ $C^{\infty}(M ; V)$ which is equivalent to the standard norm on $W^{k+m, p}(M ; V)$. Furthermore, the following operators are continuous,
\[

$$
\begin{aligned}
S P-I: W^{k, p}(M ; V) & \rightarrow W^{k+1, p}(M ; V) \\
P S-I: W^{k, p}(M ; W) & \rightarrow W^{k+1, p}(M ; W)
\end{aligned}
$$
\]

because, denoting $H^{l}(M ; V)=W^{l, 2}(M ; V)$ and $H^{l}(M ; W)=W^{l, 2}(M ; W)$, the operators

$$
\begin{gathered}
S P-I: H^{l}(M ; V) \rightarrow H^{l+n}(M ; V) \\
P S-I: H^{l}(M ; W) \rightarrow H^{l+n}(M ; W)
\end{gathered}
$$

are continuous by 48, Lemma 1.3.5], for any integers $l, n \in \mathbb{Z}$, and the standard Sobolev Embedding [4, Theorem 4.12] and duality [4, Sections 3.5-3.14], which provide continuous embeddings, $W^{k, p}(M ; V) \subset H^{l}(M ; V)$ for sufficiently small $l \leq$ $l_{0}(d, k, p)$ and $H^{l+n}(M ; W) \subset W^{k+1, p}(M ; W)$ for sufficiently large $n \geq n_{0}(d, k, p)$, where $d$ is the dimension of $M$. (The finite integers $l_{0}(d, k, p)$ and $n_{0}(d, k, p)$ can of course be determined explicitly from the full statement of the Sobolev Embedding [4. Theorem 4.12], but their precise values are unimportant here.) Thus,

$$
\begin{aligned}
\|S w\|_{W^{k+m, p}(M ; V)} & \leq C\left(\|P S w\|_{W^{k, p}(M ; W)}+\|w\|_{W^{k, p}(M ; W)}\right) \\
& \leq C\left(\|(P S-I) w\|_{W^{k, p}(M ; W)}+\|w\|_{W^{k, p}(M ; W)}\right) \\
& \left.\leq C\|w\|_{W^{k, p}(M ; W)} \quad \text { (by continuity of } P S-I \text { on } W^{k, p}(M ; W)\right)
\end{aligned}
$$

and so the operator $S: W^{k, p}(M ; V) \rightarrow W^{k+m, p}(M ; W)$ is continuous.
The embeddings

$$
W^{k+1, p}(M ; V) \Subset W^{k, p}(M ; V) \quad \text { and } \quad W^{k+1, p}(M ; W) \Subset W^{k, p}(M ; W)
$$

are compact by the Rellich-Kondrachov Theorem [4, Theorem 6.3] when $k \geq 0$ and when $k \leq-1$ using duality [14, Theorem 6.4] and compactness of the embeddings

$$
\begin{aligned}
W^{-k, p^{\prime}}\left(M ; V^{*}\right) & \Subset W^{-k-1, p^{\prime}}\left(M ; V^{*}\right) \quad \text { and } \\
W^{-k, p^{\prime}}\left(M ; W^{*}\right) & \Subset W^{-k-1, p^{\prime}}\left(M ; W^{*}\right) .
\end{aligned}
$$

Hence, the operators (now viewed as compositions with compact embeddings),

$$
\begin{gathered}
S P-I: W^{k, p}(M ; V) \rightarrow W^{k, p}(M ; V) \\
P S-I: W^{k, p}(M ; W) \rightarrow W^{k, p}(M ; W)
\end{gathered}
$$

are compact by [14, Proposition 6.3]. Thus, $P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)$ is Fredholm by 48, Section 1.4.2, Definition, p. 38].

Because $P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)$ is Fredholm, its index may be computed by [2, Lemma 4.38],

$$
\begin{aligned}
\text { Index } P= & \operatorname{dim} \operatorname{Ker}\left(P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)\right) \\
& -\operatorname{dim} \operatorname{Ker}\left(P^{*}: W^{-k, p^{\prime}}\left(M ; W^{*}\right) \rightarrow W^{-k-m, p^{\prime}}\left(M ; V^{*}\right)\right) .
\end{aligned}
$$

If $w \in\left(\operatorname{Ker} P^{*}\right) \cap W^{-k, p^{\prime}}\left(M ; W^{*}\right)$, then $w \in\left(\operatorname{Ker} P^{*}\right) \cap H^{l}\left(M ; W^{*}\right)$ for all integers $l \leq l_{0}(d, k, p)$, where $l_{0}$ is given by the Sobolev Embedding [4, Theorem 4.12], and consequently (since the Banach space dual $P^{*}$ is defined by the realization of the formal adjoint and $P^{*}$ is thus an elliptic partial differential operator of order
$m$ ) we see that $w \in\left(\operatorname{Ker} P^{*}\right) \cap C^{\infty}\left(M ; W^{*}\right)$ by elliptic regularity 48, Lemma 1.3.2 and Section 1.3.5]. Of course, if $v \in(\operatorname{Ker} P) \cap W^{k+m, p}(M ; V)$, then $v \in$ $(\operatorname{Ker} P) \cap C^{\infty}(M ; V)$ by the same argument. This yields the stated formula for the index of $P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)$.

If $K \subset W^{-k, p^{\prime}}\left(M ; W^{*}\right)$ is any subspace, we recall that the annihilator $\mathbf{1 0 1}$ Section 4.6] of $K$ in $\left(W^{-k, p^{\prime}}\left(M ; W^{*}\right)\right)^{*}=W^{k, p}(M ; W)$ is

$$
K^{\circ}=\left\{w \in W^{k, p}(M ; W):\langle\alpha, w\rangle=0, \forall \alpha \in K\right\}
$$

and $\langle\cdot, \cdot\rangle: W^{-k, p^{\prime}}\left(M ; W^{*}\right) \times\left(W^{-k, p^{\prime}}\left(M ; W^{*}\right)\right)^{*} \rightarrow \mathbb{R}$ is the canonical pairing. We can therefore identify the range of $P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)$ using

$$
\begin{aligned}
\operatorname{Ran} & \left(P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)\right) \\
= & \overline{\operatorname{Ran}}\left(P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)\right) \quad \text { (by closed range) } \\
= & \operatorname{Ker}\left(P^{*}: W^{-k, p^{\prime}}\left(M ; W^{*}\right) \rightarrow W^{-k-m, p^{\prime}}\left(M ; V^{*}\right)\right)^{\circ} \\
& (\text { by }[\mathbf{1 4}, \text { Corollary } 2.18(\text { iv }))]) .
\end{aligned}
$$

If $K:=\operatorname{Ker}\left(P^{*}: W^{-k, p^{\prime}}\left(M ; W^{*}\right) \rightarrow W^{-k-m, p^{\prime}}\left(M ; V^{*}\right)\right)$, then $K \subset C^{\infty}\left(M ; W^{*}\right)$, as we observed by elliptic regularity and

$$
\begin{aligned}
& K^{\circ}=\left\{w \in W^{k, p}(M ; W): w(\alpha)=0, \forall \alpha \in K\right\} \\
&=\left\{w \in W^{k, p}(M ; W): \alpha(w)=0, \forall \alpha \in K\right\} \\
&\left(\operatorname{by}\left(W^{k, p}(M ; W)\right)^{* *} \cong W^{k, p}(M ; W)\right) \\
&=\left\{w \in W^{k, p}(M ; W): \iota(\kappa)(w)=0, \forall \kappa \in K^{*}\right\} \\
&=\left\{w \in W^{k, p}(M ; W):(w, \kappa)_{\left.L^{2}(M ; W)=0, \forall \kappa \in K^{*}\right\}}^{=}\right. \\
&\left(K^{*}\right)^{\perp} \cap W^{k, p}(M ; W),
\end{aligned}
$$

where $K^{*} \subset C^{\infty}(M ; W)$ and $\iota: K^{*} \cong K^{* *}=K$ is the canonical isomorphism, $\kappa \mapsto \iota(\kappa)=(\cdot, \kappa)_{L^{2}(M ; W)}$. This yields the claimed identification of the range of $P$.

Noting that $\left(W^{-k, p^{\prime}}(M ; V)\right)^{*}=W^{k, p}\left(M ; V^{*}\right)$ and $\left(W^{-k-m, p^{\prime}}(M ; W)\right)^{*}=$ $W^{k+m, p}\left(M ; W^{*}\right)$, we see that the operator,

$$
P^{*}: W^{k+m, p}\left(M ; W^{*}\right) \rightarrow W^{k, p}\left(M ; V^{*}\right)
$$

is Fredholm by 48, Lemma 1.4.3], since the same is true for the operator,

$$
P: W^{-k, p^{\prime}}(M ; V) \rightarrow W^{-k-m, p^{\prime}}(M ; W)
$$

If $V=W^{*}$, we may write $P^{*}=P+\left(P^{*}-P\right)$ and observe that $P^{*}-P$ : $W^{k+m, p}(M ; V) \rightarrow W^{k+1, p}(M ; W)$ is bounded since $P^{*}-P$ has order $m-1$. The inclusion, $W^{k+1, p}(M ; W) \Subset W^{k, p}(M ; W)$, is compact by 4, Theorem 6.3], so $P^{*}-P: W^{k+m, p}(M ; V) \rightarrow W^{k, p}(M ; W)$ is compact. Hence, $\operatorname{Index}\left(P^{*}\right)=\operatorname{Index} P$ by [67, Corollary 19.1.8]. But we also have Index $P^{*}=-\operatorname{Index} P$ by 48, Lemma 1.4.4 (a)] and thus Index $P=0$.

The following corollary is partially based on Abramovich and Aliprantis [2, Theorem 4.46], Melrose [91, Lecture 9, Proposition 18] and Treves [113, Theorem 2.4].

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Corollary A. 3 (Green's operators for elliptic partial differential operators on Sobolev spaces). Assume the hypotheses of Theorem A. 1 and let $G: W^{k, p}(M ; W) \rightarrow$ $W^{k+m, p}(M ; V)$ be the Green's operator for $P$ defined by

$$
G w:= \begin{cases}P^{-1} w, & w \in \operatorname{Ran} P \subset W^{k, p}(M ; W) \\ 0, & w \in(\operatorname{Ran} P)^{\perp} .\end{cases}
$$

Then $G: C^{\infty}(M ; W) \rightarrow C^{\infty}(M ; V)$ is an elliptic pseudo-differential operator of order $-m$ with

$$
G P=\mathrm{id}-\Pi_{1} \quad \text { and } \quad P G=\mathrm{id}-\Pi_{2},
$$

where $\Pi_{1}$ is the $L^{2}$-orthogonal projection onto $\operatorname{Ker} P \subset W^{k+m, p}(M ; V)$ and $\Pi_{2}$ is the $L^{2}$-orthogonal projection onto $(\operatorname{Ran} P)^{\perp} \subset W^{k, p}(M ; W)$. Moreover, if $k \geq 0$, the following operator is bounded,

$$
G: W^{k, p}(M ; W) \rightarrow W^{k+m, p}(M ; V)
$$

Proof. From Melrose [91, Lecture 6, Proposition 11] or 91, Lecture 7, Theorem 3] or Treves [113 Theorem 2.4], there exists an elliptic pseudo-differential operator, $G: C^{\infty}(M ; W) \rightarrow C^{\infty}(M ; V)$, of integer order $-m$ such that id $-G P$ and id $-P G$ are infinitely smoothing operators. We then appeal to Melrose 91 , Lecture 9, Proposition 18] to conclude that id $-G P$ and id $-P G$ can be identified as the stated $L^{2}$-orthogonal projections.

For $k \geq 0$ and any $w \in W^{k, p}(M ; W)$, we have

$$
\begin{aligned}
\|G w\|_{W^{k+m, p}(M)} & \left.\leq C\left(\|P G w\|_{W^{k, p}(M)}+\|w\|_{W^{k, p}(M)}\right) \quad \text { (by [32, Theorem 14.60] }\right) \\
& \leq C\left(\left\|\left(\operatorname{id}-\Pi_{2}\right) w\right\|_{W^{k, p}(M)}+\|w\|_{W^{k, p}(M)}\right) \\
& \leq C\|w\|_{W^{k, p}(M)}
\end{aligned}
$$

and so the conclusion on boundedness follows.

## APPENDIX B

## Equivalence of Sobolev norms defined by Sobolev and smooth connections

Suppose that $(X, g)$ is a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ is a compact Lie group, and $P$ is a smooth principal $G$-bundle over $X$. In standard references for gauge theory [27, 43], it is generally assumed in the construction of Sobolev completions of spaces such as $\Omega^{l}(X ; \operatorname{ad} P)$ that one defines Sobolev norms using a covariant derivative $\nabla_{A}$ determined by a connection $A$ on $P$ that is smooth or of class $W^{k, p}$ for $p \geq 1$ and an integer $k \geq 1$ large enough that $k p>d$ or even $k p \gg d$. However, in this monograph, we often consider connections $A$ with more borderline regularity, for example of class $W^{1, q}$ for $q>d / 2$, and in that situation, one must exercise care in the definition of Sobolev spaces using such connections. Lemmas B. 1 and B. 2 provide some guidance.

Lemma B. 1 (Second-order Kato inequality and second-order Sobolev norms). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a Lie group, and $P$ be a smooth principal $G$-bundle over $X$, and $V=P \times_{\varrho} \mathbb{V}$ be a smooth Riemannian vector bundle over $X$ defined by a finite-dimensional, orthogonal representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{R}}(\mathbb{V})$. Then there exists a constant $C=$ $C(g, q) \in[1, \infty)$ with the following significance. If $A$ is a $W^{1, q}$ connection on $P$ with $q>d / 2$, then for all $v \in C^{\infty}(X ; V)$,

$$
\begin{equation*}
\|v\|_{C(X)} \leq C\|v\|_{W_{A}^{2, q}(X)} \tag{B.1}
\end{equation*}
$$

Proof. The first-order analogue of (B.1), namely,

$$
\|v\|_{C(X)} \leq \kappa_{1}\|v\|_{W_{A}^{1, q}(X)}
$$

when $q>d$ and $\kappa_{1}=\kappa_{1}(g) \in[1, \infty)$ is the norm of the Sobolev embedding $W^{1, q}(X) \subset C(X)$, is an immediate consequence of the pointwise first-order Kato inequality, $|\nabla| v\left|\left|\leq\left|\nabla_{A} v\right|\right.\right.$ from [43, Inequality (6.20)], in turn a consequence of the compatibility of the fiber metric on $V$ with $\nabla_{A}$.

We first note that, for $f \in C^{\infty}(X ; \mathbb{R})$, the norm

$$
\|f\|_{W^{2, q}(X)}=\left\|\nabla^{2} f\right\|_{L^{q}(X)}+\|\nabla f\|_{L^{q}(X)}+\|f\|_{L^{q}(X)}
$$

is equivalent (with respect to a constant depending at most on $(g, q)$ ) by virtue of [47, Theorem 9.11] to

$$
\|f\|_{W^{2, q}(X)}=\|\Delta f\|_{L^{q}(X)}+\|f\|_{L^{q}(X)},
$$

where $\Delta$ is the Laplace operator defined by the Riemannian metric $g$ on $X$. Now recall the pointwise identity [43, Equation (6.18)],

$$
\Delta|v|^{2}=2\left\langle\nabla_{A}^{*} \nabla_{A} v, v\right\rangle-2\left|\nabla_{A} v\right|^{2} .
$$

Hence, letting $\kappa_{2}=\kappa_{2}(g) \in[1, \infty)$ denote the norm of the Sobolev embedding $W^{2, q}(X) \subset C(X)$,

$$
\begin{aligned}
\|v\|_{C(X)}^{2} & =\left\||v|^{2}\right\|_{C(X)} \\
& \leq \kappa_{2}\left\||v|^{2}\right\|_{W^{2, q}(X)} \\
& =\kappa_{2}\left(\left\|\Delta|v|^{2}\right\|_{L^{q}(X)}+\left\||v|^{2}\right\|_{L^{q}(X)}\right) \\
& \leq \kappa_{2}\left(2\|v\|_{C(X)}\left\|\nabla_{A}^{*} \nabla_{A} v\right\|_{L^{q}(X)}+2\left\|\nabla_{A} v\right\|_{L^{2 q}(X)}^{2}+\|v\|_{C(X)}\|v\|_{L^{q}(X)}\right)
\end{aligned}
$$

Recall that $W^{1, q}(X) \subset L^{2 q}(X)$, for $q<d$, if and only if $2 q \leq q^{*}=d q /(d-q)$, that is, $2(d-q) \leq d$ or $q \geq d / 2$; the embedding is immediate from 4, Theorem 4.12] when $q \geq d$. Thus, applying the first-order Kato Inequality and the preceding Sobolev embedding for functions,

$$
\||v|\|_{L^{2 q}(X)} \leq \kappa_{1}\left(\|\nabla|v|\|_{L^{q}(X)}+\|v\|_{L^{q}(X)}\right) \leq \kappa_{1}\left(\left\|\nabla_{A} v\right\|_{L^{q}(X)}+\|v\|_{L^{q}(X)}\right)
$$

we obtain

$$
\begin{aligned}
\|v\|_{C(X)}^{2} \leq & \kappa_{2}\left(2\|v\|_{C(X)}\left\|\nabla_{A}^{*} \nabla_{A} v\right\|_{L^{q}(X)}+2 \kappa_{1}^{2}\left(\left\|\nabla_{A}^{2} v\right\|_{L^{q}(X)}+\left\|\nabla_{A} v\right\|_{L^{q}(X)}\right)^{2}\right. \\
& \left.+\|v\|_{C(X)}\|v\|_{L^{q}(X)}\right) .
\end{aligned}
$$

We now use Young's Inequality $2 a b \leq \varepsilon a^{2}+\varepsilon^{-1} b^{2}$ from [47, Inequality (7.8)] and rearrangement with a suitably small and universal $\varepsilon$ to give

$$
\|v\|_{C(X)}^{2} \leq C^{2}\left(\left\|\nabla_{A}^{*} \nabla_{A} v\right\|_{L^{q}(X)}^{2}+\left\|\nabla_{A}^{2} v\right\|_{L^{q}(X)}^{2}+\left\|\nabla_{A} v\right\|_{L^{q}(X)}^{2}+\|v\|_{L^{q}(X)}^{2}\right)
$$

where $C=C(g, q) \in[1, \infty)$. We simplify the right-hand side in the preceding inequality via

$$
\left\|\nabla_{A}^{*} \nabla_{A} v\right\|_{L^{q}(X)} \leq z\left(\left\|\nabla_{A}^{2} v\right\|_{L^{q}(X)}+\left\|\nabla_{A} v\right\|_{L^{q}(X)}+\|v\|_{L^{q}(X)}\right)
$$

where $z$ is a constant depending at most on the Riemannian metric on $X$. The desired Sobolev inequality (B.1) now follows by taking square roots.

Lemma B. 2 (Equivalence of Sobolev norms defined by Sobolev and smooth connections). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$, and $V=P \times_{\varrho} \mathbb{V}$ be a smooth Riemannian vector bundle over $X$ defined by a finite-dimensional, orthogonal representation, $\varrho: G \rightarrow \operatorname{Aut}_{\mathbb{R}}(\mathbb{V})$, and $q>d / 2$ and $p$ obey $d / 2 \leq p \leq q$. Let $A_{1}$ be a $C^{\infty}$ connection on $P$, and $A_{0}$ be a Sobolev connection on $P$, and $a_{0}:=A_{0}-A_{1}$.
(1) There exists $C=C(g, p) \in[1, \infty)$ such that, if $a_{0} \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then

$$
\|\xi\|_{L^{r}(X)} \leq C\|\xi\|_{W_{A_{0}}^{1, p}(X)}, \quad \text { for } \begin{cases}1 \leq r \leq d p /(d-p) & \text { if } p<d \\ 1 \leq r<\infty & \text { if } p=d \\ r=\infty & \text { if } p>d\end{cases}
$$

for all $\xi \in C^{\infty}(X ; V)$; moreover, there exists $C=C\left(A_{1}, g, p\right) \in[1, \infty)$ such that
$\|\xi\|_{L^{r}(X)} \leq C\|\xi\|_{W_{A_{1}}^{2, p}(X)}, \quad$ for $\begin{cases}1 \leq r \leq d p /(d-2 p) & \text { if } p<d / 2, \\ 1 \leq r<\infty & \text { if } p=d / 2, \\ r=\infty & \text { if } p>d / 2,\end{cases}$
for all $\xi \in C^{\infty}(X ; V)$.
(2) If $a_{0} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then

$$
\|\xi\|_{C(X)} \leq C\|\xi\|_{W_{A_{0}}^{2, q}(X)}, \quad \forall \xi \in C^{\infty}(X ; V)
$$

for some $C=C(g, q) \in[1, \infty)$.
(3) If $a_{0} \in W_{A_{1}}^{1, p}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then
$C^{-1}\|\xi\|_{W_{A_{1}}^{1, p}(X)} \leq\|\xi\|_{W_{A_{0}}^{1, p}(X)} \leq C\|\xi\|_{W_{A_{1}}^{1, p}(X)}, \quad \forall \xi \in C^{\infty}(X ; V)$,
for some $C=C\left(g, G, p,\left\|a_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\right) \in[1, \infty)$.
(4) If $a_{0} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then
$C^{-1}\|\xi\|_{W_{A_{1}}^{2, q}(X)} \leq\|\xi\|_{W_{A_{0}}^{2, q}(X)} \leq C\|\xi\|_{W_{A_{1}}^{2, q}(X)}, \quad \forall \xi \in C^{\infty}(X ; V)$,
for some $C=C\left(g, G, p, q,\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}\right) \in[1, \infty)$.
(5) If $a_{0} \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then

$$
\|\xi\|_{W_{A_{0}}^{2, p}(X)} \leq C\|\xi\|_{W_{A_{1}}^{2, p}(X)}, \quad \forall \xi \in C^{\infty}(X ; V)
$$

for some $C=C\left(g, G, p, q,\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}\right) \in[1, \infty)$.
(6) If $a_{0} \in W_{A_{1}}^{2, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then
$C^{-1}\|\xi\|_{W_{A_{1}}^{2, p}(X)} \leq\|\xi\|_{W_{A_{0}}^{2, p}(X)} \leq C\|\xi\|_{W_{A_{1}}^{2, p}(X)}, \quad \forall \xi \in C^{\infty}(X ; V)$,
for some $C=C\left(g, G, p, q,\left\|a_{0}\right\|_{W_{A_{1}}^{2, q}(X)}\right) \in[1, \infty)$.
Proof. Item (11) is a well-known consequence of the Sobolev Embedding 4, Theorem 4.12] for scalar functions and the Kato Inequality [43, Equation (6.20)] in the case of the embedding $W^{1, p}(X) \subset L^{r}(X)$. Item (2) restates the conclusion of Lemma B. 1 .

For Item (3), we use $\nabla_{A_{0}} \xi=\nabla_{A_{1}} \xi+\left[a_{0}, \xi\right]$ and estimate

$$
\begin{aligned}
\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)} & \leq\left\|\nabla_{A_{0}} \xi\right\|_{L^{p}(X)}+\left\|\left[a_{0}, \xi\right]\right\|_{L^{p}(X)} \\
& \leq\|\xi\|_{W_{A_{0}}^{1, p}(X)}+z\left\|a_{0}\right\|_{L^{2 p}(X)}\|\xi\|_{L^{2 p}(X)} \\
& \leq C\left(1+\left\|a_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\right)\|\xi\|_{W_{A_{0}}^{1, p}(X)}
\end{aligned}
$$

where we used the continuous Sobolev embedding $W^{1, p}(X) \subset L^{2 p}(X)$ for $p \geq d / 2$ and Item (11) to obtain the last inequality. Here, $z=z(g, G) \in[1, \infty)]$ and $C \in$ $[1, \infty)$ has the stated dependencies. The analogous estimate with the roles of $A_{0}$ and $A_{1}$ reversed follows by a symmetric argument.

For Item (4), we first write

$$
\begin{equation*}
\nabla_{A_{1}}^{2} \xi=\nabla_{A_{0}}^{2} \xi+\nabla_{A_{0}} a_{0} \times \xi+a_{0} \times \nabla_{A_{0}} \xi+a_{0} \times a_{0} \times \xi \tag{B.2}
\end{equation*}
$$

Taking $L^{q}$ norms of both sides of (B.2), we see that

$$
\begin{aligned}
\left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{q}(X)} \leq & \left\|\nabla_{A_{0}}^{2} \xi\right\|_{L^{q}(X)}+\left\|\nabla_{A_{0}} a_{0} \times \xi\right\|_{L^{q}(X)}+\left\|a_{0} \times \nabla_{A_{0}} \xi\right\|_{L^{q}(X)} \\
& +\left\|a_{0} \times a_{0} \times \xi\right\|_{L^{q}(X)}
\end{aligned}
$$

and thus, for $z=z(g, G) \in[1, \infty)$,

$$
\begin{align*}
&\left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{q}(X)} \leq\left\|\nabla_{A_{0}}^{2} \xi\right\|_{L^{q}(X)}+z\left\|\nabla_{A_{0}} a_{0}\right\|_{L^{q}(X)}\|\xi\|_{C(X)}  \tag{B.3}\\
&+z\left\|a_{0}\right\|_{L^{2 q}(X)}\left\|\nabla_{A_{0}} \xi\right\|_{L^{2 q}(X)}+z\left\|a_{0}\right\|_{L^{2 q}}^{2}\|\xi\|_{C(X)}
\end{align*}
$$

By Item (3), we have

$$
\left\|\nabla_{A_{0}} a_{0}\right\|_{L^{q}(X)} \leq\left\|a_{0}\right\|_{W_{A_{0}}^{1, q}(X)} \leq C\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}
$$

and by Item (11) and the fact that $W^{1, q}(X) \subset L^{2 q}(X)$ for $q>d / 2$, we obtain

$$
\left\|\nabla_{A_{0}} \xi\right\|_{L^{2 q}(X)} \leq C\|\xi\|_{W_{A_{0}}^{2, q}(X)}
$$

Similarly, Item (2) gives

$$
\|\xi\|_{C(X)} \leq C\|\xi\|_{W_{A_{0}}^{2, q}(X)}
$$

By substituting the preceding inequalities into ( $\overline{\mathrm{B} .3}$ ), we find that

$$
\|\xi\|_{W_{A_{1}}^{2, q}(X)} \leq C\|\xi\|_{W_{A_{0}}^{2, q}(X)}
$$

where $C \in[1, \infty)$ has the stated dependencies. The analogous inequality with the roles of $A_{0}$ and $A_{1}$ reversed follows by a symmetric argument.

For Item (5), define $r \in[p, \infty]$ by $1 / p=1 / q+1 / r$, recall that $p=d / 2<q$ or $d / 2<p \leq q$, interchange the roles of $A_{0}$ and $A_{1}$ in (B.2), and take $L^{p}$ norms to give

$$
\begin{aligned}
\left\|\nabla_{A_{0}}^{2} \xi\right\|_{L^{p}(X)} \leq & \left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)}+\left\|\nabla_{A_{1}} a_{0} \times \xi\right\|_{L^{p}(X)}+\left\|a_{0} \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \\
& +\left\|a_{0} \times a_{0} \times \xi\right\|_{L^{p}(X)} \\
\leq & \left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)}+z\left\|\nabla_{A_{1}} a_{0}\right\|_{L^{q}(X)}\|\xi\|_{L^{r}(X)}+z\left\|a_{0}\right\|_{L^{2 p}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{2 p}(X)} \\
& +z\left\|a_{0}\right\|_{L^{2 q}}^{2}\|\xi\|_{L^{r}(X)} \\
\leq & \left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)}+C\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}\|\xi\|_{W_{A_{1}}^{2, p}(X)}+C\left\|a_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{W_{A_{1}}^{1, p}(X)} \\
& +C\left\|a_{0}\right\|_{W_{A_{1}}^{1, q}(X)}^{2}\|\xi\|_{W_{A_{1}}^{2, p}(X)},
\end{aligned}
$$

where $z=z(g, G) \in[1, \infty)$ and, to obtain the last inequality, we use the continuous Sobolev embeddings $W^{1, p}(X) \subset L^{2 p}(X)$ and $W^{1, q}(X) \subset L^{2 q}(X)$ for $d / 2 \leq p \leq q$ and Item (11) together with the continuous Sobolev embedding $W^{2, p}(X) \subset L^{r}(X)$, for $r \in[1, \infty)$ if $p=d / 2$ and $r=\infty$ if $p>d / 2$. Therefore, we obtain

$$
\|\xi\|_{W_{A_{0}}^{2, p}(X)} \leq C\|\xi\|_{W_{A_{1}}^{2, p}(X)}
$$

where $C \in[1, \infty)$ has the stated dependencies.
For Item (6), we take $L^{p}$ norms of ( (B.2) and use $\nabla_{A_{0}} a_{0}=\nabla_{A_{1}} a_{0}+\left[a_{0}, a_{0}\right]$ to give

$$
\begin{aligned}
\left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)} \leq & \left\|\nabla_{A_{0}}^{2} \xi\right\|_{L^{p}(X)}+\left\|\nabla_{A_{0}} a_{0} \times \xi\right\|_{L^{p}(X)}+\left\|a_{0} \times \nabla_{A_{0}} \xi\right\|_{L^{p}(X)} \\
& +\left\|a_{0} \times a_{0} \times \xi\right\|_{L^{p}(X)} \\
\leq & \left\|\nabla_{A_{0}}^{2} \xi\right\|_{L^{p}(X)}+z\left\|\nabla_{A_{1}} a_{0}\right\|_{L^{2 p}(X)}\|\xi\|_{L^{2 p}(X)}+\left\|a_{0}\right\|_{C(X)}\left\|\nabla_{A_{0}} \xi\right\|_{L^{p}(X)} \\
& +2 z\left\|a_{0}\right\|_{C(X)}^{2}\|\xi\|_{L^{p}(X)} .
\end{aligned}
$$

Applying the continuous Sobolev embeddings $W^{1, p}(X) \subset L^{2 p}(X)$ and $W^{2, q}(X) \subset$ $C(X)$ and Items (11) and (3), we discover that

$$
\begin{aligned}
\left\|\nabla_{A_{1}}^{2} \xi\right\|_{L^{p}(X)} \leq & \left\|\nabla_{A_{0}}^{2} \xi\right\|_{L^{p}(X)}+C\left\|\nabla_{A_{1}} a_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\|\xi\|_{W_{A_{0}}^{1, p}(X)} \\
& +\left\|a_{0}\right\|_{W_{A_{1}}^{2, q}(X)}\left\|\nabla_{A_{0}} \xi\right\|_{L^{p}(X)}+2 z\left\|a_{0}\right\|_{W_{A_{1}}^{2, q}(X)}^{2}\|\xi\|_{L^{p}(X)} \\
\leq & C\|\xi\|_{W_{A_{0}}^{2, p}(X)}
\end{aligned}
$$

where $C \in[1, \infty)$ has the stated dependencies. The analogous inequality with the roles of $A_{0}$ and $A_{1}$ reversed follows by a symmetric argument.

## APPENDIX C

## Fredholm and index properties of a Hodge Laplacian with Sobolev coefficients

In this chapter we include proofs of results regarding the Fredholm properties of the Hodge Laplace operators encountered in Chapters 2 and 3 that would be standard if the operator had smooth coefficients and acted on $L^{2}$ rather than $L^{p}$ Sobolev spaces as we allow here.

When $A$ is a smooth connection, the Hodge Laplace operator, $\Delta_{A}$, in (2.2.1) is an elliptic, second-order partial differential operator with smooth coefficients that is $L^{2}$-self-adjoint and so Theorem A. 1 immediately provides the

Proposition C. 1 (Fredholm and index zero properties of a Laplace operator with smooth coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group, and $P$ be a smooth principal $G$-bundle over $X$, and $l \geq 0$ an integer. If $A$ and $A_{1}$ are $C^{\infty}$ connections on $P$ and $1<p<\infty$, then the operator,

$$
\begin{equation*}
\Delta_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{C.1}
\end{equation*}
$$

is Fredholm with index zero and closed range $K^{\perp} \cap L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, where $\perp$ denotes $L^{2}$-orthogonal complement and $K \subset W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is the kernel of $\Delta_{A}$ in (C.1).

Proposition C. 1 and a compact operator perturbation argument provides following useful generalization from the case of a $C^{\infty}$ to a $W^{1, q}$ connection $A$.

Corollary C. 2 (Fredholm and index zero properties of a Laplace operator with Sobolev coefficients). Assume the hypotheses of Proposition C.1, but allow A to be a $W^{1, q}$ connection with $d / 2<q<\infty$ and restrict $p \in(1, \infty)$ so that $d / 2 \leq p \leq q$. Then the operator,

$$
\begin{equation*}
\Delta_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{C.2}
\end{equation*}
$$

is Fredholm with index zero and closed range $K^{\perp} \cap L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, where $\perp$ denotes $L^{2}$-orthogonal complement and $K \subset W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ is the kernel of $\Delta_{A}$ in (C.2).

Proof. By hypothesis, $A_{1}$ is a $C^{\infty}$ connection on $P$. We write $A=A_{1}+a$, for $a \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ and proceed by modifying the proof of Proposition 2.2.1 to show that, for suitable $r \in(p, \infty]$ and $t \in[p, \infty)$, the operator

$$
\begin{equation*}
\Delta_{A}-\Delta_{A_{1}}: W_{A_{1}}^{1, t}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \cap L^{r}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{C.3}
\end{equation*}
$$

is bounded and, because the Sobolev embedding $W^{2, p}(X) \Subset W^{1, t}(X) \cap L^{r}(X)$ will be compact by the Rellich-Kondrachov Theorem 4, Theorem 6.3], then the following composition of that compact embedding and the preceding bounded operator,

$$
\begin{equation*}
\Delta_{A}-\Delta_{A_{1}}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \tag{C.4}
\end{equation*}
$$

is compact by [14, Proposition 6.3].
By modifying the derivation of the estimate (2.2.11) in the proof of Proposition 2.2.1. we claim that, for $r \in(p, \infty]$ defined by $1 / p=1 / q+1 / r$,

$$
\begin{align*}
\left\|\left(\Delta_{A}-\Delta_{A_{1}}\right) \xi\right\|_{L^{p}(X)} \leq z\left(\|a\|_{W_{A_{1}}^{1, q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}^{2}\right) & \|\xi\|_{L^{r}(X)}  \tag{C.5}\\
& +z\|a\|_{W_{A_{1}}^{1, q}(X)}\|\xi\|_{W_{A_{1}}^{1, t}(X)}
\end{align*}
$$

where $z=z(g, G, l) \in[1, \infty)$ and 1) $q<d$ and $t \in(p, \infty)$ defined by $1 / p=1 / q^{*}+1 / t$ with $q^{*}=d q /(d-q)$, or 2) $q=d$ and $t=p+\delta$ for a small enough $\delta \in(0,1]$, or 3) $q>d$ and $t=p$.

To see that (C.5) holds, observe that (2.2.7) (with $A_{s}$ replaced by $A_{1}$ ) gives

$$
\Delta_{A} \xi=\Delta_{A_{1}} \xi+\nabla_{A_{1}} a \times \xi+a \times \nabla_{A_{1}} \xi+a \times a \times \xi
$$

Proceeding as in the derivation of (2.2.9), we see that the preceding identity yields

$$
\begin{aligned}
&\left\|\left(\Delta_{A}-\Delta_{A_{1}}\right) \xi\right\|_{L^{p}(X)} \leq z\left\|\nabla_{A_{1}} a\right\|_{L^{q}(X)}\|\xi\|_{L^{r}(X)}+\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \\
&+z\left\|\left.a\right|^{2}\right\|_{L^{q}(X)}\|\xi\|_{L^{r}(X)}
\end{aligned}
$$

Hence, using continuity of the Sobolev multiplication map, $L^{2 q}(X) \times L^{2 q}(X) \rightarrow$ $L^{q}(X)$, and Sobolev embedding, $W^{1, q}(X) \subset L^{2 q}(X)$, by 4, Theorem 4.12, Part I, Cases A, B, and C] (using $q \geq d / 2$ ), and Kato Inequality [43, Equation (6.20)], we find that

$$
\begin{aligned}
\left\|\left(\Delta_{A}-\Delta_{A_{1}}\right) \xi\right\|_{L^{p}(X)} \leq z\left(\|a\|_{W_{A_{1}}^{1, q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}^{2}\right)\|\xi\|_{L^{r}(X)} & \\
& +\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)}
\end{aligned}
$$

Now, defining $v \in(p, \infty)$ by a choice of small $\delta \in(0,1]$ and $1 / p=1 / v+1 /(p+\delta)$ and $w \in(p, \infty)$ by $1 / p=1 / q^{*}+1 / w$ and $q^{*}=d q /(d-q)$ (when $q<d$ ), we obtain that

$$
\left\|a \times \nabla_{A_{1}} \xi\right\|_{L^{p}(X)}
$$

is bounded by

$$
z\|a\|_{L^{\infty}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \leq z\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{p}(X)} \quad \text { if } q>d
$$

or

$$
z\|a\|_{L^{v}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{p+\delta}(X)} \leq z\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{p+\delta}(X)} \quad \text { if } q=d,
$$

or

$$
z\|a\|_{L^{q^{*}}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{w}(X)} \leq z\|a\|_{W_{A_{1}}^{1, q}(X)}\left\|\nabla_{A_{1}} \xi\right\|_{L^{w}(X)} \quad \text { if } q<d
$$

The final three inequalities follow from continuity of the Sobolev embeddings, $W^{1, q}(X) \subset C(X)$ when $q>d$, and $W^{1, q}(X) \subset L^{v}(X)$ when $q=d$, and $W^{1, q}(X) \subset$ $L^{q^{*}}(X)$ when $q<d$, by [4, Theorem 4.12, Part I, Cases A, B, and C], respectively. This proves Claim (C.5).

In particular, the operator (C.3) is bounded for each of the three above cases for $(q, t)$.

When $p<q$ and thus $r<\infty$, the embedding $W^{2, p}(X) \subset L^{r}(X)$ is compact since $p \geq d / 2$ by hypothesis and so the embedding $W^{2, p}(X) \subset W^{2, d / 2}(X)$ is continuous while the embedding $W^{2, p}(X) \subset L^{r}(X)$ is compact by 4, Theorem 6.3, Part I]. When $p=q$ and thus $r=\infty$, the embedding $W^{2, p}(X) \subset C(X)$ is compact since $q>d / 2$ by hypothesis and thus $p>d / 2$ so the embedding $W^{2, p}(X) \subset C(X)$ is compact by [4, Theorem 6.3, Part II].

To see compactness of the embedding $W^{2, p}(X) \subset W^{1, t}(X)$ and hence of the embedding $W^{2, p}(X) \subset W^{1, t}(X) \cap L^{r}(X)$ and consequently the operator (C.4), we consider the three cases for $(q, t)$ separately.

Case $1\left(q<d\right.$ and $t \in(p, \infty)$ defined by $1 / p=1 / q^{*}+1 / t$ with $\left.q^{*}=d q /(d-q)\right)$. We have a compact Sobolev embedding, $W^{2, p}(X) \Subset W^{1, t}(X)$, by 4, Theorem 6.3, Part I] provided $t<p^{*}=d p /(d-p)$. To check that the strict inequality $t<p^{*}$ holds, observe that $p \leq q$ and so $1 / p \geq 1 / q$, which gives

$$
1 / t=1 / p-1 / q^{*}=1 / p-1 / q+1 / d \geq 1 / d
$$

If $p>d / 2$, then $p^{*}=d p /(d-p)>d$ and $1 / p^{*}<1 / d$ and thus $1 / t \geq 1 / d>1 / p^{*}$ and $t<p^{*}$, as desired. If $p=d / 2$, then $p<q$ by hypothesis and so $1 / p<1 / q$, which now gives

$$
1 / t>1 / d
$$

But $p=d / 2 \Longrightarrow p^{*}=d$ and thus $1 / t>1 / d=1 / p^{*}$ and $t<p^{*}$, again as desired.
Hence, the embedding $W^{2, p}(X) \subset W^{1, t}(X)$ is compact and so the operator (C.4) is compact for this case.

CASE $2(q=d$ and $t=p+\delta$ for small enough $\delta>0)$. We have a compact Sobolev embedding, $W^{2, p}(X) \Subset W^{1, p+\delta}(X)$, by [4, Theorem 6.3, Part I], where it suffices to choose $\delta \in(0,1]$ small enough that $p+\delta<p^{*}=d p /(d-p) \in[d, \infty)$ when $d / 2 \leq p<d$ and any $\delta \in(0,1]$ will do when $p=d$. Hence, the embedding $W^{2, p}(X) \subset W^{1, t}(X)$ is compact with $t=p+\delta$ and so the operator (C.4) is compact for this case.

Case $3(q>d$ and $t=p)$. We have a compact Sobolev embedding, $W^{2, p}(X) \Subset$ $W^{1, p}(X)$, by [4, Theorem 6.3, Parts I and II]. Hence, the embedding $W^{2, p}(X) \subset$ $W^{1, t}(X)$ is compact with $t=p$ and so the operator (C.4) is compact for this case.

Proposition C. 1 implies that the operator

$$
\Delta_{A_{1}}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)
$$

is Fredholm with index zero while the operator $\Delta_{A}-\Delta_{A_{1}}$ in (C.4) is compact from each of the preceding cases, so the operator (C.2) is Fredholm with index zero by 67. Corollary 19.1.8].

The identification of the range of the operator (C.2) follows mutatis mutandis the proof of the corresponding fact in the statement of Theorem A. 1 The only difference, after noting that $\Delta_{A}^{*}=\Delta_{A}$ and $p \leq q \Longrightarrow q^{\prime} \leq p^{\prime}$ and thus $L^{p^{\prime}}(X) \subset L^{q^{\prime}}(X)$, is that we appeal to the following regularity result for distributional solutions, $b \in L^{p^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$, to an elliptic linear partial differential equation, $\Delta_{A} b=0$, with Sobolev rather than $C^{\infty}$ coefficients,

$$
\begin{aligned}
& \operatorname{Ker}\left(\Delta_{A}: L^{p^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow W_{A_{1}}^{-2, p^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right) \\
&=\operatorname{Ker}\left(\Delta_{A}: W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right) \rightarrow L^{p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)\right)
\end{aligned}
$$

Indeed, because $b \in L^{q^{\prime}}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ by the preceding remarks, Lemma 2.3.1 implies that $b \in W_{A_{1}}^{2, q}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ and, in particular, $b \in W_{A_{1}}^{2, p}\left(X ; \Lambda^{l} \otimes \operatorname{ad} P\right)$ since $p \leq q$. This completes the proof of Corollary C.2,

We shall also need to consider Fredholm properties of the perturbed Laplace operator,

$$
d_{A}^{*} d_{A+a}: C^{\infty}(X ; \operatorname{ad} P) \rightarrow C^{\infty}(X ; \operatorname{ad} P),
$$

when a $C^{\infty}$ connection, $A$, one-form, $a \in \Omega^{l}(X ; \operatorname{ad} P)$, and Fréchet space, $\Omega^{l}(X ; \operatorname{ad} P)$, are replaced by suitable Sobolev counterparts. As usual, we begin with the simpler case of smooth coefficients.

Lemma C. 3 (Fredholm and index zero properties of a perturbed Laplace operator with smooth coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group and $P$ be a smooth principal $G$-bundle over $X$. If $A$ and $A_{1}$ are $C^{\infty}$ connections on $P$ and $a \in C^{\infty}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$ and $1<q<\infty$, then the operator,

$$
\begin{equation*}
d_{A}^{*} d_{A+a}: W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow L^{q}(X ; \operatorname{ad} P), \tag{C.6}
\end{equation*}
$$

is Fredholm with index zero.
Proof. We observe that

$$
d_{A}^{*} d_{A+a}: C^{\infty}(X ; \operatorname{ad} P) \rightarrow C^{\infty}(X ; \operatorname{ad} P)
$$

is an elliptic, linear, second-order partial differential operator such that

$$
d_{A}^{*} d_{A+a}-d_{A+a}^{*} d_{A}=d_{A}^{*}[a, \cdot]-[a, \cdot]^{*} d_{A}: C^{\infty}(X ; \operatorname{ad} P) \rightarrow C^{\infty}(X ; \operatorname{ad} P)
$$

is a first-order differential operator. The conclusions now follow from Theorem A. 1

Lemma C. 4 (Fredholm and index zero properties of a perturbed Laplace operator with Sobolev coefficients). Let $(X, g)$ be a closed, smooth Riemannian manifold of dimension $d \geq 2$, and $G$ be a compact Lie group and $P$ be a smooth principal $G$-bundle over $X$. If $A_{1}$ is a $C^{\infty}$ connection on $P$ and $A$ is a $W^{1, q}$ connection on $P$ with $d / 2<q<\infty$ and $a \in W_{A_{1}}^{1, q}\left(X ; \Lambda^{1} \otimes \operatorname{ad} P\right)$, then the operator,

$$
\begin{equation*}
d_{A}^{*} d_{A+a}: W_{A_{1}}^{2, q}(X ; \operatorname{ad} P) \rightarrow L^{q}(X ; \operatorname{ad} P) \tag{C.7}
\end{equation*}
$$

is Fredholm with index zero.
Proof. The argument is almost identical to the proof of Corollary C.2. Write $A=A_{1}+a_{1}$ and observe that

$$
d_{A}^{*} d_{A+a}=d_{A_{1}}^{*} d_{A_{1}}+d_{A_{1}}^{*}\left[a_{1}+a, \cdot\right]+\left[a_{1}, \cdot\right]^{*} d_{A_{1}}+\left[a_{1}, \cdot\right]^{*}\left[a_{1}+a, \cdot\right] .
$$

By again retracing the steps in the proof of Proposition 2.2.1 we find that, for $r \in(p, \infty]$ defined by $1 / p=1 / q+1 / r$ and $\xi \in W_{A_{1}}^{2, q}(X ; \operatorname{ad} P)$,

$$
\begin{align*}
& \left\|\left(d_{A}^{*} d_{A+a}-d_{A_{1}}^{*} d_{A_{1}}\right) \xi\right\|_{L^{p}(X)} \\
& \leq  \tag{C.8}\\
& \quad z\left(\|a\|_{W_{A_{1}}^{1, q}(X)}+\left\|a_{1}\right\|_{W_{A_{1}}^{1, q}(X)}+\left\|a_{1}\right\|_{W_{A_{1}}^{1, q}(X)}^{2}\right. \\
& \left.\quad+\left\|a_{1}\right\|_{W_{A_{1}}^{1, q}(X)}\|a\|_{W_{A_{1}}^{1, q}(X)}\right)\|\xi\|_{L^{r}(X)} \\
& \quad+z\left(\left\|a_{1}\right\|_{W_{A_{1}}^{1, q}(X)}+\|a\|_{W_{A_{1}}^{1, q}(X)}\right)\|\xi\|_{W_{A_{1}}^{1, t}(X)},
\end{align*}
$$

with the values of $t$ specified in the proof of Corollary C.2. The remainder of the proof of CorollaryC.2 now applies to show that the operator (C.7) is Fredholm with index zero.

## APPENDIX D

## Convergence of gradient flows under the validity of the Łojasiewicz-Simon gradient inequality

As we noted in Section 1.2, the gradient inequality for an energy function, $\mathscr{E}: \mathscr{X} \supset \mathscr{U} \rightarrow \mathbb{R}$, with gradient map, $\mathscr{M}: \mathscr{X} \supset \mathscr{U} \rightarrow \tilde{X}$, is most useful when it has the strong form (1.2.9) implied by Theorem 3, namely,

$$
\|\mathscr{M}(x)\|_{\mathscr{H}} \geq Z\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}, \quad \forall x \in \mathscr{U} \text { with }\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma,
$$

where $\mathscr{H}$ is a Hilbert space, the Banach space $\tilde{\mathscr{X}}$ is continuously embedded in $\mathscr{H}$, and the Banach space $\mathscr{X}$ is a dense subspace of $\mathscr{H}$ with continuous embedding $\mathscr{X} \subset \mathscr{H}$ and thus $\mathscr{H}^{*} \subset \mathscr{X}^{*}$ is also a continuous embedding.

In this Appendix, we briefly explain why Theorem 3is so useful in applications to questions of global existence and convergence of a strong solution, that is, $u \in$ $C([0, T) ; \mathscr{X})$ with time derivative $\dot{u} \in C((0, T) ; \mathscr{H})$ (for $T \in(0, \infty])$, to the Cauchy problem for the gradient system

$$
\begin{equation*}
\dot{u}(t)=-\mathscr{M}(u(t)) \quad \text { in } \mathscr{H}, \quad t \in(0, T), \quad u(0)=u_{0} . \tag{D.1}
\end{equation*}
$$

The importance of a geometric version of Theorem 3 to a more specific setting in geometric analysis was famously pioneered by Simon in [105], generalizing a result of Łojasiewicz for gradient flows in Euclidean spaces [88.

A weak solution to the gradient system for $\mathscr{E}$ has the form $u \in C([0, T) ; \mathscr{X})$ with time derivative $\dot{u} \in C\left((0, T) ; \mathscr{X}^{*}\right)$, obeying

$$
\begin{equation*}
\dot{u}(t)=-\mathscr{E}^{\prime}(u(t)) \quad \text { in } \mathscr{X}^{*}, \quad t \in(0, T), \quad u(0)=u_{0} . \tag{D.2}
\end{equation*}
$$

To illustrate the application of Theorem 3, we include from 32 a proof of a simplified version of our [32, Proposition 24.12] that yields convergence, $u(t) \rightarrow u_{\infty}$ in $\mathscr{H}$ as $t \rightarrow \infty$ for a global strong solution, $u \in C([0, \infty) ; \mathscr{X}) \cap C^{1}([0, \infty) ; \mathscr{H})$ to (D.1), when $\mathscr{M}$ and $\mathscr{E}$ obey the version of Łojasiewicz-Simon gradient inequality in (1.2.9).

The statement and proof of the forthcoming Proposition D. 1 are closely modeled on Huang's [68, Proposition 3.3.2], but for the gradient system (D.1) in a Hilbert space. By contrast, Huang's version allows apparently more general weak gradient-like differential inequalities in Banach spaces, namely 68, Equation (3.10a) or Equation $\left(3.10^{\prime}\right)$ ], with auxiliary conditions such as those in his 68, Equation (3.10b)] or [68, Equation (3.10')]. However, examples satisfying Huang's gradient-like differential inequalities and auxiliary conditions appear to us to be difficult to find except when they reduce to a pure gradient system (D.1) in a Hilbert space or Simon's gradient-like system [105, Equation (3.1)] in a Hilbert space,

$$
\dot{u}(t)=-\mathscr{M}(u(t))+\mathscr{R}(t) \quad \text { in } \mathscr{H}, \quad t \in(0, T), \quad u(0)=u_{0}
$$

where $\mathscr{R} \in C((0, \infty) ; \mathscr{H})$ obeys a decay condition (as $t \rightarrow \infty$ ) implying Huang's 68, Equation (3.10b)] or Simon's hypothesis in [105, Equation (3.1)],

$$
\|\mathscr{R}(t)\|_{\mathscr{H}} \leq \alpha\|\dot{u}(t)\|_{\mathscr{H}},
$$

where $\alpha \in(0,1)$ is a constant.
Proposition D. 1 (Convergence of gradient flow under the validity of the Łojasiewicz-Simon gradient inequality). Let $\mathscr{U}$ be an open subset of a Banach space, $\mathscr{X}$, that is continuously embedded and dense in a Hilbert space, $\mathscr{H}$. Let $\mathscr{E}: \mathscr{U} \subset \mathscr{X} \rightarrow \mathbb{R}$ be a $C^{1}$ function on an open subset, $\mathscr{U} \subset \mathscr{X}$, with gradient map $\mathscr{M}: \mathscr{U} \subset \mathscr{X} \rightarrow \mathscr{H}$, and $x_{\infty} \in \mathscr{U}$ be a critical point of $\mathscr{E}$, that is, $\mathscr{E}^{\prime}\left(x_{\infty}\right)=0$. Let $u \in C([0, \infty) ; \mathscr{U}) \cap C^{1}([0, \infty) ; \mathscr{H})$ be a strong solution to (D.1) such that

$$
\begin{equation*}
\inf \{|\mathscr{E}(u(t))|: t \geq 0\}>-\infty \tag{D.3}
\end{equation*}
$$

If $\mathscr{E}$ and $\mathscr{M}$ satisfy a Eojasiewicz-Simon gradient inequality (1.2.9) in the orbit $O(u)=\{u(t): t \geq 0\}$, that is,

$$
\begin{equation*}
\|\mathscr{M}(u(t))\|_{\mathscr{H}} \geq Z\left|\mathscr{E}(u(t))-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}, \quad \forall t \geq 0 \tag{D.4}
\end{equation*}
$$

for constants $Z \in(0, \infty)$ and $\theta \in[1 / 2,1)$, then

$$
\begin{equation*}
\int_{0}^{\infty}\|\dot{u}(t)\|_{\mathscr{H}} d t \leq \int_{\mathscr{E}_{\infty}}^{\mathscr{E}(u(0))} \frac{1}{Z\left|s-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}} d s<\infty \tag{D.5}
\end{equation*}
$$

where $\mathscr{E}_{\infty}:=\lim _{t \rightarrow \infty} \mathscr{E}(u(t)) \in \mathbb{R}$, and thus

$$
u(t) \rightarrow u_{\infty} \quad \text { in } \mathscr{H}, \quad \text { as } t \rightarrow \infty
$$

for some $u_{\infty} \in \mathscr{H}$.
Proof. The function $[0, \infty) \ni t \mapsto \mathscr{E}(u(t)) \in \mathbb{R}$ is $C^{1}$ by direct calculation and obeys (see [68, Proposition 3.1.2])

$$
\begin{aligned}
-\frac{d}{d t} \mathscr{E}(u(t))= & -\mathscr{E}^{\prime}(u(t))(\dot{u}(t)) \quad \text { (Chain Rule) } \\
= & -\langle\dot{u}(t), \mathscr{M}(u(t))\rangle_{\mathscr{X} \times \mathscr{X}^{*} \quad \quad(\text { by Definition 1.2.1) })}^{=} \\
= & -(\dot{u}(t), \mathscr{M}(u(t)))_{\mathscr{H}} \quad\left(\text { by } \mathscr{H} \subset \mathscr{X}^{*} \text { via } \mathscr{H} \ni h \mapsto(\cdot, h)_{\mathscr{H}} \in \mathscr{X}^{*}\right) \\
= & \|\mathscr{M}(u(t))\|_{\mathscr{H}}^{2}=\|\mathscr{M}(u(t))\|_{\mathscr{H}}\|\dot{u}(t)\|_{\mathscr{H}} \geq 0, \\
& \forall t \in(0, \infty) \quad(\text { by }(\text { D.1) }) .
\end{aligned}
$$

Hence, $\mathscr{E}(u(t))$ is a nonincreasing and uniformly bounded function of $t \in[0, \infty)$ by (D.3), so $\mathscr{E}_{\infty}=\lim _{t \rightarrow \infty} \mathscr{E}(u(t))$ exists, as asserted by the proposition. Set $H(t):=$ $\mathscr{E}(u(t))$, for all $t \in[0, \infty)$, and observe that $H(t)$ is monotone and absolutely continuous on $[0, \infty)$ and obeys, by the preceding equality,

$$
\begin{equation*}
-\frac{d}{d t} H(t)=\|\mathscr{M}(u(t))\|_{\mathscr{H}}\|\dot{u}(t)\|_{\mathscr{H}}, \quad \forall t \in[0, \infty) \tag{D.6}
\end{equation*}
$$

Let $\phi: \mathbb{R} \rightarrow \mathbb{R}$ be the function defined by $\phi(s)=Z\left|s-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}$, for all $s \in \mathbb{R}$, and let $\Phi: \mathbb{R} \rightarrow \mathbb{R}$ be the absolutely continuous function given by

$$
\Phi(x):=\int_{\mathscr{E}_{\infty}}^{x} \frac{1}{\phi(s)} d s=\int_{\mathscr{E}_{\infty}}^{x} \frac{1}{Z\left|s-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}} d s, \quad \forall x \in \mathbb{R}
$$

where $\lim _{t \rightarrow \infty} H(t)=\mathscr{E}_{\infty}$. The function $\Phi$ is differentiable a.e. on $\mathbb{R}$ with $\Phi^{\prime}(x)=$ $1 / \phi(x)$ for a.e. $x \in \mathbb{R}$. According to [68, Lemma 3.2.1], the composition $\Phi \circ H$ is absolutely continuous on $[0, \infty)$ and there holds

$$
\begin{equation*}
\frac{d}{d t} \Phi(H(t))=\frac{H^{\prime}(t)}{\phi(H(t))}, \quad \forall t \in \Lambda \tag{D.7}
\end{equation*}
$$

where $\Lambda \subset[0, \infty)$ is such that the complement, $[0, \infty) \backslash \Lambda$, has zero Lebesgue measure.

For any $t \in \Lambda$, we have two possibilities: either i) $\|\mathscr{M}(u(t))\|_{\mathscr{H}}=0$, or ii) $\|\mathscr{M}(u(t))\|_{\mathscr{H}}>0$. For Case (iii), we observe that the Łojasiewicz-Simon gradient inequality (D.4) takes the shape,

$$
\begin{equation*}
\phi(H(t))=Z\left|\mathscr{E}(u(t))-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta} \leq\|\mathscr{M}(u(t))\|_{\mathscr{H}}, \tag{D.8}
\end{equation*}
$$

and so

$$
\begin{aligned}
-\frac{d}{d t} \Phi(H(t)) & =-\frac{H^{\prime}(t)}{\phi(H(t))} \quad(\text { by }(\overline{\mathrm{D} .7)}) \\
& =\frac{\|\mathscr{M}(u(t))\|_{\mathscr{H}}\|\dot{u}(t)\|_{\mathscr{H}}}{\phi(H(t))} \quad(\text { by }(\overline{\mathrm{D} .6})) \\
& \geq \frac{\|\mathscr{M}(u(t))\|_{\mathscr{H}}\|\dot{u}(t)\|_{\mathscr{H}}}{\|\mathscr{M}(u(t))\|_{\mathscr{H}}} \quad(\text { by }(\overline{\mathrm{D} .8})) \\
& =\|\dot{u}(t)\|_{\mathscr{H}},
\end{aligned}
$$

that is,

$$
\begin{equation*}
-\frac{d}{d t} \Phi(H(t)) \geq\|\dot{u}(t)\|_{\mathscr{H}} . \tag{D.9}
\end{equation*}
$$

Therefore, by the non-negativity of the function $-d \Phi(H(t)) / d t$, combined with the fact that $[0, \infty) \backslash \Lambda$ has Lebesgue measure zero, we obtain the estimate,

$$
\begin{equation*}
-\frac{d}{d t} \Phi(H(t)) \geq\|\dot{u}(t)\|_{\mathscr{H}}, \quad \text { a.e. } t \in[0, \infty) \tag{D.10}
\end{equation*}
$$

for both Cases (ii) and(iii). Integration and the fact that $\lim _{t \rightarrow \infty} H(t)=\mathscr{E}_{\infty}$ yields

$$
\int_{0}^{\infty}\|\dot{u}(t)\|_{\mathscr{H}} d t \leq \Phi(H(0))-\lim _{t \rightarrow \infty} \Phi(H(t))=\Phi(H(0))-\Phi\left(\mathscr{E}_{\infty}\right)=\Phi(H(0))
$$

By the definitions of $\Phi(x)$ and $H(t)$, this is (D.5), since

$$
\Phi(H(0))=\int_{\mathscr{E}_{\infty}}^{H(0)} \frac{1}{\phi(s)} d s=\int_{\mathscr{E}_{\infty}}^{\mathscr{E}(u(0))} \frac{1}{Z\left|s-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}} d s
$$

The final convergence assertion follows from the fact that

$$
u\left(t_{n}\right)-u\left(t_{m}\right)=\int_{t_{m}}^{t_{n}} \dot{u}(t) d t, \quad \forall t_{m}, t_{n} \in[0, \infty)
$$

and thus, for any unbounded sequence of times, $\left\{t_{n}\right\}_{n=1}^{\infty} \subset[0, \infty)$, the sequence of points, $\left\{u\left(t_{n}\right)\right\}_{n=1}^{\infty} \subset \mathscr{H}$, is Cauchy in $\mathscr{H}$ and thus converges to a limit, $u_{\infty} \in$ $\mathscr{H}$, as $t \rightarrow \infty$, independent of the choice, $\left\{t_{n}\right\}_{n=1}^{\infty}$. This completes the proof of Proposition D. 1

The convergence of $u(t)$ in $\mathscr{H}$ provided by Proposition D.1 can be improved with the assumption of the following

Hypothesis D. 2 (Regularity and a priori interior estimate for a trajectory). (See [32, Hypothesis 24.10].) Let $C_{1}$ and $\rho$ be positive constants and let $T \in(0, \infty]$. Given $u \in C([0, \infty) ; \mathscr{X}) \cap C^{1}([0, \infty) ; \mathscr{H})$, we say that $\dot{u}:[0, T) \rightarrow \mathscr{H}$ obeys an $a$ priori interior estimate on $(0, T]$ if, for every $S \geq 0$ and $\delta>0$ obeying $S+\delta \leq T$, the map $\dot{u}:[S+\delta, T) \rightarrow \mathscr{X}$ is Bochner integrable and there holds

$$
\begin{equation*}
\int_{S+\delta}^{T}\|\dot{u}(t)\|_{\mathscr{X}} d t \leq C_{1}\left(1+\delta^{-\rho}\right) \int_{S}^{T}\|\dot{u}(t)\|_{\mathscr{H}} d t \tag{D.11}
\end{equation*}
$$

Given Hypothesis D.2 the bound (D.5) in Proposition D.1 improves to

$$
\begin{equation*}
\int_{\delta}^{\infty}\|\dot{u}\|_{\mathscr{X}} d t \leq C_{1}\left(1+\delta^{-\rho}\right) \int_{\mathscr{E}_{\infty}}^{\mathscr{E}(u(0))} \frac{1}{c\left|s-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}} d s<\infty \tag{D.12}
\end{equation*}
$$

for all $\delta>0$, and so the convergence in $\mathscr{H}$ improves to $u(t) \rightarrow u_{\infty}$ in $\mathscr{X}$ as $t \rightarrow \infty$.
The application of Proposition D. 1 and similar results to proofs of global existence, convergence, convergence rates, and stability of solutions to (D.1) are described at length in [32, Section 2.1].

## APPENDIX E

## Huang's Łojasiewcz-Simon gradient inequality for analytic functions on Banach spaces

For the convenience of the reader, we quote Huang's [68, Theorem 2.4.2 (i)] for the Łojasiewcz-Simon gradient inequality for analytic functions on Banach spaces. We first recall Huang's hypotheses for [68, Theorem 2.4.2 (i)].

Hypothesis E. 1 (Hypotheses for the abstract Łojasiewicz-Simon gradient inequality with Hilbert space gradient norm). (See [68, pp. 34-35].) Assume the following conditions.
(1) Let $\mathscr{H}$ be a Hilbert space, $\mathscr{A}: \mathscr{D}(\mathscr{A}) \subset \mathscr{H} \rightarrow \mathscr{H}$ a linear, positive definite, self-adjoint operator, and $\mathscr{H}_{\mathscr{A}}:=\left(\mathscr{D}(\mathscr{A}),(\cdot, \cdot)_{\mathscr{A}}\right)$ be the Hilbert space with inner product,

$$
(u, v)_{\mathscr{A}}:=(\mathscr{A} u, \mathscr{A} v)_{\mathscr{H}}, \quad \forall u, v \in \mathscr{D}(\mathscr{A})
$$

where $(\cdot, \cdot)_{\mathscr{H}}$ is the inner product on $\mathscr{H}$.
(2) Let $\mathscr{X} \subset \tilde{\mathscr{X}}$ be Banach spaces such the following embeddings are continuous,

$$
\mathscr{X} \subset \mathscr{H}_{\mathscr{A}}, \quad \tilde{X} \subset \mathscr{H} .
$$

(3) Let $\mathscr{E}: \mathscr{X} \rightarrow \mathbb{R}$ be a function with $C^{1}$ gradient map, $\mathscr{M}: \mathscr{U} \subset \mathscr{X} \rightarrow \tilde{\mathscr{X}}$, so that

$$
\mathscr{E}^{\prime}(x) v=(v, \mathscr{M}(x))_{\mathscr{H}}, \quad \forall x \in \mathscr{U} \text { and } v \in \mathscr{X}
$$

where $\mathscr{U} \subset \mathscr{X}$ is an open subset, and having the following properties:
(a) $\mathscr{M}$ is a Fredholm map of index zero, that is, for each $x \in \mathscr{U}$,

$$
\mathscr{M}^{\prime}(x): \mathscr{X} \rightarrow \tilde{\mathscr{X}}
$$

is a Fredholm operator of index zero.
(b) For each $x \in \mathscr{U}$, the bounded, linear operator,

$$
\mathscr{M}^{\prime}(x): \mathscr{X} \rightarrow \tilde{\mathscr{X}},
$$

has an extension

$$
\mathscr{M}_{1}(x): \mathscr{H}_{A} \rightarrow \mathscr{H}
$$

which is symmetrid and also a Fredholm operator of index zero and such that the map

$$
\mathscr{U} \ni x \mapsto \mathscr{M}_{1}(x) \in \mathscr{L}\left(\mathscr{H}_{\mathscr{A}}, \mathscr{H}\right) \quad \text { is continuous, }
$$

or, equivalently, the map $\mathscr{U} \ni x \mapsto \mathscr{M}_{1}(x) \mathscr{A}^{-1} \in \mathscr{L}(\mathscr{H})$ is continuous.

[^9]Theorem E. 2 (Łojasiewcz-Simon gradient inequality for analytic functions on Banach spaces and Hilbert space gradient). [68, Theorem 2.4.2 (i)] Assume Hypothesis E. 1 on $\mathscr{E}$, $\mathscr{M}, \mathscr{M}_{1}, \mathscr{H}, \mathscr{U}, \mathscr{X}$, and $\tilde{\mathscr{X}}$, and that $\mathscr{M}: \mathscr{U} \subset \mathscr{X} \rightarrow \tilde{\mathscr{X}}$ is real analytic. If $x_{\infty} \in \mathscr{U}$ is a critical point of $\mathscr{E}$, that is, $\mathscr{M}\left(x_{\infty}\right)=0$, then there are positive constants, $c, \sigma$, and $\theta \in[1 / 2,1)$ such that
(E.1) $\quad\|\mathscr{M}(x)\|_{\mathscr{H}} \geq c\left|\mathscr{E}(x)-\mathscr{E}\left(x_{\infty}\right)\right|^{\theta}, \quad \forall x \in \mathscr{U}$ such that $\left\|x-x_{\infty}\right\|_{\mathscr{X}}<\sigma$.

## APPENDIX F

## Quantitative implicit and inverse function theorems

Statements of the Implicit Function Theorem are often most useful in applications when equipped with explicit estimates for the radii of the balls containing the domain and range of the implied function that the statement produces, along with an explicit estimate for the Lipschitz constant of that function. We shall state and prove a version of such a result in more generality than we need in our current application, since precise statements are not easy to find in the literature.

Theorem F. 1 (Quantitative implicit function theorem for maps of Banach spaces). Let $\mathbb{K}=\mathbb{R}$ or $\mathbb{C}$, and $k \geq 1$ be an integer or $\infty$, and $\mathscr{X}, \mathscr{Y}$, and $\mathscr{Z}$ be Banach spaces over $\mathbb{K}$, and $\mathscr{U} \subset \mathscr{X}$ and $\mathscr{V} \subset \mathscr{Y}$ be open neighborhoods of points $x_{0} \in \mathscr{X}$ and $y_{0} \in \mathscr{Y}$, and $f: \mathscr{U} \times \mathscr{V} \rightarrow \mathscr{Z}$ be a $C^{k}$ (respectively, analytic) map such that $f\left(x_{0}, y_{0}\right)=0$ and the partial derivative of $f$ at $\left(x_{0}, y_{0}\right)$ with respect to the second variable,

$$
\begin{equation*}
D_{2} f\left(x_{0}, y_{0}\right) \in \mathscr{L}(\mathscr{Y}, \mathscr{Z}), \tag{F.1}
\end{equation*}
$$

is an isomorphism of Banach spaces. Let $\zeta \in(0,1]$ be small enough that $B_{\zeta}\left(x_{0}\right) \subset$ $\mathscr{U}$ and $B_{\zeta}\left(y_{0}\right) \subset \mathscr{V}$ and assume that

$$
\begin{equation*}
M:=\left\|\left(D_{2} f\left(x_{0}, y_{0}\right)\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Z}, \mathscr{Y})} \tag{F.2a}
\end{equation*}
$$

$$
\begin{equation*}
\sup _{(x, y) \in B_{\zeta}\left(x_{0}\right) \times B_{\zeta}\left(y_{0}\right)}\left\|D_{2} f(x, y)-D_{2} f\left(x_{0}, y_{0}\right)\right\|_{\mathscr{L}(\mathscr{X} \times \mathscr{Y}, \mathscr{Z})} \leq \frac{1}{2 M} \tag{F.2b}
\end{equation*}
$$

$$
\begin{equation*}
\beta:=\sup _{(x, y) \in B_{\zeta}\left(x_{0}\right) \times B_{\zeta}\left(y_{0}\right)}\left\|D_{1} f(x, y)\right\|_{\mathscr{L}(\mathscr{X}, \mathscr{Z})}<\infty . \tag{F.2c}
\end{equation*}
$$

Then there are a constant $\delta \in(0, \min \{\zeta, \zeta /(2 \beta M)\}]$ and a unique $C^{k}$ (respectively, analytic) map $g: \mathscr{X} \supset B_{\delta}\left(x_{0}\right) \rightarrow B_{\zeta}\left(y_{0}\right) \subset \mathscr{Y}$ such that $y_{0}=g\left(x_{0}\right)$ and

$$
\begin{equation*}
f(x, g(x))=0, \quad \forall x \in B_{\delta}\left(x_{0}\right), \tag{F.3a}
\end{equation*}
$$

$$
\begin{gather*}
D g(x)=-\left(D_{2} f(x, g(x))\right)^{-1} D_{1} f(x, g(x)) \in \mathscr{L}(\mathscr{X}, \mathscr{Y}), \quad \forall x \in B_{\delta}\left(x_{0}\right),  \tag{F.3b}\\
\left\|g\left(x_{1}\right)-g\left(x_{2}\right)\right\|_{\mathscr{Y}} \leq 2 \beta M\left\|x_{1}-x_{2}\right\|_{\mathscr{X}}, \quad \forall x_{1}, x_{2} \in B_{\delta}\left(x_{0}\right) .
\end{gather*}
$$

Remark F. 2 (Inverse and implicit function theorems for smooth or analytic maps of Banach spaces). Statements and proofs of the Inverse Function Theorem for $C^{k}$ maps of Banach spaces are provided by Abraham, Marsden, and Ratiu [1. Theorem 2.5.2], Apostol [5, Theorem 13.6] (Euclidean space only), Deimling [26, Theorem 4.15.2], Lang [82, Theorem XIV.1.2], and Zeidler [126, Theorem 4.F]; statements and proofs of the Inverse Function Theorem for analytic maps of Banach spaces are provided by Berger [7, Corollary 3.3.2] (complex), Deimling [26, Theorem 4.15.3] (real or complex), and Zeidler [126, Corollary 4.37] (real
or complex). The corresponding $C^{k}$ or Analytic Implicit Function Theorems are proved in the standard way as corollaries, for example [1, Theorem 2.5.7] and $\mathbf{1 2 6}$, Theorem 4.H], although the reverse strategy is used by Jost [73, Theorem 10.1] and others.

A quantitative version of the Inverse Function Theorem is given by Abraham, Marsden, and Ratiu as [1, Proposition 2.5.6], although in the latter version the conclusion (F.3C) on the Lipschitz constant for the inverse map, $g$, is omitted. See also Frøyshov 46, Proposition B.0.2], Kronheimer and Mrowka [80, Proposition 18.3.6], McDuff and Salamon [90, Proposition A.3.4], Mrowka and Rollin [93, Proposition 2.3.5], and Salamon [104, Theorem B.1] for closely-related versions of the Quantitative Inverse Function Theorem in the gauge-theory literature, though none quite as general as our Theorem F.4.

REmARK F. 3 (Alternative hypothesis for implicit function theorems for $C^{2}$ maps of Banach spaces). When $f$ is $C^{2}$, one can replace the hypothesis (F.2b) in Theorem F. 1 by the simpler condition,

$$
\begin{equation*}
K:=\sup _{(x, y) \in B_{\zeta}\left(x_{0}\right) \times B_{\zeta}\left(y_{0}\right)}\left\|D^{2} f(x, y)\right\|_{\mathscr{L}(\mathscr{X} \times \mathscr{Y}, \mathscr{Z})}<\infty \tag{F.4}
\end{equation*}
$$

The Mean Value Theorem for $C^{1}$ maps of Banach spaces (see (F.7) gives

$$
\begin{aligned}
D_{2} f(x, y)-D_{2} f\left(x_{0}, y_{0}\right)= & \int_{0}^{1} D_{21} f\left(t x+(1-t) x_{0}, t y+(1-t) y_{0}\right)\left(x-x_{0}\right) d t \\
& +\int_{0}^{1} D_{2}^{2} h\left(t x+(1-t) x_{0}, t y+(1-t) y_{0}\right)\left(y-y_{0}\right) d t \\
& \forall x \in \mathscr{U} \text { and } y \in \mathscr{V}
\end{aligned}
$$

assuming that $\mathscr{U}$ and $\mathscr{V}$ are convex without loss of generality. Therefore,

$$
\left\|D_{2} f(x, y)-D_{2}\left(x_{0}, y_{0}\right)\right\|_{\mathscr{L}(\mathscr{X} \times \mathscr{Y}, \mathscr{Z})} \leq K\left(\left\|x-x_{0}\right\|_{\mathscr{X}}+\left\|y-y_{0}\right\|_{\mathscr{Y}}\right)
$$

and thus it suffices to choose $\eta \in(0, \zeta]$ small enough that $\eta \leq 1 /(4 K M)$ in order to achieve (F.2b), with $\zeta$ replaced by $\eta$ there (so $(x, y) \in B_{\eta}\left(x_{0}\right) \times B_{\eta}\left(y_{0}\right)$ ) and in the conclusions of Theorem F. 1

Proof of Theorem F.1. When $\mathbb{K}=\mathbb{R}$ and in the absence of the supplementary hypothesis (F.2b), the existence of a unique $C^{k}$ map $g$ obeying (F.3a) is provided by standard statements of the Implicit Function Theorem (see Abraham, Marsden, and Ratiu [1, Theorem 2.5.7], Lang [82, Theorem XIV.2.1], and, in the case of Euclidean spaces, Apostol [5. Theorem 13.7]). For the cases $\mathbb{K}=\mathbb{C}$ or $f$ real or complex analytic, see the references cited in Remark F. 2 for proofs of the Inverse Function Theorem and recall that the Implicit Function Theorem can be obtained in a straightforward way as a corollary of the Inverse Function Theorem.

Hence, we need only focus on the quantitative conclusions of Theorem F. 1 and, for this purpose, we shall adapt Lang's proof of the Inverse Function Theorem (see Lang [82, Theorem XIV.1.2]) and Jost's proof of the Implicit Function Theorem [73, Theorem 10.1]. Recall that given $x$ in an open neighborhood of $x_{0}$, the existence of $y$ in an open neighborhood of $y_{0}$ solving $f(x, y)=0$ in F.3a) is equivalent to the existence of $y$ solving $h(x, y)=y$, where

$$
\begin{equation*}
h(x, y):=y-\left(D_{2} f\left(x_{0}, y_{0}\right)\right)^{-1} f(x, y) \tag{F.5}
\end{equation*}
$$

Observe that the smooth map $h: \mathscr{U} \times \mathscr{V} \rightarrow \mathscr{Y}$ has partial derivatives,

$$
\begin{aligned}
& D_{1} h(x, y)=-\left(D_{2} f\left(x_{0}, y_{0}\right)\right)^{-1} D_{1} f(x, y) \\
& D_{2} h(x, y)=\operatorname{id} \mathscr{Y}-\left(D_{2} f\left(x_{0}, y_{0}\right)\right)^{-1} D_{2} f(x, y), \quad \forall(x, y) \in \mathscr{U} \times \mathscr{V}
\end{aligned}
$$

and $D_{2} h\left(x_{0}, y_{0}\right)=0$. Moreover,

$$
\begin{aligned}
\left\|D_{2} h(x, y)\right\|_{\mathscr{L}(\mathscr{Y})} & =\left\|\left(D_{2} f\left(x_{0}, y_{0}\right)\right)^{-1}\left(D_{2} f\left(x_{0}, y_{0}\right)-D_{2} f(x, y)\right)\right\|_{\mathscr{L}(\mathscr{Y})} \\
& \leq\left\|\left(D_{2} f\left(x_{0}, y_{0}\right)\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Z}, \mathscr{Y})}\left\|D_{2} f\left(x_{0}, y_{0}\right)-D_{2} f(x, y)\right\|_{\mathscr{L}(\mathscr{Y}, \mathscr{Z})} .
\end{aligned}
$$

By hypothesis (F.2b) and the preceding bound for $\left\|D_{2} h(x, y)\right\|_{\mathscr{L}(\mathscr{Y})}$ and by hypothesis ( F .2 C$)$ and the preceding expression for $D_{1} h(x, y)$, we therefore have

$$
\begin{align*}
\left\|D_{1} h(x, y)\right\|_{\mathscr{L}(\mathscr{Y})} & \leq \beta M  \tag{F.6a}\\
\left\|D_{2} h(x, y)\right\|_{\mathscr{L}(\mathscr{Y})} & \leq \frac{1}{2}, \quad \forall(x, y) \times B_{\zeta}\left(x_{0}\right) \times B_{\zeta}\left(y_{0}\right) . \tag{F.6b}
\end{align*}
$$

The Mean Value Theorem for $C^{1}$ maps of Banach spaces [82, Theorem XIII.4.2] gives

$$
\begin{align*}
& h\left(x_{1}, y_{1}\right)-h\left(x_{2}, y_{2}\right)  \tag{F.7}\\
& =\int_{0}^{1} D_{1} h\left(t x_{1}+(1-t) x_{2}, t y_{1}+(1-t) y_{2}\right)\left(x_{1}-x_{2}\right) d t \\
& \quad+\int_{0}^{1} D_{2} h\left(t x_{1}+(1-t) x_{2}, t y_{1}+(1-t) y_{2}\right)\left(y_{1}-y_{2}\right) d t \\
& \quad \forall x_{1}, x_{2} \in \mathscr{U} \text { and } y_{1}, y_{2} \in \mathscr{V},
\end{align*}
$$

assuming that $\mathscr{U}$ and $\mathscr{V}$ are convex without loss of generality. Hence, by (F.6b) and (F.7)
(F.8) $\left\|h\left(x, y_{1}\right)-h\left(x, y_{2}\right)\right\| \mathscr{Y} \leq \frac{1}{2}\left\|y_{1}-y_{2}\right\|_{\mathscr{Y}}, \quad \forall x \in \bar{B}_{\zeta}\left(x_{0}\right)$ and $y_{1}, y_{2} \in \bar{B}_{\zeta}\left(y_{0}\right)$.

Moreover,

$$
\begin{aligned}
\left\|h(x, y)-y_{0}\right\|_{\mathscr{Y}} & =\left\|h(x, y)-h\left(x_{0}, y_{0}\right)\right\|_{\mathscr{Y}} \quad\left(\text { by } h\left(x_{0}, y_{0}\right)=y_{0}\right) \\
& \leq\left\|h(x, y)-h\left(x_{0}, y\right)\right\|_{\mathscr{Y}}+\left\|h\left(x_{0}, y\right)-h\left(x_{0}, y_{0}\right)\right\|_{\mathscr{Y}} \\
& \leq \beta M\left\|x-x_{0}\right\|_{\mathscr{X}}+\frac{1}{2}\left\|y-y_{0}\right\|_{\mathscr{Y}} \quad(\text { by (F.6a), (F.7), and (F.8)). }
\end{aligned}
$$

Hence, for $\delta \in(0, \zeta]$ obeying $\delta \leq \zeta /(2 \beta M)$, we obtain

$$
\begin{equation*}
\left\|h(x, y)-y_{0}\right\|_{\mathscr{Y}} \leq \zeta, \quad \forall(x, y) \in B_{\delta}\left(x_{0}\right) \times B_{\zeta}\left(y_{0}\right) \tag{F.9}
\end{equation*}
$$

Thus, for each $x \in \bar{B}_{\delta}\left(x_{0}\right)$, the map $h(x, \cdot): \bar{B}_{\zeta}\left(y_{0}\right) \rightarrow \bar{B}_{\zeta}\left(y_{0}\right)$ is a contraction. We can now apply the Banach Contraction Mapping Lemma (for example, see 73 Lemma 10.2] or [82, Lemma XIV.1.1]) to find a unique $y=g(x) \in \bar{B}_{\zeta}\left(y_{0}\right)$ for every $x \in \bar{B}_{\delta}\left(x_{0}\right)$ such that $h(x, y)=y$, that is, $f(x, y)=0$, and $y_{0}=g\left(x_{0}\right)$.

To prove (F.3c), we shall apply the Mean Value Theorem,

$$
\begin{equation*}
g\left(x_{1}\right)-g\left(x_{2}\right)=\int_{0}^{1} D g\left(t x_{1}+(1-t) x_{2}\right)\left(x_{1}-x_{2}\right) d t, \quad \forall x \in B_{\delta}\left(x_{0}\right) \tag{F.10}
\end{equation*}
$$

Computing the derivative of the equation (F.3a) with respect to $x$ yields

$$
D_{1} f(x, g(x))+D_{2} f(x, g(x)) D g(x)=0 \in \mathscr{L}(\mathscr{X}, \mathscr{Z}), \quad \forall x \in B_{\delta}\left(x_{0}\right),
$$

so that F.3b) holds provided $D_{2} f(x, g(x))$ is invertible. To verify invertibility, we define an operator $T_{x} \in \mathscr{L}(\mathscr{Y})$ by

$$
\left(D_{2} f\left(x_{0}, g\left(x_{0}\right)\right)\right)^{-1} D_{2} f(x, g(x))=: \mathrm{id} \mathscr{Y}-T_{x} \in \mathscr{L}(\mathscr{Y}), \quad \forall x \in B_{\delta}\left(x_{0}\right)
$$

Observe that

$$
T_{x}=\mathrm{id} \mathscr{Y}-\left(D_{2} f\left(x_{0}, g\left(x_{0}\right)\right)\right)^{-1} D_{2} f(x, g(x))=D_{2} h(x, g(x)) \quad(\text { by (F.5) })
$$

and so $\left\|T_{x}\right\|_{\mathscr{L}(\mathscr{Y})} \leq 1 / 2$ for all $x \in B_{\delta}\left(x_{0}\right)$ by (F.6b). Since $\left\|T_{x}\right\|_{\mathscr{L}(\mathscr{Y})}<1$, for all $x \in B_{\delta}\left(x_{0}\right)$, the operator id $\mathscr{y}-T_{x}$ is invertible, for all $x \in B_{\delta}\left(x_{0}\right)$. Hence, $D_{2} f(x, g(x)) \in \mathscr{L}(\mathscr{Y}, \mathscr{Z})$ is invertible for all $x \in B_{\delta}\left(x_{0}\right)$, as we had claimed in (F.3b), with

$$
\begin{equation*}
\left(D_{2} f(x, g(x))\right)^{-1}=\left(\mathrm{id}_{\mathscr{Y}}-T_{x}\right)^{-1}\left(D_{2} f\left(x_{0}, g\left(x_{0}\right)\right)\right)^{-1}, \quad \forall x \in B_{\delta}\left(x_{0}\right) \tag{F.11}
\end{equation*}
$$

We can estimate the operator norm of $\left(\text { id }_{\mathscr{Y}}-T_{x}\right)^{-1}$ (and hence that of $\left.\left(D_{2} f(x, g(x))\right)^{-1}\right)$ using the Neumann series,

$$
\left(\mathrm{id} \mathscr{Y}-T_{x}\right)^{-1}=\sum_{k=0}^{\infty} T_{x}^{k}
$$

to give

$$
\left\|\left(\mathrm{id}_{\mathscr{Y}}-T_{x}\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Y})} \leq \sum_{k=0}^{\infty}\left\|T_{x}^{k}\right\|_{\mathscr{L}(\mathscr{Y})} \leq \sum_{k=0}^{\infty}\left\|T_{x}\right\|_{\mathscr{L}(\mathscr{Y})}^{k}=\frac{1}{1-\left\|T_{x}\right\|_{\mathscr{L}(\mathscr{Y})}}
$$

and consequently,

$$
\begin{equation*}
\left\|\left(\mathrm{id} \mathscr{Y}-T_{x}\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Y})} \leq 2, \quad \forall x \in B_{\delta}\left(x_{0}\right) \tag{F.12}
\end{equation*}
$$

We can therefore estimate the operator norm of $\left(D_{2} f(x, g(x))\right)^{-1}$ by

$$
\begin{aligned}
& \left\|\left(D_{2} f(x, g(x))\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Z}, \mathscr{Y})} \\
& \quad=\left\|\left(\mathrm{id}_{\mathscr{Y}}-T_{x}\right)^{-1}\left(D_{2} f\left(x_{0}, g\left(x_{0}\right)\right)\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Z}, \mathscr{Y}) \quad(\text { by (F.11) })} \\
& \quad \leq\left\|\left(\mathrm{id}_{\mathscr{Y}}-T_{x}\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Y})}\left\|\left(D_{2} f\left(x_{0}, g\left(x_{0}\right)\right)\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Z}, \mathscr{Y})} \\
& \quad \leq 2 M \quad\left(\text { by } \quad \text { F.2a) and } y_{0}=g\left(x_{0}\right) \text { and (F.12) }\right), \quad \forall x \in B_{\delta}\left(x_{0}\right) .
\end{aligned}
$$

By hypothesis (F.2c) and the fact that $g\left(B_{\delta}\left(x_{0}\right)\right) \subset B_{\zeta}\left(x_{0}\right)$ by construction, we have

$$
\sup _{x \in B_{\delta}\left(x_{0}\right)}\left\|D_{1} f(x, g(x))\right\|_{\mathscr{L}(\mathscr{X} \times \mathscr{Y}, \mathscr{Z})} \leq \beta
$$

and so the expression (F.3b) for $D g(x)$ and the preceding two estimates give

$$
\begin{aligned}
\|D g(x)\|_{\mathscr{L}(\mathscr{X}, \mathscr{Y})} & =\left\|\left(D_{2} f(x, g(x))\right)^{-1} D_{1} f(x, g(x))\right\|_{\mathscr{L}(\mathscr{X}, \mathscr{Y})} \\
& \leq\left\|\left(D_{2} f(x, g(x))\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Z}, \mathscr{Y})}\left\|D_{1} f(x, g(x))\right\|_{\mathscr{L}(\mathscr{X}, \mathscr{Z})} \\
& \leq 2 \beta M, \quad \forall x \in B_{\delta}\left(x_{0}\right)
\end{aligned}
$$

The Lipschitz bound (F.3c) now follows from the preceding estimate and (F.10). This completes the proof of Theorem F. 1

The quantitative inverse function theorem below (see [34, Theorem 3.2]) follows in a standard way from the preceding quantitative implicit function theorem.

Theorem F. 4 (Quantitative inverse function theorem for maps of Banach spaces). Let $\mathbb{K}=\mathbb{R}$ or $\mathbb{C}$, and $k \geq 1$ be an integer or $\infty$, and $\mathscr{X}$ and $\mathscr{Y}$ be Banach spaces over $\mathbb{K}$, and $\mathscr{U} \subset \mathscr{X}^{-}$be an open neighborhood of a point $x_{0} \in \mathscr{X}$, and $f: \mathscr{U} \rightarrow \mathscr{Y}$ be a $C^{k}$ (respectively, analytic) map such that $f\left(x_{0}\right)=y_{0}$ and the derivative of $f$ at $x_{0}$,

$$
\begin{equation*}
D f\left(x_{0}\right) \in \mathscr{L}(\mathscr{X}, \mathscr{Y}), \tag{F.13}
\end{equation*}
$$

is an isomorphism of Banach spaces. Let $\zeta \in(0,1]$ be small enough that $B_{\zeta}\left(x_{0}\right) \subset$ $\mathscr{U}$ and assume that

$$
\begin{gather*}
M:=\left\|\left(D f\left(x_{0}\right)\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Y}, \mathscr{X})},  \tag{F.14a}\\
\sup _{x \in B_{\zeta}\left(x_{0}\right)}\left\|D f(x)-D f\left(x_{0}\right)\right\|_{\mathscr{L}(\mathscr{X}, \mathscr{Y})} \leq \frac{1}{2 M} .
\end{gather*}
$$

Then there are a constant $\delta \in(0, \min \{\zeta, \zeta /(2 M)\}]$ and a unique $C^{k}$ (respectively, analytic) map $g: \mathscr{Y} \supset B_{\delta}\left(y_{0}\right) \rightarrow B_{\zeta}\left(x_{0}\right) \subset \mathscr{X}$ such that $x_{0}=g\left(y_{0}\right)$ and

$$
\begin{align*}
f(g(y)) & =y, \quad \forall y \in B_{\delta}\left(y_{0}\right),  \tag{F.15a}\\
D g(x) & =(D f(x))^{-1} \in \mathscr{L}(\mathscr{Y}, \mathscr{X}), \quad \forall x \in B_{\delta}\left(x_{0}\right),  \tag{F.15b}\\
\left\|g\left(y_{1}\right)-g\left(y_{2}\right)\right\|_{\mathscr{X}} & \leq 2 M\left\|y_{1}-y_{2}\right\|_{\mathscr{Y}}, \quad \forall y_{1}, y_{2} \in B_{\delta}\left(y_{0}\right) . \tag{F.15c}
\end{align*}
$$

REMARK F. 5 (Alternative hypothesis for inverse function theorems for $C^{2}$ maps of Banach spaces). When $f$ is $C^{2}$, one can replace the hypothesis (F.14b) in Theorem F. 4 by the simpler condition,

$$
\begin{equation*}
K:=\sup _{x \in B_{\zeta}\left(x_{0}\right)}\left\|D^{2} f(x)\right\|_{\mathscr{L}(\mathscr{X}, \mathscr{Y})}<\infty \tag{F.16}
\end{equation*}
$$

Just as in Remark F.3 the Mean Value Theorem for $C^{1}$ maps of Banach spaces ensures that we can then choose $\eta \in(0, \zeta]$ small enough that $\eta \leq 1 /(2 K M)$ in order to achieve (F.14b), with $\zeta$ replaced by $\eta$ there (so $x \in B_{\eta}\left(x_{0}\right)$ ) and in the conclusions of Theorem F.4. The hypothesis (F.16) is the one used in [1, Proposition 2.5.6] and 46, Proposition B.0.2], whereas the remaining examples of quantitative inverse function theorems cited in Remark F. 2 use the hypothesis (F.14b).

Proof of Theorem F.4 Set $\mathscr{Z}:=\mathscr{Y}$ and $F(x, y):=f(x)-y$ and observe that $D_{1} F(x, y)=D f(x)$, so (F.1) is obeyed (interchanging the roles of $x$ and $y$ ). Moreover, $\left(D_{1} F\left(x_{0}, y_{0}\right)\right)^{-1}=\left(D f\left(x_{0}\right)\right)^{-1}$, so $\left\|\left(D_{1} F\left(x_{0}, y_{0}\right)\right)^{-1}\right\|_{\mathscr{L}(\mathscr{Y}, \mathscr{X})}=M$ in (F.2a) by (F.14a). Also, (F.14b) yields F.2b) for all $y \in \mathscr{Y}$. Finally, $D_{2} F(x, y)=$ -id $\mathscr{y}$ for all $x \in \mathscr{U}$ and $y \in \mathscr{Y}$, so $\beta=1$ in (F.2c). Hence, Theorem F. 1 applied to $F(x, y)$ yields $\delta \in(0, \min \{\zeta, \zeta /(2 M)\}]$ and a $C^{k}$ (respectively, analytic) map $g: \mathscr{Y} \supset B_{\delta}\left(y_{0}\right) \rightarrow B_{\zeta}\left(x_{0}\right) \subset \mathscr{X}$ such that $F(g(y), y)=0$, that is, $f(g(y))=y$ for all $y \in B_{\delta}\left(y_{0}\right)$. Lastly, (F.3b) yields (F.15b) and (F.3c) yields (F.15c).

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[^0]:    ${ }^{1}$ To appear in Memoirs of the American Mathematical Society.

[^1]:    ${ }^{1}$ Unless stated otherwise, all Banach spaces are considered to be real in this monograph.

[^2]:    ${ }^{1}$ In the sense of Remark 2.2 .2

[^3]:    ${ }^{2}$ In applications of Proposition 2.5.1 we can choose $\delta=1$ without loss of generality.
    ${ }^{3}$ The use of the constant $M$ could be avoided if we replaced the right-hand-side of Inequality (2.5.6) by $N\left(1+\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}\right)\left\|A-A_{0}\right\|_{W_{A_{1}}^{1, p}(X)}$.

[^4]:    ${ }^{4}$ The Quantitative Implicit Function Theorem F. 1 can be used to give a precise bound on $\chi_{s}$ given a bound on $b_{s}$.

[^5]:    ${ }^{5}$ For example, when $d \geq 4$, we have $1 / q^{*}+1 / q<1 / d+2 / d=3 / d<1$ for $d \geq 4$ and thus, when $d \geq 4$, it is easy to see that we can choose $r$ as stated.

[^6]:    ${ }^{1}$ Observe that $d / 2 \geq 4 d /(d+4)$ for all $d \geq 4$ and $2 \geq 4 d /(d+4)$ when $d=2,3$, so if $p$ obeys the hypotheses of Theorem 4 then it obeys the hypotheses of Proposition 3.1.1.

[^7]:    ${ }^{1}$ Gilkey and Hörmander allow $m, k \in \mathbb{R}$, as we could also in Theorem A. 1 but we shall omit this refinement.

[^8]:    ${ }^{2}$ This could be accomplished geometrically with the aid of a Dirac operator, $D: C^{\infty}(M ; V) \rightarrow$ $C^{\infty}(M ; V)$, when $V$ and $W$ are spinor bundles and $M$ admits a $\operatorname{spin}^{c}$ structure 84.

[^9]:    ${ }^{1}$ See 14 Section 7.4], 76, Section 5.3.3], or 125 Section 7.3].

