

Grafting Seiberg-Witten monopoles

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Abstract We demonstrate that the operation of taking disjoint unions of J -holomorphic curves (and thus obtaining new J -holomorphic curves) has a Seiberg-Witten counterpart. The main theorem asserts that, given two solutions $(A_i; \rho_i)$, $i = 0, 1$ of the Seiberg-Witten equations for the $Spin^c$ -structures $W_{E_i}^+ = E_i \oplus (E_i \otimes K^{-1})$ (with certain restrictions), there is a solution $(A; \rho)$ of the Seiberg-Witten equations for the $Spin^c$ -structure W_E with $E = E_0 \oplus E_1$, obtained by "grafting" the two solutions $(A_i; \rho_i)$.

AMS Classification 53D99, 57R57; 53C27, 58J05

Keywords Symplectic 4-manifolds, Seiberg-Witten gauge theory, J -holomorphic curves

1 Introduction

In his series of groundbreaking works [4], [5], [6], Taubes showed that the Seiberg-Witten invariants and the Gromov-Witten invariants (as defined in [7]) for a symplectic 4-manifold $(X; \omega)$ are the same. His results opened the door to a whole new world of interactions between the two theories that had previously only been speculations. The most spectacular outcomes of this interplay were new results that in one theory were obvious but when translated into the other theory, became highly nontrivial. An example of such a phenomenon is the simple formula relating the Seiberg-Witten invariant of a $Spin^c$ -structure W to the Seiberg-Witten invariant of its dual $Spin^c$ -structure W^+ , i.e. the one with $c_1(W^+) = -c_1(W)$. The formula reads:

$$SW_X(W^+) = SW_X(W)$$

When translated into the Gromov-Witten language, this duality becomes

$$Gr_X(E) = Gr_X(K - E) \quad (1)$$

Here K is the canonical class of $(X; \omega)$ and $E \in H^2(X; \mathbb{Z})$ is related to W as $c_1(W^+) = 2E - K$. This is a highly nonobvious result about J -holomorphic curves, even in the simplest case when $E = 0$. In that case we obtain that

$Gr_X(K) = Gr_X(0) = 1$, the latter equation simply being the definition of $Gr_X(0)$. This gives an existence result of a J -holomorphic representative for the class K , a result unknown prior to Taubes' theorem. The formula (1) has recently been proved by S. Donaldson and I. Smith [1] without any reference to Seiberg-Witten theory (but under slightly stronger restrictions on $(X; !)$ than in Taubes' theorem).

In the author's opinion, proving a result about Gromov-Witten theory which had only been known through its relation with Seiberg-Witten theory, without relying on the latter, has a number of benefits. One is to understand Gromov-Witten theory from within better. But also to possibly generalize the theorem to a broader class of manifolds. Recall that Taubes' theorem equates the two invariants only on symplectic 4-manifolds. Both Seiberg-Witten and Gromov-Witten theory are defined over larger sets of manifolds, namely all smooth 4-manifolds and all symplectic manifolds (of any dimension) respectively. On the other hand, even within the category of symplectic 4-manifolds, one can hope for more nonvanishing theorems i.e. theorems of the type $Gr_X(E) \neq 0$ for classes $E \neq 0; K$. The techniques used by Donaldson and Smith are promising in that direction.

The aim of this paper is to prove a result in the same vein but going the opposite direction. Namely, on the Gromov-Witten side, given two classes $E_i \in H^2(X; \mathbb{Z})$, $i = 0; 1$ with $E_0 \cdot E_1 = 0$ and J -holomorphic curves σ_i with $[\sigma_i] = P:D:(E_i)$, one can define a new J -holomorphic curve $\sigma = \sigma_0 \# \sigma_1$. By the assumption $E_0 \cdot E_1 = 0$, the two curves σ_i are either disjoint or share toroidal components (see [2]). In the former case, σ is simply the disjoint union of σ_0 and σ_1 and in the latter case one needs to replace the tori shared by σ_0 and σ_1 with their appropriate multiple covers. This induces a map on moduli spaces:

$$M_X^{Gr}(E_0) \times M_X^{Gr}(E_1) \xrightarrow{\#} M_X^{Gr}(E_0 + E_1) \quad (2)$$

This article describes the Seiberg-Witten counterpart of (2). That is, given two complex line bundles E_0 and E_1 (with certain restrictions, see the assumption 3.1 below for a precise statement) and two solutions $(A_i; \sigma_i)$ of the Seiberg-Witten equations for the $Spin^c$ -structures $W_{E_i} = E_i \oplus (E_i \otimes K^{-1})$, $i = 0; 1$ and with Taubes' large r perturbation, we show how to produce a solution $(A; \sigma) = (A_0; \sigma_0) \# (A_1; \sigma_1)$ for the $Spin^c$ -structure W_E with $E = E_0 \oplus E_1$. This operation induces the following commutative diagram:

$$\begin{array}{ccc}
 M_X^{SW}(E_0) & M_X^{SW}(E_1) & \xrightarrow{t} M_X^{SW}(E_0 + E_1) \\
 \downarrow \cong & \downarrow \cong & \downarrow \cong \\
 M_X^{Gr}(E_0) & M_X^{Gr}(E_1) & \xrightarrow{t} M_X^{Gr}(E_0 + E_1)
 \end{array} \tag{3}$$

Here the map $\tau : M_X^{SW}(E) \rightarrow M_X^{Gr}(E)$ is the map described in [4] that associates to each solution of the Seiberg-Witten equations an embedded J -holomorphic curve. The solution $(A; \psi)$ is obtained by "grafting" the two solutions $(A_i; \psi_i)$. The key observation here is that for the large r version of Taubes' perturbation, a solution $(B; \psi)$ of the Seiberg-Witten equations for the $Spin^c$ -structure W_E is "concentrated" near the zero set of $\rho_{\bar{r}}$, the E component of ψ . That is, the restriction of $(B; \psi)$ to the complement of a regular neighborhood of $\psi^{-1}(0)$ converges pointwise (under certain bundle identifications) to the unique solution $(A_0; \rho_{\bar{r}} u_0)$ for the anticanonical $Spin^c$ -structure $W_0 = \mathbb{C} \otimes K^{-1}$. This is used to define a first approximation of ψ by declaring it to be equal to ψ_i in a regular neighborhood V_i of $\psi_i^{-1}(0)$ and equal to $\rho_{\bar{r}} u_0$ on the complement of $V_0 \cup V_1$. Bump functions are used to produce a smooth spinor. The first approximation of A is simply the product connection $A_0 + A_1$. The contraction mapping principle is then evoked to deform this approximate solution to an honest solution of the Seiberg-Witten equations. The author has learned the techniques employed in this article from the inspiring work of Taubes on gauge theory of symplectic 4-manifolds, most notably from [5].

The article is organized as follows. In section 2 we review the needed Seiberg-Witten theory on symplectic 4-manifolds. Section 3 explains how to define an "almost" monopole $(A^0; \psi^0)$ from a pair of monopoles $(A_i; \psi_i)$, $i = 0, 1$. It also analyzes the asymptotic (as $r \rightarrow \infty$) regularity theory for the linearized operators $L_{(A_i; \psi_i)}$ and deduces a corresponding result for $L_{(A^0; \psi^0)}$. The latter is used in combination with the contraction mapping principle to obtain an "honest" monopole $(A; \psi)$. Section 4 compares the present method of grafting monopoles to the one used in exploring Seiberg-Witten theory on manifolds X which are obtained as a fiber sum: $X = X_1 \# X_2$. Section 5 proves a converse to theorem 3.11. It explains which monopoles in the $Spin^c$ -structure W_E can be obtained as products of monopoles $(A_i; \psi_i)$ in the $Spin^c$ -structures W_{E_i} , $i = 0, 1$ with $E_0 + E_1 = E$ and with the property that $(A_i; \psi_i)$ does not contain multiply covered tori.

Acknowledgment The author would like to express his gratitude to his thesis advisor, professor Ron Fintushel, for his continuing help and encouragement during the process of writing this article, the author's doctoral thesis.

2 The Seiberg-Witten equations on symplectic manifolds

2.1 Introduction

Let $(X; \omega)$ be a symplectic, smooth, compact 4-manifold with symplectic form ω . Denote by \mathcal{J} the set of all almost complex structures J on TX that are compatible with ω , i.e. the ones for which

$$g(v; w) = g_J(v; w) = \omega(v; Jw) \quad v, w \in TX$$

defines a Riemannian metric on X . Given a $J \in \mathcal{J}$, the associated metric g_J will always be assumed throughout to be the metric of choice.

On any almost complex 4-manifold there is a anticanonical $Spin^c$ -structure $W_0 = W_0^+ \oplus W_0^-$ determined by the almost complex structure as:

$$\begin{aligned} W_0^+ &= \{v \in T_x X \mid \omega(v, Jv) > 0\} \\ W_0^- &= \{v \in T_x X \mid \omega(v, Jv) < 0\} \\ v &= \frac{\rho}{2} v_{0,1} \wedge - v_{0,1} \quad v \in T_x X; \quad \rho \in \mathbb{R} \end{aligned}$$

In the above, $v_{0,1} \in \wedge^{0,1} T_x X$ denotes the $(0,1)$ projection of $v \in T_x X$, the dual of $v \in T_x X$. All other $Spin^c$ -structures can be obtained from W_0 by tensoring it with a complex line bundle E and extending Clifford multiplication trivially over the E factor, i.e.

$$W_E = E \otimes W_0$$

$$v \cdot (\cdot) = \rho \cdot (v \cdot) \quad \rho \in E_x; \quad v \in T_x X; \quad \rho \in W_{0,x}; \quad x \in X$$

The symplectic form ω induces a splitting of $\wedge^{2,+} T_x X$ as

$$\wedge^{2,+} T_x X = \mathbb{R} \omega \oplus \wedge^{0,2} T_x X \tag{4}$$

which will be used below to write the curvature component of the Seiberg-Witten equations as two equations, one for each of the summands on the right-hand side of (4).

Given a $Spin^c$ -structure W_E on X , the Seiberg-Witten equations are a coupled, elliptic system of equations for a pair (A, ψ) of a connection A on E and a positive spinor $\psi \in W_E^+ = (E \otimes (E \otimes K^{-1}))$. The connection A on E together with a fixed connection A_0 on K^{-1} (which will be made specific in a bit), induces a $Spin^c$ -connection on W_E which we will denote by ∇^A and

which in turn gives rise to the Dirac operator $D_A : (W_E)^+ \rightarrow (W_E)^-$. It proves convenient to write the spinor ψ in the form

$$\psi = \rho_{\bar{r}}(\psi_+ ; \psi_-) \quad 2 (E); \quad 2 (E \oplus K^{-1})$$

where $r = 1$ is a parameter whose significance will become clear later. With this understood, the Seiberg-Witten equations read:

$$D_A(\psi) = 0$$

$$F_A^+ = q(\psi_+ ; \psi_-) + \quad (5)$$

Here $2i^{2,+}$ is a fixed imaginary, self-dual two form on X and $q : (W_E^+)^+ \rightarrow (W_E^+)^- \oplus i^{2,+}$ is the bilinear quadratic map given explicitly by

$$q(\psi_+ ; \psi_-) = \frac{ir}{8}(j \psi_+^2 - j \psi_-^2) + \frac{ir}{4}(\psi_+ + \psi_-) \quad (6)$$

2.2 The anticanonical $Spin^c$ -structure

It is another result of Taubes' [8] that the Seiberg-Witten invariant of the anticanonical $Spin^c$ -structure on a symplectic manifold is equal to ± 1 . Furthermore, the equations have exactly one solution $(A_0, \rho_{\bar{r}} u_0)$, $u_0 \in \mathbb{C} \oplus K^{-1}$, for the choice of

$$F_{A_0}^+ = \frac{ir}{8} \quad (7)$$

in (5) and for $r = 1$. The purpose of this section is to describe the solution $(A_0, \rho_{\bar{r}} u_0)$ and its linearized operator.

The pair $(A_0, \rho_{\bar{r}} u_0)$ is characterized (up to gauge) by the condition

$$hr^0 u_0 ; u_0 i = 0 \quad (8)$$

(where r^0 is the $Spin^c$ -connection induced by A_0) and can be obtained as follows: let u_0 be any section of $\mathbb{C} \oplus K^{-1}$ with $\int u_0 j = 1$ and whose projection onto the second summand is zero. Likewise, let A be any connection on K^{-1} and let r^A be its induced $Spin^c$ -connection on $W_0^+ = \mathbb{C} \oplus K^{-1}$. Set $a = \int hu_0 ; r^A u_0 i$. This defines an imaginary valued 1-form as can easily be seen:

$$a + a = \int hr^A u_0 ; u_0 i + \int hu_0 ; r^A u_0 i = \int dj u_0 j^2 = 0$$

Define the connection A_0 on K^{-1} by $A_0 = A - a$ which induces the $Spin^c$ -connection $r^0 = r^A - a$ on W_0^+ . This connection clearly satisfies (8). With the choice of ψ as in (7), the Seiberg-Witten equations (5) take the form:

$$D_A(\psi) = 0$$

$$F_A^+ = \frac{ir}{8}(j \psi_+^2 - 1 - j \psi_-^2) + F_{A_0}^+ + \frac{ir}{4}(\psi_+ + \psi_-) \quad (9)$$

Since the $\bar{\partial}$ -component of u_0 is zero and since $\int \langle j, j \rangle = \int u_0 j = 1$, the pair $(A_0; u_0)$ clearly solves the second equation of (9). The fact that it also solves the first equation relies on the property $d! = 0$ of $!$ as well as (8). Taubes [8] showed that there are, up to gauge, no other solutions to (9) and, as we shall presently see, that the solution $(A_0; u_0)$ is a smooth solution in the sense that the linearization of (9) at $(A_0; u_0)$ has trivial cokernel. These two facts together show that $SW_X(W_0) = \mathbb{Z}$.

Define $S : L^{1,2}(j^{-1} W_0^+) \rightarrow L^2(j^{-1} W_0^-)$ to be the linearized Seiberg-Witten operator for the solution $(A_0; u_0)$. Thus, for $(b; (c; \gamma)) \in L^{1,2}(j^{-1} W_0^+) \oplus (L^2(K^{-1}))$ we have:

$$S(b; (c; \gamma)) = \begin{pmatrix} d b + i \frac{\rho_r}{2} \text{Im}(u_0 \cdot c); \\ d^* b - \frac{\rho_r}{2} r q(c; u_0) - \frac{\rho_r}{2} r q(u_0; \gamma); \\ D_{A_0}(c; \gamma) + \frac{\rho_r}{2} b \cdot u_0 \end{pmatrix}; \tag{10}$$

Let $S^* : L^2(j^{-1} W_0^-) \rightarrow L^{1,2}(j^{-1} W_0^+)$ be the formal adjoint of S . The following proposition and corollary are proved in [5], section 4.

Proposition 2.1 *Let S and S^* be as above. Then the operator SS^* on $L^2(j^{-1} W_0^-)$ is given by*

$$SS^* = \frac{1}{4} r^0; r^0 + R_0 + \frac{\rho_r}{2} R_1 + \frac{r}{8} \tag{11}$$

where $r^0;$ is the adjoint of r^0 and where $R_i; i = 0; 1$ are certain r -independent endomorphism on $L^2(j^{-1} W_0^-)$.

The proof is a straightforward calculation, terms of the form $D_{A_0} D_{A_0}$ are simplified using the Weitzenböck formula for the Dirac operator. An important consequence of (11) is the following:

Corollary 2.2 *With S and S^* as above, the smallest eigenvalue λ_1 of SS^* is bounded from below by $r=16$. In particular, S is invertible and S^{-1} satisfies the bounds*

$$\|j S^{-1} y\|_2 \leq \frac{4}{r} \|j y\|_2 \quad \text{and} \quad \|j S^{-1} y\|_{1,2} \leq C \|j y\|_2 \tag{12}$$

where C is r -independent.

2.3 The general case and bounds on $(a; \cdot)$

Consider now a $Spin^c$ -structure $W_E = E \oplus W_0$ on X . The connection A_0 on K^{-1} and a choice of a connection B_0 on E together induce a connection $B_0 \oplus A_0$ on $E \oplus K^{-1} = c_1(W_E^+)$ by the product rule:

$$B_0 \oplus A_0(\cdot_1 \oplus \cdot_2) = B_0(\cdot_1) \oplus \cdot_2 + \cdot_1 \oplus B_0(\cdot_2) + \cdot_1 \oplus \cdot_2 \oplus A_0(\cdot)$$

The space of connections on $E \oplus K^{-1}$ is an affine space with associated vector space $i \oplus \frac{1}{X}$. With the choice of a "base" connection $B_0 \oplus A_0$ in place, we will from now on regard solutions to the Seiberg-Witten equations as pairs $(a; \cdot) \in i \oplus \frac{1}{X} \oplus (W_E^+)$ rather than $(A; \cdot) \in \text{Conn}(E \oplus K^{-1}) \oplus (W_E^+)$, the relation between the two being:

$$A = B_0 \oplus A_0 + a$$

We will agree to use henceforth the choice of \cdot in (5) to be:

$$\cdot = -\frac{ir}{8} \cdot + F_{A_0}^+ \tag{13}$$

For $\cdot \in (E \oplus \mathbb{C} \oplus K^{-1})$ we will write $\cdot = \rho_{\bar{r}}(\cdot, u_0; \cdot)$ with $\cdot \in (E)$ and $\cdot \in (E \oplus K^{-1})$ and u_0 as in the previous section.

With these conventions understood and with the use of (4), the Seiberg-Witten equations (5) become:

$$\begin{aligned} D_a \cdot &= 0 \\ F_a^{1,1} &= \frac{ir}{8} (j \cdot^2 - j \cdot^2 - 1) \cdot \\ F_a^{0,2} &= \frac{ir}{4} \end{aligned} \tag{14}$$

Here F_a^{ij} is the orthogonal projection of $2 F_{B_0}^+ + d^* a$ onto $i \oplus j$. The linearized operator $L_{(a; \cdot)} : L^{1,2}(i \oplus W_E^+) \oplus L^2(i \oplus i^{2,+} \oplus W_E^-)$ of the Seiberg-Witten equations for a solution $(a; \cdot)$ of (14) is:

$$\begin{aligned} L_{(a; \cdot)}(b; (\cdot; \cdot)) &= d^* b + i \frac{\rho_{\bar{r}}}{2} \text{Im}(h; \cdot; \cdot + h; \cdot; \cdot); \\ &= d^* b - \rho_{\bar{r}} q(\cdot; \cdot; \cdot) - \rho_{\bar{r}} q(\cdot; \cdot; \cdot); \\ &= D_a(\cdot; \cdot) + \frac{\rho_{\bar{r}}}{2} b; \end{aligned} \tag{15}$$

It is another result of Taubes' that the operator $L_{(a; \cdot)}$ has Fredholm index zero on a symplectic manifold with $b_2^+ = 2$, provided that E is a basic class. As we will use this fact repeatedly throughout the paper, we give a short proof of it here:

Theorem 2.3 (Taubes) *Let X be a symplectic manifold with $b_2^+ \geq 2$ and $E \in H^2(X; \mathbb{Z})$ a basic class, i.e. $SW(W_E) \neq 0$. Let $(a; \cdot)$ be a solution of (14) and $L_{(a; \cdot)}$ be the operator defined by (15). Then the Fredholm index of $L_{(a; \cdot)}$ is equal to zero.*

Proof As E is assumed to be a Seiberg-Witten basic class, it has to also be a Gromov-Witten basic class. In particular, the dimension of the Gromov-Witten moduli space has to be non-negative:

$$\dim M^{Gr}(E) = \frac{1}{2}(E^2 - K \cdot E) \geq 0$$

Let Σ be an embedded J -holomorphic curve in X with $[\Sigma] = P \cdot D_+(E)$. Then the adjunction formula for Σ states:

$$2g - 2 = E^2 \cdot \Sigma + K \cdot \Sigma$$

Combining these last two relations we obtain two inequalities:

$$E^2 \cdot \Sigma \geq g - 1 \quad \text{and} \quad K \cdot \Sigma \leq g - 1$$

Let $n \geq 0$ be the integer such that $E^2 \cdot \Sigma = g - 1 + n$ and $K \cdot \Sigma = g - 1 - n$. Since E is a Gromov-Witten basic class, by duality, so is $K - E$. But then (by positivity of intersection of J -holomorphic curves) we must have:

$$0 \leq E \cdot (K - E) = E \cdot K - E^2 = g - 1 - n - (g - 1 + n) = -2n \leq 0$$

This forces $n = 0$ and so $E^2 = g - 1 = K \cdot \Sigma$. Using these in the index formula for $L_{(A; \cdot)}$ immediately yields the desired result:

$$\text{Ind } L_{(a; \cdot)} = \frac{1}{4} (2E - K)^2 - (3 + 2e) = \frac{1}{4} K^2 - (3 + 2e) = 0 \quad \square$$

We also use this section to remind the reader of several useful bounds that a solution $(a; \cdot)$ of the Seiberg-Witten equations satisfies. These bounds are provided courtesy of [4] and their proofs rely solely on properties of the Seiberg-Witten equations.

A solution $(a; \cdot)$ of (14) satisfies the following bounds:

$$\begin{aligned}
 |j| &\leq 1 + \frac{C}{r} \\
 |j|^2 &\leq \frac{C}{r}(1 - |j|^2) + \frac{C^3}{r^3} \\
 |j^A| &\leq \frac{C^{\rho_r}}{r} \exp\left(-\frac{\rho_r}{C} \text{dist}(x; \pi^{-1}(0))\right) \quad ; \quad x \in X \\
 |j - j(x)| &\leq C \exp\left(-\frac{\rho_r}{C} \text{dist}(x; \pi^{-1}(0))\right) \quad ; \quad x \in X \\
 |jF_a| &\leq \frac{r}{4} (1 - |j|^2) + C \\
 |jF_a(x)| &\leq C r \exp\left(-\frac{1}{C} \frac{\rho_r}{r} \text{dist}(x; \pi^{-1}(0))\right) \quad ; \quad x \in X
 \end{aligned}
 \tag{16}$$

Remark 2.4 The constant C appearing above may change its value from line to line. It is important to point out that C only depends on the $Spin^c$ -structure W_E and the Riemannian metric g but **not** on the particular choice of the parameter r . This will be the case for all the numerous constants (all labeled C) appearing subsequently and we will henceforth tacitly adopt this misuse of notation.

3 The main part

3.1 Producing the approximate solution $(a; \cdot)$ from a pair $(a_0; \cdot_0), (a_1; \cdot_1)$

Let E_0 and E_1 be two complex line bundles over X . The aim of this section is to produce an approximate solution $(a; \cdot)$ of the Seiberg-Witten equations for the $Spin^c$ -structure $W_{E_0 \oplus E_1}$ from two solutions $(a_0; \cdot_0)$ and $(a_1; \cdot_1)$ for the $Spin^c$ -structures W_{E_0} and W_{E_1} respectively. Implicit to our discussion are the choices of two "base" connections B_0 and B_1 on E_0 and E_1 and the product connection $B_0 \oplus B_1$ they determine on $E_0 \oplus E_1$. As before, we will write $\cdot_i = \rho_r(\cdot_i, u_0; \cdot_i)$, $i = 0, 1$, and $\cdot = \rho_r(\cdot, u_0; \cdot)$. We define $(a; \cdot)$ as:

$$\begin{aligned}
 a &= a_0 + a_1 \\
 &= \cdot_0 \oplus \cdot_1 \\
 &= \cdot_0 \oplus \cdot_1 + \cdot_1 \oplus \cdot_0
 \end{aligned}
 \tag{17}$$

The first task at hand is to check how close (a_i) comes to solving the Seiberg-Witten equations. We begin by calculating D_a locally at a point $x \in X$. Choose an orthonormal frame fe_jg_j in a neighborhood of x and let fe^jg_j be its dual frame.

$$\begin{aligned}
 D_a(\) &= \rho_{\bar{r}} D_a(\ 0 \ 1 \ u_0 + \ 0 \ \rho_{\bar{r}} \ 1 + \ 1 \ 0) \\
 &= \rho_{\bar{r}} D_{a_0}(\ 0 \ u_0) \ 1 + \rho_{\bar{r}} \ 0 \ D_{a_1}(\ 1 \ u_0) + \\
 &\quad + \rho_{\bar{r}} e^j : r_{e_i}^a(\ 0 \ 1 + \ 1 \ 0) \\
 &= \rho_{\bar{r}} D_{a_0}(\ 0 \ u_0) \ 1 + \rho_{\bar{r}} \ 0 \ D_{a_1}(\ 1 \ u_0) + \\
 &\quad + \rho_{\bar{r}}(\ 0 \ e^j : (r_{e_i}^{a_1} \ 1) + \ 1 \ e^j : (r_{e_i}^{a_0} \ 0) + \\
 &\quad + (r_{e_i}^{a_0} \ 0) \ e^j : \ 1 + (r_{e_i}^{a_1} \ 1) \ e^j : \ 0) \\
 &= \rho_{\bar{r}} D_{a_0}(\ 0 \ u_0) \ 1 + \rho_{\bar{r}} \ 0 \ D_{a_1}(\ 1 \ u_0) + \\
 &\quad + \rho_{\bar{r}}(\ 0 \ D_{a_1} \ 1 + \ 1 \ D_{a_0} \ 0) + \\
 &\quad + \rho_{\bar{r}}((r_{e_i}^{a_0} \ 0) \ e^j : \ 1 + (r_{e_i}^{a_1} \ 1) \ e^j : \ 0) \\
 &= (D_{a_0} \ 0) \ 1 + \ 0 \ (D_{a_1} \ 1) + \\
 &\quad + \rho_{\bar{r}}((r_{e_i}^{a_0} \ 0) \ e^j : \ 1 + (r_{e_i}^{a_1} \ 1) \ e^j : \ 0) \\
 &= \rho_{\bar{r}}(r_{e_i}^{a_0} \ 0) \ e^j : \ 1 + \rho_{\bar{r}}(r_{e_i}^{a_1} \ 1) \ e^j : \ 0
 \end{aligned} \tag{18}$$

It is easy to see, using the bounds in (16), that the first term in (18) satisfies the following pointwise estimate :

$$\begin{aligned}
 rj(r_{e_i}^{a_0} \ 0) \ e^j : \ 1 &\leq \int_X \\
 &\leq C \exp \left(-\frac{\rho_{\bar{r}}}{C} \text{dist}(x; \ 0^{-1}(0)) \right) \exp \left(-\frac{\rho_{\bar{r}}}{C} \text{dist}(x; \ 1^{-1}(0)) \right)
 \end{aligned} \tag{19}$$

The second term in (18) satisfies the same bound. In order for the right hand side of (19) to pointwise converge to zero, it is sufficient and necessary that there exist some $r_0 > 1$ such that for all $r > r_0$, the distance from $0^{-1}(0)$ to $1^{-1}(0)$ be bounded from below by some r -independent $M > 0$. This condition, under the map π from (3), is the Seiberg-Witten equivalent of the condition that $\gamma_i = (A_i; \ \gamma_i)$ be disjoint curves. Thus, from now onward we will make the following assumption.

Assumption 3.1 *As above, let $E_0; E_1 \in H^2(X; \mathbb{Z})$ be two line bundles over X . Let $(\ \gamma_i; a_i)$, $i = 0; 1$, be two solutions to the Seiberg-Witten equations (14) for the $Spin^c$ -structures W_{E_i} with $\gamma_i = \rho_{\bar{r}}(\ \gamma_i \ u_0; \ 0)$ and $\gamma_i \in (E_i)$. We henceforth make the assumption that there exists an $r_0 > 1$ and $M > 0$ such that for all $r > r_0$ the inequality*

$$\text{dist}(\ 0^{-1}(0); \ 1^{-1}(0)) \geq M \tag{20}$$

holds.

We now proceed by looking at the second equation in (14):

$$\begin{aligned}
 F_a^{1,1} - \frac{i}{8}r(j^2 - 1 - j^2)! &= \\
 &= F_{a_0}^{1,1} + F_{a_1}^{1,1} - \frac{i}{8}r(j_0^2 - j_1^2 - 1 - j_0^2 - j_1^2 - j_0^2 - j_1^2 - \\
 &\quad - 2h_{0,1}!) \\
 &= F_{a_0}^{1,1} + F_{a_1}^{1,1} - \frac{i}{8}rj_1^2(j_0^2 - 1 - j_0^2)! - \frac{i}{8}rj_0^2(j_1^2 - 1 - j_1^2)! \\
 &\quad + \frac{i}{8}r(j_0^2 - 1)(j_1^2 - 1)! + \frac{i}{4}rh_{0,1}! \\
 &= \frac{i}{8}r(1 - j_1^2)(j_0^2 - 1 - j_0^2)! - \frac{i}{8}r(1 - j_0^2)(j_1^2 - 1 - j_1^2)! + \\
 &\quad + \frac{i}{8}r(j_0^2 - 1)(j_1^2 - 1)! + \frac{i}{4}rh_{0,1}!
 \end{aligned}$$

From this last equation, and again using (16), one easily deduces that:

$$\begin{aligned}
 jF_a^{1,1} - \frac{i}{8}r(j^2 - 1 - j^2)! & \\
 &= C \exp\left(-\frac{\rho_r}{C} \text{dist}(x; \sigma_0^{-1}(0))\right) \exp\left(-\frac{\rho_r}{C} \text{dist}(x; \sigma_1^{-1}(0))\right) + \frac{C}{\rho_r}
 \end{aligned} \tag{21}$$

Finally, we consider the third equation in (14):

$$\begin{aligned}
 F_a^{0,2} - \frac{i}{4}r &= F_{a_0}^{0,2} + F_{a_1}^{0,2} - \frac{i}{4}r(0 + 1) \\
 &= \frac{i}{4}r_{0,0} + \frac{i}{4}r_{1,1} - \frac{i}{4}rj_0^2 - \frac{i}{4}rj_1^2 \\
 &= \frac{i}{4}r(1 - j_1^2)_{0,0} + \frac{i}{4}r(1 - j_0^2)_{1,1}
 \end{aligned}$$

Once again using the bounds (16), we find from this last equation:

$$\begin{aligned}
 jF_a^{0,2} - \frac{i}{4}r & \\
 &= C \exp\left(-\frac{\rho_r}{C} \text{dist}(x; \sigma_0^{-1}(0))\right) \exp\left(-\frac{\rho_r}{C} \text{dist}(x; \sigma_1^{-1}(0))\right) + \frac{C}{\rho_r}
 \end{aligned} \tag{22}$$

To summarize, we have proved the following result.

Proposition 3.2 *Let $(a; \sigma)$ be defined as in (17) and assume that there exists an $r_0 > 0$ and $M > 0$ such that for all $r > r_0$, the distance $\text{dist}(\sigma_0^{-1}(0); \sigma_1^{-1}(0))$*

is bounded from below by M . Then for large enough r and any $x \in X$ the pointwise bound below holds:

$$j(D_a(\cdot); F_a^{1,1} - \frac{i}{8}r(j^2 - 1 - j^2)) ; F_a^{0,2} - \frac{i}{4}r(\cdot))j_x \leq \frac{C}{r} \quad (23)$$

3.2 Inverting the linearized operators of $(a_i; \cdot)$

This section serves as a digression. The main result of the section is theorem 3.6, an asymptotic (as $r \rightarrow \infty$) regularity statement for the linear operators $L_{(a_i; \cdot)}$.

We start with two easy auxiliary lemmas:

Lemma 3.3 *Let $L : V \rightarrow W$ be a surjective Fredholm operator between Hilbert spaces. Then there exists a $\delta > 0$ such that for every linear operator $\phi : V \rightarrow W$ with $\langle \phi(x), \phi(x) \rangle_W \leq \delta \langle x, x \rangle_V$, the operator $L + \phi$ is still surjective.*

Proof Since L is Fredholm, we can orthogonally decompose V as $V = \text{Ker}(L) \oplus \text{Im}(L)$. Let L_1 be the restriction of L to $\text{Im}(L)$. Then $L_1 : \text{Im}(L) \rightarrow W$ is an isomorphism with bounded inverse L_1^{-1} .

If the lemma were not true then we could find for all integers $n \geq 1$ an operator $\phi_n : V \rightarrow W$ with $\langle \phi_n(x), \phi_n(x) \rangle_W \leq \frac{1}{n} \langle x, x \rangle_V$ and with $\text{Coker}(L + \phi_n) \neq \emptyset$. Let $0 \neq y_n \in \text{Coker}(L + \phi_n)$ with $\langle y_n, y_n \rangle_W = 1$ and $x_n = L_1^{-1}(y_n)$. Notice that the sequence $\langle \phi_n(x_n), y_n \rangle$ is bounded by $\frac{1}{n} \langle x_n, x_n \rangle_V$. Since $y_n \in \text{Coker}(L + \phi_n)$, y_n is orthogonal to $\text{Im}(L + \phi_n)$. In particular,

$$\langle (L + \phi_n)x_n, y_n \rangle = 0$$

This immediately leads to a contradiction for large enough n since $\langle Lx_n, y_n \rangle = 1$ and $\langle \phi_n(x_n), y_n \rangle \leq \frac{1}{n}$. □

Lemma 3.4 *Let V and W be two finite rank vector bundles over X and $L_r : L^{1,2}(V) \rightarrow L^2(W)$ a smooth one-parameter family (indexed by $r \geq 1$) of elliptic, first order, differential operators of index zero. Assume further that there exists a $\delta > 0$ and $r_0 \geq 1$ such that for any zeroth order linear operator $\phi : L^{1,2}(V) \rightarrow L^2(W)$ with $\langle \phi(x), \phi(x) \rangle_2 \leq \delta \langle x, x \rangle_{1,2}$, the operator $L_r + \phi$ is onto. Then there exists a $r_1 \geq r_0$ and a $M > 0$ such that for all $r \geq r_1$ the inverses of the operators L_r are uniformly bounded by M , i.e. $\|L_r^{-1}y\|_{1,2} \leq M\|y\|_2$.*

Proof Notice that a universal upper bound on L_r^{-1} is equivalent to a universal lower bound on L_r . Suppose the lemma were not true: then there would be a sequence $r_n \rightarrow 1$ and $x_n \in L^{1,2}(V)$ with $\|x_n\|_{L^{1,2}} = 1$ and $\|L_{r_n} x_n\|_{L^2} < 1 - \epsilon$. Choose n large enough so that $1 - \epsilon > 0$ and define the operator $\psi : L^{1,2}(V) \rightarrow L^2(W)$ by $\psi(x) = -\epsilon x_n + L_{r_n}(x_n)$. For this ψ the assumption of the lemma is met, namely

$$\|\psi(x)\|_{L^2} \leq \frac{1}{r} \|x\|_{L^{1,2}} < \|x\|_{L^{1,2}}$$

Thus the operator $L_{r_n} + \psi$ should be onto and injective (since the index of $L_r + \psi$ is zero). But x_n is clearly a nonzero kernel element. This is a contradiction. \square

Recall that the set \mathcal{J} of almost-complex structures compatible with the symplectic form ω , contains a Baire subset \mathcal{J}_0 of generic almost-complex structures in the sense of Gromov-Witten theory (see [7]). Also, as in the introduction, let

$$\mathcal{M}_X^{SW} : \mathcal{M}_X^{SW}(W_E) \rightarrow \mathcal{M}_X^{Gr}(E) \tag{24}$$

be the map introduced in [4] which associates an embedded J -holomorphic curve to a Seiberg-Witten monopole.

Proposition 3.5 *Let J be chosen from \mathcal{J}_0 and let $(a; \cdot)$ be a solution of the Seiberg-Witten equations (14) such that $(a; \cdot)$ doesn't contain any multiply covered components. Then there exists a $\delta > 0$ and an $r_0 \geq 1$ such that for all linear operators $\psi : L^{1,2}(i^{-1} E \otimes W_0^+) \rightarrow L^2(i^{-1} E \otimes W_0^-)$ with norm $\|\psi(x)\|_{L^2} < \delta \|x\|_{L^{1,2}}$, the operator $L_{(a; \cdot)} + \psi$ is surjective.*

Before proceeding to the proof, notice that proposition 3.5 and lemma 3.4 immediately imply the following theorem, the main result of this section:

Theorem 3.6 *Choose $J \in \mathcal{J}_0$ and let $(a; \cdot)$ be a solution of the Seiberg-Witten equations for the $Spin^c$ -structure W_E with parameter r . Assume that $(a; \cdot)$ contains no multiply covered components. Then there exists a r -independent $M > 0$ and $r_0 \geq 1$ such that for all $r \geq r_0$:*

$$\|L_{(a; \cdot)}^{-1} x\|_{L^{1,2}} \leq M \|x\|_{L^2} \tag{25}$$

Proof of proposition 3.5 The proof is a bit technical and relies on the even more technical account from [5] on the connection between the deformation theory of the Seiberg-Witten equations on one hand and the Gromov-Witten equation on the other. The idea is however very simple: for large $r \geq 1$, a certain perturbation of the operator L (with the size of the perturbation

getting smaller with larger r) has no cokernel if a certain perturbation of the linearization of the generalized del-bar operator has no cokernel. The latter is ensured by the choice of a generic almost complex structure J from the Baire set \mathcal{J}_0 of almost complex structures compatible with $!$.

Before proceeding, the (interested) reader is advised to familiarize him/her-self with the notation from [5], in particular, sections 4 and 6 as the remainder of the proof heavily relies on it. For convenience we restate here the parts of lemma 4.11 and a slightly modified version of lemma 6.7 from [5] relevant to our situation.

Lemma 4.11 *The equation $Lq + q = g$ is solvable if and only if, for each k*

$$c_k W^k + \partial_0^k(w) + \partial_k(w) = x(g^k) + \partial_1^k(g)$$

Notice that the assignment of ∂_k to ∂ is linear i.e. for two operators ∂ and ∂' , we have $(\partial + \partial')_k = \partial_k + \partial'_k$.

Lemma 6.7⁰ *The equation $(L_{r(y)} + \partial')\rho = g$ has an $L^{1,2}$ solution ρ if and only if there exists $u = (u^1; \dots; u^k) \in L^{1,2}(N^{(k)})$ for which*

$$yU^k + \partial_0^k(u) + \partial'_k(u) = \partial_1^{-1}x(g^k) + \partial_1^k(g)$$

holds for each k .

The proof of lemma 6.7⁰ is almost identical to that of the original lemma 6.7 in [5]. The only difference is in *Step 2* where Taubes shows that one can write the equation $L_{r(y)}\rho = g$ in the form $L\rho + \partial\rho = g$ with L as in lemma 4.11 above and with ∂ an appropriate (bounded) correction term (see (6.30) in [5] for a precise definition). The difference here is that in our case one can write $(L_{r(y)} + \partial')\rho = g$ as $L\rho + \partial\rho = g$ (with L again as in lemma 4.11 of [5]) but with $\partial = \partial + \partial'$. Since ∂' is assumed bounded, lemma 4.11 applies to ∂ in the exact same way as it applied to the original ∂ and the proof of lemma 6.7 in [5] transfers verbatim to our case. Note also that the operators ∂_i^k occurring in lemmas 6.7 and 6.7⁰ are identical so in particular they continue to satisfy the bounds asserted by lemma 6.7 of [5].

According to lemma 3.3 there exists a $\delta > 0$ such that $\partial_y + \partial'$ is still surjective if $\|\partial'\| < \delta$. Choose r large enough so that $\|\partial\| < \delta = 2k$. On the other hand, since $\partial'_k(v) = (25; k'(\partial_0^k - 100; k^0 \underline{v}^{k^0}))$ we find that $\|\partial'_k\| < C\|\partial\|$. Thus choosing $\delta = 2C$ ensures that $L_{r(y)} + \partial'$ is surjective provided that $\|\partial\| < \delta$. This finishes the proof of proposition 3.5. □

3.3 The linearized operator at $(a; \cdot)$

In order to use the contraction mapping principle to deform the approximate solution $(a; \cdot)$ to an honest solution of the Seiberg-Witten equations, we need to know that $L = L_{(a; \cdot)}$ admits an inverse whose norm is bounded independently of r . We start by exploring when the equation

$$L = g \tag{26}$$

has a solution \cdot for a given g . Here:

$$g \in L^{1,2}(i^{-1}(E_0 \oplus E_1 \oplus W_0^+)) \text{ and } g \in L^2(i^{-1}(E_0 \oplus E_1 \oplus W_0^-))$$

The idea is to restrict equation (26) first to a neighborhood of $\cdot_0^{-1}(0)$. Over such a neighborhood the bundle E_1 is trivial and, under an isomorphism trivializing E_1 , the equation (26) becomes a zero-th order perturbation of the equation $L_0 \cdot_0 = g_0$ (with \cdot_0 and g_0 being appropriately defined in terms of \cdot and g). This allows one to take advantage of the results of theorem 3.6 about the inverse of $L_0 = L_{(a_0; \cdot_0)}$. Then one restricts (26) to a neighborhood of $\cdot_1^{-1}(0)$ where the bundle E_0 trivializes and once again uses theorem 3.6, this time for the inverse of $L_1 = L_{(a_1; \cdot_1)}$. Finally, one restricts to the complement of a neighborhood of $\cdot_0^{-1}(0) \cup \cdot_1^{-1}(0)$ where both E_0 and E_1 become trivial and L becomes close to S - the linearized operator of the unique solution $(A_0; \bar{r} u_0)$ for the anticanonical $Spir^c$ -structure W_0 .

To begin this process, choose regular neighborhoods V_i of $\cdot_i^{-1}(0)$, $i = 0, 1$ subject to the condition:

$$\text{dist}(V_0; V_1) \geq M \text{ for some } M > 0$$

The existence of such neighborhoods V_i follows from our main assumption (20). A priori, as one chooses larger values of r , it seems that the sets V_i may need to be chosen anew as well. However, it was shown in [4], section 5c, that in fact this is not necessary. An initial "smart" choice of V_i for large enough r ensures that for $r' > r$, the zero sets $\cdot_i^{-1}(0)$ continue to lie inside of V_i . Choose an open set U such that $X = V_0 \cup V_1 \cup U$ and such that

$$U \cap (\cdot_0^{-1}(0) \cup \cdot_1^{-1}(0)) = \emptyset$$

Arrange the choices of V_i and U further so that ∂V_i is an embedded 3-manifold of X and so that $U \setminus V_i$ contains a collar $\partial V_i \times I$. Here I is some segment $[0; d]$ and ∂V_i corresponds to $\partial V_i \times \text{fdg}$. For the sake of simplicity of notation, we shall make the assumption that for large values of r , the sets $\cdot_i^{-1}(0)$, $i = 0, 1$, are connected. The case of disconnected zero sets of the \cdot_i 's is treated much in the same way except for that in the following, one would have to choose a

bump function χ_i (see below) for each connected component. This complicates notation to a certain degree but doesn't lead to new phenomena.

Fix once and for all a bump function $\chi : [0; 1) \rightarrow [0; 1]$ which is 1 on $[0, 1]$ and 0 on $[2; 1)$. For $0 < \epsilon < d=1000$ define $\chi_i : X \rightarrow [0; 1]$ by:

$$\chi_i(x) = \begin{cases} < 1 & x \in V_i \cap @V_i \setminus I \\ \chi(t) & x = (y; t) \in @V_i \setminus I \\ 0 & x \notin V_i \end{cases} \tag{27}$$

Set $V_0^\theta = V_0 \sqcup U$ and $V_1^\theta = V_1 \sqcup U$. Define the isomorphisms $\alpha_0 : \mathbb{C} \setminus V_0^\theta \rightarrow E_1 \setminus V_0^\theta$ and $\alpha_1 : \mathbb{C} \setminus V_1^\theta \rightarrow E_0 \setminus V_1^\theta$ as $\alpha_0(\chi_i(x)) = \alpha_1(x)$ and $\alpha_1(\chi_i(x)) = \alpha_0(x)$. Also, for $i = 0, 1$ define the operators

$$M_i : L^{1,2}(i^{-1}(E_i \setminus W_0^+); V_i^\theta) \rightarrow L^2(i^{-1}(E_i \setminus W_0^-); V_i^\theta)$$

and

$$T : L^{1,2}(i^{-1}(W_0^+; U)) \rightarrow L^2(i^{-1}(W_0^-; U))$$

by demanding the diagrams

$$\begin{array}{ccc} L^{1,2}(i^{-1}(E_0 \setminus E_1 \setminus W_0^+); V_i^\theta) & \xrightarrow{i} & L^{1,2}(i^{-1}(W_i^+; V_i^\theta)) \\ \downarrow \chi_i & & \downarrow \chi_i M_i \\ L^2(i^{-1}(E_0 \setminus E_1 \setminus W_0^-); V_i^\theta) & \xrightarrow{i} & L^2(i^{-1}(W_i^-; V_i^\theta)) \end{array}$$

and

$$\begin{array}{ccc} L^{1,2}(i^{-1}(E_0 \setminus E_1 \setminus W_0^+); U) & \xrightarrow{i} & L^{1,2}(i^{-1}(W_0^+; U)) \\ \downarrow \chi_i & & \downarrow \chi_i T \\ L^2(i^{-1}(E_0 \setminus E_1 \setminus W_0^-); U) & \xrightarrow{i} & L^2(i^{-1}(W_0^-; U)) \end{array}$$

to be commutative diagrams (with $W_i = E_i \setminus W_0$).

We now start our search for a solution g of (26) in the form:

$$g = \alpha_0(\chi_0; g_0) + \alpha_1(\chi_1; g_1) + \alpha_0^{-1}((1 - \chi_0; g_0)(1 - \chi_1; g_1)) \tag{28}$$

Here $\chi_i \in L^{1,2}(i^{-1}(E_i \setminus W_0^+))$ and $\chi_i \in L^{1,2}(i^{-1}(W_0^+))$. Given a $g \in L^2(i^{-1}(E_0 \setminus E_1 \setminus W_0^-))$, define $g_i \in L^2(i^{-1}(E_i \setminus W_0^-))$ and $\tilde{g}_i \in L^2(i^{-1}(W_0^-))$ as

$$g_i = \chi_i^{-1}(\chi_i; g) \quad \text{and} \quad \tilde{g}_i = (\chi_i)^{-1}((1 - \chi_i; g)(1 - \chi_i; g)) \tag{29}$$

It is easy to check that g, g_i and \tilde{g}_i satisfy a relation similar to (28), namely:

$$g = \alpha_0(\chi_0; g_0) + \alpha_1(\chi_1; g_1) + \alpha_0^{-1}((1 - \chi_0; g_0)(1 - \chi_1; g_1)) \tag{30}$$

Putting the form (28) of ψ and the form (30) of g into equation (26), after a few simple manipulations, yields the equation

$$\begin{aligned} & \psi_0^*(M_0(\psi_0) - \psi_1^*P(d_{4,0}^* \psi_0) - g_0) + \\ & + \psi_1^*(M_1(\psi_1) - \psi_0^*P(d_{4,1}^* \psi_1) - g_1) + \\ & + \psi_0^*(1 - \psi_4^*)\psi_1^*(1 - \psi_4^*)T + \psi_1^{-1}P(d_{100,0}^* \psi_0) + \\ & + \psi_0^{-1}P(d_{100,1}^* \psi_1) - \\ & = 0 \end{aligned} \tag{31}$$

In the above, P denotes the principal symbol of L . This last equation suggests a splitting into three equations (each corresponding to one line in (31)):

$$\begin{aligned} M_0(\psi_0) - \psi_1^*P(d_{4,0}^* \psi_0) &= g_0 \\ M_1(\psi_1) - \psi_0^*P(d_{4,1}^* \psi_1) &= g_1 \\ T + \psi_1^{-1}P(d_{100,0}^* \psi_0) + \psi_0^{-1}P(d_{100,1}^* \psi_1) &= \end{aligned} \tag{32}$$

Equation (31) (and hence also equation (26)) can be recovered from (32) by multiplying the three equations by ψ_0^* , ψ_1^* and $\psi_0^*(1 - \psi_4^*)(1 - \psi_4^*)$ respectively and then adding them. Thus, given a g and with g_i and ψ_i defined by (29), solutions ψ_i and ψ_i^* of (32) lead to a solution ψ of (26) via (28). However, the problem with (32) is that the operators M_i and T are not defined over all of X . We remedy this in the next step.

Define new operators:

$$\begin{aligned} M_i^\theta &: L^{1,2}(j^{-1}(E_i - W_0^+)) \rightarrow L^2(j^{-1}(E_i - W_0^-)) \\ &\text{and} \\ T^\theta &: L^{1,2}(j^{-1}(W_0^+)) \rightarrow L^2(j^{-1}(W_0^-)) \end{aligned}$$

by

$$\begin{aligned} M_i^\theta &= \psi_{200,i}^*M_i + (1 - \psi_{200,i}^*)L_i \\ T^\theta &= (1 - \psi_{,0}^*)(1 - \psi_{,1}^*)T + (\psi_{,0}^* + \psi_{,1}^*)S \end{aligned} \tag{33}$$

Here $L_i = L_{(a_i, \psi_i)}$. Now replace the coupled equations (32) by the following system:

$$\begin{aligned} M_0^\theta(\psi_0) - \psi_1^*P(d_{4,0}^* \psi_0) &= g_0 \\ M_1^\theta(\psi_1) - \psi_0^*P(d_{4,1}^* \psi_1) &= g_1 \\ T^\theta + \psi_1^{-1}P(d_{100,0}^* \psi_0) + \psi_0^{-1}P(d_{100,1}^* \psi_1) &= \end{aligned} \tag{34}$$

The advantage of (34) over (32) is that the former is defined over all of X (notice that the support of $P(d_{100,0}^* \psi_0)$ lies in the domain of ψ_1^{-1} and the

support of $P(d_{100,1;1})$ lies in the domain of $(\cdot)^{-1}$. On the other hand, solutions of (34) give rise to solutions of (26) in the same way as solutions of (32) did because

$$\begin{aligned} d_{100,1;1} M_i^0 &= d_{100,1;1} M_i \quad i = 0; 1 \\ (1 - d_{4,0})(1 - d_{4,1})T^0 &= (1 - d_{4,0})(1 - d_{4,1})T \end{aligned}$$

Lemma 3.7 For every $\epsilon > 0$ there exists an $r > 1$ such that for $r > r$ the following hold:

$$\begin{aligned} \|j(M_i^0 - L_i)x_i\|_2 &\leq \|jx_i\|_2 \\ \|j(T^0 - S)y\|_2 &\leq \|jy\|_2 \end{aligned}$$

Here $x_i \in L^{1,2}(i^{-1} E_i, W_0^+)$ and $y \in L^{1,2}(i^{-1} W_0^+)$.

Proof The above Sobolev inequalities are proved by first calculating pointwise bounds for $j(M_i^0 - L_i)x_i$ and $j(T^0 - S)y$, $p \in X$. Notice first that $j(M_i^0 - L_i)x_i = 0$ if $p \notin V_i$ and $j(T^0 - S)y = 0$ if $p \notin U$. For $p \in V_i$ and for $q \in U$, a straightforward but somewhat tedious calculation shows that:

$$\begin{aligned} j(M_i^0 - L_i)x_i &\leq C \left(\rho_{r-1} |j - j^2| + \rho_{r-1} |j - j^2| + jr^{a_i} |j - j^2| \right) \\ j(T^0 - S)y &\leq C \left(\rho_{r-1} |j - j^2| + \rho_{r-1} |j - j^2| + \rho_{r-1} |j - j^2| \right) \\ &\quad + \rho_{r-1} |j - j^2| + jr^{a_0} |j - j^2| + jr^{a_1} |j - j^2| \end{aligned}$$

Squaring and then integrating both sides over X together with a reference to (16) gives the desired Sobolev inequalities. \square

The lemma suggests that the system (34) can be replaced by the system:

$$\begin{aligned} L_0 \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} - P(d_{4,0}) \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} &= g_0 \\ L_1 \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} - P(d_{4,1}) \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} &= g_1 \\ S \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} + P(d_{100,0}) \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} + P(d_{100,1}) \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} &= \end{aligned} \tag{35}$$

Lemmas 3.7 and 3.3 say that for $r > 0$, (34) has a solution $(\theta_0; \theta_1)$ if (35) has a solution $(\theta_0; \theta_1)$. It is this latter set of equations that we now proceed to solve.

Since S is onto, we can solve the third equation in (35), regarding θ_0 and θ_1 as parameters. Thus:

$$\begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} = S^{-1} \left(-P(d_{100,0}) \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} - P(d_{100,1}) \begin{pmatrix} \theta_0 \\ \theta_1 \end{pmatrix} \right) \tag{36}$$

Recall that the inverse of S satisfies the bound (12):

$$\|jS^{-1}y\|_2 \leq C \|jy\|_2 \quad \text{for } y \in L^2(i^{-1} W_0^-)$$

We will solve the first two equations in (35) simultaneously by first rewriting them in the form:

$$\begin{aligned} \theta_0 &= L_0^{-1}(g_0 + \theta_1 P(d_{4,0}; \theta_0)) \\ \theta_1 &= L_1^{-1}(g_1 + \theta_0 P(d_{4,1}; \theta_1)) \end{aligned} \tag{37}$$

To solve (37) is the same as to find a fixed point of the map $Y : L^2(i^{-1} W_{E_0}^+) \times L^2(i^{-1} W_{E_1}^+) \rightarrow L^2(i^{-1} W_{E_0}^+) \times L^2(i^{-1} W_{E_1}^+)$ given by

$$\begin{aligned} Y(\theta_0; \theta_1) &= \\ &= (L_0^{-1}(g_0 + \theta_1 P(d_{4,0}; \theta_1)); L_1^{-1}(g_1 + \theta_0 P(d_{4,1}; \theta_0))) \end{aligned} \tag{38}$$

with θ_0 given by (36). The existence and uniqueness of such a fixed point will be guaranteed by the fixed point theorem for Banach spaces if we can show that Y is a contraction mapping. To see this, let $x, y \in L^2(i^{-1} W_{E_0}^+) \times L^2(i^{-1} W_{E_1}^+)$ be two arbitrary sections. Using the first bound of (12) and the result of theorem 3.6 to bound the norms of L_i^{-1} , one finds:

$$\begin{aligned} \|Y(x) - Y(y)\|_2^2 &= \\ &= \|L_0^{-1}(g_0 + \theta_1 P(d_{4,0}; x)) - L_0^{-1}(g_0 + \theta_1 P(d_{4,0}; y))\|_2^2 \\ &+ \|L_1^{-1}(g_1 + \theta_0 P(d_{4,1}; x)) - L_1^{-1}(g_1 + \theta_0 P(d_{4,1}; y))\|_2^2 \\ &+ C_0 \|x - y\|_2^2 + C_1 \|x - y\|_2^2 \\ &+ C \|S^{-1}(L_1^{-1} P(d_{100,0}; y) - L_1^{-1} P(d_{100,0}; x)) + \\ &+ L_0^{-1} P(d_{100,1}; y) - L_0^{-1} P(d_{100,1}; x))\|_2^2 \end{aligned} \tag{39}$$

$$\frac{C}{r} \|x - y\|_2^2 \tag{40}$$

Choosing $r > 2C$, where C is the constant in the last line of (40), makes Y a contraction mapping. Thus we finally arrive at an L^2 solution $(\theta_0; \theta_1)$. It is in fact an $L^{1,2}$ solution because of (37). This, together with equation (36) provides a solution $(\theta_0; \theta_1; \theta)$ of (35). As explained above, this gives rise to a solution $(\theta_0; \theta_1)$ of (34) and thus provides a solution $\theta \in L^{1,2}(i^{-1} W_0^+)$ of (26). In particular, we have proved half of the following theorem.

Theorem 3.8 *Let $(a; \cdot)$ be constructed from $(a_i; \cdot)$ as in (17). Suppose that the $(a_i; \cdot)$ meet assumption 3.1, that $(a_i; \cdot)$ contains no multiply covered tori and that J has been chosen from the Baire set \mathcal{J}_0 of compatible almost complex structures. Then $L_{(a; \cdot)} : L^{1,2}(i^{-1} E_0 \oplus E_1 \oplus W_0^+) \rightarrow L^2(i^{-1} E_0 \oplus E_1 \oplus W_0^-)$ is invertible with bounded inverse $\|L_{(a; \cdot)}^{-1}\| \leq C \|y\|_2$ for all sufficiently large r . Here C is independent of r .*

Proof It remains to prove the inequality $\|jL_{(a; \cdot)}^{-1}\|_{YJ_{1,2}} \leq C\|jYj\|_2$. Each of the two lines of (37), together with the bound (25) on L_i^{-1} , yields:

$$\|j\|_{YJ_{1,2}} \leq C(\|jg_0j\|_2 + \|j\|_{(0; 1)}\|j\|_2) \tag{41}$$

A bound for the second term on the right-hand side of (41) comes from (36) and the L^2 bound in (12):

$$\|j\|_{(0; 1)}\|j\|_2 \leq \frac{C}{r}(\|j\|_2 + \|j\|_{0j_2} + \|j\|_{1j_2}) \tag{42}$$

Adding the two inequalities (41) for $i = 0, 1$ and using (42) gives:

$$(1 - \frac{C}{r}) (\|j\|_{0j_2} + \|j\|_{1j_2}) \leq C(\|jg_0j\|_2 + \|jg_1j\|_2 + \frac{1}{r}\|j\|_2) \tag{43}$$

For large enough r , this last inequality gives a bound on the $L^{1,2}$ norm of $(\cdot; \cdot)$ in terms of an r -independent multiple of the L^2 norm of $(g_0; g_1)$. With this established, the missing piece, namely the $L^{1,2}$ bound of \cdot , comes from (36) and the $L^{1,2}$ bound in (12):

$$\|j\|_{YJ_{1,2}} \leq C(\|j\|_2 + \|j\|_{0j_2} + \|j\|_{1j_2}) \leq C(\|j\|_2 + \|jg_0j\|_2 + \|jg_1j\|_2) \tag{44}$$

It remains to relate the now established bound on $(\cdot; \cdot)$ to a bound for $(\cdot; \cdot)$. To begin doing that, write the systems (35) and (34) schematically as:

$$F(\cdot; \cdot) = (g_0; g_1) \quad \text{and} \quad G(\cdot; \cdot) = (g_0; g_1)$$

Lemma 3.7 implies that for any $\epsilon > 0$ there exists a $r_\epsilon > 1$ such that for all $r > r_\epsilon$ the inequality $\|j(F - G)\|_2 \leq \epsilon\|j\|_2$ holds. The established surjectivity of F guarantees (by means of lemma 3.3) that G is also surjective. The proof of theorem 3.11 thus far, also shows that $\|jF^{-1}j\| \leq C$ where C is r -independent. Now the standard inequality

$$\|jG^{-1}j\| \leq \|jF^{-1}j\| + \|jG^{-1} - F^{-1}j\| \leq \|jF^{-1}j\| + \|jF^{-1}j\| + \|jG^{-1}j\| \|jG - Fj\|$$

implies the r -independent bound for $\|jG^{-1}j\|$:

$$\|jG^{-1}j\| \leq \frac{\|jF^{-1}j\|}{1 - \|jF^{-1}j\| \|jG - Fj\|} \leq \frac{C}{1 - C\epsilon}$$

This last inequality provides $L^{1,2}$ bounds on $(\cdot; \cdot)$ and \cdot in terms of the L^2 norms of $(g_0; g_1)$ and \cdot which in turn imply an r -independent $L^{1,2}$ bound on $\cdot = L^{-1}g$ in terms of the L^2 norm of g through (28) and (29). This finishes the proof of theorem 3.8. \square

3.4 Deforming $(a; \cdot)$ to an honest solution

The goal of this section is to show that the approximate solution $(a; \cdot)$ can be made into an honest solution of the Seiberg-Witten equations by a deformation whose size goes to zero as r goes to infinity.

To set the stage, let $SW : L^{1,2}(i^{-1} W_E^+) \times L^2(i^{-1} i^{-2,+} W_E^-)$ denote the Seiberg-Witten operator

$$SW(b; \cdot) = (d b; F_b^+ - F_{A_0}^+ - q(\cdot; \cdot) + \frac{ir}{8} \cdot; D_b)$$

We will search for a zero of SW of the form $(a; \cdot) + (a^\theta; \cdot^\theta)$ with $(a^\theta; \cdot^\theta) \in B(\cdot)$. Here $B(\cdot)$ is the closed ball in $L^{1,2}(i^{-1} W_E^+)$ centered at zero and with radius $\epsilon > 0$ which we will choose later but which should be thought of as being small. The equation $SW((a; \cdot) + (a^\theta; \cdot^\theta)) = 0$ can be written as:

$$0 = SW(a; \cdot) + L_{(a; \cdot)}(a^\theta; \cdot^\theta) + Q(a^\theta; \cdot^\theta) \tag{45}$$

Here $Q : L^{1,2}(i^{-1} W_E^+) \times L^2(i^{-1} i^{-2,+} W_E^-)$ is the quadratic map given by:

$$Q(b; \cdot) = (b; \cdot) \cdot \frac{i}{8} (j_0^2 - j_2^2) + \frac{i}{4} (\cdot_0^2 + \cdot_2^2) \tag{46}$$

Lemma 3.9 For $x, y \in L^{1,2}(i^{-1} W_E^+)$, the map Q satisfies the inequality:

$$|Q(x) - Q(y)| \leq C(|x|_{1,2} + |y|_{1,2})|x - y|_{1,2} \tag{47}$$

Proof This is a standard inequality for quadratic maps and it can be explicitly checked using the definition of Q and the multiplication theorem for Sobolev spaces. We give the calculation for the first component of the right hand side of (46). Let $x = (b; \cdot)$ and $y = (c; \cdot')$, then we have:

$$\begin{aligned} |Q(b; \cdot) - Q(c; \cdot')| &= |Q(b; \cdot) - Q(c; \cdot) + Q(c; \cdot) - Q(c; \cdot')| \leq |Q(b; \cdot) - Q(c; \cdot)| + |Q(c; \cdot) - Q(c; \cdot')| \\ &\leq C|b - c|_{1,2} + C|c|_{1,2} |c - c'|_{1,2} \\ &\leq C(|b; \cdot|_{1,2} + |c; \cdot|_{1,2})|c - c'|_{1,2} \end{aligned}$$

The other components are checked similarly. □

Solving equation (45) for $(a^\theta; \cdot^\theta) \in L^{1,2}(i^{-1} W_E^+)$ is equivalent to finding a fixed point of the map $Y : B(\cdot) \rightarrow B(\cdot)$ given by:

$$Y(b; \cdot) = -L_{(a; \cdot)}^{-1}(SW(a; \cdot) + Q(b; \cdot)) \tag{48}$$

In order for the image of Y to lie in $B(\epsilon)$ we need to choose r large enough and ϵ small enough. To make this precise, let $(b; \gamma) \in B(\epsilon)$. Using the bounds in (23) we find that

$$\|SW(a; \gamma)\|_{1,2} \leq \frac{C}{r}$$

and so together with the results of theorem 3.8 and lemma 3.9 we get:

$$\|Y(b; \gamma)\|_{1,2} \leq \frac{C}{r} + C\epsilon^2$$

Choosing $\epsilon < 1/2C$ and $r > 4C^2\epsilon^2$ ensures that Y is well defined.

Lemma 3.10 *The map $Y : B(\epsilon) \rightarrow B(\epsilon)$ as defined by (48) is a contraction mapping for r large enough and ϵ small enough.*

Proof Let $x, y \in B(\epsilon)$, then using (47) we find:

$$\|Y(x) - Y(y)\|_{1,2} \leq C\|Q(x) - Q(y)\|_{1,2} \leq C\|x - y\|_{1,2} \quad (49)$$

Choosing $\epsilon < 1/2C$ makes $C\|x - y\|_{1,2} < 2\epsilon$ less than ϵ . □

We summarize in the following:

Theorem 3.11 *Let $(a; \gamma)$ be constructed from $(a_i; \gamma_i)$ as in (17). Suppose that the $(a_i; \gamma_i)$ meet assumption 3.1, that $(a; \gamma)$ contains no multiply covered tori and that J has been chosen from the Baire set \mathcal{J}_0 of compatible almost complex structures.*

Then there exists a $\epsilon_0 > 0$ such that for any $0 < \epsilon < \epsilon_0$ there exists an $r \geq 1$ such that for every $r \geq r$ there exists a unique solution $(a; \gamma) + (a^\theta; \gamma^\theta)$ of the Seiberg-Witten equations (with perturbation parameter r) with $(a^\theta; \gamma^\theta) \in L^{1,2}(i^{-1}W_E^+)$ satisfying the bound $\|(a^\theta; \gamma^\theta)\|_{1,2} \leq \epsilon$.

Remark 3.12 It is not known if theorem 3.11 holds under the relaxed hypothesis allowing $(a; \gamma)$ to contain multiply covered tori. The difficulty in dealing with this case stems from the fact that the operators $L_{(a; \gamma)}$ may no longer have trivial cokernel.

4 Comparison with product formulas

Before proceeding further, we would like to take a moment to point out the similarities and differences between our construction of $(A; \nu)$ from $(A_i; \nu_i)$ on one hand and product formulas for the Seiberg-Witten invariants on manifolds that are fiber sums of simpler manifolds. We begin by briefly (and with few details) recalling the scenario of the latter.

Let $X_i, i = 0; 1$ be two compact smooth 4-manifolds and $\nu_i \in \pi_1 X_i$ embedded surfaces of the same genus and with $\nu_0 \nu_1 = -\nu_1 \nu_0$. In this setup one can construct the fiber sum

$$X = X_0 \#_{\nu_i} X_1$$

by cutting out tubular neighborhoods $N(\nu_i)$ in X_i and gluing the manifolds $X_i^0 = \overline{X_i \setminus N(\nu_i)}$ along their diffeomorphic boundaries.

Under certain conditions one can calculate some of the Seiberg-Witten invariants of X in terms of the Seiberg-Witten invariants of the building blocks X_i (see e.g. [3]). One accomplishes this by showing that from solutions $(B_i; \nu_i), i = 0; 1$ on X_i one can construct a solution $(B; \nu)$ on X (this isn't possible for any pair of solutions $(B_i; \nu_i)$ but the details are not relevant to the present discussion). This is done by inserting a "neck" of length $r \gg 1$ between the X_i^0 so as to identify X with

$$X = X_0^0 \cup ([0; r] \times Y) \cup X_1^0$$

with $Y = @N(\nu_0) = @N(\nu_1)$. A partition of unity $f'_{0; 1} g$ is chosen for each value of $r \gg 1$ subject to the conditions:

$$\begin{aligned} f'_i &= 1 \text{ on } X_i^0 \\ f'_i &= 0 \text{ outside of } X_i^0 \cup [0; r] \times Y \\ f'_0 f'_1 &= \frac{C}{r} \text{ on } [0; r] \times Y \end{aligned}$$

An approximation B^0 of B is then defined to be $B^0 = f'_0 B_0 + f'_1 B_1$ (similarly for B^0 , a first approximation for B). The measure of the failure of $(B^0; \nu^0)$ to solve the Seiberg-Witten equations can be made as small as desired by making r large. The honest solution $(B; \nu)$ is then sought in the form $(B^0; \nu^0) + (b; \eta)$ with $(b; \eta)$ small. The correction term $(b; \eta)$ is found as a fixed point of the map

$$(b; \eta) \mapsto Z(b; \eta) = -L_{(B^0, \nu^0)}^{-1} (Q(b; \eta) + \text{err})$$

Here "err" is the size of $SW(B^0; \nu^0)$ and L and Q are as in the previous section. Choosing r large enough and $\|b; \eta\|$ small enough makes Z a contraction

mapping and so the familiar fixed point theorem for Banach spaces guarantees the existence of a unique fixed point.

In the case of fiber sums there are product formulas that allow one to calculate the Seiberg-Witten invariants of X in terms of the invariants of the manifolds X_j . The formulas typically have the form:

$$SW_X(W_E) = \sum_{E_0 + E_1 = E} SW_{X_0}(W_{E_0}) SW_{X_1}(W_{E_1}) \quad (50)$$

Due to the similarity of our construction of grafting monopoles to the one used to construct $(B; \cdot)$ from $(B_j; \cdot_j)$, it is natural to ask if such or similar formulas exist for the present case, that is, can one calculate $SW_X(W_{E_0 \cup E_1})$ in terms of $SW_X(W_{E_0})$ and $SW_X(W_{E_1})$? The author doesn't know the answer. However, if they do exist, they can't be expected to be as simple as (50). The reason for this can be understood by trying to take the analogy between our setup and that for fiber sums further.

In the case of fiber sums, once one has established that the two solutions $(B_j; \cdot_j)$ on X_j can be used to construct a solution $(B; \cdot)$ on X , one needs to establish a converse of sorts. That is, one needs to show that every solution $(B; \cdot)$ on X is of that form. It is at this point where the analogy between the two situations breaks down. It is conceivable in our setup, that there will be solutions for the $Spin^c$ -structure $(E_0 \cup E_1) \rightarrow W_0^+$ that can not be obtained as products of solutions for the $Spin^c$ -structures $E_j \rightarrow W_0^+$. Worse even, there might be monopoles that can not be obtained as products of solutions for any $Spin^c$ -structures $F_j \rightarrow W_0^+$ with the choice of F_j , $j = 0, 1$ such that $E = F_0 \cup F_1$ and $F_j \not\equiv 0$. Those are the monopoles where $\pi^{-1}(0)$ is connected. Thus if a product formula for our situation exists, it must in addition to a term similar to the right hand side of (50) also contain terms which count these "undecomposable" solutions. But then again, they might not exist.

The next section describes which solutions of the Seiberg-Witten equations for the $Spin^c$ -structure $(E_0 \cup E_1) \rightarrow W_0^+$ are obtained as products of solutions for the $Spin^c$ -structures $E_j \rightarrow W_0^+$, $E = E_0 \cup E_1$.

5 The image of the multiplication map

This section describes a partial converse to theorem 3.11. Recall that

$$\mu : M_X^{SW}(W_E) \rightarrow M_X^{Gr}(E)$$

is the map assigning a J -holomorphic curve to a Seiberg-Witten monopole.

Theorem 5.1 *Let $E = E_0 \cup E_1$ and let $(A; \gamma)$ be a solution of the Seiberg-Witten equations in the $Spin^c$ -structure W_E with perturbation term $\eta = F_{A_0}^+ - ir! = 8$ and with $\eta = P_{\bar{r}}(u_0; \gamma)$. Assume further that J has been chosen from the Baire set \mathcal{J}_0 and that $(A; \gamma)$ contains no multiply covered components. If there exists an r_0 such that for all $r > r_0$, $\eta^{-1}(0)$ splits into a disjoint union $\eta^{-1}(0) = \cup_{i=0}^1 \tau_i$ with $[\tau_i] = P.D.(E_i)$ then $(A; \gamma)$ lies in the image of the multiplication map*

$$M_X^{SW}(E_0) \times M_X^{SW}(E_1) \rightarrow M_X^{SW}(E_0 \cup E_1)$$

The proof of theorem 5.1 is divided into 3 sections. In section 5.1 we give the definition of $(A_i^{\theta}; \gamma_i^{\theta})$ - first approximations of Seiberg-Witten monopoles $(A_i; \gamma_i)$ for the $Spin^c$ -structure W_{E_i} which when multiplied give the monopole $(A; \gamma)$ from theorem 5.1. Section 5.2 shows that for large values of r , $(A_i^{\theta}; \gamma_i^{\theta})$ come close to solving the Seiberg-Witten equations. In the final section 5.3 we show that $L_{(A_i^{\theta}; \gamma_i^{\theta})}$ is surjective with inverse bounded independently of r . The contraction mapping principle is then used to deform the approximate solutions $(A_i^{\theta}; \gamma_i^{\theta})$ to honest solutions $(A_i; \gamma_i)$. Section 5.3 also explains why $(A_0; \gamma_0) \cup (A_1; \gamma_1) = (A; \gamma)$.

We tacitly carry the assumptions of the theorem until the end of the section.

5.1 Definition $(A_i^{\theta}; \gamma_i^{\theta})$

The basic idea behind the definition of $(A_i^{\theta}; \gamma_i^{\theta})$ is again that of grafting existing solutions. For example, one would like γ_0^{θ} to be defined as the restriction of γ_0 to a neighborhood of τ_0 (under an appropriate bundle isomorphism trivializing E_1 over that neighborhood) and to be the restriction of $P_{\bar{r}}u_0$ outside that neighborhood. This is essentially how the construction goes even though a bit more care is required, especially in splitting the connection A into A_0^{θ} and A_1^{θ} .

To begin with, choose regular neighborhoods V_0 and V_1 of τ_0 and τ_1 . Once r is large enough, these choices don't need to be readjusted for larger values of r . Choose, as in section 3.3, an open set U such that:

$$X = V_0 \cup U \cup V_1 \\ U \cap \tau_i = \emptyset$$

Also, just as in section 3.3, arrange the choices so that $U \cap V_i$ contains a collar $@V_i \times [0; d]$ (with $@V_i$ corresponding to $@V_i$ in fdg) and choose $\epsilon > 0$ smaller than $d=1000$. Assume that the curves τ_i are connected, the general case goes through with little difficulty but with a bit more complexity of notation.

Over $U \sqcup V_1$, choose a section $\sigma_0 \in \Omega^0(E_0; U \sqcup V_1)$ with $\int \sigma_0 = 1$. Choose a connection B_0 on E_0 with respect to which σ_0 is covariantly constant over $U \sqcup V_1$, i.e.

$$B_0(\sigma_0(x)) = 0 \quad \forall x \in U \sqcup V_1 \tag{51}$$

Notice that such a connection is automatically flat over $U \sqcup V_1$. Choose a connection B_1 on E_1 such that $B_0 \oplus B_1 = A$ over X . Now define $\tilde{\sigma}_1 \in \Omega^0(E_1; U \sqcup V_1)$ and $\tilde{\sigma}_1 \in \Omega^0(E_1; K^{-1}; U \sqcup V_1)$ by:

$$\tilde{\sigma}_1 = \sigma_0 \oplus \tilde{\sigma}_1^0 \tag{52}$$

$$\tilde{\sigma}_1 = \sigma_0 \oplus \tilde{\sigma}_1^0 \tag{53}$$

Proceed similarly over V_0 . However, since some of the data is now already defined, more caution is required. Choose a section $\sigma_1 \in \Omega^0(E_1; V_0)$ with:

$$\sigma_1 = \tilde{\sigma}_1^0 \quad \text{on } (U \setminus V_0) \cap (@V_0 \text{ [0;4 i)}) \tag{54}$$

$$\int \sigma_1 = 1 \quad \text{on } (V_0 \cap U) \cap (@V_0 \text{ [0;2 i)})$$

We continue by defining $\tilde{\sigma}_0^0$ and $\tilde{\sigma}_0^0$ over V_0 by:

$$\tilde{\sigma}_0^0 = \sigma_0 \oplus \tilde{\sigma}_0^0 \tag{55}$$

$$\tilde{\sigma}_0^0 = \tilde{\sigma}_0^0 \oplus \sigma_0 \tag{56}$$

Choose one forms a_0 and a_1 such that over V_0 the following two relations hold:

$$(B_1 + i a_1) \sigma_1 = 0 \tag{57}$$

$$(B_0 + i a_0) \oplus (B_1 + i a_1) = A \tag{58}$$

With these preliminaries in place, we are now ready to define $(A_i^0; \tilde{\sigma}_i^0)$:

$$\begin{aligned} \tilde{\sigma}_0^0 &= \int \sigma_0 \tilde{\sigma}_0^0 + (1 - \int \sigma_0) \sigma_0 & \tilde{\sigma}_0^0 &= \int \sigma_0 \tilde{\sigma}_0^0 \\ \tilde{\sigma}_1^0 &= (1 - \int \sigma_0) \tilde{\sigma}_1^0 + \int \sigma_0 \sigma_1 & \tilde{\sigma}_1^0 &= (1 - \int \sigma_0) \tilde{\sigma}_1^0 \\ A_0^0 &= B_0 + i \int \sigma_0 a_0 & A_1^0 &= B_1 + i \int \sigma_0 a_1 \end{aligned} \tag{59}$$

Lemma 5.2 *The $(A_i^0; \tilde{\sigma}_i^0)$ defined above, satisfy the following properties:*

- (a) $A_0^0 \oplus A_1^0 = A$ on all of X .
- (b) $\tilde{\sigma}_0^0 = \sigma_0$ on $(U \setminus V_0) \cap (@V_0 \text{ [0;4 i)})$.
- (c) $F_{B_0} = 0$ on $U \sqcup V_1$ and $F_{B_1 + i a_1} = 0$ on V_0 .
- (d) On $(U \setminus V_0) \cap (@V_0 \text{ [0;4 i)})$, $\int \sigma_0 \tilde{\sigma}_0^0$ and $\int \sigma_0 \tilde{\sigma}_1^0$ converge exponentially fast to zero as $r \rightarrow 1^-$.

Proof (a) This is trivially true everywhere except possibly on the support of $d_{4,0}$ which is contained in $U \setminus V_0$. However, on $U \setminus V_0$ we have $A = B_0 - B_1$ and $A = (B_0 + ia_0) - (B_1 + ia_1)$ and thus $a_0 + a_1 = 0$. In particular, $A_0^q - A_1^q = B_0 - B_1 + i_{4,0}(a_0 + a_1) = B_0 - B_1 = A$.

(b) Notice that on $(U \setminus V_0) \cap @V_0 \cap [0; 4; i)$, $\psi_1 = \tilde{\psi}_1^q$. Thus, $\psi_0 = \psi_1$ and $\psi_0 = \tilde{\psi}_0^q = \psi_1$ imply that $\psi_0 = \tilde{\psi}_0^q$. The claim now follows from the definition of ψ_0 .

(c) Follows from the fact that both connection annihilate nowhere vanishing sections on the said regions.

(d) On $(U \setminus V_0) \cap @V_0 \cap [0; 4; i)$ we have $\psi_0 = \psi_1$ and $r^A = r^{B_0 + ia_0} - r^{B_1 + ia_1}$. Also, recall that $r^{B_0} \psi_0 = 0$ and $r^{B_1 + ia_1} \psi_1 = 0$. Thus:

$$r^a \psi = (r^{B_0 + ia_0} - r^{B_1 + ia_1})(\psi_0 - \psi_1) = ia_0 \psi_0 - \psi_1$$

This equation yields:

$$j a_0 j = \frac{j r^A j}{j j} \tag{60}$$

The claim follows now for a_0 by evoking the bounds (16). The same result holds for a_1 by the proof of part (a) where it is shown that $a_0 + a_1 = 0$ on $U \setminus V_0$. The statement for da_i follows from part (c), the equation $F_A = F_{B_0 + ia_0} + F_{B_1 + ia_1}$ and the bounds (16) for $jF_A j$. \square

5.2 Pointwise bounds on $SW(A_i^q; \psi_i)$

Proposition 5.3 *Let $(A_i^q; \psi_i)$ be defined as above, then there exists a constant C and an $r_0 > 1$ such that for all $r > r_0$ the inequality*

$$jSW(A_i^q; \psi_i)j_X \leq \frac{C}{r}$$

holds for all $X \in \mathcal{X}$.

Proof We calculate the size of the contribution of each of the three Seiberg-Witten equations separately. The only nontrivial part of the calculation is in the region of X which contains the support of $d_{4,0}$ i.e. in $@V_0 \cap [4; 8]$. We will tacitly use the results of lemma 5.2 in the calculations below.

a) The Dirac equation

To begin with, we calculate the expression $D_A((\gamma_0 \ u_0 + \tilde{\gamma}_0) \ 1)$ in two different ways. On one hand we have:

$$D_A((\gamma_0 \ u_0 + \tilde{\gamma}_0) \ 1) = D_A(\ + \ 4,0) = (1 - \ 4,0)D_A + d \ 4,0:$$

On the other hand we get:

$$\begin{aligned} D_A((\gamma_0 \ u_0 + \tilde{\gamma}_0) \ 1) &= & (61) \\ &= \ 1 \ D_{A_0^0}(\gamma_0 \ u_0 + \tilde{\gamma}_0) + e^i:(\gamma_0 \ u_0 + \tilde{\gamma}_0) \ A_1^0(1) \\ &= \ 1 \ D_{A_0^0}(\gamma_0 \ u_0 + \tilde{\gamma}_0) + i(\ 4,0 - 1)a_{1 \ 1} \end{aligned}$$

Equating the results of the two calculations we obtain:

$$\begin{aligned} j \ j \ j D_{A_0^0}(\gamma_0 \ u_0 + \tilde{\gamma}_0)j &= j \ 1 \ D_{A_0^0}(\gamma_0 \ u_0 + \tilde{\gamma}_0)j \\ C(j \ j j a_{1j} + j D_A \ j + j \ j) &= \frac{C}{r} \end{aligned}$$

Since over $@V_0 \ [4 ; 8]$, $j \ j ! \ 1$ exponentially fast as $r ! \ 1$ we obtain that:

$$j D_{A_0^0}(\gamma_0 \ u_0 + \tilde{\gamma}_0)j = \frac{C}{r} \tag{62}$$

b) The (1;1)-component of the curvature equation

Again, we only calculate for $x \geq @V_0 \ [4 ; 8]$:

$$\begin{aligned} F_{A_0^0}^{(1;1)} - F_{A_0}^{(1;1)} - \frac{ir}{8} \ j \ \gamma_0 j^2 - 1 - j \ \tilde{\gamma}_0 j^2 \ ! &= \ 4,0 \ (da_0)^{(1;1)} + \frac{ir}{8} j \ \tilde{\gamma}_0 j^2 ! \\ &= \ 4,0 \ (da_0)^{(1;1)} + \frac{ir}{8j \ j^2} j \ 4,0 j^2 j \ j^2 ! \end{aligned}$$

Both terms in the last line converge in norm exponentially fast to zero on $@V_0 \ [4 ; 8]$ as $r ! \ 1$.

c) The (0;2)-component of the curvature equation

Similar to the calculation for the (1;1)-component of the curvature equation on $@V_0 \ [4 ; 8]$, we have for the (0;2)-component of the same equation:

$$F_{A_0^0}^{(0;2)} - F_{A_0}^{(0;2)} - \frac{ir}{4} \ \gamma_0 \ \tilde{\gamma}_0 = \ 4,0 \ (da_0)^{(0;2)} - \frac{ir}{4j \ j^2} \ 4,0 \ -$$

Once again, both terms on the right-hand side of the above equation converge in norm exponentially fast to zero as r converges to infinity. The proofs for the case of $(A_1^0; \ 0_1)$ are similar and are left to the reader. □

5.3 Surjectivity of $L_{(A_i^0; \vartheta)}$ and deforming $(A_i^0; \vartheta)$ to an exact solution

The strategy employed here is very similar to the one used in section 3.4 and we only spell out part of the details. We start by showing that $L_{(A_0^0; \vartheta)}$ is surjective, the case $L_{(A_1^0; \vartheta)}$ is identical.

We begin by asking ourselves when the equation

$$L_{(A_0^0; \vartheta)} \varphi = g_0 \tag{63}$$

has a solution $\varphi \in L^{1,2}(i^{-1} W_{E_0}^+)$ for a given $g_0 \in L^2(i^{-1} W_{E_0}^-)$. Define the analogues of the isomorphisms φ_i from section 3.3 to be:

$$\begin{aligned} \varphi_0 : \mathbb{C} \otimes (U \oplus V_1) &\rightarrow (E_0; U \oplus V_1) \text{ given by } \varphi_0(\varphi; \chi) = \varphi_0(\varphi) \text{ and} \\ \varphi_1 : \mathbb{C} \otimes V_0 &\rightarrow (E_1; V_0) \text{ given by } \varphi_1(\varphi; \chi) = \varphi_1(\varphi) \end{aligned}$$

Let $\varphi \in L^2(i^{-1} W_0^-; U \oplus V_1)$ be determined by the equation $\varphi_0 \varphi = g_0$ on $U \oplus V_1$ and $\varphi \in L^2(i^{-1} W_E^-; V_0)$ be given by the equation $\varphi_1^{-1}(\varphi) = (1 - \varphi_0) g_0$ on V_0 . Thus we can write g_0 as:

$$g_0 = \varphi_0 \varphi + (1 - \varphi_0) \varphi_1^{-1}(\varphi) \tag{64}$$

This last form suggests that, in order to split equation (63) into two components involving $L_{(A; \vartheta)}$ and S , one should search for φ in the form

$$\varphi = \varphi_0 \varphi + (1 - \varphi_0) \varphi_1^{-1}(\varphi) \tag{65}$$

with $\varphi \in L^{1,2}(i^{-1} W_0^+; U \oplus V_1)$ and $\varphi \in L^{1,2}(i^{-1} W_E^+; V_0)$. Using relations (64) and (65) in (63) one obtains the analogue of equation (31):

$$\begin{aligned} \varphi_0 \varphi_0 (T(\varphi) - \varphi_0^{-1} \varphi_1^{-1} P(d_{A_0^0; \vartheta}(\varphi) - \varphi)) + \\ + (1 - \varphi_0) \varphi_1^{-1} (M(\varphi) + \varphi_1 \varphi_0 P(d_{A_0^0; \vartheta}(\varphi) - \varphi)) = 0 \end{aligned} \tag{66}$$

The operators T^ϑ and M^ϑ are defined over $U \oplus V_1$ and V_0 respectively, through the relations:

$$\begin{aligned} L_{(A_0^0; \vartheta)} \varphi_0 &= \varphi_0 T \\ L_{(A_0^0; \vartheta)} \varphi_1^{-1} &= \varphi_1^{-1} M \end{aligned}$$

We use these operators, defined only over portions of X , to define the operators T^ϑ and M^ϑ defined on all of X by:

$$\begin{aligned} T^\vartheta &= (1 - \varphi_0) T + \varphi_0 S \\ M^\vartheta &= \varphi_0 M + (1 - \varphi_0) L_{(A; \vartheta)} \end{aligned}$$

Split equation (66) into the following two equations:

$$\begin{aligned} T^0(\cdot) - \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}^{-1} P(d_{4,0}; \cdot) &= \\ M^0(\cdot) + \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} P(d_{100,0}; \cdot) &= \& \end{aligned} \tag{67}$$

It is easy to see that solutions to the system of equations (67) provide solutions to (66) by multiplying the two lines with $\begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix}$ and $\begin{pmatrix} 1 & -4 \\ 0 & 1 \end{pmatrix}^{-1}$ respectively and adding them.

The following lemma is the analogue of lemma 3.6, its proof is identical to that of lemma 3.6 and will be skipped here.

Lemma 5.4 *For every $\epsilon > 0$ there exists an $r > 1$ such that for $r > r$ the following hold:*

$$\begin{aligned} \langle \langle M^0 - L_{(A_i; \cdot)} \rangle \rangle_{X^2} &= \langle \langle X \rangle \rangle_2 \\ \langle \langle T^0 - S \rangle \rangle_{Y^2} &= \langle \langle Y \rangle \rangle_2 \end{aligned}$$

Here $X \in L^{1,2}(i^{-1} W_E^+)$ and $Y \in L^{1,2}(i^{-1} W_0^+)$.

The lemma allows us to replace the system (67) by the system:

$$\begin{aligned} S(\cdot) - \begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}^{-1} P(d_{4,0}; \cdot) &= \\ L_{(A_i; \cdot)}(\cdot) + \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} P(d_{100,0}; \cdot) &= \& \end{aligned} \tag{68}$$

The process of solving (68) is now step by step the analogue of solving (35). In particular, we solve the first of the two equations in (68) for \cdot in terms of \cdot :

$$\cdot = (\cdot) = S^{-1} \left(\begin{pmatrix} 0 & -1 \\ 1 & 0 \end{pmatrix}^{-1} P(d_{4,0}; \cdot) + \cdot \right)$$

Use this in the second equation of (68) and rewrite it as:

$$= L_{(A_i; \cdot)}^{-1} (\& - \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} P(d_{100,0}; (\cdot)))$$

To solve this last equation is the same as to find a fixed point of the map $Y : L^2(i^{-1} W_E^+) \rightarrow L^2(i^{-1} W_E^+)$ (the analogue of the map described by (38)) given by:

$$Y(\cdot) = L_{(A_i; \cdot)}^{-1} (\& - \begin{pmatrix} 1 & 0 \\ 0 & 0 \end{pmatrix} P(d_{100,0}; (\cdot)))$$

The proof of the existence of a unique fixed point of Y follows from a word by word analogue of the proof of theorem 3.8 together with the discussion preceding the theorem.

With the surjectivity of $L_{(A_i^0; \cdot)}$ proved, the process of deforming $(A_i^0; \cdot)$ to an honest solution $(A_i; \cdot)$ is accomplished by the same method as used in section 3.4 and will be skipped here.

To finish the proof theorem 5.1, we need to show that:

$$(A_0; \rho_0) \# (A_1; \rho_1) = (A; \rho)$$

This follows from the fact that as $r \rightarrow 1$, the distance $\text{dist}((A_i; \rho_i); (A_i^l; \rho_i^l))$ converges to zero, together with the following relations which follow directly from the definitions:

$$\begin{aligned} \rho_0 &= \rho_1 = \\ \tilde{\rho}_0 &= \tilde{\rho}_1 + \rho_1 & \tilde{\rho}_0 &= \\ A_0^l & \# A_1^l = A \end{aligned}$$

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Received: 24 November 2002 Revised: 27 January 2003