# THE ADIABATIC LIMIT OF WAVE MAP FLOW ON A TWO-TORUS 

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#### Abstract

The $S^{2}$ valued wave map flow on a Lorentzian domain $\mathbb{R} \times \Sigma$, where $\Sigma$ is any flat two-torus, is studied. The Cauchy problem with initial data tangent to the moduli space of holomorphic maps $\Sigma \rightarrow S^{2}$ is considered, in the limit of small initial velocity. It is proved that wave maps, in this limit, converge in a precise sense to geodesics in the moduli space of holomorphic maps, with respect to the $L^{2}$ metric. This establishes, in a rigorous setting, a long-standing informal conjecture of Ward.


## 1. Introduction

Wave maps are the analogue of harmonic maps in the case where the domain is Lorentzian. They satisfy a semilinear wave equation which has been heavily studied, in the simplest nontrivial case of $S^{2}$ target space, as a model PDE system involving a manifold-valued field 21. The wave map equation is particularly interesting in the case where the domain is $\left(\mathbb{R} \times \Sigma, d t^{2}-g_{\Sigma}\right)$, with $\left(\Sigma, g_{\Sigma}\right)$ an oriented Riemannian two-manifold. In this case, the static wave map problem is conformally invariant, so static solutions on $\Sigma=\mathbb{R}^{2}$ in particular have no preferred scale: they can be dilated to any size without changing their energy. This suggests that timedependent solutions with initial data close to a static solution might collapse and form singularities in finite time, an issue which has been heavily studied both numerically and analytically mainly for $\Sigma=\mathbb{R}^{2}, N=S^{2}$, within a certain rotational equivariance class. Numerical studies of increasing sophistication suggested that finite-time collapse can occur and suggested formal models of the collapse process [1,11, 13, 17. The first rigorous proof of blow-up came in the work of Krieger, Schlag and Tataru [8], who proved the existence of rotationally equivariant initial data, of topological degree $n=1$, leading to finite-time collapse. Rodnianski and Sterbenz subsequently proved existence of equivariant initial data of every degree $n \geq 4$ leading to finite-time collapse, and proved that collapse is stable to small perturbations of the initial data, at least within the equivariance class [19. These results were extended to every degree $n \geq 1$ in work of Raphael and Rodnianski [18], which also established detailed asymptotics and universality properties of the collapse mechanism. For a thorough discussion of blow-up of wave maps, see [17,18.

This paper addresses a different analytic issue from singularity formation, namely the validity of the geodesic approximation to wave map flow. To motivate this, one should think of wave maps $\mathbb{R} \times \Sigma \rightarrow S^{2}$ as (formal) critical points of the action

[^0]functional
(1.1) $\quad S=\int d t(T-E), \quad$ where $\quad T=\frac{1}{2} \int_{\Sigma}\left|\phi_{t}\right|^{2}, \quad E=\frac{1}{2} \int_{\Sigma}\left|d_{\Sigma} \phi\right|^{2}$.

It follows (from Noether's Theorem) that they conserve the total energy $T+E$. A rather general argument of Lichnerowicz [12] shows that for any map $\phi: \Sigma \rightarrow$ $S^{2}$ of topological degree $n \in \mathbb{Z}$ (subject to suitable boundary conditions if $\Sigma$ is noncompact), $E \geq 4 \pi|n|$, with equality if and only if $\phi$ is $\pm$ holomorphic. So holomorphic maps, if they exist, minimize potential energy in their homotopy class. Let us denote by $\mathrm{M}_{n}$ the moduli space of degree $n$ holomorphic maps $\Sigma \rightarrow S^{2}$. Consider a wave map $\phi(t)$ with $\phi(0) \in \mathrm{M}_{n}$ and $\phi_{t}(0) \in T_{\phi(0)} \mathrm{M}_{n}$ with $\left\|\phi_{t}(0)\right\|_{L^{2}}$ small. By conservation of $E+\frac{1}{2}\left\|\phi_{t}(t)\right\|_{L^{2}}^{2}$, one expects that $\phi(t)$ will stay close to $\mathrm{M}_{n}$, on which $E$ attains its minimum value, for as long as the solution persists. This led Ward to suggest [29], in the specific case $\Sigma=\mathbb{R}^{2}$, that such wave maps should be well approximated by the dynamical system with action $S$, but with $\phi(t)$ constrained to $\mathrm{M}_{n}$ for all time. Since $E$ is constant on $\mathrm{M}_{n}$, this constrained system is equivalent to geodesic motion on $\mathrm{M}_{n}$ with respect to the $L^{2}$ metric (obtained by restricting the quadratic form $T$ to $T \mathrm{M}_{n}$ ). A similar approximation had previously been proposed by Manton [15] for low energy monopole dynamics, and the geodesic approximation is now a standard technique in the study of the dynamics of topological solitons [16.

Geodesic motion on $\mathrm{M}_{2}$ (for $\Sigma=\mathbb{R}^{2}$ ) was studied in detail in 10. There is a technical problem: the $L^{2}$ metric is only well-defined on the leaves of a foliation of $\mathrm{M}_{n}$, and one must impose by hand that $\phi(t)$ remains on a single leaf. This turns out to be ill-justified (it precludes singularity formation for $n=1$, for example, in contradiction to [8, 18). This technical deficiency is removed if we choose $\Sigma$ to be a compact Riemann surface. Here geodesic motion in $\mathrm{M}_{n}$ is globally well-defined, if incomplete [20], and the $L^{2}$ geometry of $\mathrm{M}_{n}$ is quite well understood, at least for some choices of $\Sigma$ and $n$ [14, 22, 24].

The question remains: is geodesic motion in $\mathrm{M}_{n}$ really a good approximation to wave map flow in the adiabatic (low velocity) limit? The purpose of this paper is to prove that it is for times of order (initial velocity) ${ }^{-1}$, at least in the case where $\Sigma$ is any flat two-torus. More precisely, we will prove:

Theorem 1.1 (Main Theorem). Let $\mathrm{M}_{n}$ denote the moduli space of degree $n \geq 2$ holomorphic maps from a flat two-torus $\Sigma$ to $S^{2}$. For fixed $\phi_{0} \in \mathrm{M}_{n}$ and $\phi_{1} \in$ $T_{\phi_{0}} \mathrm{M}_{n}$ consider the one parameter family of initial value problems for the wave map equation with $\phi(0)=\phi_{0}, \phi_{t}(0)=\varepsilon \phi_{1}$, parametrized by $\varepsilon>0$. There exist constants $\tau_{*}>0$ and $\varepsilon_{*}>0$, depending only on the initial data, such that for all $\varepsilon \in\left(0, \varepsilon_{*}\right]$ the problem has a unique solution for $t \in\left[0, \tau_{*} / \varepsilon\right]$. Furthermore, the time rescaled solution

$$
\phi_{\varepsilon}:\left[0, \tau_{*}\right] \times \Sigma \rightarrow S^{2}, \quad \phi_{\varepsilon}(\tau, p)=\phi(\tau / \varepsilon, p)
$$

converges uniformly in $C^{1}$ to $\psi:\left[0, \tau_{*}\right] \times \Sigma \rightarrow S^{2}$, the geodesic in $\mathrm{M}_{n}$ with the same initial data, as $\varepsilon \rightarrow 0$.

To prove this we will adapt the perturbation method devised by Stuart to prove validity of the geodesic approximation in the critically coupled abelian Higgs and Yang-Mills-Higgs models [25,26]. The wave map problem has a key similarity with these gauge-theoretic problems, namely a moduli space of static solutions which
minimize energy in their homotopy class and satisfy a system of first order "Bogomol'nyi" equations. (For wave maps, the Bogomol'nyi equation is the condition that $\phi$ be $\pm$ holomorphic, i.e. the Cauchy-Riemann equation.) Roughly, the idea is to decompose the solution $\phi(t)$ as $\phi(t)=\psi(t)+\varepsilon^{2} Y(t)$ where $\psi(t) \in \mathrm{M}_{n}$ and control the growth of a suitable Sobolev norm of the error $Y(t)$ uniformly in $\varepsilon$ by means of energy estimates. One concurrently shows that the projected trajectory $\psi(t)$ converges to a geodesic in $\mathrm{M}_{n}$.

In comparison with Stuart's work on vortices and monopoles, the situation we study is simpler in two respects: we work on a compact domain $\Sigma$ (rather than $\mathbb{R}^{2}$ or $\mathbb{R}^{3}$ ), and our system has no gauge symmetry. On the other hand, the wave map problem introduces two new challenges for the method.

First, our field is manifold-valued, so it is not clear a priori what the decomposition $\phi(t)=\psi(t)+\varepsilon^{2} Y(t)$ really means. In preliminary work on this problem, it was suggested that the correct formulation was $\phi(t)=\exp _{\psi(t)} \varepsilon^{2} Y(t)$, where $\exp : T S^{2} \rightarrow S^{2}$ is the exponential map [5. In fact, this turns out not to have the analytic properties required by Stuart's method (except for rotationally equivariant wave maps). In this paper we isometrically embed $S^{2}$ in $\mathbb{R}^{3}$ and use the ambient linear structure to project as usual, $\phi(t)=\psi(t)+\varepsilon^{2} Y(t)$. This choice is simple, but has significant repercussions: $Y$ is no longer tangent to the map $\psi$ (not a section of $\psi^{-1} T S^{2}$ ), and must satisfy a nonlinear pointwise constraint to ensure that $\phi$ is $S^{2}$ valued. The evolution of $Y$ is governed by a nonlinear wave equation whose (spatial) linear part is the Jacobi operator $J_{\psi}$ for the harmonic map $\psi: \Sigma \rightarrow S^{2}$. It turns out that $J_{\psi}$ is not self-adjoint when acting on nontangent sections (such as $Y)$. Since self-adjointness of (the analogue of) $J_{\psi}$ is crucial for Stuart's method, we must devise a way around this: we replace $J_{\psi}$ by an "improved" Jacobi operator $L_{\psi}$, which coincides with $J_{\psi}$ on tangent sections but is self-adjoint on all sections, and introduce compensating nonlinear terms into the wave equation for $Y$ using the pointwise constraint. Further difficulties result: $L_{\psi}$, unlike $J_{\psi}$, does not define a coercive quadratic form on the $L^{2}$ orthogonal complement of $T_{\psi} \mathrm{M}_{n}$. We must work instead with a weaker near-coercivity property, which turns out to suffice for our purposes.

Second, while $\Sigma$ is compact, the moduli space $\mathrm{M}_{n}$ is not. Of course, the vortex and monopole moduli spaces, dealt with by Stuart, are also noncompact, but in those cases, moving to infinity corresponds to (clusters of) solitons separating off and escaping to infinite separation, a well-controlled process. For wave maps, by contrast, approaching the boundary of $\mathrm{M}_{n}$ at infinity corresponds to one or more lumps collapsing and "bubbling off". In this process, $\psi$ becomes singular, and both geodesic motion and wave map flow become badly behaved. To handle this, we must keep careful track of the position (of $\psi \in \mathrm{M}_{n}$ ) dependence of our various estimates and modify Stuart's a priori energy bound so that we simultaneously control the error $Y(t)$, the deviation of $\psi(t)$ from the corresponding geodesic, and the distance of $\psi(t)$ from $\partial_{\infty} \mathrm{M}_{n}$.

It is interesting to speculate to what extent Theorem 1.1 can be generalized. It is clear that the proof presented here generalizes quite easily to the case of a general compact Riemann surface, provided $n$ is sufficiently large compared with the genus of $\Sigma$. The reason for restricting to the case $\Sigma=T^{2}$ is mainly one of presentation: the existence of global cartesian coordinates makes it straightforward to define the various function spaces, for example. Generalizing the target space is not so straightforward. The wave map flow $\mathbb{R} \times \Sigma \rightarrow N$ has the appropriate

Bogomol'nyi form for Stuart's method to apply whenever $\Sigma, N$ are both compact Kähler manifolds (in fact, it suffices for $\Sigma$ to be co-Kähler). The choice $N=\mathbb{C} P^{k}$, $k \geq 2$, is of some interest in mathematical physics, for example. But here the reliance on an isometric embedding $N \subset \mathbb{R}^{p}$ becomes very problematic. It seems likely that some variant of Theorem 1.1 does remain true for general compact Kähler targets, but proving it would require a rather different approach, perhaps along the lines sketched in [27.

One should note that Theorem 1.1 gives no information about singularity formation for wave maps on $\mathbb{R} \times \Sigma$ because, although there certainly are geodesics $\psi(\tau)$ which hit $\partial_{\infty} \mathrm{M}_{n}$ in time $\tau_{0}<\infty$ and the corresponding wave maps do converge uniformly to $\psi(\tau)$ on some interval $\left[0, \tau_{*}\right]$, there is no reason to expect $\tau_{*}=\tau_{0}$. In fact Raphael and Rodnianski have shown that singularity formation of equivariant wave maps on $\mathbb{R}^{2}$ deviates significantly from the dynamics predicted by the (suitably regulated) geodesic approximation [18. Since blow-up is a (spatially) local phenomenon, these results presumably apply in some form on the torus, which would imply $\tau_{*}<\tau_{0}$. Nontheless, the geodesic approximation (on compact $\Sigma$ or, regulated, on $\Sigma=\mathbb{R}^{2}$ ) predicted finite time blow-up of wave maps in $(2+1)$ dimensions, and this prediction turned out to be correct. The geodesic approximation also makes predictions about the genericity of blow-up. It is not hard to prove, for example, that generic geodesics on $\mathrm{M}_{1}$ for $\Sigma=S^{2}$ do not hit the boundary at infinity. It would be interesting to see whether the full wave map flow has this property (i.e. generic Cauchy data tangent to $\mathrm{M}_{1}$ have global smooth solutions). The analogue of Theorem 1.1 for $\Sigma=S^{2}$ could provide a starting point for proving such results.

The rest of the paper is structured as follows. In section 2 the moduli space $\mathrm{M}_{n}$ of holomorphic maps is introduced and its key property, Proposition 2.1 (existence of a smooth local parametrization about any point), established. In section 3, the projection of the wave map flow to $\mathrm{M}_{n}$ is defined, and the coupled system satisfied by $\psi$ and the error section $Y$ is derived (equation (3.15)). In section 4 some standard functional analytic definitions and results are introduced. Our aim here, and in the remainder of the paper, is to make the proof accessible to a wide mathematical physics audience, not just experts in PDE. In section 5a local existence and uniqueness theorem for the coupled system (3.15) governing ( $\psi, Y$ ) is proved (Theorem 5.1). Of course, local existence and uniqueness of wave maps in this context is not new; the extra, and new, information we obtain here is local existence and uniqueness of the projection to $\mathrm{M}_{n}$. This is the engine underlying Stuart's method, and we go through the argument in some detail, not only because there are certain new aspects we have to deal with which Stuart did not (e.g. preservation of the pointwise constraint on $Y$ ), but also because the requirements of Picard's method for this proof determine our choice of function spaces, a point which is not obvious (to the nonanalyst) in Stuart's original applications of the method [25,26]. In section 6 the key near-coercivity property of the quadratic form associated with the improved Jacobi operator is proved (Theorem 6.5) (roughly, that $\left\langle L_{\psi} Y, L_{\psi} L_{\psi} Y\right\rangle_{L^{2}}$ controls the $H^{3}$ norm of $Y$ ). In section 7, energy estimates are established which bound the growth of $Y(t)$. Finally, in section 8 , the coercivity properties and energy estimates are combined to prove long time existence of the solution ( $\psi, Y$ ) and establish convergence to the corresponding geodesic in $\mathrm{M}_{n}$. An appendix presents the proofs of some basic analytic properties of the nonlinear terms in the coupled system for $(\psi, Y)$, which are used repeatedly in section 5

## 2. The moduli space of static $n$-Lumps

Static wave maps are harmonic maps $\phi: \Sigma \rightarrow S^{2}$, that is, solutions of the harmonic map equation

$$
\begin{equation*}
\phi_{x x}+\phi_{y y}+\left(\left|\phi_{x}\right|^{2}+\left|\phi_{y}\right|^{2}\right) \phi=0 \tag{2.1}
\end{equation*}
$$

or, equivalently, critical points of the Dirichlet energy

$$
\begin{equation*}
E(\phi)=\frac{1}{2} \int_{\Sigma}\left|\phi_{x}\right|^{2}+\left|\phi_{y}\right|^{2} \tag{2.2}
\end{equation*}
$$

Maps $\Sigma \rightarrow S^{2}$ fall into disjoint homotopy classes labelled by their degree $n \in \mathbb{Z}$, which we may assume, without loss of generality, is nonnegative. An argument of Lichnerowicz [3, p. 39] shows that in the degree $n$ class, $E(\phi) \geq 4 \pi n$, with equality if and only if $\phi$ is holomorphic. Furthermore, all harmonic maps $\Sigma \rightarrow S^{2}$ of degree $n \geq 2$ are holomorphic [4. So the moduli space of interest, $\mathrm{M}_{n}$, is the space of degree $n$ holomorphic maps $\Sigma \rightarrow S^{2}$. Such maps are called " $n$-lumps" by analogy with the case $\Sigma=\mathbb{C}$, where the Dirichlet energy density typically exhibits $n$ distinct local maxima, which may loosely be thought of as smoothed out particles, or lumps of energy.

The global topology of the space $\mathrm{M}_{n}$ is quite complicated, for example, $\mathrm{M}_{2} \cong$ $[\Sigma \times \operatorname{PSL}(2, \mathbb{C})] /\left(\mathbb{Z}_{2} \times \mathbb{Z}_{2}\right)$, [23]. For our purposes local information will suffice, however. Given a smooth variation $\phi_{s}$ of $\phi=\phi_{0} \in \mathrm{M}_{n}$ we have $d E\left(\phi_{s}\right) / d s=0$ at $s=0$ (since $\phi$ is harmonic) and

$$
\begin{equation*}
\frac{d^{2} E\left(\phi_{s}\right)}{d s^{2}}=\left\langle V, J_{\phi} V\right\rangle=\int_{\Sigma} V \cdot J_{\phi} V \tag{2.3}
\end{equation*}
$$

where $V=\left.\partial_{s} \phi_{s}\right|_{s=0} \in \Gamma\left(\phi^{-1} T S^{2}\right)$ is the infinitesimal generator of the variation and

$$
\begin{equation*}
J_{\phi} V=-V_{x x}-V_{y y}-\left(\left|\phi_{x}\right|^{2}+\left|\phi_{y}\right|^{2}\right) V-2\left(\phi_{x} \cdot V_{x}+\phi_{y} \cdot V_{y}\right) \phi \tag{2.4}
\end{equation*}
$$

is the Jacobi operator at the map $\phi$ [28, p. 155]. This operator is self-adjoint and elliptic, and its spectrum determines the stability properties of $\phi$. By the Lichnerowicz argument, $\phi$ minimizes $E$ in its homotopy class, so spec $J_{\phi}$ is nonnegative. Given a variation $\phi_{s}$ through harmonic maps, that is, a curve in $\mathrm{M}_{n}$ through $\phi=\phi_{0}$, its infinitesimal generator $V=\left.\partial_{s} \phi_{s}\right|_{s=0}$ satisfies $J_{\phi} V=0$. Hence $T_{\phi} \mathrm{M}_{n} \subset \operatorname{ker} J_{\phi}$. For a general harmonic map $\phi: M \rightarrow N$ between Riemannian manifolds, the converse may be false; that is, there may be sections $V \in \operatorname{ker} J_{\phi} \subset \Gamma\left(\phi^{-1} T N\right)$ which are not tangent to any variation of $\phi$ through harmonic maps, and in this case the space of harmonic maps $M \rightarrow N$ may not be a smooth manifold around $\phi$. It is important for us to rule out this kind of bad behaviour in our case. More precisely, we need:

Proposition 2.1. Given any $\phi_{0} \in \mathrm{M}_{n}, n \geq 2$, there exists an open set $U \subset \mathbb{R}^{4 n}$ and a smooth map $\psi: U \times \Sigma \rightarrow S^{2}$ such that
(i) for each $q \in U, \psi(q, \cdot) \in \mathrm{M}_{n}$,
(ii) there exists $q_{0} \in U$ such that $\phi_{0}=\psi\left(q_{0}, \cdot\right)$, and
(iii) $\psi_{\mu}=\partial \psi / \partial q^{\mu}, \mu=1,2, \ldots, 4 n$ span $\operatorname{ker} J_{\psi(q, \cdot)}$.

Proof. Choose any $p \in S^{2}$ such that both $p$ and $-p$ are regular values of $\phi_{0}$ (such $p$ exists by Sard's Theorem). Then $\phi: \Sigma \rightarrow S^{2}$ is in $\mathrm{M}_{n}$ if and only if $s_{p} \circ \phi$, its
image under stereographic projection from $p$, is meromorphic, of degree $n$, that is, a degree $n$ elliptic function. The most general degree $n$ elliptic function is 9

$$
\begin{equation*}
\left(s_{p} \circ \phi\right)(z)=\lambda \frac{\sigma\left(z-a_{1}\right) \cdots \sigma\left(z-a_{n}\right)}{\sigma\left(z-b_{1}\right) \cdots \sigma\left(z-b_{n}\right)} \tag{2.5}
\end{equation*}
$$

where $\sigma$ is the Weierstrass sigma function, $\lambda, a_{1}, \ldots, a_{n}, b_{1}, \ldots, b_{n}$ are complex constants, $\lambda \neq 0, \sum a_{i}=\sum b_{i} \bmod \Lambda$ and $\left\{a_{i}\right\} \cap\left\{b_{j}\right\}=\emptyset$. Hence, we may parametrize a general point $\phi \in \mathrm{M}_{n}$ by $4 n$ real numbers $q^{\mu}$, for example, the real and imaginary parts of $\lambda, a_{1}, \ldots, b_{n-1}$ having set $b_{n}=a_{1}+\cdots+a_{n}-b_{1}-\cdots-b_{n-1}$. Further, $\phi$ manifestly depends smoothly on $q$ and $z$. Hence we have a smooth map $\psi: \widetilde{U} \times \Sigma \rightarrow S^{2}$ satisfying properties (i) and (ii). By our choice of $p, s_{p} \circ \phi_{0}$ has $n$ distinct zeroes and $n$ distinct poles, so $\left\{\psi_{\mu}\right\}$ at $q=q_{0}$ are linearly independent sections of $\psi(q, \cdot)^{-1} T S^{2}$, and hence, by smoothness, also linearly independent on some neighbourhood $U$ of $q_{0}$ in $\widetilde{U}$. As explained previously, $\psi_{\mu}$ span a subspace of $\operatorname{ker} J_{\psi(q, \cdot)}$, so it remains to show that ker $J_{\phi}$ has dimension $4 n$ for any $\phi \in \mathrm{M}_{n}$.

It is known [28, p. 174] that ker $J_{\phi}$ is isomorphic, as a complex vector space, to $H^{0}(\Sigma, L)$, the space of holomorphic sections of the line bundle $L=\phi^{-1} T^{\prime} S^{2}$, where $T^{\prime} S^{2}$ denotes the holomorphic tangent bundle of $S^{2}$. Since $\phi$ has degree $n$ and $T^{\prime} S^{2}$ has degree $2, L$ has degree $2 n$. Now, by the Riemann-Roch formula [6],

$$
\begin{equation*}
\operatorname{dim} H^{0}(\Sigma, L)-\operatorname{dim} H^{1}(\Sigma, L)=\operatorname{deg} L=2 n \tag{2.6}
\end{equation*}
$$

since $\Sigma$ has genus 1 . But, by Serre duality, $H^{1}(\Sigma, L) \cong H^{0}\left(\Sigma, K \otimes L^{*}\right)^{*}$ where $K$ is the canonical bundle of $\Sigma$. Now $K$ is trivial, so $K \otimes L^{*}$ has degree $-2 n$, and hence has no holomorphic sections, whence $H^{1}(\Sigma, L)=0$. It follows that ker $J_{\phi}$ has real dimension $4 n$, as required.

We can regard $q^{\mu}$ as local coordinates on $\mathrm{M}_{n}$. Given the initial data $\phi_{1} \in T_{\phi_{0}} \mathrm{M}_{n}$ of interest, we choose and fix such a $\psi: U \times \Sigma \rightarrow S^{2}$ and denote by $q_{0} \in U$ and $q_{1} \in \mathbb{R}^{4 n}$ those vectors such that $\phi_{0}=\psi\left(q_{0}, \cdot\right)$ and $\phi_{1}=q_{1}^{\mu} \psi_{\mu}\left(q_{0}, \cdot\right)$. Here, as henceforth, we use the Einstein summation convention on repeated indices. We also choose and fix a compact neighbourhood $K$ of $q_{0}$ in $U$. In a slight abuse of notation, we will also use the symbol $\psi$ to denote the associated map $U \supset K \rightarrow \mathrm{M}_{n} \subset\left(S^{2}\right)^{\Sigma}$, so $\psi(q)$ will denote the holomorphic map $z \mapsto \psi(q, z)$. We will also use $\psi(t)$ as shorthand for $\psi(q(t))$, meaning a general curve in $\mathrm{M}_{n}$.

There is a natural Riemannian metric on $\mathrm{M}_{n}$, the $L^{2}$ metric, whose components in the local coordinate system $q^{\mu}$ are

$$
\begin{equation*}
\gamma_{\mu \nu}(q)=\left\langle\psi_{\mu}, \psi_{\nu}\right\rangle=\int_{\Sigma} \frac{\partial \psi}{\partial q^{\mu}} \cdot \frac{\partial \psi}{\partial q^{\nu}} . \tag{2.7}
\end{equation*}
$$

The associated Christoffel symbol is

$$
\begin{equation*}
G_{\lambda \nu}^{\mu}(q)=\gamma^{\mu \alpha}\left\langle\psi_{\alpha}, \psi_{\lambda \nu}\right\rangle \tag{2.8}
\end{equation*}
$$

where $\gamma^{\mu \nu}$ is the inverse metric and $\psi_{\mu \nu}=\partial^{2} \psi / \partial q^{\mu} \partial q^{\nu}$. This is the metric whose geodesics approximate wave maps in the adiabatic limit.

## 3. Projection of wave map flow and the coupled system

The wave map equation for $\phi: \mathbb{R} \times \Sigma \rightarrow S^{2} \subset \mathbb{R}^{3}$ is

$$
\begin{equation*}
\phi_{t t}-\phi_{x x}-\phi_{y y}+\left(\left|\phi_{t}\right|^{2}-\left|\phi_{x}\right|^{2}-\left|\phi_{y}\right|^{2}\right) \phi=0 . \tag{3.1}
\end{equation*}
$$

Given $\varepsilon>0$, a small parameter, we decompose $\phi$ as

$$
\begin{equation*}
\phi=\psi+\varepsilon^{2} Y \tag{3.2}
\end{equation*}
$$

where, at each fixed time, $\psi(t, \cdot): \Sigma \rightarrow S^{2}$ is a degree $n$ harmonic map, and $Y: \Sigma \rightarrow \mathbb{R}^{3}$. We may think of $\psi(t)$ as a curve in $\mathrm{M}_{n}$, the moduli space of degree $n$ harmonic maps, and $Y$ as the "error" incurred by projecting $\phi(t)$ to $\psi(t)$. It is useful to think of $Y$ as a section of $\psi^{-1} \mathbb{R}^{3}$, where $\underline{\mathbb{R}}^{3}=S^{2} \times \mathbb{R}^{3}$ is the trivial $\mathbb{R}^{3}$ bundle over $S^{2}$, and $\psi^{-1} \underline{\mathbb{R}^{3}}$ is its pullback to $\Sigma$. With this in mind, we refer to $Y$ as the "error section", and to any $Z: \Sigma \rightarrow \mathbb{R}^{3}$ with $Z \cdot \psi=0$ pointwise as a "tangent section" (in bundle language, $Z$ is a section of $\psi^{-1} T S^{2} \subset \psi^{-1} \mathbb{R}^{3}$ ). Clearly, $Y$ is not a tangent section (unless $Y=0$ ). Since both $\psi$ and $\phi$ are $S^{2}$ valued, $Y$ must satisfy the pointwise constraint

$$
\begin{equation*}
\psi \cdot Y=-\frac{1}{2} \varepsilon^{2}|Y|^{2} . \tag{3.3}
\end{equation*}
$$

For a given curve $\psi(t)$, if $\phi$ is a wave map, then $Y$ must satisfy the PDE obtained by substituting (3.2) into (3.1),

$$
\begin{equation*}
Y_{t t}+J_{\psi} Y=k+\varepsilon j, \tag{3.4}
\end{equation*}
$$

where $J_{\psi}$ is the Jacobi operator associated with the harmonic map $\psi(t): \Sigma \rightarrow S^{2}$ and the terms on the right hand side are

$$
\begin{align*}
k= & -\left(\psi_{\tau \tau}+\left|\psi_{\tau}\right|^{2} \psi\right), \\
j= & 2\left(\psi_{\tau} \cdot Y_{t}\right) \psi+\varepsilon\left\{\left(\left|Y_{t}\right|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) \psi+\left(\left|\psi_{\tau}\right|^{2}-2 \psi_{x} \cdot Y_{x}-2 \psi_{y} \cdot Y_{y}\right) Y\right\} \\
(3.5) \quad & +2 \varepsilon^{2}\left(\psi_{\tau} \cdot Y_{t}\right) Y+\varepsilon^{3}\left(\left|Y_{t}\right|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) Y . \tag{3.5}
\end{align*}
$$

We have here introduced the slow time variable $\tau=\varepsilon t$ as a bookkeeping device. The precise expression for $j$ is not important. What matters is its qualitative form: it depends only on $\psi, \psi_{\tau}$ and $Y$ and its first derivatives, and the dependence is smooth (polynomial, in fact).

Superficially (3.4) looks exactly analogous to the corresponding equation in Stuart's analysis of vortex dynamics [25], but this is deceptive. As already noted, the Jacobi operator is a self-adjoint, elliptic, second order linear operator $J_{\psi}$ : $\Gamma\left(\psi^{-1} T S^{2}\right) \rightarrow \Gamma\left(\psi^{-1} T S^{2}\right)$ whose spectrum determines the stability properties of the harmonic map $\psi$ [28. It is important to realize, however, that in (3.4) $J_{\psi}$ is acting on $Y$, which is not a tangent section. So in (3.4), $J_{\psi}$ is precisely the same operator defined above (2.4), but extended to act on sections of $\psi^{-1} \mathbb{R}^{3}$. But $J_{\psi}$ is not self-adjoint (with respect to $L^{2}$ ) as an operator on $\psi^{-1} \mathbb{R}^{3}$, and self-adjointness of (the analogue of) $J_{\psi}$ is a crucial ingredient in Stuart's method. In fact

$$
\begin{equation*}
J_{\psi} Z=-\Delta Z-\left(\left|\psi_{x}\right|^{2}+\left|\psi_{y}\right|^{2}\right) Z+A_{\psi} Z \tag{3.6}
\end{equation*}
$$

where the non-self-adjoint piece and its adjoint are

$$
\begin{align*}
A_{\psi} Z & =-2\left(\psi_{x} \cdot Z_{x}+\psi_{y} \cdot Z_{y}\right) \psi \\
A_{\psi}^{\dagger} Z & =-2\left\{(\psi \cdot Z) \Delta \psi+(\psi \cdot Z)_{x} \psi_{x}+(\psi \cdot Z)_{y} \psi_{y}\right\} \tag{3.7}
\end{align*}
$$

and we have adopted the analysts' convention for the Laplacian, that is, $\Delta Z=$ $Z_{x x}+Z_{y y}$. To remedy this deficiency, we make the following definition:
Definition 3.1. Given a harmonic map $\psi: \Sigma \rightarrow S^{2}$, we define its improved Jacobi operator to be

$$
L_{\psi}: \Gamma\left(\psi^{-1} \underline{\mathbb{R}^{3}}\right) \rightarrow \Gamma\left(\psi^{-1} \underline{\mathbb{R}^{3}}\right), \quad L_{\psi}=J_{\psi}+A_{\psi}^{\dagger} .
$$

Note that $L_{\psi}$ coincides with $J_{\psi}$ on $\Gamma\left(\psi^{-1} T S^{2}\right)$, and hence $L_{\psi}$ maps tangent sections to tangent sections. Its principal part is the Laplacian, so it is elliptic, and it is manifestly self-adjoint.

Remark 3.2. Any section can be decomposed into tangent and normal components. As just observed, $L_{\psi}$ maps a tangent section $Z$ to the tangent section $J_{\psi} Z$, so an alternative way of characterizing $L_{\psi}$ is by specifying how it acts on normal sections, $\alpha \psi$, where $\alpha: \Sigma \rightarrow \mathbb{R}$. A short calculation, using harmonicity of $\psi$, yields

$$
\begin{equation*}
L_{\psi}(\alpha \psi)=-(\Delta \alpha) \psi-4\left(\alpha_{x} \psi_{x}+\alpha_{y} \psi_{y}\right) \tag{3.8}
\end{equation*}
$$

It follows immediately that $\operatorname{ker} L_{\psi}=\operatorname{ker} J_{\psi} \oplus\langle\psi\rangle$. Note that, in general, $L_{\psi}$ does not map normal sections to normal sections.

Now, for any $Y$ satisfying the pointwise constraint,

$$
\begin{equation*}
A_{\psi}^{\dagger} Y=\varepsilon^{2}\left\{|Y|^{2} \Delta \psi+2\left(Y \cdot Y_{x}\right) \psi_{x}+2\left(Y \cdot Y_{y}\right) \psi_{y}\right\}=: \widehat{j} \tag{3.9}
\end{equation*}
$$

and so, for any curve $\psi(t)$, if $\phi$ is a wave map, then $Y$ satisfies the PDE

$$
\begin{equation*}
Y_{t t}+L_{\psi} Y=k+\varepsilon j^{\prime} \tag{3.10}
\end{equation*}
$$

where $j^{\prime}=j+\widehat{j}$. Note that $j^{\prime}$ has the same qualitative analytic properties as $j$ (specifically, no higher than first derivatives of $Y$ appear).

We have yet to specify the curve $\psi(t)$ in $\mathrm{M}_{n}$. We do this by demanding that the error section $Y$ should at all times be $L^{2}$ orthogonal to $T_{\psi} \mathrm{M}_{n}$. In this case ( $\Sigma=T^{2}, n \geq 2$ ), $T_{\psi} \mathrm{M}_{n}=\operatorname{ker} J_{\psi}$, so, in terms of the local coordinate system $q$ on $\mathrm{M}_{n}$ provided by Proposition [2.1, this amounts to requiring

$$
\begin{equation*}
\left\langle Y, \psi_{\mu}\right\rangle=0, \quad \mu=1,2, \ldots, 4 n \tag{3.11}
\end{equation*}
$$

We convert this into an evolution equation for $q$ by differentiating the orthogonality constraint (3.11) twice with respect to time and using (3.10),

$$
\begin{equation*}
\left\langle-L_{\psi} Y+k+\varepsilon j^{\prime}, \psi_{\mu}\right\rangle+2 \varepsilon Y_{t}, \psi_{\mu \nu} \dot{q}^{\nu}+\varepsilon^{2}\left\langle Y, \psi_{\mu \nu \lambda}\right\rangle \dot{q}^{\nu} \dot{q}^{\lambda}+\varepsilon^{2}\left\langle Y, \psi_{\mu \nu}\right\rangle \ddot{q}^{\nu}=0 \tag{3.12}
\end{equation*}
$$

where an overdot denotes differentiation with respect to $\tau$. Now $L_{\psi}$ is self-adjoint and $\psi_{\mu} \in \operatorname{ker} J_{\psi} \subset \operatorname{ker} L_{\psi}$, so this equation simplifies to $\left\langle k, \psi_{\mu}\right\rangle=O(\varepsilon)$ or, more explicitly,

$$
\begin{equation*}
\ddot{q}^{\mu}+G_{\nu \lambda}^{\mu}(q) \dot{q}^{\nu} \dot{q}^{\lambda}=\varepsilon h^{\mu}\left(\varepsilon, q, \dot{q}, Y, Y_{t}\right)+\varepsilon^{2} \gamma^{\mu \nu}\left\langle Y, \psi_{\nu \lambda}\right\rangle \ddot{q}^{\lambda} \tag{3.13}
\end{equation*}
$$

where $\gamma$ and $G$ are the $L^{2}$ metric and its Christoffel symbol in the coordinate system $q$, (2.7), (2.8), and the function $h$ is

$$
\begin{align*}
h^{\mu}= & \gamma^{\mu \nu}\left\{\left\langle Y_{t}, \psi_{\nu \lambda}\right\rangle \dot{q}^{\lambda}+\varepsilon\left\langle Y, \psi_{\nu \lambda \rho}\right\rangle \dot{q}^{\lambda} \dot{q}^{\rho}+\varepsilon\left\langle\left(\psi_{\lambda} \cdot \psi_{\rho}\right) Y, \psi_{\nu}\right\rangle \dot{q}^{\lambda} \dot{q}^{\rho}\right. \\
& -2 \varepsilon\left\langle\left(\psi_{x} \cdot Y_{x}+\psi_{y} \cdot Y_{y}\right) Y, \psi_{\nu}\right\rangle+2 \varepsilon^{2}\left\langle\left(\psi_{\lambda} \cdot Y_{t}\right) Y, \psi_{\nu}\right\rangle \dot{q}^{\lambda} \\
& \left.+\varepsilon^{3}\left\langle\left(\left|Y_{t}\right|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) Y, \psi_{\nu}\right\rangle\right\} . \tag{3.14}
\end{align*}
$$

Taking the formal limit $\varepsilon \rightarrow 0$, (3.13) reduces to the geodesic equation on $\left(\mathrm{M}_{n}, \gamma\right)$, as one would hope.

To summarize, if $\phi$ is a wave map and $q(t)$ is a curve in $\mathrm{M}_{n}$ such that $Y=$ $\varepsilon^{-2}(\phi-\psi(q))$ satisfies the orthogonality constraint (3.11) at all times, then $(Y, q)$ satisfies the coupled system

$$
\begin{equation*}
Y_{t t}+L Y=k+\varepsilon j^{\prime}, \quad \ddot{q}^{\mu}+G_{\nu \lambda}^{\mu} \dot{q}^{\nu} \dot{q}^{\lambda}=\varepsilon h^{\mu}+\varepsilon^{2} \gamma^{\mu \nu}\left\langle Y, \psi_{\nu \lambda}\right\rangle \ddot{q}^{\lambda} \tag{3.15}
\end{equation*}
$$

and the pointwise constraint (3.3). Conversely, if $(Y, q)$ satisfies the constraints (3.3) and (3.11) and the coupled system (3.15), then $\phi=\psi(q)+\varepsilon^{2} Y$ is a wave map. Our goal is to prove that (3.15) with fixed initial data $q(0)=q_{0}, \dot{q}(0)=q_{1}$, $Y(0)=0, Y_{t}(0)=0$ has solutions with $\|Y\|_{C^{1}}$ bounded uniformly in $\varepsilon$ for times of order $\varepsilon^{-1}$. It follows immediately that, in the limit $\varepsilon \rightarrow 0, \phi(\tau / \varepsilon)$ converges uniformly to a curve $\psi(\tau)$ in $\mathrm{M}_{n}$. In the course of the proof, we will simultaneously show that $\psi(\tau)$ is the geodesic with initial data $q_{0}, q_{1}$.

## 4. Analytic preliminaries

In this section we set up the function spaces we will use and collect some standard functional analytic results which we will appeal to repeatedly. More details can be found in [2] and references therein. Let $\mathrm{H}^{k}$ denote the set of real-valued functions on $\Sigma$ whose partial derivatives up to order $k$ are square integrable. This is a Hilbert space with respect to the inner product

$$
\begin{equation*}
\langle f, g\rangle_{k}=\sum_{|\alpha| \leq k} \int_{\Sigma} D_{\alpha} f D_{\alpha} g \tag{4.1}
\end{equation*}
$$

where $\alpha$ is a multi-index taking values from $\{x, y\},|\alpha|$ is its length and $D_{\alpha}=$ $\partial_{\alpha_{1}} \partial_{\alpha_{2}} \cdots \partial_{\alpha_{|\alpha|}}$, so $D_{(x, x, y)}=\partial_{x}^{2} \partial_{y}$, for example. We denote the corresponding norm by $\|\cdot\|_{k}$,

$$
\begin{equation*}
\|f\|_{k}^{2}=\langle f, f\rangle_{k} \tag{4.2}
\end{equation*}
$$

Let $H^{k}=\mathrm{H}^{k} \oplus \mathrm{H}^{k} \oplus \mathrm{H}^{k}$, the space of $\mathbb{R}^{3}$-valued functions on $\Sigma$ whose components are in $\mathrm{H}^{k}$. This is a Hilbert space with respect to the inner product

$$
\begin{equation*}
\langle Y, Z\rangle_{k}=\left\langle Y_{1}, Z_{1}\right\rangle_{k}+\left\langle Y_{2}, Z_{2}\right\rangle_{k}+\left\langle Y_{3}, Z_{3}\right\rangle_{k} \tag{4.3}
\end{equation*}
$$

whose norm will again be denoted $\|\cdot\|_{k}$. We adopt the convention that $\|\cdot\|=\|\cdot\|_{0}$ and $\langle\cdot, \cdot\rangle=\langle\cdot, \cdot\rangle_{0}$; that is, undecorated norms and inner products refer to $L^{2}$. We will frequently, and without further comment, use the Cauchy-Schwarz inequality

$$
\begin{equation*}
\langle Y, Z\rangle_{k} \leq\|Y\|_{k}\|Z\|_{k} \tag{4.4}
\end{equation*}
$$

and the trivial bound $\|Y\|_{k} \geq\|Y\|_{k^{\prime}}$ if $k \geq k^{\prime}$.
In the sequel, we will prove existence of a solution of the coupled system (3.15) with $Y \in H^{k}, Y_{t} \in H^{k-1}$, for $k=3$. This choice of $k$ is motivated by the following fundamental fact about $\mathrm{H}^{k}$ on a compact 2-manifold.

Proposition 4.1 (Algebra property of $\mathrm{H}^{k}, k \geq 2$ ). The Banach space $\left(\mathrm{H}^{k},\|\cdot\|_{k}\right)$ is a Banach algebra for all $k \geq 2$. That is, if $f \in \mathrm{H}^{k}$ and $g \in \mathrm{H}^{k}$, then $f g \in \mathrm{H}^{k}$, and there exists a constant $\alpha_{k}>0$, depending only on $\Sigma$ and $k$, such that $\|f g\|_{k} \leq$ $\alpha_{k}\|f\|_{k}\|g\|_{k}$.

It follows directly from this that if $\left(Y, Y_{t}\right) \in H^{3} \oplus H^{2}$, then the nonlinear term $j^{\prime}$ in the coupled system is in $H^{2}$, and $\left\|j^{\prime}\right\|_{2}$ can be bounded by a polynomial in $\|Y\|_{3},\left\|Y_{t}\right\|_{2}$ (the point being that $j^{\prime}$ contains no derivatives of $Y$ higher than first, and $Y_{x}, Y_{y}, Y_{t}$ are all in $H^{2}$ ). This is crucial, not only for proving the local existence result for (3.15), but also in later sections where we prove that $\|Y\|_{3}$ is controlled by $\langle L Y, L L Y\rangle$ and make energy estimates for the solution. So the fact that we choose $k=3$ is not just motivated by a desire to get strong bounds on the error section $Y$; the method will not work for any lower $k$. Indeed, to uniformly bound $Y$ on $\Sigma$, it would suffice to control $\|Y\|_{2}$, as we have the following Sobolev inequality.

Proposition 4.2 (Sobolev inequalities). Let $C^{k}$ denote the Banach space of continuous maps $\Sigma \rightarrow \mathbb{R}^{3}$ with the usual norm $\|Y\|_{C^{k}}=\sup \left\{\left|D_{\alpha} Y(p)\right|:|\alpha| \leq k, p \in \Sigma\right\}$. Then $H^{2} \subset C^{0}, H^{3} \subset C^{1}$, and there is a constant $\alpha>0$, depending only on $\Sigma$, such that $\|Y\|_{C^{k}} \leq \alpha\|Y\|_{k+2}$ for all $Y \in H^{2}, k=0,1$. (More briefly, the inclusions $\iota: H^{2} \rightarrow C^{0}$ and $\iota: H^{3} \rightarrow C^{1}$ are continuous.)

In later sections we will need to bound $\left\|L_{\psi} Y\right\|_{k}$ in terms of $\|Y\|_{k+2}$. Of course, since $L_{\psi}$ is a linear second order operator we have trivially, for all $q \in K$, the upper bound

$$
\begin{equation*}
\left\|L_{\psi} Y\right\|_{k} \leq C\|Y\|_{k+2} \tag{4.5}
\end{equation*}
$$

where $C$ is a constant depending only on $\Sigma$ and $K$. For a lower bound, we use the fact that $L_{\psi}$ is elliptic.

Proposition 4.3 (Standard elliptic estimate). Let $D$ be an elliptic linear differential operator of order $r$ acting on sections of a vector bundle $V$ over $\Sigma$. Then there exist constants $\alpha_{k}, \beta_{k}$ depending only on $\Sigma$ and $k$, such that

$$
\|D Y\|_{k}+\alpha_{k}\|Y\|_{0} \geq \beta_{k}\|Y\|_{k+r}
$$

If we consider only sections which are $L^{2}$ orthogonal to $\operatorname{ker} D$, the same inequality holds with $\alpha_{k}=0$.

The reason for quoting this result in the context of a general vector bundle $V$ over $\Sigma$ is that we will want to apply it to both the ordinary Laplacian on $V=\underline{\mathbb{R}^{3}}$ and the classical Jacobi operator $J_{\psi}$ on $V=\psi^{-1} T S^{2} \subset \underline{\mathbb{R}}^{3}$, where the $H^{k}$ norm is defined by inclusion.

## 5. Local existence theorem

Theorem 5.1 (Local existence for the coupled system). Consider the coupled system (3.15) with initial data $q(0)=q_{0} \in K, q_{t}(0)=\varepsilon q_{1} \in \mathbb{R}^{4 n}, Y(0)=Y_{0} \in H^{3}$, $Y_{t}(0)=Y_{1} \in H^{2}$ such that

$$
\operatorname{dist}\left(q_{0}, \partial K\right)>d, \quad\left|q_{1}\right|,\left\|Y_{0}\right\|_{3},\left\|Y_{1}\right\|_{2}<\Gamma
$$

where $\Gamma, d$ are positive constants. Then there exist constants $C(K)>0$ and $T(K, \Gamma, d)>0$ such that for all $\varepsilon \in(0, C(K) / \sqrt{\Gamma})$, this initial value problem has a unique solution on $[0, T]$ with

$$
\begin{aligned}
q & \in C^{3}([0, T], K) \\
Y & \in C^{0}\left([0, T], H^{3}\right) \cap C^{1}\left([0, T], H^{2}\right) \cap C^{2}\left([0, T], H^{1}\right) .
\end{aligned}
$$

If the initial data are tangent to the $L^{2}$ orthogonality constraint (3.11) and the pointwise constraint (3.3), then the solution preserves these constraints.

We will prove this using Picard's method: we iteratively define a sequence $\left(q^{i}, Y^{i}\right) \in C^{0}\left([0, T], K \times H^{3}\right)$ which converges to a solution of the initial value problem. To establish that the iteration scheme is well-defined and convergent, the following standard energy estimate for the driven wave equation is key:

Theorem 5.2 (Existence and energy estimate for the wave equation). The driven wave equation on $[0, T] \times \Sigma$,

$$
Y_{t t}-\Delta Y=\Xi
$$

with $\Xi:[0, T] \times \Sigma \rightarrow \mathbb{R}^{3}$ smooth and smooth initial data $Y_{0}=Y(0)$ and $Y_{1}=Y_{t}(0)$ has a unique global solution. The solution is smooth, and there exists an absolute constant $c(\Sigma) \geq 1$, depending only on the choice of torus $\Sigma$, such that

$$
\max \left\{\left\|Y_{t}(t)\right\|_{2},\|Y(t)\|_{3}\right\} \leq c(\Sigma) e^{t}\left\{\left\|Y_{1}\right\|_{2}+\left\|Y_{0}\right\|_{3}+\left(\int_{0}^{t}\|\Xi(s)\|_{2}^{2} d s\right)^{\frac{1}{2}}\right\}
$$

Proof. Existence, uniqueness and smoothness follow from [7]. Let $E(t)=\left\|Y_{t}(t)\right\|_{0}^{2}+$ $\|Y(t)\|_{1}^{2}$. Then

$$
\begin{align*}
E^{\prime}(t) & =2\left\langle Y_{t}, \Delta Y+\Xi\right\rangle_{0}+2\left\langle Y, Y_{t}\right\rangle_{1}=2\left\langle Y_{t}, \Xi\right\rangle_{0}+2\left\langle Y_{t}, Y\right\rangle_{0} \\
& \leq 2\left\|Y_{t}\right\|_{0}^{2}+\|Y\|_{0}^{2}+\|\Xi\|_{0}^{2} \leq 2 E(t)+\|\Xi(t)\|_{0}^{2}  \tag{5.1}\\
\Rightarrow \quad \frac{d}{d t}\left(e^{-2 t} E(t)\right) & \leq e^{-2 t}\|\Xi(t)\|_{0}^{2} \leq\|\Xi(t)\|_{0}^{2} \\
\text { 1) } & \\
\Rightarrow E(t) & \leq e^{2 t}\left[E(0)+\int_{0}^{t}\|\Xi\|_{0}^{2}\right] .
\end{align*}
$$

Now consider $Z=\Delta Y$. This is also smooth and satisfies the wave equation with source $\Delta \Xi$. Applying the above estimate to $Z$ yields

$$
\begin{aligned}
\left\|\Delta Y_{t}\right\|_{0}^{2}+\|\Delta Y\|_{1}^{2} & \leq e^{2 t}\left[\left\|\Delta Y_{1}\right\|_{0}^{2}+\left\|\Delta Y_{0}\right\|_{1}^{2}+\int_{0}^{t}\|\Delta \Xi\|_{0}^{2}\right] \\
& \leq 2 e^{2 t}\left[\left\|Y_{1}\right\|_{2}^{2}+\left\|Y_{0}\right\|_{3}^{2}+\int_{0}^{t}\|\Xi\|_{2}^{2}\right]
\end{aligned}
$$

Since $\Delta$ is an elliptic operator, there exist positive constants $\alpha_{k}, \beta_{k}$ depending only on $k$ and $\Sigma$ such that

$$
\|\Delta Y\|_{k}^{2}+\alpha_{k}\|Y\|_{0}^{2} \geq \beta_{k}\|Y\|_{k+2}^{2}
$$

by the standard elliptic estimate, Proposition 4.3. The result immediately follows.

To prove Theorem5.1 we must first write the coupled system (3.15) as an explicit evolution system (note that both equations have $\ddot{q}$ on the right hand side). Let $X=$ $\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{3} \times H^{2}$ given the norm $\|(q, p, Y, Z)\|_{X}=\max \left\{|q|, \varepsilon^{-1}|p|,\|Y\|_{3},\|Z\|_{2}\right\}$. The $\varepsilon$ dependence of the norm is chosen so that $\left\|\left(0, q_{t}, 0,0\right)\right\|_{X}=\dot{q}$. Given any $\Gamma>0$ let $X_{\Gamma}=\{(q, p, Y, Z) \in X: q \in K,\|(0, p, Y, Z)\| \leq 8 c(\Sigma) \Gamma\}$ where $c(\Sigma) \geq 1$ is the absolute constant obtained from Theorem 5.2. Note that $X_{\Gamma}$ is a closed subset of a Banach space and hence is a complete metric space with respect to the metric induced by $\|\cdot\|_{X}$. Consider the matrix valued function $M: \mathbb{R} \times X_{\Gamma} \rightarrow \operatorname{End}\left(\mathbb{R}^{4 n}\right)$,

$$
\begin{equation*}
M(\varepsilon, q, Y)_{\nu}^{\mu}=\delta_{\nu}^{\mu}-\varepsilon^{2} \gamma^{\mu \lambda}\left\langle Y, \psi_{\lambda \nu}\right\rangle \tag{5.3}
\end{equation*}
$$

Since the matrix $\left(\gamma^{\mu \lambda}\right)$ is positive definite, $K$ is compact, all $q$-dependence is smooth, and $\|Y\|_{0} \leq 8 c(\Sigma) \Gamma$, there exists a constant $c(K)>0$ such that $M$ : $[0, c(K) / \sqrt{\Gamma}] \times X_{\Gamma} \rightarrow G L(4 n, \mathbb{R})$ and $M^{-1}:[0, c(K) / \sqrt{\Gamma}] \times X_{\Gamma} \rightarrow G L(4 n, \mathbb{R})$ is $C^{1}$ and bounded. Hence, for all $\varepsilon \in\left[0, \varepsilon_{*}(K, \Gamma)\right]$, where $\varepsilon_{*}=c(K) / \sqrt{\Gamma}$ the coupled system can be rewritten as

$$
\begin{align*}
q_{t t} & =\varepsilon^{2} f\left(\varepsilon, q, q_{t}, Y, Y_{t}\right)  \tag{5.4}\\
Y_{t t}-\Delta Y & =g\left(\varepsilon, q, q_{t}, Y, Y_{t}\right) \tag{5.5}
\end{align*}
$$

where

$$
\begin{align*}
f\left(\varepsilon, q, q_{t}, Y, Y_{t}\right) & =M^{-1}(\varepsilon, q, Y)\left(-G\left(q, q_{t}, q_{t}\right)+\varepsilon h\left(\varepsilon, q, \varepsilon^{-1} q_{t}, Y, Y_{t}\right)\right) \\
\left(5(q, u, v)^{\mu}\right. & =G_{\nu \lambda}^{\mu}(q) u^{\nu} v^{\lambda}  \tag{5.6}\\
(5.7) g\left(\varepsilon, q, q_{t}, Y, Y_{t}\right) & =-B_{\psi} Y-\psi_{\mu} f^{\mu}-\psi_{\mu \nu} \frac{q_{t}^{\mu}}{\varepsilon} \frac{q_{t}^{\nu}}{\varepsilon}+\varepsilon j^{\prime}\left(\varepsilon, q, \varepsilon^{-1} q_{t}, Y, Y_{t}\right)
\end{align*}
$$

$B_{\psi}$ denotes the first and zeroth order piece of $L_{\psi}$, so $L_{\psi}=-\Delta+B_{\psi}$, explicitly

$$
\begin{align*}
B_{\psi} Y= & -\left(\left|\psi_{x}\right|^{2}+\left|\psi_{y}\right|^{2}\right) Y-2\left(\psi_{x} \cdot Y_{x}+\psi_{y} \cdot Y_{y}\right) \psi-2(\psi \cdot Y) \Delta \psi \\
& -2(\psi \cdot Y)_{x} \psi_{x}-2(\psi \cdot Y)_{y} \psi_{y} \tag{5.8}
\end{align*}
$$

and $h$ and $j^{\prime}=j+\widehat{j}$ are as defined in (3.14), (3.5), (3.9).
It is convenient henceforth to consider $\varepsilon$ as a fixed parameter in $\left[0, \varepsilon_{*}(K, \Gamma)\right]$ and suppress the dependence of $f, g$ on $\varepsilon$. The proof of existence will use Picard's method, which requires that $f, g$ be bounded and Lipschitz on $X_{\Gamma}$. This follows quickly from the following proposition, whose proof is straightforward but lengthy, and so is deferred to the appendix.

Proposition 5.3. The functions $f, g$ are continuously differentiable maps $f$ : $X_{\Gamma} \rightarrow \mathbb{R}^{4 n}$ and $g: X_{\Gamma} \rightarrow H^{2}$. Their differentials $\mathrm{d} f: X_{\Gamma} \rightarrow \mathcal{L}\left(X, \mathbb{R}^{4 n}\right)$, $\mathrm{d} g: X_{\Gamma} \rightarrow \mathcal{L}\left(X, H^{2}\right)$ are bounded, uniformly in $\varepsilon$. That is, there exist constants $\Lambda_{f}(K, \Gamma), \Lambda_{g}(K, \Gamma)>0$ such that

$$
\left|\mathrm{d} f_{x} \omega\right| \leq \Lambda_{f}\|\omega\|_{X}, \quad\left\|\mathrm{~d} f_{x} \omega\right\|_{2} \leq \Lambda_{g}\|\omega\|_{X}
$$

for all $x \in X_{\Gamma}, \omega \in X$.
Note that for Banach spaces $B, C, \mathcal{L}(B, C)$ denotes the space of bounded linear maps $B \rightarrow C$, which is itself a Banach space with respect to the norm $\|S\|_{\mathcal{L}(B, C)}=$ $\sup \left\{\|S(x)\|_{C} /\|x\|_{B}: x \in B, x \neq 0\right\}$.
Corollary 5.4. The functions $f: X_{\Gamma} \rightarrow \mathbb{R}^{4 n}$ and $g: X_{\Gamma} \rightarrow H^{2}$ are Lipschitz and bounded, uniformly in $\varepsilon$. That is, there exist constants $\Lambda_{f}, \Lambda_{g}, C_{f}, C_{g}$, depending only on $K$ and $\Gamma$, such that for all $x, x^{\prime} \in X_{\Gamma}$,

$$
\begin{array}{rlrl}
\left|f(x)-f\left(x^{\prime}\right)\right| & \leq \Lambda_{f}\left\|x-x^{\prime}\right\|_{X}, & \left\|g(x)-g\left(x^{\prime}\right)\right\|_{2} & \leq \Lambda_{g}\left\|x-x^{\prime}\right\|_{X} \\
|f(x)| & \leq C_{f}, & \|g(x)\|_{2} \leq C_{g}
\end{array}
$$

Proof. Let $x_{1}, x_{2} \in X_{\Gamma}$. Since $X_{\Gamma}$ is convex, the curve $x(t)=x_{1}+t\left(x_{2}-x_{1}\right)$, $0 \leq t \leq 1$ remains in $X_{\Gamma}$. Hence

$$
\begin{align*}
\left|f\left(x_{1}\right)-f\left(x_{2}\right)\right| & =\left|\int_{0}^{1} \mathrm{~d} f_{x(t)}\left(x_{2}-x_{1}\right) d t\right| \\
& \leq \int_{0}^{1} \Lambda_{f}\left\|x_{2}-x_{1}\right\|_{X} d t=\Lambda_{f}\left\|x_{1}-x_{2}\right\|_{X} \tag{5.9}
\end{align*}
$$

From the definition of $f$ one sees that $f(q, 0,0,0)=0$ for all $q \in K$. Hence, for all $x=(q, p, Y, Z) \in X_{\Gamma}$,

$$
\begin{equation*}
|f(x)|=|f(q, p, Y, Z)-f(q, 0,0,0)| \leq \Lambda_{f}\|(0, p, Y, Z)\|_{X} \leq 8 c(\Sigma) \Gamma \Lambda_{f} \tag{5.10}
\end{equation*}
$$

The proof for $g$ follows mutatis mutandis.
To establish uniqueness of the solution and to show that $q$ is three times continuously differentiable, we will need the following extension property of $\mathrm{d} f$ and $\mathrm{d} g$, whose proof is also deferred to the appendix.

Proposition 5.5. The differentials $\mathrm{d} f: X_{\Gamma} \rightarrow \mathcal{L}\left(X, \mathbb{R}^{4 n}\right)$ and $\mathrm{d} g: X_{\Gamma} \rightarrow \mathcal{L}\left(X, H^{2}\right)$ of $f$ and $g$ extend continuously to maps $\mathrm{d} f^{e x t}: X_{\Gamma} \rightarrow \mathcal{L}\left(\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{1} \times L^{2}, \mathbb{R}^{4 n}\right)$ and $\mathrm{d} g^{e x t}: X_{\Gamma} \rightarrow \mathcal{L}\left(\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{1} \times L^{2}, L^{2}\right)$, bounded by $\Lambda_{f}$ and $\Lambda_{g}$ respectively.

The rest of this section is devoted to the proof of Theorem 5.1. Let $T>0$ be chosen such that

$$
\begin{align*}
& T \leq \log 2, \quad T \leq \frac{1}{\varepsilon}, \quad T \leq \frac{d}{\varepsilon\left(2 \Gamma+C_{f}(K, \Gamma)\right)}, \quad T \leq \frac{\Gamma}{\varepsilon C_{f}(K, \Gamma)}, \\
& \sqrt{T} \leq \frac{\Gamma}{C_{g}(K, \Gamma)}, \quad T \leq \frac{1}{4 \varepsilon \Lambda_{f}(K, \Gamma)}, \quad \sqrt{T} \leq \frac{1}{8 c(\Sigma) \Lambda_{g}} \tag{5.11}
\end{align*}
$$

Given a complete subset $B$ of a Banach space with norm $\|\cdot\|_{B}$, denote by $C_{T} B$ the space of continuous maps $[0, T] \rightarrow B$ equipped with the sup norm $\|b\| \|=$ $\sup \left\{\|b(t)\|_{B}: t \in[0, T]\right\} .\left(C_{T} B,\| \| \cdot \| \mid\right)$ is itself a complete subset of a Banach space.
5.1. Definition of the iteration scheme. We will produce a sequence $\omega^{i} \in$ $C_{T} X_{\Gamma}$ converging to a solution of the initial value problem for the coupled system. Choose and fix $\delta \in(0, \Gamma / 4)$, and let $Y_{0}^{i}, Y_{1}^{i} \in C^{\infty}\left(\Sigma, \mathbb{R}^{3}\right)$ be sequences such that

$$
\begin{equation*}
\left\|Y_{0}^{i}-Y_{0}\right\|_{3}<\frac{\delta}{2^{i}}, \quad\left\|Y_{1}^{i}-Y_{1}\right\|_{2}<\frac{\delta}{2^{i}} \tag{5.12}
\end{equation*}
$$

Such sequences exist since $C^{\infty}$ is dense in $H^{k}$ for all $k \geq 0$. Let $\omega^{0}=\left(q_{0}, \varepsilon q_{1}, Y_{0}^{0}, Y_{1}^{0}\right)$, which is constant in $t$ and smooth on $\Sigma$ and trivially lies in $C_{T} X_{\Gamma}$. Given $\omega^{i}$, we define the next iterate to be the solution of the initial value problem $q^{i+1}(0)=$ $q_{0}, q_{t}^{i+1}(0)=\varepsilon q_{1}, Y^{i+1}(0)=Y_{0}^{i+1}, Y_{t}^{i+1}(0)=Y_{1}^{i+1}$ for

$$
\begin{align*}
q_{t t}^{i+1} & =\varepsilon^{2} f\left(\omega^{i}\right)  \tag{5.13}\\
Y_{t t}^{i+1}-\Delta Y^{i+1} & =g\left(\omega^{i}\right) \tag{5.14}
\end{align*}
$$

We must first check that the sequence $\omega^{i}$ is well-defined. So, assume that $\omega^{i}$ is smooth and lies in $C_{T} X_{\Gamma}$. Then

$$
\begin{equation*}
q^{i+1}(t)=q_{0}+\varepsilon t q_{1}+\varepsilon^{2} \int_{0}^{t}\left(\int_{0}^{s} f\left(\omega^{i}(r)\right) d r\right) d s \tag{5.15}
\end{equation*}
$$

which exists since $f \circ \omega^{i}$ is continuous. Now

$$
\begin{equation*}
\left|q^{i+1}(t)-q_{0}\right| \leq \varepsilon T\left|q_{1}\right|+\frac{1}{2} \varepsilon^{2} T^{2} C_{f}(K, \Gamma) \leq \varepsilon T\left(\Gamma+\frac{C_{f}(K, \Gamma)}{2}\right) \leq \frac{d}{2} \tag{5.16}
\end{equation*}
$$

by our choice of $T$, so $q^{i+1}(t)$ remains in $K$. Further,

$$
\begin{equation*}
\varepsilon^{-1}\left|q_{t}^{i+1}(t)\right| \leq\left|q_{1}\right|+\varepsilon \int_{0}^{t}\left|f\left(\omega^{i}(s)\right)\right| d s \leq \Gamma+T \varepsilon C_{f}(K, \Gamma) \leq 2 \Gamma<8 c(\Sigma) \Gamma \tag{5.17}
\end{equation*}
$$

by our choice of $T$. Turning to $Y^{i+1}$, we see by inspection that if $Y^{i}, q^{i}$ are smooth, then $g\left(\omega^{i}\right)$ is smooth, so the solution $Y^{i+1}(t)$ exists, is unique and smooth, by Theorem 5.2] which also yields the energy estimate

$$
\max \left\{\left\|Y_{t}^{i+1}(t)\right\|_{2},\left\|Y^{i+1}(t)\right\|_{3}\right\} \leq C(\Sigma) e^{t}\left\{\left\|Y_{0}^{i+1}\right\|_{3}+\left\|Y_{1}^{i+1}\right\|_{2}+\sqrt{t} C_{g}(K, \Gamma)\right\}
$$

$$
\begin{equation*}
\leq 2 C(\Sigma)\left\{2(\Gamma+\delta)+\sqrt{T} C_{g}(K, \Gamma)\right\}<8 C(\Sigma) \Gamma \tag{5.18}
\end{equation*}
$$

by our choice of $T$ and $\delta$. Hence, if $\omega^{i}$ is smooth and in $C_{T} X_{\Gamma}$, so is $\omega^{i+1}$. We have already observed that $\omega^{0}$ is smooth and in $C_{T} X_{\Gamma}$, so, by induction, the sequence $\omega^{i} \in C_{T} X_{\Gamma}$ is well-defined.
5.2. Convergence of the iteration scheme. We will now show that $\omega^{i}$ is Cauchy, and hence converges in $C_{T} X_{\Gamma}$. From (5.15) one has

$$
\begin{align*}
\left|q^{i+1}(t)-q^{i}(t)\right| & =\varepsilon^{2}\left|\int_{0}^{t} \int_{0}^{s}\left(f\left(\omega^{i}(r)\right)-f\left(\omega^{i-1}(r)\right)\right) d r d s\right| \\
& \leq \frac{\varepsilon^{2}}{2} T^{2} \Lambda_{f}(K, \Gamma)\left|\| \omega ^ { i } - \omega ^ { i - 1 } \| \left\|\leq \frac{1}{8}\left|\left\|\omega^{i}-\omega^{i-1}\right\|\right|\right.\right. \tag{5.19}
\end{align*}
$$

by our choice of $T$. Similarly

$$
\begin{align*}
\varepsilon^{-1}\left|q_{t}^{i+1}(t)-q_{t}^{i}(t)\right| & =\varepsilon\left|\int_{0}^{t}\left(f\left(\omega^{i}(r)\right)-f\left(\omega^{i-1}(r)\right)\right) d r\right| \\
& \leq \varepsilon T \Lambda_{f}(K, \Gamma)\left|\| \omega ^ { i } - \omega ^ { i - 1 } \| \left\|\leq \frac{1}{4}\left|\left\|\omega^{i}-\omega^{i-1}\right\|\right|\right.\right. \tag{5.20}
\end{align*}
$$

Now $Z=Y^{i+1}-Y^{i}$ satisfies the wave equation with source $g\left(\omega^{i}\right)-g\left(\omega^{i-1}\right)$ and small smooth initial data $\|Z(0)\|_{3},\left\|Z_{t}(0)\right\|_{2} \leq \delta / 2^{i-1}$. Hence, by Theorem 5.2, for each $t \in[0, T]$,

$$
\begin{aligned}
\max \left\{\|Z(t)\|_{3},\left\|Z_{t}(t)\right\|_{2}\right\} & \leq c(\Sigma) e^{t}\left\{\frac{\delta}{2^{i-2}}+\left(\int_{0}^{t}\left\|g\left(\omega^{i}(s)\right)-g\left(\omega^{i-1}(s)\right)\right\|_{2}^{2} d s\right)^{\frac{1}{2}}\right\} \\
& \leq 2 c(\Sigma)\left\{\left.\frac{\delta}{2^{i-2}}+\sqrt{T} \Lambda_{g}(K, \Gamma) \right\rvert\,\left\|\omega^{i}-\omega^{i-1}\right\| \|\right\} \\
(5.21) & \left.\leq \frac{c(\Sigma) \delta}{2^{i-3}}+\frac{1}{4} \right\rvert\,\left\|\omega^{i}-\omega^{i-1}\right\| \|
\end{aligned}
$$

Assembling these inequalities, one sees that

$$
\begin{equation*}
\left|\left\|\omega^{i+1}-\omega^{i}\right\|\right| \leq \frac{1}{4}\left|\left\|\omega^{i}-\omega^{i-1}\right\|\right|+\frac{\alpha}{2^{i}} \tag{5.22}
\end{equation*}
$$

where $\alpha=c(\Sigma) \delta / 8$. It follows that

$$
\begin{equation*}
\left|\left\|\omega^{i+1}-\omega^{i}\right\|\right| \leq \frac{1}{4^{i}}\left|\left\|\omega^{1}-\omega^{0}\right\|\right|+\frac{\alpha}{2^{i-1}} \tag{5.23}
\end{equation*}
$$

and hence, for all $k \geq 1$,

$$
\begin{align*}
\left|\left\|\omega^{i+k}-\omega^{i}\right\|\right| & \leq \sum_{j=1}^{k}\left|\left\|\omega^{i+j}-\omega^{i+j-1}\right\|\right| \\
& \leq \frac{1}{4^{i}}\left|\left\|\omega^{1}-\omega^{0}\right\|\right| \sum_{j=1}^{\infty} \frac{1}{4^{j}}+\frac{\alpha}{2^{i}} \sum_{j=0}^{\infty} \frac{1}{2^{j}} . \tag{5.24}
\end{align*}
$$

Hence $\omega^{i}$ is Cauchy with respect to $\left\|\|\cdot\| \mid\right.$, so $\omega^{i} \rightarrow \omega=(q, p, Y, Z) \in C_{T} X_{\Gamma}$.
5.3. The limit solves the initial value problem. We have established that

$$
\begin{array}{rllll}
q^{i} & \rightarrow q \quad \text { in } & C_{T} K \\
q_{t}^{i} & \rightarrow p \text { in } & C_{T} \mathbb{R}^{4 n} \\
Y^{i} & \rightarrow Y \text { in } & C_{T} H_{3} \\
Y_{t}^{i} & \rightarrow Z & \text { in } & C_{T} H^{2} \tag{5.28}
\end{array}
$$

Now, for all $i$,

$$
\begin{align*}
\left\|\omega(0)-\left(q_{0}, \varepsilon q_{1}, Y_{0}, Y_{1}\right)\right\|_{X} & \leq\left\|\omega(0)-\omega^{i}(0)\right\|_{X}+\left\|\omega^{i}(0)-\left(q_{0}, \varepsilon q_{1}, Y_{0}, Y_{1}\right)\right\|_{X} \\
& \leq\left\|\omega-\omega^{i}\right\| \left\lvert\,+\frac{\delta}{2^{i}} \rightarrow 0\right. \tag{5.29}
\end{align*}
$$

as $i \rightarrow \infty$. Hence $\omega(0)=\left(q_{0}, \varepsilon q_{1}, Y_{0}, Y_{1}\right)$; that is, the limit has the correct initial data.

We will now show that the limit solves the coupled system and has the differentiability properties claimed. Let $\tilde{Y}(t)=Y_{0}+\int_{0}^{t} Z(s) d s$. Note that $\widetilde{Y}$ is manifestly in $C^{1}\left([0, T], H^{2}\right)$, with derivative $\widetilde{Y}_{t}=Z$. Now

$$
\begin{align*}
\mid\left\|Y^{i}-\widetilde{Y}\right\|_{C_{T} H^{2}} & =\mid\left\|Y^{i}(0)-Y_{0}+\int_{0}^{t}\left(Y_{t}^{i}(s)-Z(s)\right) d s\right\| \|_{C_{T} H^{2}} \\
& \leq \frac{\delta}{2^{i}}+T\left|\left\|Y_{t}^{i}-Z\right\|\right|_{C_{T} H^{2}} \rightarrow 0 \tag{5.30}
\end{align*}
$$

as $i \rightarrow \infty$. Hence, $Y^{i} \rightarrow \widetilde{Y}$ in $C_{T} H^{2}$. But $Y^{i} \rightarrow Y$ in $C_{T} H^{3}$, hence also in $C_{T} H^{2}$, so $Y=\widetilde{Y}$. Hence, $Y \in C^{1}\left([0, T], H^{2}\right)$ and $Y_{t}=Z$.

Consider $Y_{t t}^{i+k}-Y_{t t}^{i}$. This is smooth and satisfies the wave equation with source $g\left(\omega^{i+k-1}\right)-g\left(\omega^{i-1}\right)$. Hence

$$
\begin{aligned}
\left\|Y_{t t}^{i+k}-Y_{t t}^{i}\right\|_{1} & =\left\|\Delta\left(Y^{i+k}-Y^{i}\right)+g\left(\omega^{i+k-1}\right)-g\left(\omega^{i-1}\right)\right\|_{1} \\
& \leq 2\left\|Y^{i+k}-Y^{i}\right\|_{3}+\left\|g\left(\omega^{i+k-1}\right)-g\left(\omega^{i-1}\right)\right\|_{2} \\
& \left.\leq 2\left\|Y^{i+k}-Y^{i}\right\|_{3}+\Lambda_{g}(K, \Gamma) \| \omega^{i+k-1}-\omega^{i-1}\right) \|_{X} \\
(5.31) \Rightarrow \mid\left\|Y_{t t}^{i+k}-Y_{t t}^{i}\right\|_{C_{T} H^{1}} & \leq\left.\left[2+\Lambda_{g}(K, \Gamma)\right]\left\|\omega^{i+k-1}-\omega^{i-1}\right\|\right|_{C_{T} X_{\Gamma}} .
\end{aligned}
$$

Since $\omega^{i}$ is Cauchy in $C_{T} X_{\Gamma}$, it follows that $Y_{t t}^{i}$ is Cauchy in $C_{T} H^{1}$. Hence $Y_{t t}^{i} \rightarrow W$ in $C_{T} H^{1}$. Let $\widetilde{Z}(t)=Y_{1}+\int_{0}^{t} W$. Note that $\widetilde{Z}$ is manifestly in $C^{1}\left([0, T], H^{1}\right)$ and $\widetilde{Z}_{t}=W$. Now

$$
\begin{align*}
\left|\left\|Y_{t}^{i}-\widetilde{Z}\right\|\right|_{C_{T} H^{1}} & \leq\left|\left\|Y_{1}^{i}-Y_{1}\right\|\right|_{C_{T} H^{1}}+\left|\left\|\int_{0}^{t}\left(Y_{t t}^{i}-W\right)\right\|\right|_{C_{T} H^{1}} \\
& \leq \frac{\delta}{2^{i}}+T\left|\left\|Y_{t t}^{i}-W\right\|\right|_{C_{T} H^{1}} \rightarrow 0 \tag{5.32}
\end{align*}
$$

as $i \rightarrow \infty$. Hence $Y_{t}^{i} \rightarrow \widetilde{Z}$ in $C_{T} H^{1}$. But $Y_{t}^{i} \rightarrow Z$ in $C_{T} H^{2}$, hence also in $C_{T} H^{1}$, so $Z=\widetilde{Z}$. But $Y_{t}=Z$. Hence, $Y \in C^{2}\left([0, T], H^{1}\right)$ and $Y_{t t}=W$.

By similar reasoning, $q_{t t}^{i} \rightarrow m$ in $C_{T} \mathbb{R}^{4 n}$ and $q \in C^{2}\left([0, T], \mathbb{R}^{4 n}\right)$ with $q_{t}=p$ and $q_{t t}=m$.

We can now show that $\omega$ solves the coupled system:

$$
\begin{align*}
\mid \| Y_{t t} & -\Delta Y-g(\omega) \|\left.\right|_{C_{T} H^{1}} \\
& \leq\left|\left\|Y_{t t}-Y_{t t}^{i}\right\|\right|_{C_{T} H^{1}}+\left|\left\|\Delta\left(Y-Y^{i}\right)\right\|\right|_{C_{T} H^{1}}+\mid\left\|g(\omega)-g\left(\omega^{i-1}\right)\right\| \|_{C_{T} H^{1}} \\
& \leq\left|\left\|Y_{t t}-Y_{t t}^{i}\right\|\right|_{C_{T} H^{1}}+2\left|\| Y - Y ^ { i } \| \left\|_{C_{T} H^{3}}+\left|\left\|g(\omega)-g\left(\omega^{i-1}\right)\right\|\right|_{C_{T} H^{2}} .\right.\right. \tag{5.33}
\end{align*}
$$

Now $Y_{t t}^{i} \rightarrow Y_{t t}$ in $C_{T} H^{1}, Y^{i} \rightarrow Y$ in $C_{T} H^{3}, g: X_{\Gamma} \rightarrow H^{2}$ is continuous, and $\omega^{i} \rightarrow \omega$ in $C_{T} X_{\Gamma}$, so $g\left(\omega^{i}\right) \rightarrow g(\omega)$ in $C_{T} H^{2}$. Hence

$$
\begin{equation*}
\left|\left\|Y_{t t}-\Delta Y-g(\omega)\right\|\right|_{C_{T} H^{1}}=0 \tag{5.34}
\end{equation*}
$$

Similarly $\left\|q_{t t}-\varepsilon^{2} f(\omega)\right\|_{C_{T} \mathbb{R}^{4 n}}=0$, that is, $q_{t t}=\varepsilon^{2} f(\omega)$.

It remains to establish the higher differentiability of $q$. Differentiating the equation for $q_{t t}^{i+1}$ gives

$$
\begin{equation*}
q_{t t t}^{i+1}=\varepsilon^{2} d f_{\omega^{i}}^{e x t} \omega_{t}^{i} \tag{5.35}
\end{equation*}
$$

Now $\omega^{i} \rightarrow \omega$ in $C_{T} X_{\Gamma}, \omega_{t}^{i} \rightarrow \omega_{t}$ in $C_{T}\left(\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{2} \times H^{1}\right)$, and $d f^{e x t}$ is continuous, so $q_{t t t}^{i+1} \rightarrow \ell$, say, in $C_{T} \mathbb{R}^{4 n}$. Let $\widetilde{m}(t)=\varepsilon^{2} f\left(q_{0}, \varepsilon q_{1}, Y_{0}, Y_{1}\right)+\int_{0}^{t} \ell=q_{t t}(0)+\int_{0}^{t} \ell$. Note that $\widetilde{m} \in C^{1}\left([0, T], \mathbb{R}^{4 n}\right)$ and $\widetilde{m}_{t}=\ell$. Then

$$
\begin{equation*}
\left\|q_{t t}^{i}-\widetilde{m}\right\|_{C_{T} \mathbb{R}^{4 n}}=\left\|\int_{0}^{t}\left(q_{t t t}^{i}-\ell\right)\right\|_{C_{T} \mathbb{R}^{4 n}} \leq T\left\|q_{t t t}^{i}-\ell\right\|_{C_{T} \mathbb{R}^{4 n}} \rightarrow 0 \tag{5.36}
\end{equation*}
$$

so $q_{t t}^{i} \rightarrow \widetilde{m}$ in $C_{T} \mathbb{R}^{4 n}$. But $q_{t t}^{i} \rightarrow q_{t t}$ in $C_{T} \mathbb{R}^{4 n}$, so $q_{t t}=\widetilde{m}$. Hence, $q_{t t} \in$ $C^{1}\left([0, T], \mathbb{R}^{4 n}\right)$, as claimed.
5.4. Uniqueness of the solution. Assume that $(\widetilde{q}, \widetilde{Y})$ is another solution of (5.4), (5.5) with the same initial data and regularity as $(q, Y)$, and let $(p, Z)=(q-\widetilde{q}$, $Y-\widetilde{Y})$. Then $(p, Z)$ satisfies the system

$$
\begin{equation*}
Z_{t t}-\Delta Z=\Xi(t), \quad p_{t t}=\varepsilon^{2} \Upsilon(t) \tag{5.37}
\end{equation*}
$$

with initial data $Z(0)=Z_{t}(0)=0, p(0)=p_{t}(0)=0$, where

$$
\begin{equation*}
\Xi(t)=g(\omega)-g(\widetilde{\omega}), \quad \Upsilon(t)=f(\omega)-f(\widetilde{\omega}) \tag{5.38}
\end{equation*}
$$

and $\omega=\left(q, q_{t}, Y, Y_{t}\right), \widetilde{\omega}=\left(\widetilde{q}, \widetilde{q} t, \widetilde{Y}, \widetilde{Y}_{t}\right)$. Define

$$
\begin{equation*}
E(t)=\|Z(t)\|_{1}^{2}+\left\|Z_{t}(t)\right\|_{0}^{2}+|p|^{2}+\frac{1}{\varepsilon^{2}}\left|p_{t}\right|^{2} \tag{5.39}
\end{equation*}
$$

which, by the regularity properties of $(p, Z)$, is continuously differentiable and has $E(0)=0$. Reprising the argument in (5.1), which requires only that $Z \in H^{2}$ and $Z_{t} \in H^{1}$, one sees that

$$
\begin{equation*}
E^{\prime}(t) \leq 2 E(t)+\|\Xi(t)\|_{0}^{2}+|\Upsilon(t)|^{2} \tag{5.40}
\end{equation*}
$$

Now, arguing as in the proof of Corollary 5.4, with $\omega(s)=\omega+s(\widetilde{\omega}-\omega)$,

$$
\begin{align*}
\|\Xi(t)\|_{0} & =\left\|g(\omega)-\left.g(\widetilde{\omega})\right|_{0}=\right\| \int_{0}^{1} \mathrm{~d} g_{\omega(s)}(\widetilde{\omega}-\omega) d s \|_{0} \\
& \leq \Lambda_{g} \max \left\{|p|, \varepsilon^{-1}\left|p_{t}\right|,\|Z\|_{1},\left\|Z_{t}\right\|_{0}\right\} \tag{5.41}
\end{align*}
$$

by Proposition 5.5. Similarly $|\Upsilon(t)| \leq \Lambda_{f} \max \left\{|p|, \varepsilon^{-1}\left|p_{t}\right|,\|Z\|_{1},\left\|Z_{t}\right\|_{0}\right\}$. Hence

$$
\begin{equation*}
E^{\prime}(t) \leq \kappa E(t) \tag{5.42}
\end{equation*}
$$

where $\kappa=2+\Lambda_{f}+\Lambda_{g}$, whence it follows that

$$
\begin{equation*}
\frac{d}{d t} e^{-\kappa t} E(t) \leq 0 \tag{5.43}
\end{equation*}
$$

So $e^{-\kappa t} E(t)$ is a nonincreasing, nonnegative function which is zero at $t=0$. Hence $E(t)=0$ for all $t$, and we conclude that $(p, Z)=(0,0)$ for all $t$, that is, $(q, Y)=$ $(\widetilde{q}, \widetilde{Y})$.
5.5. Preservation of constraints. Given the solution $(q, Y)$ produced above, define for each $\mu \in\{1,2, \ldots, 4 n\}$,

$$
\begin{equation*}
a_{\mu}(t)=\left\langle Y, \frac{\partial \psi}{\partial q^{\mu}}\right\rangle . \tag{5.44}
\end{equation*}
$$

The $L^{2}$ orthogonality constraint is that $a_{\mu}(t)=0$ for all $\mu, t$. By construction, the coupled system implies that $\ddot{a}_{\mu}=0$. If the initial data are tangent to the constraint, then $a_{\mu}(0)=\dot{a}_{\mu}(0)$, and hence $a_{\mu}(t)=0$ for all $t \in[0, T]$.

Similarly, given the solution $(q, Y)$ produced above, define $\chi:[0, T] \times \Sigma \rightarrow \mathbb{R}$ by

$$
\begin{equation*}
\chi=Y \cdot \psi(q)+\frac{1}{2} \varepsilon^{2}|Y|^{2} . \tag{5.45}
\end{equation*}
$$

The pointwise constraint is that $\chi=0$ everywhere on $\Sigma$. Note that $\chi(t) \in \mathrm{H}^{3}$ for all $t \in[0, T]$ so $\chi(t): \Sigma \rightarrow \mathbb{R}$ is continuous. Assume that $\chi(0)=0$ and $\chi_{t}(0)=0$; that is, the initial data are tangent to the constraint. A straightforward, if lengthy, calculation using the coupled system and the harmonic map equation for $\psi$ shows that $\chi$ satisfies the linear PDE

$$
\begin{align*}
\chi_{t t}-\Delta \chi= & 2 \varepsilon^{2}\left\{\left[2 \psi_{x} \cdot Y_{x}+2 \psi_{y} \cdot Y_{y}-\left|\psi_{\tau}\right|^{2}+Y \cdot \Delta \psi-2 \varepsilon \psi_{\tau} \cdot Y_{t}\right.\right. \\
& \left.\left.-\varepsilon^{4}\left(\left|Y_{t}\right|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right)\right] \chi+\left(Y \cdot \psi_{x}\right) \chi_{x}+\left(Y \cdot \psi_{y}\right) \chi_{y}\right\} \\
6) & a \chi+b_{1} \chi_{x}+b_{2} \chi_{y} \tag{5.46}
\end{align*}
$$

where $a(t) \in H^{2}, b_{1}(t), b_{2}(t) \in H^{3}$ for all $t$. Let $E(t)=\|\chi(t)\|_{1}^{2}+\left\|\chi_{t}(t)\right\|_{0}^{2}$. Then

$$
\begin{align*}
E^{\prime}(t) & =2\left\langle\chi_{t}, \Delta \chi+a \chi+b_{1} \chi_{x}+b_{2} \chi_{y}\right\rangle_{0}+2\left\langle\chi_{t}, \chi\right\rangle_{1} \\
& =2\left\langle\chi_{t}, \chi+a \chi+b_{1} \chi_{x}+b_{2} \chi_{y}\right\rangle_{0} \\
& \leq\left\|\chi_{t}\right\|_{0}^{2}+\left\|\chi+a \chi+b_{1} \chi_{x}+b_{2} \chi_{y}\right\|_{0}^{2} \leq\left\|\chi_{t}\right\|_{0}^{2}+\kappa\|\chi\|_{1}^{2} \leq \kappa E(t) \tag{5.47}
\end{align*}
$$

where

$$
\begin{equation*}
\kappa=\max _{0 \leq t \leq T} 8\left(1+\|a(t)\|_{2}^{2}+\left\|b_{1}(t)\right\|_{2}^{2}+\left\|b_{2}(t)\right\|_{2}^{2}\right) \tag{5.48}
\end{equation*}
$$

Hence $e^{-\kappa t} E(t)$ is a nonincreasing, nonnegative function which is zero at $t=0$, so $E(t)=0$, whence $\|\chi(t)\|_{1}=0$ for all $t$. Since we already know that $\chi(t): \Sigma \rightarrow$ $\mathbb{R}$ is continuous, it follows that $\chi=0$ everywhere. This completes the proof of Theorem 5.1

## 6. Near coercivity of the improved Hessian

By repeatedly applying the local existence theorem, we can extend the solution of the coupled system whenever $q$ remains in $K$ and $\varepsilon^{-1} q_{t}$, and $\|Y\|_{3}$ and $\left\|Y_{t}\right\|_{2}$ remain bounded. So to prove long time existence, we must, among other things, bound the growth of $\|Y\|_{3}$. The first step is to show that $\|Y\|_{3}$ is controlled by the quadratic form $\langle L Y, L L Y\rangle$ or, more precisely, by the quadratic form $Q_{2}: H^{3} \rightarrow \mathbb{R}$ defined next.

Definition 6.1. For a fixed harmonic map $\psi(q)$ we denote by $Q_{1}, Q_{2}$ the quadratic forms

$$
\begin{aligned}
& Q_{1}: H^{1} \rightarrow \mathbb{R}, \quad Q_{1}(Y)=\int_{\Sigma}\left\{\left|Y_{x}\right|^{2}+\left|Y_{y}\right|^{2}-\left(\left|\psi_{x}\right|^{2}+\left|\psi_{y}\right|^{2}\right)|Y|^{2}\right. \\
&\left.-4\left(\psi_{x} \cdot Y_{x}+\psi_{y} \cdot Y_{y}\right) \psi \cdot Y\right\}
\end{aligned}
$$

$Q_{2}: H^{3} \rightarrow \mathbb{R}, \quad Q_{2}(Y)=Q_{1}(L Y)$.
Note that both $Q_{1}$ and $Q_{2}$ are continuous and that $Q_{1}(Y)=\langle Y, L Y\rangle$ for all $Y \in H^{2}$. It is also convenient to define the projection map $P: H^{k} \rightarrow H^{k}$,

$$
P(Y)=Y-(\psi \cdot Y) \psi
$$

which pointwise orthogonally projects $Y(p)$ to $T_{\psi(p)} S^{2}$.
Lemma 6.2. For all $Y \in H^{2}, Q_{1}(Y) \geq Q_{1}(P(Y))$.
Proof. $Y=P(Y)+f \psi$, where $f=-\psi \cdot Y \in \mathrm{H}^{2}$. Now

$$
\begin{align*}
Q_{1}(Y) & =\langle Y, L Y\rangle=\langle P(Y), L P(Y)\rangle+2\langle f \psi, L P(Y)\rangle+\langle f \psi, L(f \psi)\rangle \\
& =\langle P(Y), L P(Y)\rangle+\langle f \psi, L(f \psi)\rangle \tag{6.1}
\end{align*}
$$

since $L$ is self-adjoint and maps tangent sections to tangent sections. But, as we saw in Remark 3.2

$$
\begin{equation*}
L(f \psi)=-(\Delta f) \psi-4\left(f_{x} \psi_{x}+f_{y} \psi_{y}\right) \tag{6.2}
\end{equation*}
$$

so $\langle f \psi, L(f \psi)\rangle=-\langle f, \Delta f\rangle \geq 0$.
If our error section $Y$ were a tangent section, the results of [5] would immediately imply that $Q_{1}$ is coercive, that is, $Q_{1}(Y) \geq c(q)\|Y\|_{1}^{2}$, orthogonal to ker $J$ :
Theorem 6.3 (Haskins-Speight [5]). There exists a constant $c(q)>0$, depending continuously on $q$, such that

$$
Q_{1}(Y) \geq c(q)\|Y\|_{1}^{2}
$$

for all $Y \in H^{1}$ satisfying $\psi \cdot Y=0, L^{2}$ orthogonal to $\operatorname{ker} J_{\psi(q)}$.
Unfortunately, $\psi \cdot Y \neq 0$ in our set-up, but is small (of order $\varepsilon^{2}$ ). This means we can only establish the following "near coercivity" property for $Q_{1}$. This will suffice for our purposes, however.
Theorem 6.4 (Near coercivity of $Q_{1}$ ). There exist constants $c(q), \widetilde{c}(q)>0$, depending continuously on $q$, such that

$$
Q_{1}(Y) \geq c(q)\|Y\|_{1}^{2}-\varepsilon^{2} \widetilde{c}(q)\|Y\|_{1}\|Y\|_{2}^{2}
$$

for all $Y \in H^{2}$ satisfying the pointwise constraint (3.3), $L^{2}$ orthogonal to $\operatorname{ker} J_{\psi(q)}$.
Proof. By Lemma 6.2, $Q_{1}(Y) \geq\langle P(Y), L P(Y)\rangle=\langle P(Y), J P(Y)\rangle$ since $L \equiv J$ on tangent sections. Given any $Z \in \operatorname{ker} J,\langle Z, P(Y)\rangle=\langle Z, Y+f \psi\rangle=\langle Z, Y\rangle=0$, since $Z$ is pointwise orthogonal to $\psi$ and $Y$ is $L^{2}$ orthogonal to ker $J$. Hence $P(Y)$ is $L^{2}$ orthogonal to ker $J$, and so, by Theorem 6.3, there exists a constant $\widetilde{c}(q)>0$ such that

$$
\begin{equation*}
Q_{1}(Y) \geq \widetilde{c}(q)\|P(Y)\|_{1}^{2} \tag{6.3}
\end{equation*}
$$

Now, since $Y$ satisfies (3.3),

$$
\begin{equation*}
\left.\|P(Y)\|_{1}^{2}=\left\|Y+\frac{1}{2} \varepsilon^{2}|Y|^{2} \psi\right\|_{1}^{2} \geq\|Y\|_{1}^{2}-\left.\varepsilon^{2}\langle Y,| Y\right|^{2} \psi\right\rangle_{1} \tag{6.4}
\end{equation*}
$$

and, by the algebra property of $\mathrm{H}^{2}$ (Proposition 4.1),

$$
\begin{equation*}
\left.\left.\langle Y,| Y\right|^{2} \psi\right\rangle_{1} \leq\|Y\|_{1}\left\||Y|^{2} \psi\right\|_{2} \leq C\|Y\|_{1}\|Y\|_{2}^{2}\|\psi\|_{2} \tag{6.5}
\end{equation*}
$$

where $C>0$ is a constant depending only on $\Sigma$. Combining (6.3), (6.4) and (6.5), and noting that $\|\psi\|_{2}$ depends continuously (in fact smoothly) on $q$, the result immediately follows.

Theorem 6.5 (Near coercivity of $Q_{2}$ ). There exist constants $c(q), \widetilde{c}(q)>0$, depending continuously on $q$, such that

$$
Q_{2}(Y) \geq c(q)\|Y\|_{3}^{2}-\varepsilon^{2} \widetilde{c}(q)\left(\|Y\|_{3}^{3}+\varepsilon^{2}\|Y\|_{3}^{4}\right)
$$

for all $Y \in H^{3}$ satisfying the pointwise constraint (3.3), $L^{2}$ orthogonal to $\operatorname{ker} J_{\psi(q)}$.
Proof. Recall that the pointwise constraint (3.3) is equivalent to $\left|\psi+\varepsilon^{2} Y\right| \equiv 1$. The set of smooth maps $\Sigma \rightarrow S^{2}$ is dense in the Banach manifold of $H^{3}$ maps $\Sigma \rightarrow S^{2}$, $\psi$ and ker $J$ are smooth, and $Q_{2}: H^{3} \rightarrow \mathbb{R}$ is continuous, so it suffices to prove the inequality in the case that $Y$ is smooth. So, let $Y$ be smooth, $L^{2}$ orthogonal to $\operatorname{ker} J$ and satisfying the pointwise constraint (3.3). Define the smooth section $Z=L Y$ and the smooth real functions $\alpha=|Y|^{2}$ and $\beta=Z \cdot Y$. Since $Y$ satisfies (3.3) and $\psi$ is harmonic, it follows that

$$
\begin{equation*}
\beta=\varepsilon^{2}\left\{Y \cdot \Delta Y+\left|Y_{x}\right|^{2}+\left|Y_{y}\right|^{2}+2\left(\left|\psi_{x}\right|^{2}+\left|\psi_{y}\right|^{2}\right)|Y|^{2}\right\} \tag{6.6}
\end{equation*}
$$

In the following, $c_{1}(q), c_{2}(q), \ldots$ denotes positive functions depending continuously on $q$. We have the following elementary estimate:

$$
\begin{equation*}
\|\beta \psi\|_{1}^{2} \leq c_{1}(q)\|\beta\|_{1}^{2} \leq \varepsilon^{4} c_{2}(q)\left\{\|Y\|_{C^{1}}^{2}\|Y\|_{2}^{2}+\|Y\|_{C^{0}}^{2}\|Y\|_{3}^{2}\right\} \tag{6.7}
\end{equation*}
$$

Applying the Sobolev inequalities (Proposition 4.2) gives

$$
\begin{equation*}
\|\beta \psi\|_{1} \leq \varepsilon^{2} c_{3}(q)\|Y\|_{3}^{2} \tag{6.8}
\end{equation*}
$$

We will also need to estimate $\|\alpha\|_{3}$. Again, we have an elementary estimate:

$$
\begin{equation*}
\|\alpha\|_{3}^{2} \leq c\left\{\|Y\|_{C^{1}}^{2}\|Y\|_{2}^{2}+\|Y\|_{C^{0}}^{2}\|Y\|_{3}^{2}\right\} \tag{6.9}
\end{equation*}
$$

which, on appealing to Proposition 4.2, yields

$$
\begin{equation*}
\|\alpha\|_{3} \leq c\|Y\|_{3}^{2} \tag{6.10}
\end{equation*}
$$

By Lemma 6.2

$$
\begin{equation*}
Q_{2}(Y)=Q_{1}(Z) \leq Q_{1}(P(Z))=\langle P(Z), J P(Z)\rangle \tag{6.11}
\end{equation*}
$$

It is in the last step that we have used the smoothness of $Y$ (assuming only $Y \in H^{3}$ gives $P(Z) \in H^{1}$, which is not sufficiently regular to make sense of $J P(Z)$ ). Since $L$ is self-adjoint and $\operatorname{ker} J \subset \operatorname{ker} L, Z=L Y$ is automatically $L^{2}$ orthogonal to ker $J$, as is $P(Z)=Z+\beta \psi$ (since $\psi$ is pointwise orthogonal to anything in ker $J$ ). Hence, by Theorem 6.3 and the estimate (6.8),
$Q_{2}(Y) \geq c_{4}(q)\|P(Z)\|_{1}^{2}=c_{4}(q)\|Z-\beta \psi\|_{1}^{2} \geq c_{4}(q)\left\{\|Z\|_{1}^{2}-\|Z\|_{1}\|\beta \psi\|_{1}\right\}$
(6.12) $\geq c_{4}(q)\left\{\|Z\|_{1}^{2}-\varepsilon^{2} c_{3}(q)\|Y\|_{3}^{2}\|Z\|_{1}\right\}$.

We next estimate $\|Z\|_{1}=\|L Y\|_{1}$ in terms of $\|Y\|_{3}$. Note that $Y$ is not $L^{2}$ orthogonal to $\operatorname{ker} L$, since it has a component in the direction of $\psi$, so we cannot apply the standard elliptic estimate for $L$ directly. We must decompose

$$
\begin{equation*}
Y=P(Y)+(\psi \cdot Y) \psi=P(Y)-\frac{1}{2} \varepsilon^{2} \alpha \psi \tag{6.13}
\end{equation*}
$$

using (3.3), and handle the two terms separately. Then, by Proposition 4.3 (for the lower bound on $\|J P(Y)\|_{1}$ ) and an elementary estimate (for the upper bound on $\left.\|J P(Y)\|_{1}\right)$,

$$
\begin{align*}
\|J P(Y)\|_{1}-\frac{1}{2} \varepsilon^{2}\|L(\alpha \psi)\|_{1} \leq\|Z\|_{1} & \leq\|J P(Y)\|_{1}+\frac{1}{2} \varepsilon^{2}\|L(\alpha \psi)\|_{1}, \\
c_{5}(q)\|P(Y)\|_{3}-\varepsilon^{2} c_{6}(q)\|\alpha\|_{3} \leq\|Z\|_{1} & \leq c_{7}(q)\|P(Y)\|_{3}+\varepsilon^{2} c_{6}(q)\|\alpha\|_{3}, \\
c_{5}(q)\|Y\|_{3}-\varepsilon^{2} c_{8}(q)\|\alpha\|_{3} \leq\|Z\|_{1} & \leq c_{7}(q)\|Y\|_{3}+\varepsilon^{2} c_{8}(q)\|\alpha\|_{3}, \\
\text { 4) } \quad c_{5}(q)\|Y\|_{3}-\varepsilon^{2} c_{9}(q)\|Y\|_{3}^{2} \leq\|Z\|_{1} & \leq c_{7}(q)\|Y\|_{3}+\varepsilon^{2} c_{9}(q)\|Y\|_{3}^{2} \tag{6.14}
\end{align*}
$$

where we have used (6.10) in the last line. Combining (6.12) and (6.14), the result immediately follows.

Remark 6.6. Since $K$ is compact, we can replace $c(q), \widetilde{c}(q)$ in Theorems 6.4, 6.5 by global constants $C, \widetilde{C}>0$ under the extra assumption that $q \in K$.

## 7. Energy estimates for the coupled system

Having shown that $Q_{2}(Y)$ controls $\|Y\|_{3}^{2}$, for small $\varepsilon$, we must now bound the growth of $Q_{2}(Y)$ for a solution $(q, Y)$ of the coupled system. We do this by establishing quasi-conservation of energies $E_{1}, E_{2}$, related to $Q_{1}, Q_{2}$ :

Definition 7.1. Let $(q, Y):[0, T] \rightarrow K \times H^{3}$ with the regularity of Theorem 5.1. Associated to ( $q, Y$ ) we define the energies $E_{1}, E_{2}:[0, T] \rightarrow \mathbb{R}$,

$$
E_{1}(t)=\frac{1}{2}\left\|Y_{t}\right\|_{0}^{2}+\frac{1}{2} Q_{1}(Y), \quad E_{2}(t)=\frac{1}{2}\left\|(L Y)_{t}\right\|_{0}^{2}+\frac{1}{2} Q_{2}(Y) .
$$

Note that $E_{1}$ is $C^{1}$ and $E_{2}$ is continuous.
Throughout this section we will use the following:
Convention 7.2. $C$ will denote a positive constant depending (at most) on the choice of $\Sigma$ and $K . c\left(a_{1}, a_{2}, \ldots, a_{p}\right)$ will denote a smooth positive bounding function of $p$ nonnegative real arguments, which may also depend (implicitly) on $\Sigma, K$, and which is increasing in each of its arguments. $C_{0}, c_{0}$ will denote that the constant or bounding function depends, in addition, on the initial data $q_{0}, q_{1}, Y_{0}, Y_{1}$. The value of $C, C_{0}, c, c_{0}$ may vary from line to line.

Theorem 7.3 (Quasi-conservation of $E_{1}$ ). Let $(q, Y):[0, T] \rightarrow K \times H^{3}$ be a solution of the coupled system with the initial data and regularity of Theorem 5.1. Then

$$
E_{1}(t) \leq C_{0}+C\left(|\ddot{q}|+|\dot{q}|^{2}\right)\|Y(t)\|_{0}+\varepsilon \int_{0}^{t} c\left(|\dot{q}|,|\ddot{q}|,|\dddot{q}|,\|Y\|_{3},\left\|Y_{t}\right\|_{2}\right)
$$

Proof. The solution satisfies (3.15) and has $Y \in H^{3}, Y_{t} \in H^{2}, Y_{t t} \in H^{1}$, so

$$
\begin{align*}
\frac{d E_{1}}{d t}= & \left\langle Y_{t}, Y_{t t}\right\rangle+\left\langle Y_{t}, L Y\right\rangle+\frac{1}{2} \varepsilon\left\langle Y, L_{\tau} Y\right\rangle=\left\langle Y_{t}, k+\varepsilon j^{\prime}\right\rangle+\frac{1}{2} \varepsilon\left\langle Y, L_{\tau} Y\right\rangle \\
= & \frac{d}{d t}\langle Y, k\rangle-\varepsilon\left\langle Y, k_{\tau}-\frac{1}{2} L_{\tau} Y\right\rangle+\varepsilon\left\langle Y_{t}, j^{\prime}\right\rangle \\
\Rightarrow E_{1}(t)= & E_{1}(0)-\langle Y(0), k(0)\rangle+\langle Y(t), k(t)\rangle+\varepsilon \int_{0}^{t}\left\{\left\langle Y, \frac{1}{2} L_{\tau} Y-k_{\tau}\right\rangle+\left\langle Y_{t}, j^{\prime}\right\rangle\right\} \\
\leq & C_{0}+\|k(t)\|_{0}\|Y(t)\|_{0} \\
& +2 \varepsilon \int_{0}^{t}\left\{\|Y\|_{0}^{2}+\left\|L_{\tau} Y\right\|_{0}^{2}+\left\|k_{\tau}\right\|_{0}^{2}+\left\|Y_{t}\right\|_{0}^{2}+\left\|j^{\prime}\right\|_{0}^{2}\right\} . \tag{7.1}
\end{align*}
$$

$$
\begin{equation*}
\left\|L_{\tau} Y\right\|_{0} \leq C|\dot{q}|\|Y\|_{1} \tag{7.2}
\end{equation*}
$$

Recall that $k=-\psi_{\tau \tau}$, so

$$
\begin{equation*}
\|k(t)\|_{0} \leq C\|k(t)\|_{C^{0}} \leq C\left(|\ddot{q}|+|\dot{q}|^{2}\right) \tag{7.3}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\|k_{\tau}(t)\right\|_{0} \leq C\left\|k_{\tau}(t)\right\|_{C^{0}} \leq C\left(|\dddot{q}|+|\ddot{q}|^{2}+|\dot{q}|^{3}+|\dot{q}|^{2}\right) . \tag{7.4}
\end{equation*}
$$

Finally, it follows immediately from the algebra property of $\mathrm{H}^{2}$ (Proposition 4.1) that $\left\|j^{\prime}\right\|_{0} \leq\left\|j^{\prime}\right\|_{2} \leq c\left(|\dot{q}|,\|Y\|_{3},\left\|Y_{t}\right\|_{2}\right)$, and the result directly follows.

We will need a similar result bounding the growth of $E_{2}(t)$. Formally, this is obtained by applying the argument above with $Y$ replaced by $L Y$ (which formally solves a PDE of the form $\left.(L Y)_{t t}+L(L Y)=L k+O(\varepsilon)\right)$. Unfortunately, this argument is not rigorous since $Y$ is insufficiently regular to make sense of expressions like $L Y_{t t}$ (recall $Y_{t t}$ is only $H^{1}$ ).
Theorem 7.4 (Quasi-conservation of $E_{2}$ ). Let $(q, Y):[0, T] \rightarrow K \times H^{3}$ be a solution of the coupled system with the initial data and regularity of Theorem 5.1. Then

$$
E_{2}(t) \leq C_{0}+C\left(|\ddot{q}|+|\dot{q}|^{2}\right)\|Y(t)\|_{2}+\varepsilon \int_{0}^{t} c\left(|\dot{q}|,|\ddot{q}|,|\dddot{q}|,\|Y\|_{3},\left\|Y_{t}\right\|_{2}\right)
$$

Proof. By uniqueness, $(q, Y)$ must arise as the limit of an iteratively defined sequence of smooth functions $\left(q^{i}, Y^{i}\right)$, as constructed in the proof of Theorem 5.1. Recall that the sections $Y^{i}$ satisfy the PDEs

$$
\begin{equation*}
Y_{t t}^{i+1}-\Delta Y^{i+1}+B^{i} Y^{i}=k^{i}+\varepsilon\left(j^{\prime}\right)^{i} \tag{7.5}
\end{equation*}
$$

where $B$ denotes the first and zeroth order pieces of $L$ (so $L=-\Delta+B$ ) and the superscript $i$ on $B^{i}, k^{i},\left(j^{\prime}\right)^{i}$ denotes that the quantity is evaluated on the iterate $\left(q^{i}, Y^{i}\right)$. Recall also that $\left(Y^{i}, Y_{t}^{i}, Y_{t t}^{i}\right) \rightarrow\left(Y, Y_{t}, Y_{t t}\right)$ in $C^{0}\left(H^{3} \oplus H^{2} \oplus H^{1}\right)$. Now, for each $i$ define

$$
\begin{align*}
Z^{i} & =-\Delta Y^{i+1}+B^{i} Y^{i} \\
E^{i}(t) & =\frac{1}{2}\left\|Z_{t}^{i}\right\|^{2}+\frac{1}{2}\left\langle Z^{i},-\Delta Z^{i}+B^{i} Z^{i}\right\rangle \tag{7.6}
\end{align*}
$$

Each $E^{i}:[0, T] \rightarrow \mathbb{R}$ is smooth, $E^{i} \rightarrow E_{2}$ uniformly on $[0, T], Z^{i} \rightarrow Z=L Y$ in $H^{1}$ and $Z_{t}^{i} \rightarrow(L Y)_{t}$ in $L^{2}$. Note that $Z^{i}$ satisfies the PDE

$$
\begin{equation*}
Z_{t t}^{i}-\Delta Z^{i}+B^{i} Z^{i}=B^{i}\left(Z^{i}-Z^{i-1}\right)+\hat{k}^{i}+\varepsilon \hat{j}^{i} \tag{7.7}
\end{equation*}
$$

where $\hat{k}^{i}=-\Delta k^{i}+B^{i} k^{i-1}$ and $\hat{j}^{i}=-\Delta\left(j^{\prime}\right)^{i}+B^{i}\left(j^{\prime}\right)^{i-1}$. Since $B^{i}$ is self-adjoint, one has

$$
\begin{align*}
\frac{d E^{i}}{d t}= & \left\langle Z_{t}^{i}, Z_{t t}^{i}-\Delta Z^{i}+B^{i} Z^{i}\right\rangle+\frac{\varepsilon}{2}\left\langle Z^{i}, B_{\tau}^{i} Z^{i}\right\rangle \\
= & \left\langle Z_{t}^{i}, B^{i}\left(Z^{i}-Z^{i-1}\right)+\hat{k}^{i}\right\rangle+\varepsilon\left\{\left\langle Z_{t}^{i}, \hat{j}^{i}\right\rangle+\frac{1}{2}\left\langle Z^{i}, B_{\tau}^{i} Z^{i}\right\rangle\right\} \\
= & \frac{d}{d t}\left\langle Z^{i}, B^{i}\left(Z^{i}-Z^{i-1}\right)+\hat{k}^{i}\right\rangle-\left\langle B^{i} Z^{i}, Z_{t}^{i}-Z_{t}^{i-1}\right\rangle \\
& +\varepsilon\left\{\left\langle Z^{i}, B_{\tau}^{i}\left(\frac{1}{2} Z^{i}-Z^{i-1}\right)+\hat{k}_{\tau}^{i}\right\rangle+\left\langle Z_{t}^{i}, \hat{j}^{i}\right\rangle\right\} . \tag{7.8}
\end{align*}
$$

Integrating this from 0 to $t$ and taking the limit $i \rightarrow \infty$ yield

$$
\begin{align*}
E_{2}(t)-E_{2}(0) \leq & C_{0}+\|Z(t)\|_{0}\|\hat{k}(t)\|_{0} \\
& +C \varepsilon \int_{0}^{t}\left\{\|Z\|_{0}\left(\left\|\hat{k}_{\tau}\right\|_{0}+\mid \dot{q}\|Z Z\|_{1}\right)+\left\|Z_{t}\right\|_{0}\|\hat{j}\|_{0}\right\} \tag{7.9}
\end{align*}
$$

where we have used the facts that $Z_{t}^{i}-Z_{t}^{i-1} \rightarrow 0$ in $L^{2}, Z^{i}-Z^{i-1} \rightarrow 0$ in $H^{1}$, $\hat{k}^{i} \rightarrow \hat{k}=L k$ in $H^{3}, \hat{k}_{\tau}^{i} \rightarrow \hat{k}_{\tau}$ in $H^{3}$ and $\hat{j}^{i} \rightarrow \hat{j}=L j^{\prime}$ in $L^{2}$ (since $L: H^{2} \rightarrow L^{2}$ and $j: K \times \mathbb{R}^{4 n} \times H^{3} \times H^{2} \rightarrow H^{2}$ are continuous). Now, we have the elementary estimates

$$
\begin{aligned}
\|Z\|_{k} & =\|L Y\|_{k} \leq C\|Y\|_{k+2}, \quad k=1,2, \\
\left\|Z_{t}\right\|_{0} & =\left\|L Y_{t}+\varepsilon B_{\tau} Y\right\|_{0} \leq C\left(\left\|Y_{t}\right\|_{2}+\varepsilon \mid \dot{q}\| \| Y \|_{1}\right), \\
\|\hat{k}\|_{0} & \leq C\left(|\ddot{q}|+|\dot{q}|^{2}\right), \\
\left\|\hat{k}_{\tau}\right\|_{0} & \leq C\left(|\dddot{q}|+|\ddot{q}|^{2}+|\dot{q}|^{2}+|\dot{q}|^{3}\right) \\
\|\hat{j}\|_{0} & \leq C\left\|j^{\prime}\right\|_{2} \leq c\left(|\dot{q}|,\|Y\|_{3},\left\|Y_{t}\right\|_{2}\right) .
\end{aligned}
$$

Combining these with (7.9), the result follows.

## 8. Long time existence and proof of the Main Theorem

Throughout this section we choose and fix $q_{0} \in K$ and $q_{1} \in \mathbb{R}^{4 n}$, and denote by $(q, Y)$ the solution of the coupled system (3.15) with initial data $q(0)=q_{0}$, $\dot{q}(0)=q_{1}, Y(0)=Y_{t}(0)=0$. By Theorem 5.1 provided $\varepsilon<C(K) / \sqrt{\left|q_{1}\right|}$, this solution exists at least for time $t \in\left[0, T_{0}\right]$, where $T_{0}$ depends on the initial data, but is independent of $\varepsilon$. Moreover, the solution is unique, has the advertised regularity, satisfies the pointwise and $L^{2}$ orthogonality constraints (3.3), (3.11), and obeys the energy estimates of section 7 . Denote by $q_{*}(\tau)$ the geodesic in $\left(\mathrm{M}_{n}, \gamma\right)$ with the same initial data, $q_{*}(0)=q_{0}, \dot{q}_{*}(0)=q_{1}$. Note that $q_{*}$, considered as a function of rescaled time $\tau$, is independent of $\varepsilon$. Since geodesic flow conserves speed $\gamma\left(\dot{q}_{*}, \dot{q}_{*}\right)$, which uniformly bounds $\left|\dot{q}_{*}\right|^{2}$ on $K$, there exist $\tau_{0}>0, \alpha_{0}>0$, depending only on the initial data, such that $q_{*}$ exists and has

$$
\begin{equation*}
\left.\left|\dot{q}_{*}\right| \leq \alpha_{0}, \quad\left|\ddot{q}_{*}\right| \leq \alpha_{0}, \quad \operatorname{dist}(q(\tau), \partial K)\right)<d / 2, \tag{8.1}
\end{equation*}
$$

for all $\tau \in\left[0, \tau_{0}\right]$, where $d=\operatorname{dist}\left(q_{0}, \partial K\right)$. Hence, the geodesic $q_{*}$ exists for time $t \in\left[0, \varepsilon^{-1} \tau_{0}\right]$, which, for $\varepsilon$ small, exceeds $T_{0}$. Whenever $q, q_{*}$ both exist, we define $\varepsilon^{2} \widetilde{q}(t)$ to be the error between them, that is,

$$
\begin{equation*}
q=q_{*}+\varepsilon^{2} \widetilde{q} \tag{8.2}
\end{equation*}
$$

and

$$
\begin{equation*}
M(s)=\max _{0 \leq t \leq s}\left\{\varepsilon^{2}|\widetilde{q}(t)|^{2}+\left|\widetilde{q}^{\prime}(t)\right|^{2}+\left|\widetilde{q}^{\prime \prime}(t)\right|^{2}+\|Y(t)\|_{3}^{2}+\left\|Y_{t}(t)\right\|_{2}^{2}\right\} \tag{8.3}
\end{equation*}
$$

where primes denote differentiation with respect to $t$. This function, which measures the total error in replacing the wave map $\phi=\psi(q)+\varepsilon^{2} Y$ with the geodesic $\psi\left(q_{*}\right)$, is continuous, manifestly increasing, and has initial value $M(0)=0$. Our next task is to bound its growth. Before doing so, we define another absolute constant (depending only on $K$ and $\Sigma$ ), which will appear frequently in this section:

$$
\begin{equation*}
\alpha_{a}=\sup \left\{\left\|\gamma^{\mu \nu} \psi_{\nu \lambda}\right\|: q \in K, 0 \leq \mu, \lambda \leq 4 n\right\} . \tag{8.4}
\end{equation*}
$$

We will again use Convention 7.2 regarding bounding constants and functions.
Theorem 8.1 (A priori bound). Whenever $(q, Y)$ exists and $t<\tau_{0} / \varepsilon$, and $M(t)<$ $\varepsilon^{-4} \alpha_{a}^{-2}$, then

$$
M(t) \leq C_{0}+C_{0} M(t)^{\frac{1}{2}}+\left(\varepsilon^{2}+\varepsilon t+\varepsilon^{2} t^{2}+\varepsilon^{4} t^{4}\right) \frac{c_{0}(M(t))}{1-\varepsilon^{2} \alpha_{a} M(t)^{\frac{1}{2}}}
$$

Proof. We first derive the ODE satisfied by $\widetilde{q}$. The curve $q(\tau)$ satisfies the ODE

$$
\begin{equation*}
\ddot{q}+G(q, \dot{q}, \dot{q})=\varepsilon h\left(\varepsilon, q, \dot{q}, Y, Y_{t}\right)+\varepsilon^{2} a(q, \ddot{q}, Y) \tag{8.5}
\end{equation*}
$$

where $G, h$ are defined in (3.14), (5.7), and

$$
\begin{equation*}
a: K \times \mathbb{R}^{4 n} \times L^{2} \rightarrow \mathbb{R}^{4 n}, \quad a(q, v, Y)^{\mu}=\gamma^{\mu \nu}\left\langle\psi_{\nu \lambda}, Y\right\rangle v^{\lambda} \tag{8.6}
\end{equation*}
$$

which is smooth with respect to $q$ and linear with respect to $v$ and $Y$. The geodesic satisfies the ODE

$$
\begin{equation*}
\ddot{q}_{*}+G\left(q_{*}, \dot{q}_{*}, \dot{q}_{*}\right)=0 \tag{8.7}
\end{equation*}
$$

with the same initial data. Substituting $q=q_{*}+\varepsilon^{2} \widetilde{q}$ into (8.5) and using (8.7), we see that $\widetilde{q}$ satisfies

$$
\begin{aligned}
\widetilde{q}^{\prime \prime}= & {\left[G\left(q_{*}, \dot{q}_{*}, \dot{q}_{*}\right)-G\left(q_{*}+\varepsilon^{2} \widetilde{q}, \dot{q}_{*}, \dot{q}_{*}\right)\right]+\left[G\left(q, \dot{q}_{*}, \dot{q}_{*}-G\left(q, \dot{q}_{*}+\varepsilon \widetilde{q}, \dot{q}_{*}+\varepsilon q^{\prime}\right)\right]\right.} \\
(8.8) \quad & +\varepsilon h\left(\varepsilon, q, \dot{q}, Y, Y_{t}\right)+\varepsilon^{2} a\left(q, Y, \ddot{q}_{*}+\widetilde{q}^{\prime \prime}\right) .
\end{aligned}
$$

Now, $q$, by assumption, remains in $K$, and for all $t \in\left[0, \tau_{0} / \varepsilon\right], q_{*}$ remains in $K$ and $\left|\dot{q}_{*}\right| \leq \alpha_{0}$, so, since $G(q, u, v)$ is Lipshitz with respect to $q$ (on $K$ ) and bilinear in $(u, v)$, we have

$$
\begin{align*}
\left|\widetilde{q}^{\prime \prime}\right| & \leq C \alpha_{0}^{2} \varepsilon^{2}|\widetilde{q}|+C\left[\varepsilon^{2}|\widetilde{q}|^{2}+\varepsilon \alpha_{0}|\widetilde{q}|+\varepsilon|\widetilde{q}|\right] \varepsilon|h|+\varepsilon^{2}|a| \\
& \leq C_{0} \varepsilon M^{\frac{1}{2}}+C \varepsilon^{2} M+\varepsilon|h|+\varepsilon^{2}|a| . \tag{8.9}
\end{align*}
$$

To estimate the $h$ term, we note that it is smooth in $q$ and polynomial in $\dot{q}$ and $Y$ and its (first) derivatives, so in light of Proposition 4.1 we have (for $\varepsilon \leq 1$ ) the crude bound

$$
\begin{equation*}
\left|h\left(\varepsilon, q, \dot{q}, Y, Y_{t}\right)\right| \leq c\left(|\dot{q}|,\|Y\|_{3},\left\|Y_{t}\right\|_{2}\right) \tag{8.10}
\end{equation*}
$$

Now, by the definition of $M$,

$$
\begin{equation*}
|\dot{q}| \leq\left|\dot{q}_{*}\right|+\varepsilon\left|\widetilde{q}^{\prime}\right| \leq \alpha_{0}+\varepsilon M^{\frac{1}{2}} \tag{8.11}
\end{equation*}
$$

so $|h| \leq c_{0}(M)$. Turning to $|a|$, we have by linearity,

$$
\begin{equation*}
\left|a\left(q, Y, \ddot{q}_{*}+\widetilde{q}^{\prime \prime}\right)\right| \leq \alpha_{a}\|Y\|\left(\left|\ddot{q}_{*}\right|+\left|\widetilde{q}^{\prime \prime}\right|\right) \leq \alpha_{a} M^{\frac{1}{2}}\left(\alpha_{0}+M^{\frac{1}{2}}\right) \leq c_{0}(M) \tag{8.12}
\end{equation*}
$$

Hence,

$$
\begin{equation*}
\left|\widetilde{q}^{\prime \prime}(t)\right| \leq \varepsilon c_{0}(M(t)) . \tag{8.13}
\end{equation*}
$$

Now $\widetilde{q}^{\prime}(t)=\widetilde{q}^{\prime}(0)+\int_{0}^{t} \widetilde{q}^{\prime \prime}=\int_{0}^{t} \widetilde{q}^{\prime \prime}$, so

$$
\begin{equation*}
\left|\widetilde{q}^{\prime}(t)\right| \leq \varepsilon \int_{0}^{t} c_{0}(M(s)) d s \leq \varepsilon t c_{0}(M(t)) \tag{8.14}
\end{equation*}
$$

since $c_{0}$ and $M$ are, by definition, increasing. Similarly

$$
\begin{equation*}
\varepsilon|\widetilde{q}(t)| \leq \varepsilon \int_{0}^{t}|\widetilde{q}| \leq \varepsilon^{2} t^{2} c_{0}(M(t)) \tag{8.15}
\end{equation*}
$$

We have now bounded the growth of all the $\widetilde{q}$ terms in $M$. To bound the growth of $\|Y\|_{3}$ and $\left\|Y_{t}\right\|_{2}$, we will use the energy estimates of section 7 and the near coercivity property of $Q_{2}$ (Theorem 6.5). But to do this, we need to control $|\ddot{q}|$ and $|\dddot{q}|$, which appear in the energy estimates for $E_{1}(t)$ and $E_{2}(t)$, so we have not yet finished with the ODE for $q$. For $\ddot{q}$ we have from (8.13) the obvious bounds

$$
\begin{equation*}
|\ddot{q}| \leq\left|\ddot{q}_{*}\right|+\left|\widetilde{q}^{\prime \prime}\right| \leq \alpha_{0}+\varepsilon c_{0}(M(t)) . \tag{8.16}
\end{equation*}
$$

For $\dddot{q}$ we must work harder. So, for fixed $\varepsilon$, let $h_{i}, i=1,2,3,4$, denote the partial derivatives of $h: K \times \mathbb{R}^{4 n} \times H^{3} \times H^{2} \rightarrow \mathbb{R}^{4 n}$ with respect to each of its four entries. Similarly, let $G_{1}, a_{1}$ denote the derivatives of $G, a$ with respect to their first entries. Differentiating (8.5) with respect to $\tau$ one finds

$$
\begin{align*}
\begin{array}{r}
q
\end{array}+G_{1}(q, \dot{q}, \dot{q}) \dot{q}+2 G(q, \dot{q}, \dot{q}) \ddot{q}= & \varepsilon h_{1} \dot{q}+h_{2} \ddot{q}+h_{3} Y_{t}+h_{4} Y_{t t} \\
& +\varepsilon^{2}\left(a_{1}(q, \ddot{q}, Y) \dot{q}+a(q, \dddot{q}, Y)\right)+\varepsilon a\left(q, \ddot{q}, Y_{t}\right) \tag{8.17}
\end{align*}
$$

where we have used the bilinearity properties of $G$ and $a$. Inspecting the formula for $h$ (3.14), we obtain, using (8.16) and (8.11), a crude bound

$$
\begin{aligned}
|\dddot{q}| & \leq c_{0}(M)+c_{0}(M)\left\|Y_{t t}\right\|+\varepsilon^{2} \alpha_{a}|\dddot{q}|\|Y\| \\
& \leq c_{0}(M)+c_{0}(M)\left(\|L Y\|+\|k\|+\varepsilon\left\|j^{\prime}\right\|\right)+\varepsilon^{2} \alpha_{a} M^{\frac{1}{2}}|\dddot{q}| \\
& \leq c_{0}(M)+c_{0}(M)\left(\|Y\|_{2}+|\ddot{q}|+|\dot{q}|^{2}+c\left(|\dot{q}|,\|Y\|_{3},\left\|Y_{t}\right\|_{2}\right)\right)+\varepsilon^{2} \alpha_{a} M^{\frac{1}{2}}|\dddot{q}| \\
(8.18) & \leq c_{0}(M)+\varepsilon^{2} \alpha_{a} M^{\frac{1}{2}}|\dddot{q}| .
\end{aligned}
$$

Hence

$$
\begin{equation*}
|\dddot{q}| \leq \frac{c_{0}(M)}{1-\varepsilon^{2} \alpha_{a} M^{\frac{1}{2}}} \tag{8.19}
\end{equation*}
$$

whenever $M(t)<\varepsilon^{-4} \alpha_{a}^{-2}$. We can now turn to bounding $Y$.
Since $Q_{2}$ is nearly coercive (Theorem 6.5),

$$
\begin{equation*}
\|Y\|_{3}^{2} \leq C\left\{Q_{2}(Y)+\varepsilon^{2} c\left(\|Y\|_{3}\right)\right\}, \tag{8.20}
\end{equation*}
$$

and, by the standard elliptic estimate for $L$ (Proposition 4.3),

$$
\begin{align*}
\left\|Y_{t}\right\|_{2}^{2} & \leq C\left\{\left\|L Y_{t}\right\|^{2}+\left\|Y_{t}\right\|^{2}\right\} \leq C\left\{\left\|(L Y)_{t}\right\|^{2}+\varepsilon^{2}\left\|L_{\tau} Y\right\|^{2}+\left\|Y_{t}\right\|^{2}\right\} \\
& \leq C\left\{\left\|(L Y)_{t}\right\|^{2}+\varepsilon^{2}|\dot{q}|^{2}\|Y\|_{1}^{2}+\left\|Y_{t}\right\|^{2}\right\} \tag{8.21}
\end{align*}
$$

since the principal part of $L$ does not depend on time. Adding (8.20) and (8.21), one sees that

$$
\begin{align*}
\|Y\|_{3}^{2}+\left\|Y_{t}\right\|_{2}^{2} & \leq C\left\{E_{1}(t)+E_{2}(t)+\varepsilon^{2} c\left(\|Y\|_{3}\right)+\varepsilon^{2}|\dot{q}|^{2}\|Y\|_{1}^{2}\right\} \\
& \leq C\left\{E_{1}(t)+E_{2}(t)+\varepsilon^{2} c_{0}(M)\right\}, \tag{8.22}
\end{align*}
$$

where $E_{1}, E_{2}$ are as in Definition 7.1. Then, by Theorems 7.3 and 7.4 and the estimates (8.11), (8.16), (8.19),

$$
\begin{align*}
\|Y\|_{3}^{2}+\left\|Y_{t}\right\|_{2}^{2} \leq & C_{0}+\left(C_{0}+\varepsilon c_{0}(M)\right) M^{\frac{1}{2}} \\
& +\varepsilon \int_{0}^{t}\left(1+\frac{c_{0}(M)}{1-\varepsilon^{2} \alpha_{a} M^{\frac{1}{2}}}\right) c_{0}(M)+\varepsilon^{2} c_{0}(M) \\
\leq & C_{0}+C_{0} M^{\frac{1}{2}}+\varepsilon c_{0}(M)+\frac{\varepsilon t c_{0}(M)}{1-\varepsilon^{2} \alpha_{a} M^{\frac{1}{2}}} \tag{8.23}
\end{align*}
$$

provided $\varepsilon \leq 1$ and $M(t)<\varepsilon^{-4} \alpha_{a}^{-2}$. Combining (8.23), (8.15), (8.14) and (8.13) gives the a priori bound claimed.

Theorem 8.2 (Long time existence). There exist $\varepsilon_{*}>0$ and $\tau_{*}>0$, depending only on the initial data, such that for all $\varepsilon \in\left(0, \varepsilon_{*}\right)$ the solution $(q, Y)$ persists for all $t \in\left[0, \tau_{*} / \varepsilon\right]$ and has $M(t)$ bounded, independent of $\varepsilon$.

Proof. Choose and fix a constant $M_{*}>4 \alpha_{0}^{2}$ so large that $M_{*}>C_{0}+C_{0} M_{*}^{\frac{1}{2}}$, where $C_{0}$ is the specific constant, depending only on initial data, appearing in the a priori bound. Assume $\varepsilon>0$ is so small that

$$
\begin{equation*}
\varepsilon^{2}<\frac{1}{\alpha_{a} M_{*}^{\frac{1}{2}}}, \quad \varepsilon<\frac{d}{4 M_{*}^{\frac{1}{2}}}, \quad \varepsilon \leq \frac{1}{2}, \quad \varepsilon<\frac{C(K)}{M_{*}^{\frac{1}{4}}} \tag{8.24}
\end{equation*}
$$

where, as before, $d=\operatorname{dist}\left(q_{0}, \partial K\right)$ and $C(K)>0$ is the constant quoted in Theorem 5.1. Then, whenever $M(t) \leq M_{*},\|Y\|_{3}<M_{*}^{\frac{1}{2}},\left\|Y_{t}\right\|_{2} \leq M_{*}^{\frac{1}{2}}$,

$$
\begin{equation*}
|\dot{q}| \leq \alpha_{0}+\varepsilon M_{*}^{\frac{1}{2}}<M_{*}^{\frac{1}{2}} \tag{8.25}
\end{equation*}
$$

and

$$
\begin{equation*}
\varepsilon^{2}\|\widetilde{q}\| \leq \varepsilon M^{\frac{1}{2}}<\varepsilon M_{*}^{\frac{1}{2}}<\frac{d}{4} \tag{8.26}
\end{equation*}
$$

Hence, whenever $M(t) \leq M_{*}$, the value of the solution $\left(q, q_{t}, Y, Y_{t}\right)(t)$ satisfies the conditions of the initial data for the local existence theorem, Theorem 5.1, with $\Gamma=M_{*}^{\frac{1}{2}}$ and $\operatorname{dist}(q(t), \partial K)<\frac{d}{4}$. Given the last condition on $\varepsilon$, [8.24), it follows that we may apply Theorem 5.1 and extend the solution for a time $\delta T>0$ depending only on $M_{*}$ and $d$, independent of $\varepsilon$. It follows that the solution persists for as long as $M(t) \leq M_{*}$. Furthermore, by the first condition on $\varepsilon$, (8.24), whenever $M(t) \leq M_{*}$ the solution obeys the a priori bound, Theorem 8.1.

For each $\varepsilon>0$ let $t_{\varepsilon}=\sup \left\{t: M(t) \leq M_{*}\right\}$. We claim that there exists $\varepsilon_{*}>0$ such that $\varepsilon t_{\varepsilon}$ is bounded away from zero on $\left[0, \varepsilon_{*}\right]$. Note that this immediately implies the statement in the theorem since then there exists $\tau_{*}>0$ such that $t_{\varepsilon} \geq \tau_{*} / \varepsilon$ for all $\varepsilon \in\left(0, \varepsilon_{*}\right)$, and the solution exists, with $M(t) \leq M_{*}$ on $\left[0, \tau_{*} / \varepsilon\right]$. Assume, towards a contradiction, that no such $\varepsilon_{*}$ exists. Then there is a positive sequence $\varepsilon_{i} \rightarrow 0$ such that $\varepsilon_{i} t_{\varepsilon_{i}} \rightarrow 0$. But then the a priori bound at time $t_{\varepsilon_{i}}$ gives (recall $M$ is continuous, so $M\left(t_{\varepsilon_{i}}\right)=M_{*}$ ), in the limit $i \rightarrow \infty$,

$$
\begin{equation*}
M_{*}<C_{0}+C_{0} M_{*}^{\frac{1}{2}} \tag{8.27}
\end{equation*}
$$

a contradiction by our choice of $M_{*}$.
Proof of Main Theorem. By Theorem 8.2, for all $\varepsilon \in\left(0, \varepsilon_{*}\right]$ the solution exists for $t \in\left[0, \tau_{*} / \varepsilon\right]$ and coincides with $\psi\left(q_{*}(\tau)+\varepsilon^{2} \widetilde{q}(t)\right)+\varepsilon^{2} Y(t)$, with $\varepsilon|\widetilde{q}(t)|,\|Y(t)\|_{3}$, and
hence (by Proposition 4.2) $\|Y(t)\|_{C^{0}}$ uniformly bounded in $t$ and $\varepsilon$. The rescaled solution is
$\phi^{\varepsilon}:\left[0, \tau_{*}\right] \times \Sigma \rightarrow S^{2} \subset \mathbb{R}^{3}, \quad \phi^{\varepsilon}(\tau, x, y)=\psi\left(q_{*}(\tau)+\varepsilon^{2} \widetilde{q}(\tau / \varepsilon), x, y\right)+\varepsilon^{2} Y(\tau / \varepsilon, x, y)$, and the geodesic with the same initial data is

$$
\begin{equation*}
\psi_{*}:\left[0, \tau_{*}\right] \times \Sigma, \quad \psi_{*}(\tau, x, y)=\psi\left(q_{*}(\tau), x, y\right) \tag{8.29}
\end{equation*}
$$

Now $\psi: K \times \Sigma \rightarrow S^{2}$ is smooth, hence uniformly continuous (since $K \times \Sigma$ is compact). Hence, as $\varepsilon \rightarrow 0, \phi^{\varepsilon}$ converges uniformly on $\left[0, \tau_{*}\right] \times \Sigma$ to $\psi_{*}$. Furthermore,

$$
\begin{align*}
\phi_{x}^{\varepsilon}(\tau, x, y) & =\psi_{x}\left(q_{*}(\tau)+\varepsilon^{2} \widetilde{q}(\tau / \varepsilon), x, y\right)+\varepsilon^{2} Y_{x} \\
\phi_{y}^{\varepsilon}(\tau, x, y) & =\psi_{y}\left(q_{*}(\tau)+\varepsilon^{2} \widetilde{q}(\tau / \varepsilon), x, y\right)+\varepsilon^{2} Y_{y} \\
\phi_{\tau}^{\varepsilon}(\tau, x, y) & =\left(\dot{q}_{*}^{\mu}+\varepsilon \widetilde{q}^{\prime}(\tau / \varepsilon)\right) \psi_{\mu}\left(q_{*}(\tau)+\varepsilon^{2} \widetilde{q}(\tau / \varepsilon), x, y\right)+\varepsilon Y_{t} \tag{8.30}
\end{align*}
$$

and $\left|\widetilde{q}^{\prime}(t)\right|,\left\|Y_{x}(t)\right\|_{C^{0}},\left\|Y_{y}(t)\right\|_{C^{0}},\left\|Y_{t}(t)\right\|_{C^{0}}$ are bounded uniformly in $t$ and $\varepsilon$ (again using Proposition 4.2), so $\phi_{x}^{\varepsilon}, \phi_{y}^{\varepsilon}, \phi_{\tau}^{\varepsilon}$ converge uniformly on $\left[0, \tau_{*}\right] \times \Sigma$ to $\psi_{* x}, \psi_{* y}, \psi_{* \tau}$. Hence, $\phi^{\varepsilon}$ converges to $\psi_{*}$ in $C^{1}$.

## Appendix: Analytic properties of the nonlinear terms

The proof of the local existence theorem, Theorem 5.1 makes fundamental use of certain basic analytic properties (Propositions 5.3 and 5.5) of the right hand sides $f, g$ of the evolution system (5.4), (5.5). These properties are established by a long chain of elementary arguments which we sketch in this appendix. We regard $\varepsilon$ as a fixed parameter in $(0,1)$ and choose a local parametrization $\psi: U \times \Sigma \rightarrow S^{2}$ of $\mathrm{M}_{n}$ as given by Proposition 2.1 and a compact convex neighbourhood $K \subset U$. We begin by showing that the associated maps $\Psi_{k}: U \rightarrow H^{k}, q \mapsto \psi(q, \cdot)$ are smooth for all $k \in \mathbb{N}$.

Lemma A.1. Let $k \in \mathbb{N}$ and $f: U \times \Sigma \rightarrow \mathbb{R}$ be smooth. Then $F: U \rightarrow \mathbf{H}^{k}$, $F(q)=f(q, \cdot)$, is smooth.

Proof. As usual, we will denote partial derivatives with respect to $q^{\mu}$ by a subscript $\mu$. It suffices to show that $F$ is everywhere differentiable, since all partial derivatives $f_{\mu_{1} \mu_{2} \cdots \mu_{r}}$ are, like $f$, smooth maps $U \times \Sigma \rightarrow \mathbb{R}$. By the mean value theorem there exist $t_{1}, \ldots, t_{r} \in[0,1]$ such that

$$
\begin{aligned}
\left\|F(q+p)-F(q)-f_{\mu}(q, \cdot) p^{\mu}\right\|_{k}^{2} \leq & |p|^{2} \sum_{\mu=1}^{4 n} \int_{\Sigma}\left\{\left(f_{\mu}\left(q+t_{1} p\right)-f_{\mu}(q)\right)^{2}\right. \\
& +\left(f_{\mu x}\left(q+t_{2} p\right)-f_{\mu x}(q)\right)^{2} \\
& \left.+\cdots+\left(f_{\mu y \cdots y}\left(q+t_{r} p\right)-f_{\mu y \cdots y}(q)\right)^{2}\right\}
\end{aligned}
$$

since $f$ and all its partial derivatives are $C^{1}$ functions of $q$. But $f_{\mu}, \ldots, f_{\mu y \cdots y}$ are uniformly continuous on $B_{\delta}(q) \times \Sigma$, for $\delta>0$ sufficiently small, so

$$
\begin{equation*}
\lim _{p \rightarrow 0} \frac{1}{|p|}\left\|F(q+p)-F(q)-\left.\frac{\partial f}{\partial q^{\mu}}\right|_{q} p^{\mu}\right\|_{k}=0 \tag{A.2}
\end{equation*}
$$

as was to be proved.

Now, by the definition of $H^{k}, \psi: U \rightarrow H^{k}$ is differentiable if and only if $\psi_{i}$ : $U \rightarrow \mathrm{H}^{k}$ is differentiable for $i=1,2,3$, so we immediately obtain:

Corollary A.2. $\Psi_{k}: U \rightarrow H^{k}, q \mapsto \psi(q, \cdot)$, is smooth for all $k \in \mathbb{N}$.
The error terms are $j^{\prime}\left(q, \dot{q}, Y, Y_{t}\right)$ and $h\left(q, \dot{q}, Y, Y_{t}\right)$ where

$$
\begin{aligned}
j^{\prime}(q, p, Y, Z)= & 2 p^{\mu}\left(\psi_{\mu} \cdot Z\right) \psi+\varepsilon\left(|Z|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) \psi+\varepsilon p^{\mu} p^{\nu}\left(\psi_{\mu} \cdot \psi_{\nu}\right) Y \\
& -2 \varepsilon\left(\psi_{x} \cdot Y_{x}+\psi_{y} \cdot Y_{y}\right) Y \\
& +\varepsilon^{2}\left\{|Y|^{2} \Delta \psi+2\left(Y \cdot Y_{x}\right) \psi_{x}+2\left(Y \cdot Y_{y}\right) \psi_{y}\right\} \\
& +2 \varepsilon^{2} p^{\mu}\left(\psi_{\mu} \cdot Z\right) Y+\varepsilon^{3}\left(|Z|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) Y, \\
h(q, p, Y, Z)^{\mu}= & \gamma^{\mu \nu}\left\{\left\langle Z, \psi_{\nu \lambda}\right\rangle p^{\lambda}+\varepsilon\left\langle Y, \psi_{\lambda \nu \rho}\right\rangle p^{\lambda} p^{\rho}+\varepsilon\left\langle\psi_{\lambda} \cdot \psi_{\rho} Y, \psi_{\nu}\right\rangle p^{\lambda} p^{\rho}\right. \\
& -2 \varepsilon\left\langle\left(\psi_{x} \cdot Y_{x}+\psi_{y} \cdot Y_{y}\right) Y, \psi_{\nu}\right\rangle+2 \varepsilon^{2}\left\langle\left(\psi_{\lambda} \cdot Z\right) Y, \psi_{\nu}\right\rangle p^{\lambda} \\
& \left.+\varepsilon^{3}\left\langle\left(|Z|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) Y, \psi_{\nu}\right\rangle\right\} .
\end{aligned}
$$

We will also need to consider the quantities
(A.4) $M^{\mu}{ }_{\nu}(q, Y)=\delta^{\mu}{ }_{\nu}-\varepsilon^{2} \gamma^{\mu \lambda}(q)\left\langle Y, \psi_{\lambda \nu}\right\rangle$,

$$
\begin{align*}
A(q, Y)= & -\left(\left|\psi_{x}\right|^{2}+\left|\psi_{y}\right|^{2}\right) Y-2\left(\psi_{x} \cdot Y_{x}+\psi_{y} \cdot Y_{y}\right) \psi-2(\psi \cdot Y) \Delta \psi \\
& -2(\psi \cdot Y)_{x} \psi_{x}-2(\psi \cdot Y)_{y} \psi_{y} . \tag{A.5}
\end{align*}
$$

Let $B$ denote the Banach space $\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{3} \times H^{2}$ with norm

$$
\begin{equation*}
\|(q, p, Y, Z)\|_{B}=\max \left\{|q|,|p|,\|Y\|_{3},\|Z\|_{2}\right\} \tag{A.6}
\end{equation*}
$$

let $B_{U} \subset B$ denote the open set on which $q \in U$ and for each $\Gamma \geq 0$, let $B_{\Gamma} \subset B_{U}$ denote the closed convex subset on which $q \in K$ and $\|(0, p, Y, Z)\|_{B} \leq \Gamma$.

Proposition A.3. The quantities defined above are smooth maps $j^{\prime}: B_{U} \rightarrow H^{2}$, $h: B_{U} \rightarrow \mathbb{R}^{4 n}, M^{\mu}{ }_{\nu}: B_{U} \rightarrow \mathbb{R}, A: B_{U} \rightarrow H^{2}$.

Proof. That $j^{\prime}, h$ define maps $B_{U} \rightarrow H^{2}$ and $B_{U} \rightarrow \mathbb{R}^{4 n}$ follows immediately from the algebra property of $\mathrm{H}^{2}$ (Proposition 4.1). Now $j^{\prime}$ is a linear combination of terms formed by composing the maps

$$
\begin{align*}
\Psi_{k}: U & \rightarrow H^{k}, & & q \mapsto \psi(q, \cdot), \\
\mathrm{d} \Psi_{k}: U \times \mathbb{R}^{4 n} & \rightarrow H^{k}, & & (q, p) \mapsto p^{\mu} \psi_{\mu}(q, \cdot), \tag{A.7}
\end{align*}
$$

which are smooth by Corollary A.2, and the manifestly smooth maps

$$
\begin{align*}
H^{k} \rightarrow \mathrm{H}^{k}, & Y \mapsto Y, \\
\mathrm{H}^{k} \rightarrow H^{k}, & f \mapsto f e_{i}, \\
H^{k} \rightarrow H^{k-1}, & Y \mapsto Y, Y \mapsto Y_{x}, Y \mapsto Y_{y}, \\
\mathrm{H}^{2} \times \mathrm{H}^{2} \rightarrow \mathrm{H}^{2}, & (f, g) \rightarrow f g, \tag{A.8}
\end{align*}
$$

where $e_{1}=(1,0,0), e_{2}=(0,1,0)$ and $e_{3}=(0,0,1)$. Hence $j^{\prime}$ is smooth.
The map $h$ is handled similarly after noting that the inverse metric coefficients $\gamma^{\mu \nu}$ are smooth $U \rightarrow \mathbb{R}$, the higher derivatives

$$
\mathrm{d}^{2} \Psi_{k}: U \times \mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \rightarrow H^{k}, \quad\left(q, p_{1}, p_{2}\right) \mapsto p_{1}^{\mu} p_{2}^{\nu} \psi_{\mu \nu}(q, \cdot)
$$

(A.9) $\mathrm{d}^{3} \Psi_{k}: U \times \mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \rightarrow H^{k}, \quad\left(q, p_{1}, p_{2}, p_{3}\right) \mapsto p_{1}^{\mu} p_{2}^{\nu} p_{3}^{\lambda} \psi_{\mu \nu \lambda}(q, \cdot)$
are smooth by Lemma A.2, and in addition to the maps in (A.8) above, the map

$$
\begin{equation*}
H^{0} \times H^{0} \rightarrow \mathbb{R}, \quad(Y, Z) \mapsto\langle Y, Z\rangle \tag{A.10}
\end{equation*}
$$

is manifestly smooth.
That $M^{\mu}{ }_{\nu}$ and $A$ define smooth maps on $B_{U}$ is clear, since their $q$-dependence is smooth and they depend linearly on $Y$ (and are independent of $p$ and $Z$ ).

We can now assemble these pieces to show that $f$ and $g$ are smooth functions on $X_{\Gamma}$, the space defined in section [5. To do so, we note that $f=\widehat{f} \circ \iota$ and $g=\widehat{g} \circ \iota$ where $\iota: X_{\Gamma} \rightarrow B_{8 c(\Sigma) \Gamma}$ is the linear isometry

$$
\iota:(q, p, Y, Z) \rightarrow\left(q, \varepsilon^{-1} p, Y, Z\right)
$$

and $\widehat{f}: B_{\Gamma^{\prime}} \rightarrow \mathbb{R}^{4 n}, \widehat{g}: B_{\Gamma^{\prime}} \rightarrow H^{2}$ are
(A.11) $\widehat{f}(q, p, Y, Z)=M^{-1}(q, Y)\left(-\varepsilon^{2} G(q, p, p)+\varepsilon h(q, p, Y, Z)\right)$,
(A.12) $\widehat{g}(q, p, Y, Z)=-A(q, Y)-\psi_{\mu} \widehat{f}^{\mu}(q, p, Y, Z)-\psi_{\mu \nu} p^{\mu} p^{\nu}+\varepsilon j^{\prime}(q, p, Y, Z)$, and $G$ is defined in (5.6). The point is that $f, g$ are defined as functions of $q_{t}$ (and $q, Y, Y_{t}$ ) on a space ( $X_{\Gamma}$ ) with $\varepsilon$-dependent norm, but it is more convenient here to think of them as functions of $\dot{q}=\varepsilon^{-1} q_{t}$, on a space ( $B_{\Gamma^{\prime}}, \Gamma^{\prime}=8 c(\Sigma) \Gamma$ ) with fixed norm. Since $\iota$ is a linear isometry, $f, g$ are smooth, bounded, Lipschitz, etc., if and only if $\widehat{f}, \widehat{g}$ are.
Proposition A.4. There exists $\varepsilon_{*}=O(1 / \sqrt{\Gamma})$ such that for all $\varepsilon \in\left(0, \varepsilon_{*}\right)$, $\widehat{f}$ : $B_{\Gamma} \rightarrow \mathbb{R}^{4 n}$ and $\widehat{g}: B_{\Gamma} \rightarrow H^{2}$ are smooth.

Proof. Since $\gamma^{\mu \lambda}(q), \psi_{\lambda \nu}(q)$ are smooth and $K$ is compact, there exists $\varepsilon_{*}>0$ such that the matrix $M(q, Y)$ is uniformly invertible on $B_{\Gamma}$ for all $\varepsilon \in\left(0, \varepsilon_{*}\right)$. The components of $M^{-1}$ are rational in $M^{\mu}{ }_{\nu}$ with denominator $\operatorname{det} M$, which is bounded away from 0 on $B_{\Gamma}$. Hence, $M^{-1}$ is smooth on $B_{\Gamma}$, and the proposition follows immediately from the Leibniz rule and Proposition A. 3 ,

Since $\widehat{f}$ and $\widehat{g}$ are smooth, they are certainly continuously differentiable. To complete the proof of Proposition 5.3 it remains to show that their differentials are bounded.

Proposition A.5. For all $\varepsilon \in\left(0, \varepsilon_{*}\right)$, the derivatives $\mathrm{d} \widehat{g}: B_{\Gamma} \rightarrow \mathcal{L}\left(B, H^{2}\right)$ and $\mathrm{d} \widehat{f}: B_{\Gamma} \rightarrow \mathcal{L}\left(B, \mathbb{R}^{4 n}\right)$ are bounded, independent of $\varepsilon$.
Proof. This is established by estimating the operator norm of the derivatives termwise. For example, the first term of $j^{\prime}$,

$$
\begin{equation*}
J(q, p, Y, Z)=2 p^{\mu}\left(\psi_{\mu} \cdot Z\right) \psi \tag{A.13}
\end{equation*}
$$

has derivative
$d J_{(q, p, Y, Z)}:(\hat{q}, \hat{p}, \hat{Y}, \hat{Z}) \mapsto 2 \hat{q}^{\nu} p^{\mu}\left[\left(\psi_{\nu \mu} \cdot Z\right) \psi+\left(\psi_{\mu} \cdot Z\right) \psi_{\nu}\right]+2 \hat{p}^{\mu}\left(\psi_{\mu} \cdot Z\right) \psi+2 p^{\mu}\left(\psi_{\mu} \cdot \hat{Z}\right) \psi$, and so, by the algebra property of $\mathrm{H}^{2}$, for all $(q, p, Y, Z) \in B_{\Gamma}$,

$$
\begin{align*}
\left\|d J_{(q, p, Y, Z)}(\hat{q}, \hat{p}, \hat{Y}, \hat{Z})\right\|_{2} & \leq C(K)\left\{\left|\hat{q}\left\|p\left|\|Z\|_{2}+|\hat{p}|\|Z\|_{2}+\|Z\|_{2}\right| p \mid\right\}\right.\right. \\
& \leq C(K)\left\{2 \Gamma+\Gamma^{2}\right\}\|(\hat{q}, \hat{p}, \hat{Y}, \hat{Z})\|_{B} \tag{A.15}
\end{align*}
$$

where $C(K)>0$ is a constant depending only on $K$. Hence, for all $(q, p, Y, Z) \in B_{\Gamma}$,

$$
\begin{equation*}
\left\|d J_{(p, q, Y, Z)}\right\|_{\mathcal{L}\left(B, H^{2}\right)} \leq C(K)\left\{2 \Gamma+\Gamma^{2}\right\} \tag{A.16}
\end{equation*}
$$

The other terms of $j^{\prime}, h$ and $A$ are handled similarly. To bound $\mathrm{d} M^{-1}$ we bound $\mathrm{d} M$ and appeal to the Leibniz rule and uniform invertibility of $M$. Boundedness of the differentials of functions depending only on $q$ and $p$ is immediate by compactness and finiteness of dimension.

Finally, we turn to Proposition 5.5, which is equivalent to:
Proposition A.6. The differentials of the maps $\widehat{f}: B_{\Gamma} \rightarrow \mathbb{R}^{4 n} \widehat{g}: B_{\Gamma} \rightarrow H^{2}$ extend to maps $\mathrm{d} \widehat{f}^{e x t}: B_{\Gamma} \rightarrow \mathcal{L}\left(\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{1} \times L^{2}, \mathbb{R}^{4 n}\right)$ and $\mathrm{d} \widehat{g}^{e x t} g: B_{\Gamma} \rightarrow$ $\mathcal{L}\left(\mathbb{R}^{4 n} \times \mathbb{R}^{4 n} \times H^{1} \times L^{2}, L^{2}\right)$ bounded by $\Lambda_{f}, \Lambda_{g}$ respectively.

Proof. This follows from explicit termwise computation. For example, the last term of $h$ is

$$
m_{\mu}(q, Y, Z)=\left\langle\left(|Z|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) Y, \psi_{\mu}\right\rangle
$$

whose differential at $(q, Y, Z)$ is the linear map $\mathbb{R}^{4 n} \times H^{3} \times H^{2} \rightarrow \mathbb{R}$,

$$
\begin{aligned}
\mathrm{d} m_{\mu}:(\hat{q}, \hat{Y}, \hat{Z}) \mapsto & \left\langle\left(|Z|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) Y, \psi_{\mu \nu}\right\rangle \hat{q}^{\nu}+\left\langle\left(|Z|^{2}-\left|Y_{x}\right|^{2}-\left|Y_{y}\right|^{2}\right) \psi_{\mu}, \hat{Y}\right\rangle \\
& -2\left\langle\left(Y \cdot \psi_{\mu}\right) Y_{x}, \hat{Y}_{x}\right\rangle-2\left\langle\left(Y \cdot \psi_{\mu}\right) Y_{y}, \hat{Y}_{y}\right\rangle+2\left\langle\left(Y \cdot \psi_{\mu}\right) Z, \hat{Z}\right\rangle,
\end{aligned}
$$

which clearly extends to a bounded linear map $\mathbb{R}^{4 n} \times H^{1} \times L^{2} \rightarrow \mathbb{R}$.

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