# Bifurcation of Solutions to the Allen-Cahn Equation 

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#### Abstract

We use Morse Homology to study bifurcation of the solution sets of the AllenCahn Equation.


Key Words: Bifurcation, Allen-Cahn Equation, Morse Homology.
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## 1-Introduction.

1.1 - Main Results. Let $M:=\left(M^{n}, g\right)$ be a compact, $n$-dimensional Riemannian manifold. For $\epsilon>0$, we define the Allen-Cahn Operator over $C^{\infty}(M)$ with parameter $\epsilon$ by:

$$
\mathrm{AC}_{\epsilon, g}(u):=-\epsilon \Delta_{g} u+u^{3}-u
$$

where $\Delta_{g}$ is the Laplacian operator of $g$. The Allen-Cahn Operator appears in mathematical physics to describe the process of phase separation in metal allows (c.f. [2]), and its interesting properties have already made it the object of various mathematical studies (c.f., for example, [4], [5] and [7]). In particular, the Allen-Cahn Operator is variational, arising as the Euler-Lagrange Equations (that is, the $L^{2}$-gradient) of the Ginzburg-Landau-Wilson Free Energy Functional:

$$
\mathcal{E}_{\epsilon, g}(u)=\int_{M} \epsilon\left\|\nabla^{g} u\right\|^{2}+\frac{1}{4}\left(u^{2}-1\right)^{2} \mathrm{dVol},
$$

and for this reason naturally lends itself to analysis by Morse theoretical techniques. In this note, we use Morse Homology, we show how the space of solutions to the Allen-Cahn Equation bifurcates as $\epsilon$ becomes small. Indeed, the Morse Homology yields a lower bound for the number of solutions, which increases discretely as $\epsilon^{-1}$ crosses points of the spectrum of $-\Delta_{g}$, tending to infinity as $\epsilon$ tends to zero.
We denote by $\mathcal{M}$ the space of smooth Riemannian metrics over $M$, which we furnish with the topology of $C^{\infty}$ convergence. For any smooth metric $g$, we define the solution space $\left.\mathcal{Z}_{g} \subseteq\right] 0, \infty\left[\times C^{\infty}(M)\right.$ by:

$$
\mathcal{Z}_{g}=\left\{(\epsilon, u) \mid \mathrm{AC}_{\epsilon, g}(u)=0\right\}
$$

so that $\mathcal{Z}_{g}$ is the union of all solution sets for the metric $g$ and all parameters $\left.\epsilon \in\right] 0, \infty[$. We denote by $\left.e_{g}: \mathcal{Z}_{g} \rightarrow\right] 0, \infty[$ the projection onto the second factor, and we will see presently (c.f. Proposition 2.1.3) that this is a proper map. For all $\epsilon \in] 0, \infty[$, we denote:

$$
\mathcal{Z}_{\epsilon, g}=e_{g}^{-1}(\{\epsilon\})=\left\{u \mid \mathrm{AC}_{\epsilon, g}(u)=0\right\},
$$

so that $\mathcal{Z}_{\epsilon, g}$ is the solution set for the metric $g$ and the parameter $\epsilon$. We aim to study the manner in which $\mathcal{Z}_{\epsilon, g}$ bifurcates as $\epsilon$ tends to 0 , and the simplest way to do so is to study the geometry of $\left(\mathcal{Z}_{g}, e_{g}\right)$. We first show that, upon perturbing the metric by an arbitrarily small amount, we may suppose that this geometry is relatively straightforward but for a countable set of singularities determined by the Laplacian of $g$. Indeed, we denote by $\operatorname{Spec}\left(-\Delta_{g}\right)$ the set of all eigenvalues of $-\Delta_{g}$ (which, by convention, is non-negative), and we define the singular set, $\left.\operatorname{Sing}_{g} \subseteq\right] 0, \infty\left[\times C^{\infty}(M)\right.$ by:

$$
\operatorname{Sing}_{g}=\left\{(\epsilon, 0) \mid \epsilon^{-1} \in \operatorname{Spec}\left(-\Delta_{g}\right)\right\}
$$

We recall that a subset $X$ of $\mathcal{M}$ is said to be generic (or equivalently, in the second category in the sense of Baire), whenever it contains a countable intersection of dense open sets. A given property is then said to hold for generic elements whenever it holds for all elements of some generic set. We recall that, by the Baire Category Theorem, any property that holds for generic elements of $\mathcal{M}$ in particular holds over a dense subset of $\mathcal{M}$. Using transversality techniques, we to show:

## Theorem 1.1.1

For generic $g \in \mathcal{M}, \mathcal{Z}_{g} \backslash \operatorname{Sing}_{g}$ is a smooth, 1-dimensional submanifold of $] 0, \infty\left[\times C^{\infty}(M)\right.$. Moreover, if $\operatorname{Dim}(M) \geqslant 3$, then we may assume in addition that all critical points of $e_{g}$ are non-degenerate. In particular:
(1) if $\epsilon^{-1} \notin \operatorname{Spec}\left(-\Delta_{g}\right)$, then $\mathcal{Z}_{\epsilon, g}$ is finite; and
(2) there exists a discrete subset $X$ of the complement of $\operatorname{Spec}\left(-\Delta_{g}\right)$ such that if $\epsilon^{-1} \notin$ $X \cup \operatorname{Spec}_{g}$, then $\mathcal{Z}_{\epsilon, g}$ only consists of non-degenerate solutions of the Allen-Cahn Equation.
Remark: We recall that a solution to an elliptic partial differential equation is said to be non-degenerate whenever the linearisation of the operator about that solution is invertible.

When $u$ is a solution of the Allen-Cahn Equation with parameter $\epsilon$, we denote by $L \mathrm{AC}_{\epsilon, g}(u)$ the linearisation of the Allen-Cahn Operator about $u$. When $u$ is non-degenerate, we denote by Index $(u)$ its Morse Index, which we recall is defined to be equal to the number of strictly negative eigenvalues of $L \mathrm{AC}_{\epsilon, g}(u)$ counted with geometric multiplicity. Observe that the constant function $u=0$ is a solution of $\mathrm{AC}_{\epsilon, g}$ for all $g$ and for all $\epsilon$. Moreover, as we will see presently (c.f. Proposition 2.2.1), $u=0$ is non-degenerate if and only if $\epsilon^{-1} \notin \operatorname{Spec}\left(-\Delta_{g}\right)$, and the index of this solution is given by:

$$
\operatorname{Index}(0)=\#\left\{\lambda \in \operatorname{Spec}\left(-\Delta_{g}\right) \mid \lambda<\epsilon^{-1}\right\}
$$

Observe that Index (0) tends to infinity as $\epsilon$ tends to zero. Our second result now describes in terms of the number of solutions of a given Morse Index how $\mathcal{Z}_{\epsilon, g}$ bifurcates as $\epsilon$ tends to 0 :

## Theorem 1.1.2

For generic $g \in \mathcal{M}$, if $\epsilon^{-1} \notin \operatorname{Spec}\left(-\Delta_{g}\right)$, then for all $0 \leqslant k<\operatorname{Index}(0)$, there exist at least two non-degenerate solutions $u$ and $-u$ of the Allen-Cahn Equation such that $\operatorname{Index}(u)=\operatorname{Index}(-u)=k$.
Remark: In particular, we do not require Part 2 of Theorem 1.1.1. For countably many values of $\epsilon$, there may exist finitely many degenerate solutions. We simply ignore them.
Theorem 1.1.2 follows from Theorem 1.1.1 in a straightforward manner from standard Morse homological techniques. Indeed, for all $g$ and for all $\epsilon$, the constant functions $u= \pm 1$ are also solutions of the Allen-Cahn Equation, this time with Morse Index equal to 0 . Denoting $l=\operatorname{Index}(0)$, it follows that the chain groups $C_{0}$ and $C_{l}$ of the MorseComplex of the Ginzburg-Landau-Wilson Free Energy Functional are at least 2- and 1dimensional respectively. Since the underlying space (that is, $C^{\infty}(M)$ ) is contractible, the Morse-Homology, $H_{k}$, is non-trivial only for $k=0$. Finally, as the Allen-Cahn Operator is an odd operator, all the intermediate chain groups $C_{1}, C_{2}, \ldots, C_{l-1}$ have even dimension, and this fact, used together with the algebraic relations of Morse Homology allows us to deduce that they are non-trivial, thus proving the theorem.

Theorem 1.1.1 is proven using the Sard-Smale Theorem, and this paper is therefore mostly devoted to obtaining the requisite surjectivity results. We draw the reader's attention

## The Allen-Cahn Equation

to the fact that our usage of the Sard-Smale Theorem differs from standard approaches in one subtle but interesting respect. Indeed, whilst any application of the Sard-Smale Theorem is generally considered to require the separability of the function spaces used, we replace this condition by one that we call "paraproperness", which is to properness as paracompactness is to compactness. In Proposition 2.2.2, we make paraproperness into a useful concept by showing that it is preserved by restriction to both closed and open subsets, and in Theorem 2.2.4, we reprove the Sard-Smale Theorem in this new context. This will be of particular use in the forthcoming paper [14] where it makes possible the construction of a working Morse Homology theory in the Hölder space framework.

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## 2 - The Solution Space.

2.1-Preliminaries and Compactness. For $\lambda \in[0, \infty] \backslash \mathbb{N}$, that is, for $\lambda=+\infty$, or for $\lambda=k+\alpha$, where $k \in \mathbb{N}$ and $\alpha \in] 0,1\left[\right.$, we denote by $C^{\lambda}:=C^{\lambda}(M)$ the space of $\lambda$-times Hölder differentiable functions over $M$, and when $\lambda<\infty$, we denote by $\|\cdot\|_{\lambda}$ the corresponding $C^{\lambda}$-Hölder norm. For $\mu \in[0, \infty] \backslash \mathbb{N}$, we likewise denote by $\mathcal{M}^{\mu}$ the space of $C^{\mu}$-Riemannian metrics over $M$. It is well known that these spaces are non-separable, but as indicated in the introduction, this is of no consequence to us, and is satisfactorily treated by the concept of paraproperness (c.f. Section 2.2, below).

We consider a slightly more general problem than that discussed in the introduction. Let $f: \mathbb{R} \rightarrow \mathbb{R}$ be a smooth function such that $f$ is not linear over any interval, both $f$ and $f^{\prime}$ have non-degenerate zeroes, and:

$$
\begin{equation*}
\operatorname{LimSup}_{t \rightarrow-\infty} f(t)<0, \quad \operatorname{LimInf}_{t \rightarrow+\infty} f(t)>0 . \tag{A}
\end{equation*}
$$

As we will see presently (c.f. Proposition 2.1.2, below) our theory only depends on the restriction of $f$ to the smallest interval containing all its zeroes. We therefore modify $f$ outside this interval, and replace (A) with the following technically more convenient property:

$$
\begin{equation*}
\operatorname{Lim}_{t \rightarrow \pm \infty} f(t) / t|t|=+\infty \tag{B}
\end{equation*}
$$

For $\mu>\lambda \in[0, \infty] \backslash \mathbb{N}$, we define the Allen-Cahn Operator, $\mathrm{AC}:] 0, \infty\left[\times \mathcal{M}^{\mu+1} \times C^{\lambda+2} \rightarrow\right.$ $C^{\lambda}$ by:

$$
\begin{equation*}
\mathrm{AC}_{\epsilon, g}(u):=\mathrm{AC}(\epsilon, g, u)=\epsilon \Delta_{g} u-f(u) \tag{C}
\end{equation*}
$$

where $\Delta_{g}$ is the Laplacian operator of $g$. Since AC is constructed via a finite combination of multiplication, addition, differentiation and post-composition by smooth functions, it defines a smooth mapping between Banach manifolds. Importantly, the Allen-Cahn Operator arises as the $L^{2}$-gradient of the Ginzburg-Landau-Wilson Free Energy Functional:

$$
\mathcal{E}_{\epsilon, g}=\int_{M} \epsilon\left\|\nabla^{g} u\right\|^{2}+F(u) \mathrm{dVol},
$$

where $F$ is any primitive of $u$. In particular, solutions of the Allen-Cahn Equation are critical points of $\mathcal{E}_{\epsilon, g}$.
We define the solution space $\mathcal{Z} \subseteq] 0, \infty\left[\times \mathcal{M}^{\mu+1} \times C^{\lambda+2}\right.$ by:

$$
\mathcal{Z}=\mathrm{AC}^{-1}(0)
$$

Let $\Pi: \mathcal{Z} \rightarrow] 0, \infty\left[\times \mathcal{M}^{\mu+1}\right.$ be the projection onto the first two factors and let $\Pi_{g}: \mathcal{Z} \rightarrow$ $\mathcal{M}^{\mu+1}$ and $\Pi_{u}: \mathcal{Z} \rightarrow C^{\lambda+2}$ be the projection onto the second and third factor, respectively. For all $(\epsilon, g) \in] 0, \infty\left[\times \mathcal{M}^{\mu+1}\right.$, we define $\mathcal{Z}_{\epsilon, g} \subseteq \mathcal{Z}$, the solution space for the data $(\epsilon, g)$ by:

$$
\mathcal{Z}_{\epsilon, g}=\Pi^{-1}((\epsilon, g)) .
$$

We study the bifurcations of $\mathcal{Z}_{\epsilon, g}$ as $\epsilon$ varies. For this reason, we prefer to study all values of $\epsilon$ simultaneously and thus define $\mathcal{Z}_{g} \subseteq \mathcal{Z}$ by:

$$
\mathcal{Z}_{g}=\Pi_{g}^{-1}(g)
$$

The main results of this paper follow from the differential topological properties of $\mathcal{Z}, \Pi$ and $\Pi_{g}$, which we now proceed to study.
We first review the analytic properties of the Allen-Cahn Operator. Elements of $\mathcal{Z}$ have the following regularity properties:

## Proposition 2.1.1

Given $\mu>\lambda \in[0, \infty] \backslash \mathbb{N}$, if $(\epsilon, g, u) \in \mathcal{Z}$, then $u \in C^{\mu+2}$.
Proof: Observe that $\epsilon \Delta_{g}$ is a second-order elliptic partial differential operator with coefficients in $C^{\mu}$. Thus, if $u$ lies in $C^{\mu+2(1-k)}$ for some positive integer $k$ with $\mu+2(1-k)>0$, then, since $f$ is smooth:

$$
\epsilon \Delta_{g} u=f \circ u \in C^{\mu+2(1-k)}
$$

and by elliptic regularity (c.f. [6]), $u \in C^{\mu+2(2-k)}$. Observe that since $u \in C^{\lambda+2}$, there exists $k$ such that $u \in C^{\mu+2(1-k)}$, and it follows by induction that $u \in C^{\mu+2}$, as desired.
In order to obtain a-priori estimates, we define $T_{0}>0$ by:

$$
T_{0}=\operatorname{Sup}\{|t| \mid f(t)=0\} .
$$

It follows from (B) that $T_{0}$ is finite. We have:

## Proposition 2.1.2

For all $(\epsilon, g, u) \in \mathcal{Z}$ :

$$
\|u\|_{L^{\infty}} \leqslant T_{0}
$$

Proof: Suppose the contrary, that is, $\|u\|_{L^{\infty}}>T_{0}$. Since $M$ is compact, there exists $p \in M$ such that $|u(p)|=\|u\|_{L^{\infty}}$. If $u(p) \geqslant 0$, then $u(p)=\|u\|_{L^{\infty}}$, and since $p$ is a maximum of $u,\left(\Delta_{g} u\right)(p) \leqslant 0$, so that:

$$
f\left(\|u\|_{L^{\infty}}\right)=\epsilon\left(\Delta_{g} u\right)(p) \leqslant 0
$$

On the other hand, if $u(p)<0$, then $u(p)=-\|u\|_{L^{\infty}}$ and $\left(\Delta_{g} u\right)(p) \geqslant 0$ so that:

$$
f\left(-\|u\|_{L^{\infty}}\right)=\epsilon\left(\Delta_{g} u\right)(p) \geqslant 0
$$

In each case, this is absurd by definition of $T_{0}$ and Property (B) of $f$, and the result follows.

## Proposition 2.1.3

$\Pi$ defines a proper map from $\mathcal{Z}$ into $] 0, \infty\left[\times \mathcal{M}^{\mu+1}\right.$.
Proof: Let $\left(\epsilon_{m}, g_{m}, u_{m}\right)_{m \in \mathbb{N}}$ be a sequence in $\mathcal{Z}$ and suppose that $\left(\epsilon_{m}, g_{m}\right)_{m \in \mathbb{N}}$ converges to $\left.\left(\epsilon_{\infty}, g_{\infty}\right) \in\right] 0, \infty\left[\times \mathcal{M}^{\mu+1}\right.$, say. By Proposition 2.1.1, $u_{m} \in C^{\mu+2}(M)$ for all $M$. By the Schauder estimates (c.f. [6]), there exists $B_{1}>0$ such that for all $m$ :

$$
\begin{aligned}
\left\|u_{m}\right\|_{\mu+2} & \leqslant B_{1}\left(\left\|u_{m}\right\|_{L^{\infty}}+\left\|\epsilon_{m} \Delta_{g_{m}} u_{m}\right\|_{\mu}\right) \\
& =B_{1}\left(\left\|u_{m}\right\|_{L^{\infty}}+\left\|f \circ u_{m}\right\|_{\mu}\right)
\end{aligned}
$$

By Proposition 2.1.2, for all $m, u_{m}$ takes values in the compact set $\left[-T_{0}, T_{0}\right]$. Since $f$ is smooth, it follows from the chain rule and Gagliardo-Nirenberg-Moser type interpolation estimates (c.f. [15]) that there exists $B_{2}>0$ such that for all $m$ :

$$
\left\|f \circ u_{m}\right\|_{\mu} \leqslant B_{2}\left(\left\|u_{m}\right\|_{L^{\infty}}+\left\|u_{m}\right\|_{\mu}\right)
$$

By standard interpolation inequalities (c.f. [6]), there exists $B_{3} \geqslant 0$ such that for all $m$ :

$$
\left\|u_{m}\right\|_{\mu} \leqslant B_{3}\left\|u_{m}\right\|_{L^{\infty}}+\frac{1}{2 B_{1} B_{2}}\left\|u_{m}\right\|_{\mu+2} .
$$

Combining these estimates yields, for all $m$ :

$$
\begin{aligned}
\left\|u_{m}\right\|_{\mu+2} & \leqslant 2 B_{1}\left(1+B_{2}\right)\left(1+B_{3}\right)\left\|u_{m}\right\|_{L^{\infty}} \\
& \leqslant 2 B_{1}\left(1+B_{2}\right)\left(1+B_{3}\right) T_{0} .
\end{aligned}
$$

It now follows by the Arzela-Ascoli Theorem that there exists $u_{\infty} \in C^{\lambda+2}(M)$ towards which $\left(u_{m}\right)_{m \in \mathbb{N}}$ subconverges, and this completes the proof.
2.2 - The Regular Solution Space and Paraproperness. For all $(\epsilon, g, u) \in \mathcal{Z}$, we denote by $L A C$ the linearisation of $\mathrm{AC}_{\epsilon, g}$ about $u$. By definition, $L \mathrm{AC}=D_{3} \mathrm{AC}(\epsilon, g, u)$, where $D_{3} \mathrm{AC}$ denotes the partial derivative of AC with respect to the third component. In particular, for all $(\epsilon, g, u) \in \mathcal{Z}$ and for all $\varphi \in C^{\lambda+2}(M)$ :

$$
\begin{equation*}
L \mathrm{AC} \varphi=\epsilon \Delta_{g} \varphi-f^{\prime}(u) \varphi \tag{D}
\end{equation*}
$$

so that LAC is a self-adjoint second-order elliptic linear operator. In particular, it is Fredholm of index zero and by classical spectral theory, its spectrum is discrete, real and bounded above, and all of its eigenvalues have finite multiplicity. We say that the solution $u$ is non-degenerate whenever $L A C$ is invertible, and we define the Morse Index of $u$, which we denote by $\operatorname{Index}(u)$ by:

$$
\operatorname{Index}(u):=\operatorname{Index}(L A C)=\sum_{\lambda \in \operatorname{Spec}(L A C), \lambda>0} \operatorname{Mult}(\lambda)
$$

where, for all $\lambda \in \operatorname{Spec}(L A C)$, $\operatorname{Mult}(\lambda)$ is its multiplicity.

## Proposition 2.2.1

For all $(\epsilon, g) \in] 0, \infty\left[\times \mathcal{M}^{\mu+1}\right.$, the constant function $u=c$ is a solution to the Allen-Cahn Equation $A C_{\epsilon, g}(u)=0$ if and only if $f(c)=0$. Moreover, this solution is non-degenerate if and only if:

$$
\epsilon^{-1} f^{\prime}(c) \notin \operatorname{Spec}\left(\Delta_{g}\right),
$$

in which case its Morse Index is given by:

$$
\operatorname{Index}(c)=\sum_{\lambda \in \operatorname{Spec}\left(\Delta_{g}\right), \lambda>\epsilon^{-1} f^{\prime}(c)} \operatorname{Mult}(\lambda) .
$$

Proof: The first assertion follows immediately from (C). By definition, $u$ is non-degenerate if and only if $0 \notin \operatorname{Spec}(L A C)$. By (D), this holds if and only if $\epsilon^{-1} f^{\prime}(c) \notin \operatorname{Spec}(L A C)$ and the second assertion follows. Finally, by (D), $\lambda>0$ is an eigenvalue of LAC at $c$ if and only if $\mu:=\lambda+\epsilon^{-1} f(c)$ is an eigenvalue of $\Delta_{g}$. The third assertion then follows, and this completes the proof.
Since $f$ has non-degenerate zeroes, if $f(c)=0, f^{\prime}(c)$ is either positive or negative. When $f^{\prime}(c)$ is positive, it follows from Proposition 2.2 .1 that $u=c$ is always non-degenerate with Morse Index zero. On the other hand, if $f^{\prime}(c)$ is negative, then $u=c$ is degenerate for countably many values of $c$ and its Morse Index tends to $+\infty$ as $\epsilon$ tends to 0 . It follows that zeroes of $f$ with negative derivative behave qualitatively differently from zeroes of $f$ with positive derivative. In fact, they yield singularities which are fundamental in the sense that they cannot be removed by perturbations of the metric, and this will be key to the bifurcation theory that follows. We therefore define the singular set, Sing $\subseteq$ $] 0, \infty\left[\times \mathcal{M}^{\mu+1} \times \mathcal{C}^{\lambda+2}\right.$ by:

$$
\operatorname{Sing}=\left\{(\epsilon, g, c) \mid f(c)=0, \epsilon^{-1} f^{\prime}(c) \in \operatorname{Spec}\left(\Delta_{g}\right)\right\}
$$

and we define the regular solution space, $\mathcal{Z}^{*} \subseteq \mathcal{Z}$ by:

$$
\mathcal{Z}^{*}=\mathcal{Z} \backslash \text { Sing. }
$$

We now construct a countable exhaustion of $\mathcal{Z}^{*}$ by closed sets. For all $g \in \mathcal{M}^{\mu+1}$, we define:

$$
\operatorname{Sing}_{g}=\left\{(\epsilon, c) \mid f(c)=0, \epsilon^{-1} f^{\prime}(c) \in \operatorname{Spec}\left(\Delta_{g}\right)\right\} .
$$

By classical perturbation theory (c.f. [8]), $\operatorname{Sing}_{g}$ varies continuously with $g$ in the Hausdorff sense. For all $m \in \mathbb{N}$, we define $\mathcal{Z}_{m} \subseteq \mathcal{Z}$ by:

$$
\mathcal{Z}_{m}=\left\{(\epsilon, g, u) \in \mathcal{Z} \mid 1 / m \leqslant \epsilon \leqslant n, d\left((\epsilon, u), \operatorname{Sing}_{g}\right) \geqslant 1 / m\right\}
$$

and it follows from the continuous dependence of $\operatorname{Sing}_{g}$ on $g$ that $\mathcal{Z}_{m}$ is closed. Moreover:

$$
\mathcal{Z}^{*}=\cup_{m \in \mathbb{N}} \mathcal{Z}_{m}
$$

We now say that a continuous mapping $\Phi: X \rightarrow Y$ between two topological spaces is paraproper whenever there exists a countable exhaustion $\left(X_{m}\right)_{m \in \mathbb{N}}$ of $X$ by closed sets such that for all $m$, the restriction of $\Phi$ to $X_{m}$ is proper. Paraproperness is made workable as a concept by the following restriction property:

## Proposition 2.2.2

Let $X$ and $Y$ be topological spaces and let $\Phi: X \rightarrow Y$ be paraproper.
(1) if $K \subseteq X$ is closed, then the restriction of $\Phi$ to $K$ is paraproper; and
(2) if $X$ is metrisable and if $\Omega \subseteq X$ is open, then the restriction of $\Phi$ to $\Omega$ is paraproper.

Proof: Indeed, $\left(X_{m} \cap K\right)_{m \in \mathbb{N}}$ is a countable exhaustion of $K$ by closed sets and for all $m$, the restriction of $\Phi$ to $X_{m} \cap K$ is proper, which proves (1). Now let $d$ be a distance function over $X$. For all $n \in \mathbb{N}$, we define $\Omega_{n} \subseteq X$ by:

$$
\Omega_{n}=\left\{x \in X \mid d\left(x, \Omega^{c}\right) \geqslant(1 / n)\right\} .
$$

For all $n, \Omega_{n}$ is closed, and since $\Omega$ is open:

$$
\Omega=\bigcup_{n \in \mathbb{N}} \Omega_{n}
$$

$\left(\Omega_{n} \cap X_{m}\right)_{m, n \in \mathbb{N}}$ therefore constitutes a covering of $\Omega$ by closed sets. Moreover, for all $m, n \in \mathbb{N}$, since $\Omega_{n} \cap X_{m}$ is a closed subset of $X_{m}$, the restriction of $\Phi$ to this set is proper, and the restriction of $\Phi$ to $\Omega$ is therefore paraproper, which proves (2).
In particular $\Pi_{g}$ defines a para-proper map from $\mathcal{Z}$ into $\mathcal{M}^{\mu+1}$ :

## Proposition 2.2.3

For all $n, \Pi_{g}$ defines a proper map from $\mathcal{Z}_{n}$ into $\mathcal{M}^{\mu+1}$.
Proof: Let $\left(\epsilon_{m}, g_{m}, u_{m}\right)_{m \in \mathbb{N}}$ be a sequence in $\mathcal{Z}_{n}$ and suppose that $\left(g_{m}\right)_{m \in \mathbb{N}}$ converges to $g_{\infty} \in \mathcal{M}^{\mu+1}$. Since $\epsilon_{m} \in[1 / n, n]$ for all $n$, by the Heine-Borel Theorem, we may suppose that there exists $\epsilon_{\infty} \in[1 / n, n]$ towards which $\left(\epsilon_{m}\right)_{m \in \mathbb{N}}$ converges. By Proposition 2.1.3, there exists $u_{\infty} \in C^{\lambda+2}(M)$ towards which $\left(u_{m}\right)_{m \in \mathbb{N}}$ subconverges. Since $\mathcal{Z}_{n}$ is closed, $\left(\epsilon_{\infty}, g_{\infty}, u_{\infty}\right) \in \mathcal{Z}_{n}$, and the result follows.

Paraproperness now substitutes separability in our version of the Sard-Smale Theorem (c.f. [13]):

## Theorem 2.2.4, Sard-Smale

If $X$ and $Y$ are smooth Banach manifolds, and if $\Phi: X \rightarrow Y$ is a smooth, paraproper Fredholm map, then the set of regular values of $\Phi$ is generic in $Y$.

Remark: As Smale's result often mystifies, it is worth underlining the straightforward idea behind it. Using Fredholm Theory and the Implicit Function Theorem for Banach manifolds we reduce the problem to one of smooth maps between finite dimensional manifolds, and the result then follows by the classical Sard Theorem.
Proof: Let $\left(X_{n}\right)_{n \in \mathbb{N}}$ be a countable exhaustion of $X$ by closed sets such that for all $n$, the restriction of $\Phi$ to $X_{n}$ is proper. For all $n$, we denote the restriction of $\Phi$ to $X_{n}$ by $\Phi_{n}$, and we denote the set of regular values of $\Phi_{n}$ in $Y$ by $Y_{n}$. Since $\Phi_{n}$ is proper, and since surjectivity of Fredholm maps is an open property, $Y_{n}$ is open for all $n$.

## The Allen-Cahn Equation

We now show that $Y_{n}$ is dense in $Y$. Indeed, choose $y \in Y$. Since we are only concerned with a neighbourhood of $y$ in $Y$, without loss of generality, we may suppose that $Y$ is a Banach space and that $y=0$. Define $\Psi: X \times Y \rightarrow Y$ by $\Psi(\tilde{x}, \tilde{y})=\Phi(\tilde{x})+\tilde{y}$. Now choose $x \in \Phi_{n}^{-1}(0)$. Since $\Phi$ is Fredholm, $D \Phi(x)$ is closed and has finite dimensional cokernel, which we denote by $E_{x}$. In particular, the restriction of $D \Psi(x, 0)$ to $T_{x} X \times E_{x}$ is surjective, and since surjectivity of Fredholm maps is an open property, there exists a neighbourhood $U_{x}$ of $x$ in $X_{n}$ such that the restriction of $D \Psi(\tilde{x}, 0)$ to $T_{\tilde{x}} X \times E_{x}$ is surjective for all $\tilde{x} \in U_{x}$. Since $\Phi_{n}^{-1}(0)$ is compact, it may be covered by finitely many such open sets, and there therefore exists a finite-dimensional subspace $E \subseteq Y$ such that the restriction of $D \Psi(\tilde{x}, 0)$ to $T_{\tilde{x}} X \times E$ is surjective for all $\tilde{x} \in \Phi_{n}^{-1}(y)$. We now consider the restriction of $\Psi$ to $X \times E$ and we denote $Z=\Psi^{-1}(0)$. By the Implicit Function Theorem for Banach manifolds, there exists a neighbourhood $\Omega$ of $\Phi_{n}^{-1}(0) \times\{0\}$ in $Z$ which is a smooth finite-dimensional submanifold of $X \times E$. Moreover, since $\Phi_{n}^{-1}(0)$ is compact, upon reducing $\Omega$ is necessary, we may suppose that this submanifold is separable. Let $\pi: \Omega \rightarrow E$ be the projection onto the first factor. Observe that if $\tilde{y} \in E$ is a regular value of $\pi$, then it is also a regular value of $\Phi_{n}$. However, by Sard's Theorem, regular values of $\pi$ are dense in $E$. It follows that $y=0$ is a concentration point of regular values of $\Phi_{n}$, and $Y_{n}$ is therefore a dense subset of $Y$ as asserted.

Since the set of regular values of $\Phi$ coincides with $\cap_{n \in \mathbb{N}} Y_{n}$, it follows that this set is generic, which completes the proof.
2.3 - The Regular Solution Space. In this section, we prove the following:

## Proposition 2.3.1

If $\operatorname{Dim}(M) \geqslant 2$, then for all $\mu>\lambda \in\left[0, \infty\left[\backslash \mathbb{N}, \mathcal{Z}^{*}\right.\right.$ is a smooth Banach manifold modelled on $\mathbb{R} \times \mathcal{M}^{\mu+1}$. Moreover, $\Pi_{g}$ defines a smooth, paraproper Fredholm map from $\mathcal{Z}^{*}$ into $\mathcal{M}^{\mu}$ of Fredholm index equal to 1.
We prove this result using the Implicit Function Theorem for Banach manifolds. It is thus necessary to show that the derivative of AC is surjective at every point of $\mathcal{Z}^{*}$. We denote by $D_{1} \mathrm{AC}, D_{2} \mathrm{AC}$ and $D_{3} \mathrm{AC}$ the partial derivatives of AC with respect to the first, second and third components in $] 0, \infty\left[\times \mathcal{M}^{\mu+1} \times C^{\lambda+2}\right.$ respectively. We are interested in particular in $D_{2} \mathrm{AC}$. The tangent space of $\mathcal{M}^{\mu+1}$ at any point canonically identifies with the space of $C^{\mu+1}$ sections of $\operatorname{Symm}(T M)$. We denote this space by $\Gamma^{\mu+1}:=\Gamma^{\mu+1}(\operatorname{Symm}(T M))$ and we refer to elements therin as first order perturbations of the metric. We then identify $C^{\mu+1}$ with a subspace of $\Gamma^{\mu+1}$ by identifying every $C^{\mu+1}$ function $f$ with the $C^{\mu+1}$ section $f g$, and this induces the orthogonal splitting, $\Gamma^{\mu+1}=\Gamma_{0}^{\mu+1} \oplus C^{\mu+1}$, where $\Gamma_{0}^{\mu+1}:=\Gamma^{\mu+1}\left(\operatorname{Symm}_{0}(T M)\right)$ is the space of trace-free sections of $\operatorname{Symm}(T M)$. The first order perturbations arising from sections of $C^{\mu+1}$ are precisely the conformal perturbations of the metric. However, it turns out that the useful perturbations for us are those whose trace vanishes. Indeed, for $g \in \mathcal{M}^{\mu+1}$, and for any first order perturbation $A$ of $g$, denoting by $\delta_{A} \Delta_{g}$ the resulting first order perturbation of $\Delta_{g}$, we obtain:

## Proposition 2.3.2

If $\operatorname{Tr}(A)=0$, then, viewing $A$ as a section of $\operatorname{End}(T M)$, for all $\varphi \in C^{\lambda+2}$ :

$$
\left(\delta_{A} \Delta_{g}\right) \varphi=-\nabla \cdot(A \nabla \varphi)
$$

## The Allen-Cahn Equation

where $\nabla$ and $\nabla$. are the gradient and divergence operators of $g$ respectively.
Proof: We denote respectively by $\delta_{A} \Omega_{g}$ and $\delta_{A} \operatorname{Hess}_{g}$ the first order perturbations resulting from $A$ of the Levi-Civita covariant derivative and the Hessian operator of $g$. The Koszul formula yields:

$$
\left(\delta_{A} \Omega_{g}\right)^{k}{ }_{; i j}=\frac{1}{2}\left(A_{i ; j}^{k}+A_{j ; i}^{k}-A_{i j ;}^{k}\right),
$$

where indices are raised and lowered with respect to $g$. Thus:

$$
\begin{aligned}
\delta_{A} \operatorname{Hess}_{g}(u)_{i j} & =-\frac{1}{2}\left(A^{k}{ }_{i ; j}+A^{k}{ }_{j ; i}-A_{i j ;}{ }^{k}\right) u_{; k} \\
\Rightarrow \quad \delta_{A} \Delta_{g}(u) & =-A_{j}{ }^{i} u_{; i}{ }^{j}-A_{j}{ }^{i j j} u_{; i}+\frac{1}{2} \operatorname{Tr}(A)^{; k} u_{; k} \\
& =-\left(A_{j}{ }^{i} u_{; i}\right)^{; j}+\frac{1}{2} \operatorname{Tr}(A)^{; k} u_{; k},
\end{aligned}
$$

and since $\operatorname{Tr}(A)=0$, the result follows.
We recall the following straightforward result:

## Proposition 2.3.3

Let $X$ be a set consisting of at least $n$ distinct points. Let $E$ be an n-dimensional subset of the space of real-valued functions over $X$. Then there exist $n$ points $p_{1}, \ldots, p_{n} \in X$ such that the mapping Eval: $E \rightarrow \mathbb{R}^{n}$ given by:

$$
\operatorname{Eval}(f)_{k}=f\left(p_{k}\right)
$$

is a linear isomorphism.
This allows us to prove the required surjectivity result:

## Proposition 2.3.4

If $\operatorname{Dim}(M) \geqslant 2$, if $(u, g, \epsilon) \in \mathcal{Z}$ and if $u$ is non-constant, then $D A C$ is surjective at $(u, g, \epsilon)$.
Proof: Since $D_{3} \mathrm{AC}=L A C$ is elliptic, it has finite-dimensional cokernel, which we denote by $E$. Since $D_{3} A C$ is self-adjoint with respect to the $L^{2}$-inner-product of $g$, for all $\varphi \in E$ :

$$
D_{3} \mathrm{AC}(\varphi)=-\epsilon \Delta_{g} \varphi+f^{\prime}(u) \varphi=0 .
$$

Since $u$ is non-constant, there exists $p \in M$ such that $\nabla u(p) \neq 0$. Let $\Omega$ be a neighbourhood of $p$ in $M$ diffeomorphic to the unit ball in Euclidean space over which $\nabla u$ doesn't vanish. By Aronszajn's unique continuation theorem (c.f. [3]), no non-trivial element of $E$ vanishes over $\Omega$. Furthermore, since $f^{\prime}$ has non-degenerate zeroes, $f^{\prime}(u)$ does not vanish identically over $\Omega$, and therefore no non-zero element of $E$ restricts to a constant map over this set. Thus, by Proposition 2.3.3, there exist $p_{1}, \ldots, p_{m} \in \Omega \backslash\{p\}$ such that the mapping $\alpha: C^{\lambda} \rightarrow \mathbb{R}^{m}$ given by:

$$
\alpha(\varphi)_{k}=\varphi(p)-\varphi\left(p_{k}\right),
$$

restricts to a bijection on $E$.

For any vector $\xi:=\left(\xi_{0}, \ldots, \xi_{m}\right)$ of functions in $C_{0}^{\infty}(\Omega)$ we define $\alpha_{\xi}: C^{\lambda} \rightarrow \mathbb{R}^{m}$ by:

$$
\alpha_{\xi}(\varphi)_{i}=\int_{M}\left(\xi_{0}-\xi_{i}\right) \varphi \mathrm{dVol}_{g}
$$

If $\xi_{0}-\xi_{k}$ is sufficiently close to $\delta_{p}-\delta_{p_{k}}$ in the weak sense for all $k$, where $\delta_{p}$ and $\delta_{p_{k}}$ are the Dirac delta functions supported at $p$ and $p_{k}$ respectively, then $\alpha_{\xi}$ is close to $\alpha$ and, in particular, is invertible. It follows that if $F$ is the linear span of $\left(\xi_{0}-\xi_{k}\right)_{1 \leqslant k \leqslant m}$, then the $L^{2}$-inner-product restricts to a non-degenerate bilinear form over $E \times F$. In particular, $\operatorname{Dim}(F)=\operatorname{Dim}(E)$ and:

$$
F \cap \operatorname{Im}\left(D_{3} \mathrm{AC}\right)=F \cap E^{\perp}=\{0\}
$$

$F$ is therefore complementary to $\operatorname{Im}\left(D_{3} \mathrm{AC}\right)$ in $C^{\lambda}$. That is:

$$
C^{\lambda}=F \oplus \operatorname{Im}\left(D_{3} \mathrm{AC}\right)
$$

However, we may suppose in addition that for all $1 \leqslant k \leqslant m$ :

$$
\int_{M}\left(\xi_{0}-\xi_{k}\right) \mathrm{dVol}_{g}=0
$$

It then follows from classical de-Rham cohomology theory that for all $k$ there exists a smooth vector field $X_{k}$ supported in $\Omega$ such that:

$$
\nabla \cdot X_{k}=\xi_{0}-\xi_{k}
$$

By Proposition 2.1.1, $\nabla u$ is of class $C^{\mu+1}$, and thus, since it is non-vanishing over $\Omega$, there exists for all $k$ a $C^{\mu+1}$ field $A_{k}$ of symmetric matrices such that $A_{k} \nabla u=X_{k}$. In addition, since $M$ has dimension at least 2, we may assume moreover that $\operatorname{Tr}\left(A_{k}\right)=0$ for all $k$, and it follows from Proposition 2.3.2 that:

$$
D_{2} \mathrm{AC} \cdot A_{k}=-\epsilon \nabla \cdot\left(A_{k} \nabla u\right)=\epsilon \nabla \cdot X_{k}=\epsilon\left(\xi_{0}-\xi_{k}\right) .
$$

It follows that $F \subseteq \operatorname{Im}\left(D_{2} \mathrm{AC}\right)$ and so $C^{\lambda} \subseteq \operatorname{Im}(D A C)$ and surjectivity follows.
Proposition 2.3.1 follows readily:
Proof of Proposition 2.3.1: By Propositions 2.2.1 and 2.3.4, DAC is surjective at every point of $\mathcal{Z}^{*}$. Since $D_{3} A C$ is self-adjoint and elliptic, it is Fredholm of index zero, and it follows from the Implicit Function Theorem for Banach manifolds that $\mathcal{Z}^{*}$ is a smooth Banach manifold modelled on $\mathbb{R} \times \mathcal{M}^{\mu+1}$ and $\Pi_{g}$ is a smooth Fredholm map of Fredholm index equal to 1 . Finally, by Proposition $2.2 .3, \Pi_{g}$ is para-proper, and this completes the proof.

Applying the Sard/Smale Theorem, now yields:

## Proposition 2.3.5

If $\operatorname{Dim}(M) \geqslant 2$, then for generic $g \in \mathcal{M}^{\mu+1}, \mathcal{Z}_{g}^{*}$ is a smooth, 1-dimensional submanifold of $] 0, \infty\left[\times C^{\lambda+2}\right.$. Moreover, if we denote by $\left.\epsilon_{g}: \mathcal{Z}_{g}^{*} \rightarrow\right] 0, \infty[$ the projection onto the first factor, then $\epsilon_{g}$ is proper.

Remark: Observe that paraproperness allows us to show that $\mathcal{Z}_{g}^{*}$ is separable even though $\mathcal{Z}^{*}$ isn't.
Proof: By the Sard-Smale Theorem, the set of regular values of $\Pi_{g}$ is generic in $\mathcal{M}^{\mu+1}$. Let $g \in \mathcal{M}^{\mu+1}$ be a regular value of $\Pi_{g}$. By definition, $D \Pi_{g}(\epsilon, u, g)$ is surjective for all $(\epsilon, u) \in \mathcal{Z}_{g}^{*}$. Since $\Pi_{g}$ is a smooth Fredholm map of Fredholm index equal to 1, it follows from the Implicit Function Theorem for Banach manifolds that $\mathcal{Z}_{g}^{*}$ is a (not necessarily separable) smooth, 1-dimensional submanifold of $] 0, \infty\left[\times C^{\lambda+2}\right.$. By Proposition 2.1.3, $\left.\epsilon_{g}: \mathcal{Z}_{g} \rightarrow\right] 0, \infty[$ is proper, and since $] 0, \infty\left[\right.$ has a compact exhaustion, so too does $\mathcal{Z}_{g}$. In particular, $\mathcal{Z}_{g}$ is separable, and therefore so too is $\mathcal{Z}_{g}^{*}$, which completes the proof.
2.4 - Non-Degeneracy of Critical Points. Let $\epsilon:] 0, \infty\left[\times \mathcal{M}^{\mu+1} \times C^{\lambda+2} \rightarrow\right] 0, \infty[$ be the projection onto the first factor, and denote its restriction to $\mathcal{Z}_{g}^{*}$ by $\epsilon_{g}$. We now aim to show that for generic $g \in \mathcal{M}^{\mu+1}$, every critical point $\epsilon_{g}$ is non-degenerate. We first characterise those points where $d \epsilon_{g}$ vanishes:

## Proposition 2.4.1

At every point of $\mathcal{Z}^{*}$ :

$$
\operatorname{Ker}(D \epsilon) \cap \operatorname{Ker}\left(D \Pi_{g}\right) \cap T \mathcal{Z}^{*}=\{0\} \times\{0\} \times \operatorname{Ker}(L A C)
$$

Proof: Indeed, by definition:

$$
\begin{aligned}
T \mathcal{Z}^{*} & =\operatorname{Ker}(D \mathrm{AC}) \\
& =\operatorname{Ker}\left(D_{1} \mathrm{AC} \circ D \epsilon+D_{2} \mathrm{AC} \circ D \Pi_{g}+D_{3} \mathrm{AC} \circ D \Pi_{u}\right),
\end{aligned}
$$

where $D_{1} \mathrm{AC}, D_{2} \mathrm{AC}$ and $D_{3} \mathrm{AC}$ represent the partial derivatives of AC with respect to the first, second and third components respectively. Thus:

$$
\begin{aligned}
\operatorname{Ker}(D \epsilon) \cap \operatorname{Ker}\left(D \Pi_{g}\right) \cap T \mathcal{Z}^{*} & =\operatorname{Ker}(D \epsilon) \cap \operatorname{Ker}\left(D \Pi_{g}\right) \cap \operatorname{Ker}\left(D_{3} \mathrm{AC} \circ D \Pi_{u}\right) \\
& =\{0\} \times\{0\} \times \operatorname{Ker}(L A C),
\end{aligned}
$$

as desired.

## Proposition 2.4.2

If $g$ is a regular value of $\Pi_{g}$, then at every point of $\mathcal{Z}_{g}^{*}$ :

$$
\operatorname{Ker}\left(d \epsilon_{g}\right)=\operatorname{Ker}(D \epsilon) \cap T \mathcal{Z}_{g}^{*}=\{0\} \times\{0\} \times \operatorname{Ker}(L A C)
$$

In particular, LAC has nullity at most 1 .
Proof: If $g$ is a regular value of $\Pi_{g}$, then:

$$
\operatorname{Ker}\left(D \Pi_{g}\right) \cap T \mathcal{Z}^{*}=T \mathcal{Z}_{g}^{*}
$$

and the result now follows by Proposition 2.4.1.
By Proposition 2.4.2, if $g$ is a regular value of $\Pi_{g}$ and if $p \in \mathcal{Z}_{g}$ is such that $d \epsilon_{g}=0$, then $\operatorname{Ker}(L A C)$ is 1 -dimensional. In particular, we may split $C^{\lambda+2}$ as the direct sum of $\operatorname{Ker}(L \mathrm{AC})$ and $\operatorname{Ker}(L \mathrm{AC})^{\perp}$ where $\operatorname{Ker}(L A C)^{\perp}$ is the orthogonal complement of $\operatorname{Ker}(L A C)$ in $C^{\lambda+2}$ with respect to the $L^{2}$-inner-product.

## Proposition 2.4.3

If $g$ is a regular value of $\Pi_{g}$, and if $p \in \mathcal{Z}_{g}$ is such that $d \epsilon_{g}(p)=0$, then there is a neighbourhood $\Omega$ of $p$ in $\mathcal{Z}^{*}$ which is a graph over $\mathcal{M}^{\mu+1} \times \operatorname{Ker}(L A C)$.
Proof: Let $\pi: C^{\lambda+2} \rightarrow \operatorname{Ker}(L A C)$ be the orthogonal projection. Consider the restriction of the mapping $\left(\Pi_{g}, \pi \circ \Pi_{u}\right)$ to $\mathcal{Z}^{*}$. Since $d \epsilon_{g}(p)=0$, bearing in mind Proposition 2.4.2, at $p$ :

$$
\operatorname{Ker}\left(D \Pi_{g}\right) \cap T \mathcal{Z}^{*}=T \mathcal{Z}_{g}^{*}=\operatorname{Ker}\left(d \epsilon_{g}\right) \cap T \mathcal{Z}_{g}^{*}=\{0\} \times\{0\} \times \operatorname{Ker}(L \mathrm{AC})
$$

In particular, the restriction of $\pi \circ \Pi_{u}$ to $\operatorname{Ker}\left(D \Pi_{g}\right)$ is a linear isomorphism. The restriction of $\left(\Pi_{g}, \pi \circ \Pi_{u}\right)$ to $T \mathcal{Z}^{*}$ is therefore also a linear isomorphism at $p$ and the result now follows by the Inverse Function Theorem for smooth maps between Banach manifolds.
Let $\Omega \subseteq \mathcal{Z}^{*}$ be as in Proposition 2.4.3. We construct a non-vanishing vector field, $X$ over $\Omega$ which is always tangent to $\mathcal{Z}_{g}$ as follows. Choose $\varphi_{0} \in \operatorname{Ker}(L A C)$ such that $\|\varphi\|_{L^{2}}^{2}=1$. Let $X$ be the unique, smooth vector field over $\Omega$ which projects down to $\varphi$. There exist smooth functions $s: \Omega \rightarrow \mathbb{R}$ and $\varphi: \Omega \rightarrow \varphi_{0}+\operatorname{Ker}(L A C)^{\perp}$ such that, throughout $\Omega$ :

$$
X=(s, 0, \varphi)
$$

Trivially:

$$
D \Pi_{g} \cdot X=0
$$

so that $X$ is always tangent to $\mathcal{Z}_{g}$, as desired.
We now recall the following formula for the variation of a non-degenerate eigenvalue. Let $E \subseteq F \subseteq L^{2}(M)$ be Banach spaces and let $i: E \rightarrow F$ be a continuous embedding with dense image. It is normal to suppress $i$ and identify elements of $E$ with their image in $F$. Let $A \in \operatorname{Lin}(E, F)$ be a bounded, linear map. We recall that $E$ is said to be self-adjoint if and only if for all $u, v \in E$ :

$$
\langle u, A(v)\rangle=\langle A(u), v\rangle
$$

The Implicit Function Theorem for Banach manifolds readily yields:

## Proposition 2.4.4

Let $X, E$ and $F$ be Banach spaces. Let $A: X \rightarrow \operatorname{Lin}(E, F)$ be a smooth mapping such that for all $x \in X, A_{x}:=A(x)$ is self-adjoint and Fredholm of index zero. Suppose that $\operatorname{Null}\left(A_{0}\right)=1$ and let $\varphi_{0}$ be a non-zero element of $\operatorname{Ker}\left(A_{0}\right)$. Then there exists a neighbourhood $U$ of 0 in $X$ and smooth maps $\lambda: X \rightarrow \mathbb{R}$ and $\varphi: X \rightarrow \varphi+\operatorname{Ker}\left(A_{0}\right)^{\perp}$ such that $\lambda(0)=0, \varphi(0)=\varphi_{0}$ and for all $x \in X$ :

$$
A(x) \varphi(x)=\lambda(x) \varphi(x)
$$

Moreover, for any tangent vector $\xi$ to $X$ at 0 :

$$
d \lambda(\xi)=\left\langle D A_{0}(\xi) \varphi_{0}, \varphi_{0}\right\rangle
$$

By Proposition 2.4.4, upon reducing $\Omega$ if necessary, there exist smooth functions $\lambda: \Omega \rightarrow \mathbb{R}$ and $\tilde{\varphi}: \Omega \rightarrow \varphi_{0}+\operatorname{Ker}\left(L A C\left(p_{0}\right)\right)^{\perp}$ such that $\lambda\left(p_{0}\right)=0, \tilde{\varphi}\left(p_{0}\right)=\varphi_{0}$ and throughout $\Omega$ :

$$
L \mathrm{AC} \tilde{\varphi}=\lambda \tilde{\varphi}
$$

The role played by $\lambda$ is revealed by the following quantitative analogue of Proposition 2.4.2:

## Proposition 2.4.5

Let $g \in \mathcal{M}^{\mu+1}$ be a regular value of $\Pi_{g}$. If $p \in \mathcal{Z}_{g}$ is such that $d \epsilon_{g}(p)=0$, then:

$$
\left\langle D_{1} A C \cdot D^{2} \epsilon_{g}\left(X_{p}, X_{p}\right), \varphi_{0}\right\rangle=-d \lambda\left(X_{p}\right)
$$

In particular, if $p$ is a non-degenerate zero of $\lambda$, then it is also a non-degenerate zero of $d \epsilon_{g}$.
Proof: By definition, AC vanishes over $\mathcal{Z}^{*}$ and so:

$$
D_{1} \mathrm{AC} \circ D \epsilon(X)=-D_{2} \mathrm{AC} \circ D \Pi_{g}(X)-D_{3} \mathrm{AC} \circ D \Pi_{u}(X)=-L \mathrm{AC} \circ D \Pi_{u}(X)
$$

Since $D \epsilon\left(X_{p}\right)=0$, differentiating a second time yields:

$$
D_{1} \mathrm{AC} \circ D^{2} \epsilon\left(X_{p}, X_{p}\right)=-D_{X_{p}} L \mathrm{AC} \circ D \Pi_{u}\left(X_{p}\right)-L \mathrm{AC} \circ D^{2} \Pi_{u}\left(X_{p}, X_{p}\right)
$$

Observe that $D^{2} \Pi_{u}\left(X_{p}, X_{p}\right)$ takes values in $\operatorname{Ker}(L A C)^{\perp}$. Moreover, $\varphi_{0} \in \operatorname{Ker}(L A C)$, and since $L A C$ preserves both $\operatorname{Ker}(L A C)$ and $\operatorname{Ker}(L A C)^{\perp}$, taking the inner-product with $\varphi_{0}$ yields:

$$
\begin{aligned}
\left\langle D_{1} \mathrm{AC} \circ D^{2} \epsilon\left(X_{p}, X_{p}\right), \varphi_{0}\right\rangle & =-\left\langle D_{X_{p}} L \mathrm{AC} \circ D \Pi_{u}\left(X_{p}\right), \varphi_{0}\right\rangle \\
& =-\left\langle D_{X_{p}} L \mathrm{AC} \varphi_{0}, \varphi_{0}\right\rangle .
\end{aligned}
$$

Thus, by Proposition 2.4.4:

$$
\left\langle D_{1} \mathrm{AC} \circ D^{2} \epsilon_{g}\left(X_{p}, X_{p}\right), \varphi_{0}\right\rangle=-d \lambda\left(X_{p}\right),
$$

as desired.
The above discussion is most usefully summarised as follows:

## Proposition 2.4.6

There exists an open subset $\Omega \subseteq \mathcal{Z}^{*}$ and a smooth function $\lambda: \Omega \rightarrow \mathbb{R}$ with the following properties:
(1) if $p \in \mathcal{Z}^{*}$ is such that $D \Pi_{g}(p)$ is surjective and $d \epsilon_{g}(p)=0$, then $p \in \Omega$ and $\lambda=0$; and (2) for all $p \in \Omega, \lambda(p)$ is an eigenvalue of $L A C(p)$.

Proof: Let $p \in \mathcal{Z}^{*}$ be such that $D \Pi_{g}(p)$ is surjective and $D \epsilon_{g}(p)=0$. Let $\Omega_{p}$ and $\lambda: \Omega_{p} \rightarrow \mathbb{R}$ be as in the preceeding discussion. Upon reducing $\Omega_{p}$ if necessary, we may assume that $\lambda$ is the eigenvalue of $L A C(p)$ with least absolute value. It follows then that $\lambda$ is uniquely defined, and taking the union over all such $\Omega_{p}$ yields the desired open set and smooth function.

## Proposition 2.4.7

Suppose that $\operatorname{Dim}(M) \geqslant 3$. Choose $p \in \Omega$ and let $\varphi \in C^{\lambda+2}$ be an element of $\operatorname{Ker}(L A C(p))$. If there exists a point $x \in M$ such that $d u(x)$ and $d \varphi(x)$ are both non-vanishing and noncolinear, then $d \lambda$ is non-zero at $p$.

Proof: Let $A$ be a trace-free first order perturbation of $g$ and let $\delta_{A} \mathrm{dVol}_{g}, \delta_{A} L A C$ and $\delta_{A} \lambda$ denote the resulting first order perturbations of $\mathrm{dVol}_{g}, L A C$ and $\lambda$ respectively. Then:

$$
\delta_{A} \mathrm{dVol}_{g}=\operatorname{Tr}(A) \mathrm{dVol}_{g}=0
$$

Thus, by Proposition 2.4.4:

$$
\delta_{A} \lambda=\int \varphi\left(\delta_{A} L \mathrm{AC}\right) \varphi \mathrm{dVol}_{g}
$$

However, by Proposition 2.3.2:

$$
\delta_{A} L \mathrm{AC} \varphi=-\epsilon \nabla \cdot(A \nabla \varphi),
$$

and so:

$$
\begin{aligned}
\delta_{A} \lambda & =-\epsilon \int \varphi \nabla \cdot(A \nabla \varphi) \mathrm{dVol}_{g} \\
& =\epsilon \int\langle A, \nabla \varphi \otimes \nabla \varphi\rangle \mathrm{dVol}_{g} .
\end{aligned}
$$

Let $p \in M$ be such that $d u(p)$ and $d \varphi(p)$ are non-vanishing and non-colinear. Since $M$ is 3 -dimensional, there exists a first-order perturbation $A$ of $g$, supported near $p$ such that:

$$
\operatorname{Tr}(A)=0, \quad\left(\delta_{A} \Delta_{g}\right) u=0, \quad \delta_{A} \lambda \neq 0
$$

It follows from the first two relations that the vector $(0, A, 0)$ is tangent to $\mathcal{Z}^{*}$, and it follows from the third relation that $d \lambda(0, A, 0) \neq 0$, which completes the proof.

Applying the Sard/Smale Theorem, we now obtain:

## Proposition 2.4.8

If $\operatorname{Dim}(M) \geqslant 3$, then for generic $g \in \mathcal{M}^{\mu+1}$ and for $p \in \mathcal{Z}_{g}$, if $d \epsilon_{g}(p)=0$, then either:
(1) $D^{2} \epsilon_{g}(p) \neq 0$; or
(2) if $\varphi \in \operatorname{Ker}(L A C(p))$, then for all $Y \in T M$, if $d u(Y)=0$, then $d \varphi(Y)=0$.

Proof: Let $X \subseteq \Omega$ be the set of all points $p$ such that $d \epsilon_{g}(p)=0$ and (2) is satisfied. Observe that (2) implies that $d u$ and $d \varphi$ are everywhere colinear. Since this is a closed condition, it follows that $X$ is a closed subset of $\Omega$, and $\tilde{\Omega}:=\Omega \backslash X$ is therefore open. Let $Y \subseteq \tilde{\Omega}$ be the set of all points where $\lambda$ vanishes. By Proposition 2.4.7 and the Implicit Function Theorem for Banach manifolds, $Y$ is a smooth codimension-1 Banach submanifold of $\tilde{\Omega}$. Observe that the restriction of $\Pi_{g}$ to $Y$ is a smooth Fredholm map of Fredholm index 0 . Moreover, by Proposition 2.2.2, this restriction is paraproper. It therefore follows from Theorem 2.2.4 that for generic $g \in \mathcal{M}^{\mu+1}, g$ is a regular value of this restriction. Moreover, since the intersection of two generic sets is also generic, we may assume that $g$ is also a regular value of $\Pi_{g}$. For such a $g, \mathcal{Z}_{g}$ is a smooth 1-dimensional manifold and the restriction of $\lambda$ to $\mathcal{Z}_{g}$ has non-degenerate zeroes at all points where (2) is satisfied, and the result now follows by Proposition 2.4.5.
2.5 - The Degenerate Case. We now eliminate Case (2) of Proposition 2.4.8. We begin by characterising its geometry:

## Proposition 2.5.1

Let $u, \varphi \in C^{\lambda+2}(M)$ be such that $u$ is non-constant, $\varphi$ is non-zero, and:

$$
\epsilon \Delta_{g} u=f(u), \quad \epsilon \Delta_{g} \varphi=f^{\prime}(u) \varphi
$$

If for all vectors $X \in T M$ such that $d u(X)=0$ we have $d \varphi(X)=0$, then $\|d u\|_{g}$ is constant over each connected component of every level set of $u$.
Remark: In fact, we can prove more: the complement of the vanishing set of $d u$ in $M$ is foliated by compact hypersurfaces of constant mean curvature. This property is interesting, as it is independent of the parameter $\epsilon$.
Proof: Observe that if $u$ is constant over any non-trivial neighbourhood, then it is equal to a zero of $f, c$ say over this neighbourhood. Since $u=c$ is also a solution of $\mathrm{AC}(u)=0$, it follows from Aronszajn's unique continuation theorem (c.f. [3]) that $f=c$ over the whole of $M$, which is absurd, and it follows that $d u$ is almost everywhere non-vanishing.
Now choose $p \in M$ such that $d u(p) \neq 0$. Let $\Omega$ be a neighbourhood of $p$ over which $d u$ does not vanish. Observe that the image of the restriction of $u$ to $\Omega$ is an open interval, $I$, say. Moreover, by Aronszajn's unique continuation theorem again, the restriction of $\varphi$ to $\Omega$ is non-zero. Let $\mathcal{F}$ denote the foliation of $\Omega$ by level hypersurfaces of $u$. By hypothesis, $\varphi$ is constant over each leaf of $\mathcal{F}$. Thus, upon reducing $\Omega$ if necessary, there exists a non-zero $C^{\lambda+2}$-function $\Phi: I \rightarrow \mathbb{R}$ such that, over $\Omega, \varphi=\Phi(u)$. Taking the Laplacian of both sides of this relation yields:

$$
f^{\prime}(u) \Phi(u)=\epsilon \Phi^{\prime \prime}(u)\|d u\|_{g}^{2}+f(u) \Phi^{\prime}(u)
$$

We claim that $\Phi^{\prime \prime}$ is almost everywhere non-vanishing. Indeed, otherwise, upon reducing $\Omega$ further if necessary, we may suppose that $\Phi$ is linear and that $f^{\prime} \Phi-f \Phi^{\prime}=0$. The restriction of $f$ to $I$ is therefore also linear, which is absurd by the hypothesis on $f$, and $\Phi^{\prime \prime}$ is therefore almost everywhere non-vanishing, as asserted. However, whenever $\Phi^{\prime \prime}(u) \neq 0$, we have:

$$
\|d u\|_{g}^{2}=\frac{1}{\epsilon \Phi^{\prime \prime}(u)}\left(f^{\prime}(u) \Phi(u)-f(u) \Phi^{\prime}(u)\right)
$$

from which it follows that $\|d u\|_{g}^{2}$ is constant over every leaf of $\mathcal{F}$ where $\Phi^{\prime \prime}(u)$ does not vanish. Since the set of all such leaves is dense, it follows that $\|d u\|_{g}$ is constant over every leaf of $\mathcal{F}$.
Choose $t \in \mathbb{R}$ and denote $X=u^{-1}(t)$. Let $X_{0}, X_{1} \subseteq X$ be respectively the subset of $X$ consisting of those points where $d u$ vanishes, and the subset of $X$ consisting of those points where it does not vanish. Trivially, $\|d u\|_{g}$ is constant over $X_{0}$. Observe that $X_{1}$ is a submanifold of $M$. Moreover, by the above discussion, $\|d u\|_{g}$ is constant over every connected component of $X_{1}$. Every connected component of $X_{1}$ is therefore a closed submanifold, and, in particular, is disjoint from $X_{0}$. It follows that if $X^{\prime}$ is a connected component of $X$, then $X^{\prime}$ is either contained wholly in $X_{0}$ or wholly in $X_{1}$. In either case, $\|d u\|_{g}$ is constant over $X^{\prime}$, and this completes the proof.
The following refinement of Proposition 2.5.1 is easier to work with:

## Proposition 2.5.2

Under the same hypotheses as Proposition 2.5.1, if $X_{p} \in T M$ is such that $d u\left(X_{p}\right)=0$, then:

$$
\operatorname{Hess}^{g}(u)\left(\nabla^{g} u, X_{p}\right)=0 .
$$

Proof: If $d u(p)=0$, then $\nabla^{g} u(p)=0$, and the result follows trivially. Otherwise, $d u(p) \neq$ 0 , and, by Proposition 2.5.1, $\left\|\nabla^{g} u\right\|=\|d u\|_{g}$ is constant over the level hypersurface of $u$ passing through $p$. Since $d u\left(X_{p}\right)=0, X_{p}$ is tangent to $S$, and so:

$$
\operatorname{Hess}(u)\left(\nabla^{g} u, X_{p}\right)=\left\langle\nabla_{X_{p}} \nabla^{g} u, \nabla^{g} u\right\rangle=\frac{1}{2} X_{p}\left\|\nabla^{g} u\right\|^{2}=0,
$$

as desired.

## Proposition 2.5.3

Suppose that $\operatorname{Dim}(M) \geqslant 2$ and let $u \in C^{\lambda+2}$ be a non-constant function such that $\epsilon \Delta_{g} u=$ $f(u)$. Choose $p \in M$ such that $\nabla^{g} u(p) \neq 0$ and $Y \in T_{p} M$ such that $d u(Y)=0$. There exists a $C^{\mu+1}$ first order perturbation $A$ of the metric supported in an arbitrarily small neighbourhood of $p$ such that such that:
(1) $A(p)=0$;
(2) $\left(\delta_{A} \Delta_{g}\right) u=0$; and
(3) $\delta_{A} \operatorname{Hess}(u)(\nabla u, Y)(p) \neq 0$.

Proof: Let $\Omega$ be a neighbourhood of $p$ diffeomorphic to the unit ball. Let $X$ be a smooth divergence-free vector field supported in $\Omega$ such that $X(p)=0$. Since $M$ is at least two dimensional, and since $\nabla u$ does not vanish over $\Omega$, there exists a $C^{\mu+2}$ section $A$ of $\operatorname{Symm}(T M)$ supported in $\Omega$ such that $A \cdot \nabla u=X$ and $\operatorname{Tr}(A)=0$. In particular, we may suppose that $A(p)=0$. By Proposition 2.3.2, for any such $A$ :

$$
\left(\delta_{A} \Delta_{g}\right) u=0
$$

As in the proof of Proposition 2.3.2, the first order perturbation of the Hessian of $u$ is given by:

$$
\delta_{A} \operatorname{Hess}_{g}(u)(X, Y)=\frac{1}{2}\left(A_{i j}{ }^{k}-A_{i ; j}^{k}-A_{j ; i}^{k}\right) u_{k} X^{i} Y^{j}
$$

Thus, bearing in mind that $\left(\delta_{A} \nabla^{g}\right) u(p)=X(p)=0$ :

$$
\begin{aligned}
\delta_{A} \operatorname{Hess}_{g}(u)\left(\nabla^{g} u, Y\right) & =\frac{1}{2}\left(A_{i j ;}{ }^{k}-A^{k}{ }_{i ; j}-A_{j ; i}^{k} u_{k} u^{i} Y^{j}\right. \\
& =-\frac{1}{2} A_{k i ; j} u^{k} u^{i} Y^{j} \\
& =-\frac{1}{2} Y\left\langle A \nabla^{g} u, \nabla^{g} u\right\rangle+\left\langle A \nabla_{Y}^{g} \nabla^{g} u, \nabla^{g} u\right\rangle \\
& =-\frac{1}{2} Y\left\langle A \nabla^{g} u, \nabla^{g} u\right\rangle+\operatorname{Hess}_{g}(u)(Y, X) \\
& =-\frac{1}{2} Y\left\langle X, \nabla^{g} u\right\rangle \\
& =-\frac{1}{2} Y d u(X) .
\end{aligned}
$$

Since $X$ is divergence free and compactly supported in $\Omega$, it follows from classical de-Rham cohomology theory that there exists a 2 -form $Z$ supported in $\Omega$ such that:

$$
X=\nabla^{g} \cdot Z
$$

We choose exponential coordinates about $p$, and write $Z$ as:

$$
Z=\sum_{i<j} Z^{i j} \partial_{i} \wedge \partial_{j}
$$

so that:

$$
X^{k}=\sum_{i>k} \partial_{i} Z^{i k}-\sum_{i<k} \partial_{i} Z^{k i}
$$

We choose the basis at $p$ such that $\nabla^{g} u$ and $Y$ are colinear with $\partial_{1}$ and $\partial_{2}$ respectively. Thus, at the origin, bearing in mind that $X(p)=0$ :

$$
-\frac{1}{2} Y d u(X)=\frac{1}{2}\left\|\nabla^{g} u\right\|_{g}\|Y\|_{g} \sum_{i>1} \partial_{2} \partial_{i} Z^{1 i}
$$

We choose $Z$ such that $\partial_{i} Z^{j k}=0$ for all $i, j$ and $k, \partial_{2} \partial_{1} Z^{11}=1$ and $\partial_{2} \partial_{i} Z^{1 i}=0$ for all $i>1$. Then, if $X=\nabla^{g} \cdot Z$ :

$$
X(p)=0, \quad-\frac{1}{2} Y d u(X)=\frac{1}{2}\left\|\nabla^{g} u\right\|_{g}\|Y\|_{g}
$$

$X=\nabla^{g} \cdot Z$ is the desired vector field, and this completes the proof.

## Proposition 2.5.4

If $\operatorname{Dim}(M) \geqslant 2$, then for generic $g \in \mathcal{M}^{\mu+1}$, if $(\epsilon, g, u) \in \mathcal{Z}_{g}$, if $u$ is non-constant and if $\varphi \in \operatorname{Ker}(L A C(u))$ is non-zero, then there exists a point $p \in M$ such that $d u(p)$ and $d \varphi(p)$ are both non-zero and non-colinear.

Proof: Let $X_{p}$ and $Y_{q}$ be unit vectors over distinct points of $M$. Let $\Omega:=\Omega\left(X_{p}, Y_{q}\right) \subseteq \mathcal{Z}^{*}$ be the open set of all $(\epsilon, g, u)$ such that $\nabla^{g} u(p)$ and $\nabla^{g} u(q)$ are both non-zero and noncolinear with $X_{p}$ and $Y_{q}$ respectively. We define the functions $\Phi_{p}, \Phi_{q}: \Omega \rightarrow \mathbb{R}$ by:

$$
\Phi_{p}(\epsilon, g, u)=\operatorname{Hess}_{g}(u)\left(X_{p}^{\perp}, \nabla^{g} u(p)\right), \quad \Phi_{q}(\epsilon, g, u)=\operatorname{Hess}_{g}(u)\left(Y_{q}^{\perp}, \nabla^{g} u(q)\right)
$$

where $X_{p}^{\perp}$ and $Y_{q}^{\perp}$ are the orthogonal projections of $X_{p}$ and $Y_{q}$ respectively onto the normal hyperplanes to $\nabla^{g} u(p)$ and $\nabla^{g} u(q)$ respectively. Observe that both $\Phi_{p}$ and $\Phi_{q}$ define smooth functions over $\Omega$. Moreover, it follows from Proposition 2.5.3 that $D\left(\Phi_{p}, \Phi_{q}\right)$ is surjective at every point of $\Omega$. Thus, if $Z:=Z\left(X_{p}, Y_{q}\right)$ is the zero set of this functional then it is a smooth, codimension 2 submanifold of $\Omega$. In particular, the restriction of $\Pi_{g}$ to $Z$ is a smooth Fredholm map of index -1 . Thus, if $g \in \mathcal{M}^{-1}$ is a regular value of the restriction of $\Pi_{g}$ to $Z$, then $\Pi_{g}^{-1}(g) \cap Z$ is a smooth submanifold of $Z$ of dimension equal to -1 , that is, it is empty. However, by Proposition 2.2.2, the restriction of $\Pi_{g}$ to $\Omega$, and
therefore also to $Z$, is para-proper, and it follows by Theorem 2.2.4 that the set of regular values of this restriction is generic in $\mathcal{M}^{\mu+1}$.

Let $\mathcal{X} \subseteq(U M \times U M) \backslash \pi^{-1}$ (Diag) be a countable dense family of pairs ( $X_{p}, Y_{q}$ ) of unit vectors above distinct points of $M$. Since the intersection of a countable family of generic sets is generic, it follows that for generic $g \in \mathcal{M}^{\mu+1}$, and for all $\left(X_{p}, Y_{q}\right) \in \mathcal{X}, \Pi_{g}^{-1}(g) \cap Z\left(X_{p}, Y_{q}\right)$ is empty. For such a $g$, choose $(\epsilon, g, u) \in \mathcal{Z}_{g}$. Let $\tilde{p}, \tilde{q} \in M$ be distinct points such that both $d u(\tilde{p})$ and $d u(\tilde{q})$ are non-zero, and let $\tilde{X}_{\tilde{p}}$ and $\tilde{Y}_{\tilde{q}}$ be unit vectors in $U M$ normal to $\nabla^{g} u(\tilde{p})$ and $\nabla^{g} u(\tilde{q})$ respectively. Since $\mathcal{X}$ is dense, there exists a pair $\left(X_{p}, Y_{q}\right) \in \mathcal{X}$ such that $d u(p)$ and $d u(q)$ are non-zero and $X_{p}$ and $Y_{q}$ are non-colinear with $\nabla^{g} u(p)$ and $\nabla^{g} u(q)$ respectively. However, by definition of $g,(\epsilon, g, u) \notin Z\left(X_{p}, Y_{q}\right)$, from which it follows that one of $\Phi_{p}(\epsilon, g, u)$ and $\Phi_{q}(\epsilon, g, u)$ is non-zero. In other words, without loss of generality:

$$
\operatorname{Hess}_{g}(u)\left(X_{p}^{\perp}, \nabla^{g} u(p)\right) \neq 0
$$

and it now follows from Proposition 2.5.2 that there exists at least one point in $M$ where $d u$ and $d \varphi$ are non-zero and non-colinear, as desired.

Combining these relations, we obtain Theorem 1.1.1:
Proof of Theorem 1.1.1: Since the intersection of finitely many generic sets is generic, this follows from Propositions 2.3.5, 2.4.8 and 2.5.4.
2.6 - The Solution Space at Infinity. Now fix $g \in \mathcal{M}^{\mu+1}$. We show that for $\epsilon$ sufficiently large, the only elements of $\mathcal{Z}_{\epsilon, g}$ are the constant solutions. We recall that for all $g, \operatorname{Ker}\left(\Delta_{g}\right)^{\perp}$ coincides with the space of functions whose integral with respect to the volume form of $g$ vanishes.

## Proposition 2.6.1

Let $c \in \mathbb{R}$ be such that $f(c)=0$. There exist $B>0$ and $\delta>0$ such that if $\epsilon>B$, if $v \in C^{\lambda+2}$ and $t \in \mathbb{R}$ are such that:

$$
\int_{M} v d V o l_{g}=0, \quad\|v\|_{\lambda+2}<\delta \epsilon^{-1}, \quad|t|<\delta
$$

and if $A C(\epsilon, g, c+v+t)=0$, then $(v, t)=(0,0)$.
Proof: Define $\mathcal{F}: \operatorname{Ker}\left(\Delta_{g}\right)^{\perp} \times \mathbb{R}^{2} \rightarrow C^{\lambda}$ by:

$$
\mathcal{F}(v, t, \eta)=\Delta_{g} v-f(c+\eta v+t)
$$

Observe that $\mathcal{F}$ is a smooth function between Banach manifolds. Moreover, if we denote by $D_{1} \mathcal{F}$ and $D_{2} \mathcal{F}$ its partial derivatives with respect to the first and second factors respectively, then since $f^{\prime}(c) \neq 0, D_{1} \mathcal{F}+D_{2} \mathcal{F}$ is surjective at $(0,0,0)$. It follows from the Implicit Function Theorem for Banach manifolds that there exists $b>0$ and a neighbourhood $W$ of $(0,0)$ in $\operatorname{Ker}\left(\Delta_{g}\right)^{\perp} \times \mathbb{R}$ such that if $\eta<b$ then there exists a unique point $\left(v_{\eta}, t_{\eta}\right) \in W$ such that $\mathcal{F}\left(v_{\eta}, t_{\eta}, \eta\right)=0$. Since, in particular, $\mathcal{F}(0,0, \eta)=0$ for all $\eta$, it follows that if
$(v, t) \in W$ is such that $\mathcal{F}(v, t, \eta)=0$, then $(v, t)=(0,0)$. Let $B=1 / b$ and let $\delta>0$ be such that:

$$
\left\{(v, t)\left|\|v\|_{\lambda+2}<\delta,|t|<\delta\right\} \subseteq W\right.
$$

We claim that $B$ and $\delta$ have the desired properties. Indeed, let $\epsilon>B, v \in C^{\lambda+2}$ and $t \in \mathbb{R}$ be such that $v \in \operatorname{Ker}\left(\Delta_{g}\right)^{\perp},\|v\|_{\lambda+2}<\delta \epsilon^{-1},|t|<\delta$ and $\mathrm{AC}(\epsilon, g, c+v+t)=0$. Then, denoting $\eta=1 / \epsilon$ :

$$
\mathcal{F}(\epsilon v, t, \eta)=\Delta_{g}(\epsilon v)-f(c+\eta(\epsilon v)+t)=0 .
$$

Since $\|\epsilon v\|_{\lambda+2},|t|<\delta$, it follows from the preceeding discussion that $(v, t)=(0,0)$, as desired.

## Proposition 2.6.2

There exists $B>0$ such that if $\epsilon>B$, then $\mathcal{Z}_{\epsilon, g}$ only consists of constant solutions.
Proof: Suppose the contrary. There exists a sequence $\left(u_{n}, t_{n}, \epsilon_{n}\right)_{n \in \mathbb{N}} \in \operatorname{Ker}\left(\Delta_{g}\right)^{\perp} \times \mathbb{R}^{2}$ such that $\left(\epsilon_{n}\right)_{n \in \mathbb{N}}$ tends to $+\infty, u_{n}$ is non-zero, and for all $n$ :

$$
\mathrm{AC}\left(\epsilon_{n}, g, u_{n}+t_{n}\right)=0
$$

For all $n$, denote $v_{n}=u_{n}+t_{n}$. Observe that the argument of Proposition 2.1.3 is uniform in $\epsilon$ as $\epsilon$ tends to $+\infty$, and there therefore exists $v_{\infty} \in C^{\lambda+2}$ towards which $\left(v_{n}\right)_{n \in \mathbb{N}}$ subconverges. For all $n$ :

$$
\Delta_{g} v_{n}-\epsilon_{n}^{-1} f\left(v_{n}\right)=0
$$

Upon taking limits, it follows that $\Delta_{g} v_{\infty}=0$, and so $v_{\infty}$ is equal to a constant, $c$, say. On the other hand, for all $n$ :

$$
\int f\left(v_{n}\right) \mathrm{dVol}=\int \epsilon_{n} \Delta_{g} v_{n} \mathrm{dVol}=0
$$

and upon taking limits, it follows that:

$$
f(c) \operatorname{Vol}(M)=\int f(c) \mathrm{d} \operatorname{Vol}=0
$$

and so $c$ is a zero of $f$. In particular, $\Delta_{g}\left(\epsilon_{n} u_{n}\right)=\left(f\left(v_{n}\right)\right)_{n \in \mathbb{N}}$ converges to 0 in the $C^{\lambda_{-}}$ topology. However, by the Closed Graph Theorem, the restriction of $\Delta_{g}$ to $\operatorname{Ker}\left(\Delta_{g}\right)^{\perp}$ is a linear isomorphism onto its image, and it follows that $\left(\epsilon_{n}\left\|u_{n}\right\|_{\lambda+2}\right)_{n \in \mathbb{N}}$ converges to 0 . Finally, for all $n$ :

$$
\operatorname{Vol}(M) t_{n}=\int v_{n} \mathrm{dVol}
$$

from which it follows that $\left(\left|t_{n}-c\right|\right)_{n \in \mathbb{N}}$ converges to 0 . It now follows from Proposition 2.6.1 that for sufficiently large $n, u_{n}=0$. This absurd by hypothesis, and the result follows.
2.7 - Morse Homology. We now study the Morse Homology of the Allen Cahn Equation. The construction is fairly standard, and we refer the reader to our forthcoming paper [14] for a detailled outline in the Hölder space framework. We assume henceforth that $\operatorname{Dim}(M) \geqslant 3$. Let $g$ be as in Theorem 1.1.1 and let $\epsilon$ be such that $\epsilon^{-1} \notin \operatorname{Spec}\left(-\Delta_{g}\right)$.

For all $k \in \mathbb{N}$, we define $\mathcal{Z}_{\epsilon, g, k} \subseteq \mathcal{Z}_{\epsilon, g}$ by:

$$
\mathcal{Z}_{\epsilon, g, k}=\left\{u \in \mathcal{Z}_{\epsilon, g} \mid \operatorname{Index}(u)=k\right\},
$$

and for all $k \in \mathbb{N}$, we define the chain group $C_{k}$ by:

$$
C_{k}=\mathbb{Z}_{2}\left[\mathcal{Z}_{\epsilon, g, k}\right]=\left\{f: \mathcal{Z}_{\epsilon, g, k} \rightarrow \mathbb{Z}_{2}\right\}
$$

Morse Homology theory defines a canonical chain mapping $\partial_{k}: C_{k} \rightarrow C_{k-1}$ in terms of solutions to the parabolic Allen-Cahn Equation, $\mathrm{pAC}_{\epsilon, g}:=\partial_{t}-\mathrm{AC}_{\epsilon, g}$, over the space $\mathbb{R} \times M$. The Morse Homology of the Allen-Cahn Equation is then defined to be the homology of the chain complex $\left(C_{*}, \partial_{*}\right)$. That is, for all $k$ :

$$
\mathrm{HAC}_{k}=\frac{\operatorname{Ker}\left(\partial_{k}\right)}{\operatorname{Im}\left(\partial_{k+1}\right)}
$$

Importantly, $\mathrm{HAC}_{*}$ is independant, up to isomorphism, of the pair $(\epsilon, g)$ used to define it. In actual fact, the preceeding construction would require that all elements of $\mathcal{Z}_{\epsilon, g}$ be non-degenerate. However, since all critical points of $e_{g}$ are themselves non-degenerate, we use a perturbation argument to show that degenerate elements of $\mathcal{Z}_{\epsilon, g}$ do not contribute to the homology: in other words, we simply ignore them. The justification is analogous to the manner in which the function $F_{\epsilon}(t):=t^{3}+\epsilon t$ has a degenerate critical point at 0 when $\epsilon=0$, and no critical points for $\epsilon>0$, in contrast to the function $G_{\epsilon}(t)=t^{4}+\epsilon t^{2}$, which has a critical point at 0 for all $\epsilon$.

In order to calculate the Morse Homology, we suppose that $\epsilon \gg 0$. By Proposition 2.6.2, we may suppose that $\mathcal{Z}_{\epsilon, g}$ only consists of constant solutions, and furthermore, by Proposition 2.2.1, we may suppose that the Morse Index of the constant solution $u=c$ is equal to 0 or 1 according as $f^{\prime}(c)$ is positive or negative respectively. Let $F$ be any primitive of $f$, let $c_{ \pm}$be zeroes of $f$, and let $w: \mathbb{R} \rightarrow \mathbb{R}$ be such that:

$$
\partial_{t} w=-f \circ w, \quad \operatorname{Lim}_{t \rightarrow \pm \infty}=c_{ \pm}
$$

That is, $w$ is a gradient flow of $F$ from $c_{-}$to $c_{+}$. We extend $w$ to a function from $\mathbb{R} \times M$ into $\mathbb{R}$ by setting it to be constant in the $x$ direction. Observe that $w$ is then a bounded solution to the parabolic Allen-Cahn Equation. That is:

$$
\mathrm{pAC}_{\epsilon, g} w=\left(\partial_{t}-\mathrm{AC}_{\epsilon, g}\right) w=0
$$

We therefore refer to such a function $w$ as a space-constant trajectory. As in the elliptic case, we say that $w$ is non-degenerate whenever the linearisation of $\mathrm{pAC}_{\epsilon, g}$ around $w$ defines a surjective mapping from the inhomogeneous Sobolev space $H^{1,2}(\mathbb{R} \times \mathbb{M})$ into $L^{2}(\mathbb{R} \times \mathbb{M})$. In order to correctly calculate the Morse Homology, we have to show that all trajectories that we study are non-degenerate. However:

## Proposition 2.7.1

For all $g \in \mathcal{M}^{\mu+1}$, there exists $B>0$ such that for $\epsilon>B$, every space-constant trajectory is non-degenerate.
Proof: Let $w: \mathbb{R} \times M \rightarrow \mathbb{R}$ be a space constant trajectory, and let $L$ be the linearisation of $\mathrm{pAC}_{\epsilon, g}$ about $w$. For all $\varphi: \mathbb{R} \times M \rightarrow M$ :

$$
L \varphi=\left(\partial_{t}-\epsilon \Delta_{g}\right) \varphi-\left(f^{\prime} \circ w\right)(t) \varphi .
$$

By the Sturm-Liouville Theorem, there exists an orthonormal basis $\left(\psi_{n}\right)_{n \in \mathbb{N}}$ of $L^{2}(M)$ consisting of eigenfunctions of $-\Delta_{g}$. Let $0=\lambda_{0}<\lambda_{1} \leqslant \ldots$ be the corresponding eigenvalues. Define $B>0$ such that $B>\left\|f^{\prime}\right\|_{L^{\infty}} / \lambda_{1}$. We claim that $B$ has the desired properties. Indeed, choose $\epsilon>B$. By Proposition 2.2.1, both $c_{-}$and $c_{+}$are non-degenerate with Morse Indices equal either to 0 or 1 . Observe, moreover, that $\operatorname{Index}\left(c_{-}\right)=1$ and $\operatorname{Index}\left(c_{+}\right)=0$. By the Atiyah-Patodi-Singer Index Theorem (c.f. [10]), $L$ defines a Fredholm mapping from $H^{1,2}(\mathbb{R} \times M)$ into $L^{2}(\mathbb{R} \times M)$ of Fredholm index equal to 1 . Thus, in order to show that $w$ is non-degenerate, it suffices to show that $\operatorname{Dim}(\operatorname{Ker}(L)) \leqslant 1$. However, choose $\varphi \in \operatorname{Ker}(L)$. For $k \geqslant 1$, define $\varphi_{k}: \mathbb{R} \rightarrow \mathbb{R}$ by $\varphi_{k}(t)=\left\langle\varphi_{t}, \psi_{n}\right\rangle$. Observe that $\varphi_{k} \in L^{2}(\mathbb{R})$. However:

$$
\dot{\varphi}_{k}=\left(\epsilon \lambda_{k}+\left(f^{\prime} \circ w\right)(t)\right) \varphi_{k}
$$

Since $\epsilon>B$, there exists $\delta>0$ such that $\left(\epsilon \lambda_{k}+\left(f^{\prime} \circ w\right)(t)\right)>\delta$. Thus, over any interval in which $\varphi_{n}$ is non-vanishing, we have:

$$
\partial_{t}\left(\log \left(\left|\varphi_{n}\right|\right)\right) \geqslant \delta
$$

and since $\psi_{k} \in L^{2}(\mathbb{R})$, it must therefore vanish identically. $\varphi_{t}$ therefore lies in the linear span of $\psi_{0}$ for all $t$. That is, it is constant in space. However, since the space of solutions to a first order ODE is at most 1-dimensional, it follows that $\operatorname{Ker}(L)$ is also at most one-dimensional, and we conclude that $w$ is non-degenerate, as desired.
This allows us to calculate the Morse-Homology:

## Proposition 2.7.2

The Morse Homology of the Allen-Cahn Operator is given by:

$$
H A C_{k}=\left\{\begin{array}{l}
\mathbb{Z}_{2} \text { if } k=0 \\
0 \text { otherwise }
\end{array}\right.
$$

Proof: Let $B$ be as in Proposition 2.7.1 and choose $\epsilon>B$. Upon increasing $B$ is necessary, it follows from Propositions 2.2 .1 and 2.6 .2 that $\mathcal{Z}_{\epsilon, g}$ only consists of constant solutions and, moreover, that if $u=c$ is a constant solution, then it is non-degenerate and its Morse Index is equal to 0 or 1 according as $f^{\prime}(c)$ is positive or negative respectively. By Property (B) of $f, f$ has an odd number of zeroes, $c_{1}<\ldots<c_{2 n+1}$. Moreover, if $k$ is odd, then $f^{\prime}\left(c_{k}\right)>0$, and if $k$ is even, then $f^{\prime}\left(c_{k}\right)<0$. Consequently:

$$
\mathcal{Z}_{d, 0}=\left\{c_{1}, c_{3}, \ldots, c_{2 n+1}\right\}, \quad \mathcal{Z}_{d, 1}=\left\{c_{2}, c_{4}, \ldots, c_{2 n}\right\}
$$

and $\mathcal{Z}_{d, p}$ is empty for all $p \geqslant 2$. In particular, for all $p \geqslant 2, C_{p}=0$ and so $H \mathrm{AC}_{p}=0$. For $1 \leqslant k \leqslant n$, there are two space-constant trajectories leaving $c_{2 k}$, terminating in $c_{2 k-1}$ and $c_{2 k+1}$ respectively. Moreover, by Proposition 2.7.1, these space-constant trajectories are non-degenerate, and by the unstable manifold theorem (c.f. [16]), up to reparametrisation in time, there are no other bounded solutions $w_{t}(\cdot):=w(t, \cdot)$ to the parabolic Allen-Cahn Equation which converge to $c_{2 k}$ as $t$ tends to minus infinity. It follows from the definition of the chain map (c.f. [14]) that:

$$
\partial_{1} c_{2 k}=c_{2 k-1}+c_{2 k+1} .
$$

In particular, $\left\{\partial_{1} c \mid c \in \mathcal{Z}_{d, 1}\right\}$ is a linearly independent subset of $C_{0}$, and so:

$$
\operatorname{Dim}\left(H \mathrm{AC}_{1}\right)=\operatorname{Dim}\left(\operatorname{Ker}\left(\partial_{1}\right)\right)=0
$$

Finally, by the Rank-Nullity Theorem, $\operatorname{Dim}\left(\operatorname{Im}\left(\partial_{1}\right)\right)=n$, and so $\operatorname{Dim}\left(H \mathrm{AC}_{0}\right)=1$, and the result now follows.

We now return to the specific case studied in the introduction where $f(u)=u^{3}-u$, and we prove Theorem 1.1.2:

Proof of Theorem 1.1.2: For all $k$, we define $X_{k} \subseteq \mathcal{Z}_{\epsilon, g}$ to be the set of all stationary solutions of Morse-Index equal to $k$, and we define $C_{k}$ and $\partial_{k}$ as outlined above. Denote $l=\operatorname{Index}(0)$. Since $f$ is odd, multiplication by -1 maps $\mathcal{Z}_{\epsilon, g}$ to itself, and all solutions of $\mathrm{AC}_{\epsilon, g} u=0$ which are different to 0 therefore exist in pairs. The set $X_{k}$ therefore has even cardinality for all $k \neq l$ and odd cardinality when $k=l$. In other words, $C_{k}$ has odd dimension for $k \neq l$ and even dimension for $k=l$. For all $k$, let $K_{k}$ be the kernel of $\partial_{k}$. We claim that $K_{k}$ is odd-dimensional for all $0<k<l$. Indeed, choose $0<k<l-1$ and suppose that $K_{k}$ is odd-dimensional. Then, since $\mathrm{HAC}_{k}=0$, it follows that the image of $\partial_{k+1}$ is also odd-dimensional, and since $C_{k+1}$ is even-dimensional, it follows by the Rank-Nullity Theorem that $K_{k+1}$ is odd-dimensional. However, since $K_{0}=C_{0}$ and since $H \mathrm{AC}_{0}=\mathbb{Z}_{2}, \operatorname{Im}\left(\partial_{1}\right)$ is also odd-dimensional, and it follows by the Rank-Nullity Theorem that $K_{1}$ is also odd-dimensional. We conclude by induction that $K_{k}$ is odd dimensional for all $0<k<l$ as asserted. In particular, for all $0<k<l, K_{k}$ is non-trivial, and thus so too is $C_{k}$, from which it follows that $X_{k}$ is non-empty, as desired.

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