# Energy concentration for the Landau-Lifshitz equation 

Roger Moser<br>Department of Mathematical Sciences, University of Bath, Bath BA2 7AY, United Kingdom

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#### Abstract

For the Landau-Lifshitz equation on a domain with three space dimensions, we consider energy concentration phenomena arising in the context of weakly convergent sequences of solutions. The concentration measure can be interpreted as a family of generalized curves. We establish a connection to a geometric flow. © 2007 Elsevier Masson SAS. All rights reserved.


## Résumé

Pour l'équation de Landau-Lifshitz sur un domaine en trois dimensions d'espace, nous considérons des phénomènes de concentration d'énergie survenant dans le contexte des suites faiblement convergentes de solutions. La mesure de concentration peut être interprétée comme famille de courbes généralisées. Nous établissons une connexion avec un flot géométrique. © 2007 Elsevier Masson SAS. All rights reserved.

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## 1. Introduction

For an open set $\Omega \subset \mathbb{R}^{3}$ and $T>0$, define $\Omega_{T}=\Omega \times(0, T)$. We consider the Landau-Lifshitz equation

$$
\begin{equation*}
\frac{\partial u}{\partial t}+a u \wedge(u \wedge \Delta u)+b u \wedge \Delta u=0 \quad \text { in } \Omega_{T} \tag{1}
\end{equation*}
$$

for a map $u: \Omega_{T} \rightarrow \mathbb{S}^{2}$, where $\mathbb{S}^{2} \subset \mathbb{R}^{3}$ is the unit 2-sphere. Here $\wedge$ denotes the vector product in $\mathbb{R}^{3}$ and $a, b \in \mathbb{R}$ are fixed constants. We assume that $a>0$, which makes the equation parabolic.

For a map $u \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$, the Laplacian $\Delta u$ has the orthogonal decomposition

$$
\begin{equation*}
-\Delta u=|\nabla u|^{2} u+u \wedge(u \wedge \Delta u) \tag{2}
\end{equation*}
$$

The expression $\Delta u+|\nabla u|^{2} u$ is called the tension field of $u$; it is minus the $L^{2}$-gradient of the Dirichlet energy

$$
E(u)=\frac{1}{2} \int_{\Omega}|\nabla u|^{2} d x
$$

[^0]under the constraint that $u$ takes values in $\mathbb{S}^{2}$. Using the decomposition (2), we find that (1) is equivalent to
\[

$$
\begin{equation*}
\tilde{a} \frac{\partial u}{\partial t}+\tilde{b} u \wedge \frac{\partial u}{\partial t}=\Delta u+|\nabla u|^{2} u \tag{3}
\end{equation*}
$$

\]

as well as

$$
\begin{equation*}
\tilde{a} u \wedge \frac{\partial u}{\partial t}-\tilde{b} \frac{\partial u}{\partial t}=\operatorname{div}(u \wedge \nabla u), \tag{4}
\end{equation*}
$$

where

$$
\tilde{a}=\frac{a}{a^{2}+b^{2}} \quad \text { and } \quad \tilde{b}=\frac{b}{a^{2}+b^{2}} .
$$

Both (3) and (4) can also be interpreted in the weak sense for a map in the Sobolev space

$$
H_{\mathrm{loc}}^{1}\left(\Omega_{T}, \mathbb{S}^{2}\right)=\left\{u \in H_{\mathrm{loc}}^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right):|u|=1 \text { almost everywhere }\right\} .
$$

Eq. (1) stems from a model in micromagnetics which describes the dynamics of the magnetization vector field of a ferromagnetic body. The equation is also of geometrical interest, however. It can be thought of as a hybrid of two other geometric evolution equations: the so-called harmonic map heat flow

$$
\begin{equation*}
\frac{\partial u}{\partial t}=\Delta u+|\nabla u|^{2} u \tag{5}
\end{equation*}
$$

(which is the negative $L^{2}$-gradient flow for the functional $E$ ) on the one hand and the equation

$$
\begin{equation*}
\frac{\partial u}{\partial t}=u \wedge \Delta u \tag{6}
\end{equation*}
$$

on the other hand. Solutions of (6) are often called Schrödinger maps, because the equation is of the type of a nonlinear Schrödinger equation (which is most obvious when $u$ is composed with the stereographic projection).

Suppose that we have a sequence of solutions $u_{k} \in C^{\infty}\left(\bar{\Omega} \times[0, T), \mathbb{S}^{2}\right)$ of (1) such that

$$
\sup _{k \in \mathbb{N}} E\left(u_{k}(\cdot, 0)\right)<\infty .
$$

Also assume for the moment that $\Omega$ is bounded and has a smooth boundary with outer normal vector $v$, and that $v \cdot \nabla u_{k}=0$ on $\partial \Omega \times[0, T)$. Taking the scalar product with $\Delta u_{k}$ on both sides of (1), integrating over $\Omega \times\{t\}$, and performing an integration by parts, we see that

$$
\begin{equation*}
\frac{d}{d t} E\left(u_{k}(\cdot, t)\right)+a \int_{\Omega \times\{t\}}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x=0 . \tag{7}
\end{equation*}
$$

Hence

$$
E\left(u_{k}\left(\cdot, t_{0}\right)\right)+a \int_{0}^{t_{0}} \int_{\Omega}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x d t=E\left(u_{k}(\cdot, 0)\right)
$$

for every $t_{0} \in[0, T)$. In particular the sequence $\left\{u_{k}\right\}$ is bounded in $H^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right)$ and there exists a weakly convergent subsequence $\left\{u_{k_{i}}\right\}$. Using the representation (4) for the Landau-Lifshitz equation, it is readily verified that the limit map is a weak solution. The observed convergence, however, is not strong in $H^{1}\left(\Omega_{T}, \mathbb{S}^{2}\right)$ in general, since the energy density $\frac{1}{2}\left|\nabla u_{k_{i}}\right|^{2}$ may concentrate near a certain subset of $\Omega_{T}$. Suppose that $\mathcal{L}^{4}$ is the Lebesgue measure on $\Omega_{T}$. Then a way to describe this energy concentration is to consider the Radon measures

$$
m_{k}=\frac{1}{2} \mathcal{L}^{4}\left\llcorner\left|\nabla u_{k}\right|^{2} .\right.
$$

For a suitable choice of the above subsequence, we have convergence of $m_{k_{i}}$ to a Radon measure $m$ on $\Omega_{T}$, which need not be absolutely continuous with respect to $\mathcal{L}^{4}$. The singular part of $m$ measures the energy concentration, therefore we call it the concentration measure. In the situation studied here, it turns out that the concentration measure can be interpreted as a geometric object. Namely, for almost every $t \in[0, T)$, it gives rise to a type of generalized curve in $\Omega$. Our aim is to study the evolution of these generalized curves.

Without the above boundary conditions, the energy identity (7) is no longer valid, but it can be replaced by a localized version. For any $\eta \in C_{0}^{1}\left(\Omega_{T}\right)$, a smooth solution of the Landau-Lifshitz equation satisfies

$$
\begin{equation*}
\frac{1}{2} \frac{d}{d t} \int_{\Omega \times\{t\}} \eta|\nabla u|^{2} d x=\int_{\Omega \times\{t\}}\left(\frac{1}{2} \frac{\partial \eta}{\partial t}|\nabla u|^{2}-a \eta|u \wedge \Delta u|^{2}-a \nabla \eta \cdot\langle\Delta u, \nabla u\rangle+b \nabla \eta \cdot\langle u \wedge \Delta u, \nabla u\rangle\right) d x \tag{8}
\end{equation*}
$$

Here and throughout the paper, we denote scalar products in $T \Omega$ by a dot and in $T \mathbb{S}^{2}$ by $\langle\cdot, \cdot\rangle$ for clarity. Using (8) instead of (7), we now obtain local estimates for the energy under suitable conditions, and the questions we study remain the same.

For the harmonic map heat flow (5), the corresponding problem has been studied by Li and Tian [12] and by Lin and Wang [18] (also for higher-dimensional domains and other target manifolds). These papers give a connection between the energy concentration measure and the mean curvature flow. We will show a similar connection for the Landau-Lifshitz equation, but the mean curvature flow has to be replaced by another geometric flow.

We briefly recall the definition of the mean curvature flow, before we modify it in order to obtain the equation that is relevant in the context of this paper. For simplicity, we restrict our attention to the flow for closed curves in $\Omega$ here. In this situation, the mean curvature flow is also known as the curve shortening flow. Suppose that $F: \mathbb{S}^{1} \times[0, T) \rightarrow \Omega$ is a smooth map such that $\Sigma_{t}=F\left(\mathbb{S}^{1} \times\{t\}\right)$ is an embedded curve for every $t \in[0, T)$. Let $H_{t}: \Sigma_{t} \rightarrow \mathbb{R}^{3}$ be the curvature vector for $\Sigma_{t}$. We say that $F$ is a solution of the mean curvature flow (curve shortening flow) if it satisfies the equation

$$
\frac{\partial F}{\partial t}(s, t)=H_{t}(F(s, t)) \quad \text { for } s \in \mathbb{S}^{1}, t \in[0, T)
$$

or in shorter notation,

$$
\begin{equation*}
\frac{\partial F}{\partial t}=H_{t} \tag{9}
\end{equation*}
$$

This is the negative $L^{2}$-gradient flow for the length functional. We write $\mathscr{H}^{1}$ for the 1-dimensional Hausdorff measure. If we have a smooth solution of (9), then for any $\eta \in C_{0}^{1}\left(\Omega_{T}\right)$, we can compute

$$
\begin{equation*}
\frac{d}{d t} \int_{\Sigma_{t}} \eta d \mathscr{H}^{1}=\int_{\Sigma_{t}}\left(\frac{\partial \eta}{\partial t}-\eta\left|H_{t}\right|^{2}+\nabla \eta \cdot H_{t}\right) d \mathscr{H}^{1} \tag{10}
\end{equation*}
$$

A relaxed version of this identity (with the equality replaced by an inequality for test functions with non-negative values) was used by Brakke [3] to define a generalization of the mean curvature flow. This generalization is also the formulation that is used in [18], and with a modification in [12], to describe the evolution of the energy concentration set for the harmonic map heat flow.

Next we regard $\Sigma_{t}$ as oriented curves, and we choose a unit tangent vector $\tau_{t}: \Sigma_{t} \rightarrow \mathbb{S}^{2}$ that is continuous (also with respect to $t$ ). We replace (9) by the equation

$$
\begin{equation*}
\frac{\partial F}{\partial t}=a H_{t}-b \tau_{t} \wedge H_{t} \tag{11}
\end{equation*}
$$

or, equivalently,

$$
\tilde{a} \frac{\partial F}{\partial t}+\tilde{b} \tau_{t} \wedge \frac{\partial F}{\partial t}=H_{t}
$$

Then instead of (10) we obtain

$$
\begin{equation*}
\frac{d}{d t} \int_{\Sigma_{t}} \eta d \mathscr{H}^{1}=\int_{\Sigma_{t}}\left(\frac{\partial \eta}{\partial t}-a \eta\left|H_{t}\right|^{2}+a \nabla \eta \cdot H_{t}-b \nabla \eta \cdot \tau_{t} \wedge H_{t}\right) d \mathscr{H}^{1} \tag{12}
\end{equation*}
$$

A certain formal similarity between (8) and (12) can immediately be seen. We will establish a rigorous connection in this paper and show that (a generalized version of) the flow (11) describes the behaviour of the concentration measure for the Landau-Lifshitz equation in a certain sense.

Before we can state our main result, we need some notation. We assume that the reader is familiar with the notion of countable rectifiability. If not, we refer to Federer [7] or Simon [24]. Suppose $\Sigma \subset \Omega$ is a countably 1-rectifiable
set. Then we write $T_{x} \Sigma$ for the approximate tangent space at every point $x \in \Sigma$ where it exists. If $\phi \in C^{1}\left(\Omega, \mathbb{R}^{3}\right)$ and $L$ is a linear subspace of $\mathbb{R}^{3}$, then $\operatorname{div}_{L} \phi(x)$ is the divergence of $\phi$ at $x$ with respect to $L$. That is, if $\Pi_{L}$ denotes the matrix of the orthogonal projection onto $L$, then

$$
\operatorname{div}_{L} \phi(x)=\operatorname{tr}\left(\Pi_{L} \nabla \phi(x)\right)
$$

Similarly, we write
$\operatorname{div}_{\Sigma} \phi(x)=\operatorname{div}_{T_{x} \Sigma} \phi(x)$.
Suppose that $p \in \mathbb{S}^{2}$ and $X=\left(X_{1}, X_{2}, X_{3}\right), Y=\left(Y_{1}, Y_{2}, Y_{3}\right) \in \mathbb{R}^{3} \otimes T_{p} \mathbb{S}^{2}$; then we write $X \otimes Y$ for the (3 $\times 3$ )-matrix with $(\alpha, \beta)$-th component $\left\langle X_{\alpha}, Y_{\beta}\right\rangle$.

Theorem 1.1. Suppose $u_{k} \in C^{\infty}\left(\Omega_{T}, \mathbb{S}^{2}\right)$ are solutions of the Landau-Lifshitz equation (1) such that
$\sup _{k \in \mathbb{N}} \int_{\Omega_{T}}\left|\nabla u_{k}\right|^{2} d z<\infty$.
Then there exists a subsequence $\left\{u_{k_{i}}\right\}$ such that the following is true.
(i) There exists a weak solution $u \in H_{\mathrm{loc}}^{1}\left(\Omega_{T}, \mathbb{S}^{2}\right)$ of (4) such that $u_{k_{i}} \rightharpoonup u$ weakly in $H_{\mathrm{loc}}^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right)$ and $\Delta u+$ $|\nabla u|^{2} u \in L_{\mathrm{loc}}^{2}\left(\Omega_{T}, \mathbb{R}^{3}\right)$.
(ii) There exists a function $\theta: \Omega_{T} \rightarrow \mathbb{N} \cup\{0\}$ such that for almost every $t \in[0, T)$,

- $\theta_{t}=\theta(\cdot, t)$ is locally $\mathscr{H}^{1}$-integrable, and
- $\Sigma_{t}=\theta_{t}^{-1}(\mathbb{N})$ is closed relative to $\Omega$ and countably 1-rectifiable,
and for every $\eta \in C_{0}^{0}\left(\Omega_{T}\right)$,

$$
\begin{equation*}
\lim _{i \rightarrow \infty} \int_{0}^{T} \int_{\Omega} \eta\left|\nabla u_{k_{i}}\right|^{2} d x d t=\int_{0}^{T} \int_{\Omega} \eta|\nabla u|^{2} d x d t+8 \pi \int_{0}^{T} \int_{\Omega} \eta \theta d \mathscr{H}^{1} d t \tag{13}
\end{equation*}
$$

(iii) There exists a vector field $H: \Omega_{T} \rightarrow \mathbb{R}^{3}$ such that for almost every $t \in[0, T)$,

- $H_{t}=H(\cdot, t) \in L_{\text {loc }}^{2}\left(\mathscr{H}^{1} L \theta_{t}, \mathbb{R}^{3}\right)$, and
- $H_{t}(x) \perp T_{x} \Sigma_{t}$ at $\mathscr{H}^{1}$-almost every $x \in \Sigma_{t}$,
and for every $\phi \in C_{0}^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right)$,

$$
\begin{align*}
& \int_{0}^{T} \int_{\Omega}\left(\frac{1}{2}|\nabla u|^{2} \operatorname{div} \phi-\operatorname{tr}((\nabla u \otimes \nabla u) \nabla \phi)-\langle\Delta u, \nabla u\rangle \cdot \phi\right) d x d t \\
& \quad+4 \pi \int_{0}^{T} \int_{\Omega}\left(\operatorname{div}_{\Sigma_{t}} \phi+H \cdot \phi\right) \theta d \mathscr{H}^{1} d t=0 \tag{14}
\end{align*}
$$

(iv) There exist three vector fields $\tau: \Omega_{T} \rightarrow \mathbb{S}^{2}$ and $H^{+}, H^{-}: \Omega_{T} \rightarrow \mathbb{R}^{3}$ with $H^{+}+H^{-}=H$, such that for almost every $t \in[0, T)$,

- $\tau_{t}=\tau(\cdot, t)$ is $\mathscr{H}^{1}$-measurable,
- $H_{t}^{ \pm}=H^{ \pm}(\cdot, t) \in L_{\text {loc }}^{2}\left(\mathscr{H}^{1}\left\llcorner\theta_{t}, \mathbb{R}^{3}\right)\right.$, and
- $\tau_{t}(x) \in T_{x} \Sigma_{t}$ and $H_{t}^{ \pm}(x) \perp T_{x} \Sigma_{t}$ for $\mathscr{H}^{1}$-almost every $x \in \Sigma_{t}$,
and for every $\eta \in C_{0}^{1}\left(\Omega_{T},[0, \infty)\right)$,

$$
\begin{align*}
& \int_{0}^{T} \int_{\Omega}\left(\frac{1}{2}|\nabla u|^{2} \frac{\partial \eta}{\partial t}-a \eta|u \wedge \Delta u|^{2}-a \nabla \eta \cdot\langle\Delta u, \nabla u\rangle+b \nabla \eta \cdot\langle u \wedge \Delta u, \nabla u\rangle\right) d x d t \\
& \quad+4 \pi \int_{0}^{T} \int_{\Omega}\left(\frac{\partial \eta}{\partial t}-a \eta|H|^{2}+a \nabla \eta \cdot H-b \nabla \eta \cdot\left(\tau \wedge\left(H^{+}-H^{-}\right)\right)\right) \theta d \mathscr{H}^{1} d t \geqslant 0 \tag{15}
\end{align*}
$$

Next we give a few remarks on how the quantities in this theorem can be interpreted. First note that the functions $\theta_{t}$, which are supported on the countably 1-rectifiable sets $\Sigma_{t}$, may be regarded as multiplicity functions for a generalized type of curves in $\Omega$. By (13), these generalized curves describe the energy concentration for the subsequence $\left\{u_{k_{i}}\right\}$. If $u$ happens to be sufficiently smooth, then it is readily checked that the first integral in (14) vanishes. Hence in this case we have

$$
\int_{\Sigma_{t}}\left(\operatorname{div}_{\Sigma_{t}} \phi+H_{t} \cdot \phi\right) \theta_{t} d \mathscr{H}^{1}=0
$$

for every $\phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$ and almost every $t \in[0, T)$. In other words, the vector field $H_{t}$ is the curvature of the generalized curve given by $\theta_{t}$. In general, however, the above identity may be false. The pair consisting of the map $u(\cdot, t)$ and $\theta_{t}$ should then be thought of as a geometrical object which has a "curvature" given jointly by $\Delta u+|\nabla u|^{2} u$ and $H_{t}$.

Inequality (15), finally, describes (partially) the evolution of the generalized curves given by $\theta_{t}$. Again, if $u$ is smooth, then the first integral vanishes and we have

$$
\int_{0}^{T} \int_{\Sigma_{t}}\left(\frac{\partial \eta}{\partial t}-a \eta\left|H_{t}\right|^{2}+a \nabla \eta \cdot H_{t}-b \nabla \eta \cdot\left(\tau_{t} \wedge\left(H_{t}^{+}-H_{t}^{-}\right)\right)\right) \theta_{t} d \mathscr{H}^{1} d t \geqslant 0
$$

for every $\eta \in C_{0}^{1}\left(\Omega_{T}\right)$ with $\eta \geqslant 0$. This is a relaxed version of (12), integrated over [0,T)—up to one discrepancy in the last term of the integrand. The fact that the quantity $\tau_{t} \wedge H_{t}$ in (12) splits up into two parts is not surprising, since the multiplicity function $\theta_{t}$ may encode a piece of a curve several times with opposite orientations. If this is the reason for the appearance of $H_{t}^{+}$and $H_{t}^{-}$, one would suspect that they must be parallel to $H_{t}$ and the ratios of the lengths must be fractions with the denominator $\theta_{t}$ almost everywhere. We are unable, however, to prove this, owing to the lack of a sufficiently strong convergence of the relevant quantities.

For the proof of Theorem 1.1, we use similar arguments as in [12] and [18] up to a certain point. In contrast to the harmonic map heat flow, however, the Landau-Lifshitz equation gives rise to an additional difficulty. The terms

$$
b \int_{\Omega \times\{t\}} \nabla \eta \cdot\langle u \wedge \Delta u, \nabla u\rangle d x \quad \text { and } \quad b \int_{\Sigma_{t}} \nabla \eta \cdot \tau_{t} \wedge H_{t} d \mathscr{H}^{1}
$$

in (8) and (12), respectively, are harder to control than the other terms. Unlike the terms with the coefficient $a$, they cannot be simplified by a mere integration by parts. This is the reason for most of the new concepts and arguments that we use.

Several generalizations of the theorem are conceivable. On the one hand, we may consider a sequence of weak solutions of the Landau-Lifshitz equation. In order to apply the tools that we need for the proof, we have to impose certain additional conditions (such as the stability hypothesis introduced by Feldman [8] for the harmonic map heat flow; cf. [20] for a version for the Landau-Lifshitz equation). But then a similar result follows with the same methods. On the other hand, we may change the dimension of $\Omega$. The expected energy concentration set is always of codimension 2 . Thus in dimensions 1 and 2 , it is clearly not described by a curvature driven flow. No result of the type of Theorem 1.1 can then be expected. In dimension 4, the same arguments that we use in this paper still work with only minor modifications. In order to keep the presentation simple, we leave it to the reader to verify this (as well as the results for weak solutions). In dimension 5 and higher, however, several of the tools that we use are no longer available. The most important obstacle to proving results similar to Theorem 1.1 in higher dimensions is the lack of a so-called monotonicity formula for the Landau-Lifshitz equation. Such a formula exists, e.g., for the harmonic map heat flow (cf. Struwe [26]) and is used extensively in [12] and [18].

We close this section with the introduction of some more notation. We write $B_{r}\left(x_{0}\right)$ for an open ball in $\mathbb{R}^{3}$ with centre $x_{0}$ and radius $r$. Sometimes we work with two-dimensional balls, and then we normally use the notation $B_{r}^{2}\left(x_{0}^{\prime}\right)$ to avoid confusion (where $x_{0}^{\prime} \in \mathbb{R}^{2}$ ). At one point, however, we work exclusively in $\mathbb{R}^{2}$, and then we drop the superscript. Furthermore, we use the abbreviations $B_{r}=B_{r}(0)$ and $B_{r}^{2}=B_{r}^{2}(0)$. We denote the $j$-dimensional Lebesgue measure by $\mathcal{L}^{j}$ and the $j$-dimensional Hausdorff measure by $\mathscr{H}^{j}$ (as we have already done for certain dimensions).

We write $z=(x, t)$ for a generic point in $\mathbb{R}^{3} \times \mathbb{R}$. In space-time, it is natural to use the parabolic metric

$$
d((x, t),(y, s))=\max \{|x-y|, \sqrt{|t-s|}\}
$$

rather than the Euclidean metric in the context of a parabolic problem such as (1). Therefore, we also consider "balls" in this metric, for which we use the notation

$$
P_{r}\left(x_{0}, t_{0}\right)=B_{r}\left(x_{0}\right) \times\left(t_{0}-r^{2}, t_{0}+r^{2}\right) .
$$

Moreover, $P_{r}=P_{r}(0)$. The $j$-dimensional Hausdorff measure with respect to $d$ is denoted by $\mathscr{H}_{d}^{j}$.
We often work with Radon measures, and we use the fact that a Radon measure (on $\Omega$, say) can be identified with a functional in the dual space of $C_{0}^{0}(\Omega)$. When we indicate a convergence of a sequence of Radon measures, we always mean weak* convergence in $\left(C_{0}^{0}(\Omega)\right)^{*}$.

## 2. A few tools from geometric measure theory

To study a blow-up measure as in Theorem 1.1, the notion of generalized varifolds that has been introduced by Ambrosio and Soner [2], and independently by Lin [14,15], is very useful. For the problem that we study in this paper, we have to modify the theory somewhat, in order to make sense of the notion of an orientation in the generalized setting. For convenience, we also use a different representation for the generalized varifolds. Before we give the details, we recall the basic definitions for (ordinary) varifolds. For further details, see Allard [1] or Simon [24]. A good source for other information on geometric measure theory is the book by Federer [7].

Let $G(3,1)$ be the Grassmann manifold of all 1-dimensional linear subspaces of $\mathbb{R}^{3}$ (which can of course be identified with the real projective plane). Moreover, let $G_{1}(\Omega)=\Omega \times G(3,1)$. A 1-varifold on $\Omega$ is a Radon measure on $G_{1}(\Omega)$. Important examples are the so-called integral varifolds. Suppose $\theta: \Omega \rightarrow \mathbb{N} \cup\{0\}$ is a locally $\mathscr{H}^{1}$-integrable function such that $\Sigma=\theta^{-1}(\mathbb{N})$ is countably 1-rectifiable. Then the varifold $V$ on $\Omega$ such that

$$
\int_{G_{1}(\Omega)} \psi d V=\int_{\Sigma} \psi\left(x, T_{x} \Sigma\right) \theta(x) d \mathscr{H}^{1}(x)
$$

for every $\psi \in C_{0}^{0}\left(G_{1}(\Omega)\right)$, is called an integral 1-varifold. (Thus in particular the functions $\theta_{t}$ in Theorem 1.1 give rise to integral varifolds.)

We also consider the Grassmann manifold $G^{0}(3,1)$ consisting of all oriented 1-dimensional subspaces of $\mathbb{R}^{3}$ (in other words, the 2 -sphere). We set $G_{1}^{0}(\Omega)=\Omega \times G^{0}(3,1)$. An oriented 1 -varifold on $\Omega$ is a Radon measure on $G_{1}^{0}(\Omega)$. If $V$ is an oriented 1 -varifold on $\Omega$, then the projection $P: G^{0}(3,1) \rightarrow G(3,1)$ induces naturally a 1-varifold on $\Omega$ by the push-forward $\left(\operatorname{id}_{\Omega} \times P\right)_{\#} V$.

Suppose $\theta$ and $\Sigma$ are as above and $\tau: \Omega \rightarrow G^{0}(3,1)$ is an $\mathscr{H}^{1}$-measurable function such that $P(\tau(x))=T_{x} \Sigma$ for $\mathscr{H}^{1}$-almost every $x \in \Sigma$. Furthermore, let $\theta_{1}, \theta_{2}: \Omega \rightarrow \mathbb{N} \cup\{0\}$ be $\mathcal{H}^{1}$-measurable functions such that $\theta=\theta_{1}+\theta_{2}$. Then the oriented 1-varifold $W$ on $\Omega$ defined by the condition that

$$
\int_{G_{1}^{0}(\Omega)} \psi d W=\int_{\Sigma}\left(\psi(x, \tau(x)) \theta_{1}(x)+\psi(x,-\tau(x)) \theta_{2}(x)\right) d \mathscr{H}^{1}(x)
$$

for every $\psi \in C_{0}^{0}\left(G_{1}(\Omega)\right)$, is called an oriented integral 1-varifold.
For a 1 -varifold $V$ on $\Omega$, the first variation $\delta V$ is the functional on $C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$ given by

$$
\delta V(\phi)=\int_{G_{1}(\Omega)} \operatorname{div}_{L} \phi(x) d V(x, L) .
$$

For an oriented 1 -varifold $W$, we set $\delta W=\delta\left(\operatorname{id}_{\Omega} \times P\right)_{\#} W$.
We now give another representation of the same concepts, before we finally generalize them. Note that any $L \in G(3,1)$ can be identified with the $(3 \times 3)$-matrix belonging to the orthogonal projection onto the orthogonal complement of $L$. We denote this matrix by $\Pi_{L}^{\perp}$, and similarly we write $\Pi_{L}$ for the matrix of the orthogonal projection onto $L$. Any point $\tau \in \mathbb{S}^{2}=G^{0}(3,1)$ can be identified with the pair $\left(\Pi_{P}^{\perp}(\tau), \Lambda_{\tau}\right)$, where $\Lambda_{\tau} \in \mathbb{R}^{3 \times 3}$ is the matrix such that $\Lambda_{\tau} \xi=\tau \wedge \xi$ for all $\xi \in \mathbb{R}^{3}$. Let $\tilde{F}$ be the space of all pairs $(A, B) \in \mathbb{R}^{3 \times 3} \times \mathbb{R}^{3 \times 3}$ such that
(a) $A^{T}=A$,
(b) $A^{2}=A$,
(c) $\operatorname{tr} A=2$,
(d) $B^{T}=-B$,
(e) $B^{2}=-A$.

Then we have a natural diffeomorphism $\Phi: G^{0}(3,1) \rightarrow \tilde{F}$. If $V$ is an oriented 1 -varifold on $\Omega$, then $\left(\operatorname{id}_{\Omega} \times \Phi\right)_{\#} V$ is a Radon measure on $\Omega \times \tilde{F}$, and this mapping provides a natural identification of the space of oriented 1 -varifolds on $\Omega$ with the space of Radon measures on $\Omega \times \tilde{F}$. If $\mu=\left(\operatorname{id}_{\Omega} \times \Phi\right)_{\#} V$, then we have

$$
\begin{equation*}
\delta V(\phi)=\int_{\Omega \times \tilde{F}}(\operatorname{div} \phi(x)-\operatorname{tr}(A \nabla \phi(x))) d \mu(x, A, B) . \tag{16}
\end{equation*}
$$

Next we relax the conditions (a)-(e) as follows. We replace (b) and (e) by
(b') $A$ is positive semidefinite and $|A|^{2} \leqslant 4$,
(e) $|B|^{2} \leqslant 2$.

Here $|\cdot|$ denotes the Hilbert-Schmidt norm. Let $F$ be the space of all pairs $(A, B) \in \mathbb{R}^{3 \times 3} \times \mathbb{R}^{3 \times 3}$ such that (a), (b'), (c), (d), and (e') are satisfied. We now consider Radon measures on $\Omega \times F$.

We have the inclusion map $\iota: \tilde{F} \rightarrow F$, thus any oriented 1 -varifold $V$ on $\Omega$ induces a Radon measure (id ${ }_{\Omega} \times$ $(\iota \circ \Phi))_{\#} V$ on $\Omega \times F$. Moreover, we will see that as a result of the relaxation, a map $u \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$ also induces a Radon measure on $\Omega \times F$ in a natural way. This is the reason why we consider this space.

For any $u \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$, we define the functions $A_{u}, B_{u}: \Omega \rightarrow \mathbb{R}^{3 \times 3}$ with

$$
\begin{equation*}
A_{u}(x)=2 \frac{\nabla u(x) \otimes \nabla u(x)}{|\nabla u(x)|^{2}} \quad \text { and } \quad B_{u}(x)=2 \frac{(u(x) \wedge \nabla u(x)) \otimes \nabla u(x)}{|\nabla u(x)|^{2}} \tag{17}
\end{equation*}
$$

for every $x \in \Omega$ with $\nabla u(x) \neq 0$, and $A_{u}(x)=A_{0}, B_{u}(x)=B_{0}$ if $\nabla u(x)=0$ for a fixed (but arbitrary) point $\left(A_{0}, B_{0}\right) \in F$. It is readily checked that $\left(A_{u}(x), B_{u}(x)\right) \in F$ at every $x \in \Omega$. Now consider the Radon measure $\mu$ on $\Omega$ such that for every $\psi \in C_{0}^{0}(\Omega \times F)$,

$$
\begin{equation*}
\int_{\Omega \times F} \psi d \mu=\frac{1}{2} \int_{\Omega} \psi\left(x, A_{u}(x), B_{u}(x)\right)|\nabla u(x)|^{2} d x . \tag{18}
\end{equation*}
$$

Measures of this type will play an important role in the proof of Theorem 1.1.
Motivated by (16), we define

$$
\delta \mu(\phi)=\int_{\Omega \times F}(\operatorname{div} \phi(x)-\operatorname{tr}(A \nabla \phi(x))) d \mu(x, A, B)
$$

for $\phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$ whenever $\mu$ is a Radon measure on $\Omega \times F$. In order to simplify the notation, we often write such an integral in the form

$$
\int_{\Omega \times F}(\operatorname{div} \phi-\operatorname{tr}(A \nabla \phi)) d \mu .
$$

That is, the symbols $A$ and $B$ represent the standard coordinate functions on $F$. Moreover, whenever it is convenient, we identify a function on $\Omega$ with a function on $\Omega \times F$ that depends only on the first argument. For a measure of the form (18), we then compute

$$
\begin{equation*}
\delta \mu(\phi)=\int_{\Omega}\left(\frac{1}{2}|\nabla u|^{2} \operatorname{div} \phi-\operatorname{tr}((\nabla u \otimes \nabla u) \nabla \phi)\right) d x=\int_{\Omega}\langle\Delta u, \nabla u\rangle \cdot \phi d x \tag{19}
\end{equation*}
$$

by an integration by parts.

For any Radon measure $\mu$ on $\Omega \times F$, we define

$$
W(\mu)=\sup \left\{(\delta \mu(\phi))^{2}: \phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right) \text { with } \int_{\Omega \times F} \phi \cdot A \phi d \mu \leqslant 1\right\} .
$$

Suppose $\mu=\left(\operatorname{id}_{\Omega} \times(\iota \circ \Phi)\right)_{\#} V$ for an integral 1-varifold $V$ given by the function $\theta: \Omega \rightarrow \mathbb{N} \cup\{0\}$. Let $\Sigma=\theta^{-1}(\mathbb{N})$, as before. If there exists a vector field $H: \Omega \rightarrow \mathbb{R}^{3}$ such that

$$
\delta V(\mu)+\int_{\Sigma} \phi \cdot H \theta d \mathscr{H}^{1}=0
$$

for every $\phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$ and $H(x) \perp T_{x} \Sigma$ for $\mathscr{H}^{1}$-almost every $x \in \Sigma$, then we easily calculate

$$
W(\mu)=\int_{\Sigma}|H|^{2} \theta d \mathscr{H}^{1} .
$$

If no such $H$ exists, then we have $W(\mu)=\infty$. Thus the functional $W$ can be regarded as a generalization of the Willmore functional. (See Willmore [28] for a definition and basic properties of the Willmore functional.)

If $\mu$ is defined by (18) for a map $u \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$, then we obtain

$$
W(\mu) \leqslant \int_{\Omega}|u \wedge \Delta u|^{2} d x
$$

from (19). If the rank of $\nabla u$ is either 2 or 0 at almost every point of $\Omega$, then we have even equality here.
Consider again a general Radon measure $\mu$ on $\Omega \times F$. We use the projection $Q: \Omega \times F \rightarrow \Omega$ to define the weight measure $\|\mu\|=Q_{\#} \mu$, which is a Radon measure on $\Omega$. The fibre measure $\mu^{\left(x_{0}\right)}$ on $F$, which is defined by

$$
\int_{F} \zeta d \mu^{\left(x_{0}\right)}=\lim _{r \searrow 0} \frac{\int_{B_{r}\left(x_{0}\right) \times F} \zeta(A, B) d \mu(x, A, B)}{\|\mu\|\left(B_{r}\left(x_{0}\right)\right)} \quad \text { for } \zeta \in C^{0}(F),
$$

then exists for $\|\mu\|$-almost every $x_{0} \in \Omega$. Moreover, we have

$$
\int_{\Omega \times F} \psi d \mu=\int_{\Omega} \int_{F} \psi(x, A, B) d \mu^{(x)}(A, B) d\|\mu\|(x)
$$

for every $\psi \in C_{0}^{0}(\Omega \times F)$ (cf. Allard [1], Section 3.3). Finally, we define the functions $A^{(\mu)}, B^{(\mu)}: \Omega \rightarrow \mathbb{R}^{3 \times 3}$ by

$$
\begin{aligned}
& A^{(\mu)}(x)=\int_{F} A d \mu^{(x)}, \\
& B^{(\mu)}(x)=\int_{F} B d \mu^{(x)}
\end{aligned}
$$

at every point where this exists, and $A^{(\mu)}=A_{0}, B^{(\mu)}=B_{0}$ elsewhere. Since $F$ is convex, we have $\left(A^{(\mu)}(x)\right.$, $\left.B^{(\mu)}(x)\right) \in F$ for every $x \in \Omega$.

If $W(\mu)<\infty$, then Proposition 2.1 in [19], applied to the measure

$$
\|\mu\|\left\llcorner A^{(\mu)}\right.
$$

implies that there exists a unique $\|\mu\|$-measurable function $H: \Omega \rightarrow \mathbb{R}^{3}$ with

$$
\begin{equation*}
H(x) \perp \operatorname{ker} A^{(\mu)}(x) \tag{20}
\end{equation*}
$$

for $\|\mu\|$-almost every $x \in \Omega$, such that

$$
\begin{equation*}
\delta \mu(\phi)+\int_{\Omega \times F} \phi \cdot A H d \mu=0 \tag{21}
\end{equation*}
$$

for every $\phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$ and

$$
W(\mu)=\int_{\Omega \times F} H \cdot A H d \mu
$$

If $\mu$ belongs to an integral varifold, then this $H$ coincides of course with the vector field considered earlier.
If $\mu$ belongs to a map $u \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$-that is, if it is given by (18)—then we have

$$
\|\mu\|=\frac{1}{2} \mathcal{L}^{3}\left\llcorner|\nabla u|^{2}\right.
$$

We denote the Dirac measure centred at a point $p$ (in $\Omega, F$, or $\Omega \times F$ ) by $\delta_{p}$. Then we have

$$
\mu^{(x)}=\delta_{A_{u}(x)} \times \delta_{B_{u}(x)}
$$

and

$$
A^{(\mu)}(x)=A_{u}(x), \quad B^{(\mu)}(x)=B_{u}(x)
$$

for every $x \in \Omega$ where $\nabla u(x) \neq 0$. According to (19) and (21), we have

$$
\int_{\Omega \times F} \phi \cdot A H d \mu=-\int_{\Omega} \phi \cdot\langle\Delta u, \nabla u\rangle d x
$$

for every $\phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$. Hence

$$
(\nabla u \otimes \nabla u) H=-\langle\Delta u, \nabla u\rangle
$$

almost everywhere in $\Omega$. We set

$$
Z=\{x \in \Omega: \text { dim ker } \nabla u(x) \neq 2\}
$$

We claim that

$$
((u \wedge \nabla u) \otimes \nabla u) H=-\langle u \wedge \Delta u, \nabla u\rangle
$$

almost everywhere in $Z$. That is,

$$
\begin{equation*}
\int_{Z \times F} \phi \cdot B H d \mu=-\int_{Z} \phi \cdot\langle u \wedge \Delta u, \nabla u\rangle d x \tag{22}
\end{equation*}
$$

for every $\phi \in C_{0}^{0}\left(\Omega, \mathbb{R}^{3}\right)$.
This is in fact quite easy to verify if the right coordinates are used. Suppose $x \in Z$ with $\nabla u(x) \neq 0$. Since $\nabla u(x) \otimes$ $\nabla u(x)$ is positive semidefinite and of rank 2 , there exists some $R \in \mathrm{SO}(3)$ such that

$$
R^{T}(\nabla u(x) \otimes \nabla u(x)) R=\left(\begin{array}{ccc}
c_{1}^{2} & 0 & 0  \tag{23}\\
0 & c_{2}^{2} & 0 \\
0 & 0 & 0
\end{array}\right)
$$

for certain numbers $c_{1}, c_{2}>0$. Replace $u$ by the map $\tilde{u}(y)=u(R y+x)$, then $\nabla \tilde{u}(0) \otimes \nabla \tilde{u}(0)$ is the matrix on the right hand side of (23). We work with $\tilde{u}$ instead of $u$ now (and we drop the tilde again). We also choose coordinates in the target space such that $u(x)=(0,0,1)$ and

$$
\frac{\partial u}{\partial x^{1}}(x)=\left(c_{1}, 0,0\right) \quad \text { and } \quad \frac{\partial u}{\partial x^{2}}(x)=\left(0, \pm c_{2}, 0\right)
$$

Then

$$
(u(x) \wedge \nabla u(x)) \otimes \nabla u(x)= \pm\left(\begin{array}{ccc}
0 & c_{1} c_{2} & 0 \\
-c_{1} c_{2} & 0 & 0 \\
0 & 0 & 0
\end{array}\right)
$$

Let

$$
\Delta u(x)+|\nabla u(x)|^{2} u(x)=\left(d_{1}, d_{2}, 0\right)
$$

Then

$$
\langle\Delta u, \nabla u\rangle=\left(\begin{array}{c}
c_{1} d_{1} \\
\pm c_{2} d_{2} \\
0
\end{array}\right),
$$

hence

$$
H(x)=\left(\begin{array}{c}
d_{1} / c_{1} \\
\pm d_{2} / c_{2} \\
0
\end{array}\right) .
$$

Thus we have

$$
((u(x) \wedge \nabla u(x)) \otimes \nabla u(x)) H(x)=\left(\begin{array}{c}
c_{1} d_{2} \\
\mp c_{2} d_{1} \\
0
\end{array}\right)=-\langle u(x) \wedge \Delta u(x), \nabla u(x)\rangle,
$$

which proves (22).
If we assume that $Z=\Omega$, then the energy identities (8) and (12) become formally the same in the framework of generalized varifolds. This is one of the reasons why this is a useful tool for our problem.

When we work with a map $u \in C^{\infty}\left(\Omega_{T}, \mathbb{S}^{2}\right)$, then we define the functions $A_{u}, B_{u}: \Omega_{T} \rightarrow \mathbb{R}^{3 \times 3}$ similarly as above. They give rise to a Radon measure $\mu$ on $\Omega_{T} \times F$. The quantities $\|\mu\|, \delta \mu, \mu^{(z)}, A^{(\mu)}, B^{(\mu)}$, etc., are then defined similarly as on $\Omega$.

We now briefly discuss another tool from geometric measure theory, namely the measure-function pairs introduced by Hutchinson [10].

Let $M$ be a manifold, with or without boundary, which is embedded in $\mathbb{R}^{m}$. If $\mu$ is a Radon measure on $M$ and $f: M \rightarrow \mathbb{R}^{n}$ a locally $\mu$-integrable function on $M$, then we say that $(\mu, f)$ is a measure-function pair over $M$ (with values in $\left.\mathbb{R}^{n}\right)$. We are mainly interested in the case $f \in L^{2}\left(\mu, \mathbb{R}^{n}\right)$.

To a measure-function pair $(\mu, f)$ we can assign the graph measure $[\mu, f]$ on $M \times \mathbb{R}^{n}$, which is a Radon measure defined by the condition

$$
\int_{M \times \mathbb{R}^{n}} \xi d[\mu, f]=\int_{M} \xi(x, f(x)) d \mu(x) \quad \text { for } \xi \in C_{0}^{0}\left(M \times \mathbb{R}^{n}\right) .
$$

Obviously we can represent the generalized varifold belonging to a map in $C^{\infty}\left(\Omega, S^{2}\right)$ as a graph measure of this type. More important for our purpose, however, are measure-function pairs of the form $(\mu, H)$, where $\mu$ is a generalized varifold with $W(\mu)<\infty$ and $H$ is the function characterized by (20) and (21), interpreted as a function on $\Omega \times F$ which depends only on the first argument. In general, this function is not necessarily in $L_{\mathrm{loc}}^{1}\left(\mu, \mathbb{R}^{3}\right)$. However, if $\mu$ belongs to a map $u \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$ such that $\nabla u$ is of rank 2 in $\Omega$, then we have even $H \in L_{\text {loc }}^{\infty}\left(\mu, \mathbb{R}^{3}\right)$.

Now suppose that we have a sequence of measure-function pairs ( $\mu_{k}, f_{k}$ ) over $M$ with values in $\mathbb{R}^{n}$ such that $f_{k} \in L^{2}\left(\mu_{k}, \mathbb{R}^{n}\right)$. Furthermore, we assume that $(\mu, f)$ is another measure-function pair of this type, such that $\mu_{k} \rightarrow \mu$. We say that the sequence $\left\{\left(\mu_{k}, f_{k}\right)\right\}$ converges weakly to $(\mu, f)$ if we have $\mu_{k} L f_{k} \rightarrow \mu\llcorner f$. We say that the convergence is strong if $\left[\mu_{k}, f_{k}\right] \rightarrow[\mu, f]$ and

$$
\int_{\left.\mathbb{R}^{n}:|x|^{2}+|y|^{2} \geqslant R\right\}} y^{2} d\left[\mu_{k}, f_{k}\right](x, y) \rightarrow 0 \quad \text { as } R \rightarrow \infty
$$

uniformly in $k$. Other characterizations of these (and similar) notions of convergence are given in [10]. The proof of the following result is also to be found there.

## Proposition 2.1.

(i) Suppose $\left(\mu_{k}, f_{k}\right)$ are measure-function pairs over $M$ such that

$$
\sup _{k \in \mathbb{N}}\left\|f_{k}\right\|_{L^{2}\left(\mu_{k}\right)}<\infty
$$

If $\mu$ is a Radon measure on $M$ with $\mu_{k} \rightarrow \mu$, then there exists a function $f \in L^{2}(\mu)$ such that a subsequence of $\left\{\left(\mu_{k}, f_{k}\right)\right\}$ converges weakly to $(\mu, f)$.
(ii) Suppose that the sequence $\left\{\left(\mu_{k}, f_{k}\right)\right\}$ converges weakly to $(\mu, f)$. Then

$$
\|f\|_{L^{2}(\mu)} \leqslant \liminf _{k \rightarrow \infty}\left\|f_{k}\right\|_{L^{2}\left(\mu_{k}\right)}
$$

The convergence is strong if and only if

$$
\|f\|_{L^{2}(\mu)}=\lim _{k \rightarrow \infty}\left\|f_{k}\right\|_{L^{2}\left(\mu_{k}\right)}
$$

## 3. Other tools

In this section we collect a few facts that are either proved elsewhere or that follow from known arguments. First we have an estimate for the derivatives of a solution of the Landau-Lifshitz equation with small energy. The arguments that prove this result can be found in Sections 5.2 and 5.5 of [20] (although the result is not stated in the same form in that work).

Lemma 3.1. There exist two numbers $C, \epsilon_{0}>0$ such that every solution $u \in C^{\infty}\left(P_{1}, \mathbb{S}^{2}\right)$ of (1) with

$$
\frac{1}{2} \int_{P_{1}}|\nabla u|^{2} d z=\epsilon \leqslant \epsilon_{0}
$$

satisfies

$$
\sup _{P_{1 / 2}}\left(|\nabla u|^{2}+\left|\nabla^{2} u\right|^{2}\right) \leqslant C \epsilon
$$

The next statement gives a version of the so-called monotonicity formula for harmonic maps. A proof for this formula is given in [20], Section 4.1.

Lemma 3.2. Suppose $u \in C^{\infty}\left(B_{1}, \mathbb{S}^{2}\right)$. For $0<\rho \leqslant 1$, set

$$
\Phi(\rho)=\frac{1}{\rho} \int_{B_{\rho}}\left(\frac{1}{2}|\nabla u|^{2}-x \cdot\langle\Delta u, \nabla u\rangle\right) d x
$$

Then

$$
\Phi(r)-\Phi(s)=\int_{B_{r} \backslash B_{s}}\left(\frac{|x \cdot \nabla u|^{2}}{|x|^{3}}-\frac{x \cdot\langle\Delta u, \nabla u\rangle}{|x|}\right) d x
$$

for $0<s \leqslant r \leqslant 1$.
We will also need the following estimates. Here we write $x=\left(x^{\prime}, x^{3}\right)$ for a point in $\mathbb{R}^{3}$, where $x^{\prime}=\left(x^{1}, x^{2}\right)$.
Lemma 3.3. Suppose $\xi \in C_{0}^{1}\left(-\frac{1}{2}, \frac{1}{2}\right)$ and $\zeta \in C_{0}^{1}(-1,1)$ satisfy

$$
\int_{-1 / 2}^{1 / 2} \xi(s) d s=\int_{-1}^{1} \zeta(t) d t=1
$$

Then there exists a constant $C$ with the following properties.
(i) For every $u \in C^{\infty}\left(B_{1}, \mathbb{S}^{2}\right)$ and every $\eta \in C_{0}^{1}\left(B_{1 / 2}^{2}\right)$, the inequality

$$
\begin{aligned}
& \left.\left|\int_{B_{1}} \eta\left(x^{\prime}\right) \xi\left(x^{3}\right)\right| \nabla u\left(x^{\prime}, x^{3}\right)\right|^{2} d x-\int_{B_{1 / 2}^{2}} \eta\left(x^{\prime}\right)\left|\nabla u\left(x^{\prime}, 0\right)\right|^{2} d x^{\prime} \mid \\
& \quad \leqslant C\|\eta\|_{C^{1}\left(B_{1 / 2}^{2}\right)}\|\nabla u\|_{L^{2}\left(B_{1}\right)}\left(\left\|\frac{\partial u}{\partial x^{3}}\right\|_{L^{2}\left(B_{1}\right)}+\|u \wedge \Delta u\|_{L^{2}\left(B_{1}\right)}\right)+C\|\eta\|_{C^{0}\left(B_{1 / 2}^{2}\right)}\left\|\frac{\partial u}{\partial x^{3}}(\cdot, 0)\right\|_{L^{2}\left(B_{1}^{2}\right)}^{2}
\end{aligned}
$$

holds.
(ii) For every solution $u \in C^{\infty}\left(P_{1}, \mathbb{S}^{2}\right)$ of (1) and every $\eta \in C_{0}^{1}\left(B_{1}\right)$, the inequality

$$
\begin{aligned}
& \left.\left|\int_{P_{1}} \eta(x) \zeta(t)\right| \nabla u(x, t)\right|^{2} d z-\int_{B_{1}} \eta(x)|\nabla u(x, 0)|^{2} d x \mid \\
& \quad \leqslant C\|\eta\|_{C^{1}\left(B_{1}\right)}\left\|\frac{\partial u}{\partial t}\right\|_{L^{2}\left(P_{1}\right)}\left(\|\nabla u\|_{L^{2}\left(P_{1}\right)}+\left\|\frac{\partial u}{\partial t}\right\|_{L^{2}\left(P_{1}\right)}\right)
\end{aligned}
$$

holds.
(iii) For every solution $u \in C^{\infty}\left(P_{1}, \mathbb{S}^{2}\right)$ of (1) and every $\eta \in C_{0}^{1}\left(B_{1 / 2}\right)$, the inequality

$$
\begin{aligned}
& \left.\left|\int_{P_{1}} \eta(x) \zeta(t)\right| \nabla u(x, t)\right|^{2} d z-\int_{B_{1}} \eta(x)|\nabla u(x, 0)|^{2} d x \mid \\
& \quad \leqslant C\|\eta\|_{C^{1}\left(B_{1 / 2}\right)}\left\|\frac{\partial u}{\partial t}\right\|_{L^{2}\left(P_{1}\right)}\left(\|\nabla u(\cdot, 0)\|_{L^{2}\left(B_{1}\right)}+\left\|\frac{\partial u}{\partial t}\right\|_{L^{2}\left(P_{1}\right)}\right)
\end{aligned}
$$

holds.
Proof. For part (i) of the lemma, we define $\tilde{\eta}\left(x^{\prime}, x^{3}\right)=\eta\left(x^{\prime}\right)$ and $\tilde{\xi}\left(x^{\prime}, x^{3}\right)=\xi\left(x^{3}\right)$. We consider the vector field

$$
A=\tilde{\eta}\left(\frac{1}{2}|\nabla u|^{2} e_{3}-\left\langle\frac{\partial u}{\partial x^{3}}, \nabla u\right\rangle\right),
$$

where $e_{3}=(0,0,1)$. We compute

$$
\operatorname{div} A=-\tilde{\eta}\left\langle\frac{\partial u}{\partial x^{3}}, \Delta u\right\rangle-\nabla \tilde{\eta} \cdot\left\langle\frac{\partial u}{\partial x^{3}}, \nabla u\right\rangle .
$$

Thus we find

$$
\begin{aligned}
& \int_{B_{1 / 2}^{2} \times\{s\}} \tilde{\eta}\left(\frac{1}{2}|\nabla u|^{2}-\left|\frac{\partial u}{\partial x^{3}}\right|^{2}\right) d x^{\prime}-\int_{B_{1 / 2}^{2} \times\{0\}} \tilde{\eta}\left(\frac{1}{2}|\nabla u|^{2}-\left|\frac{\partial u}{\partial x^{3}}\right|^{2}\right) d x^{\prime} \\
& \quad=-\int_{B_{1 / 2}^{2} \times(0, s)}\left(\tilde{\eta}\left\langle\frac{\partial u}{\partial x^{3}}, \Delta u\right\rangle+\nabla \tilde{\eta} \cdot\left\langle\frac{\partial u}{\partial x^{3}}, \nabla u\right\rangle\right) d x
\end{aligned}
$$

for $s \in\left(0, \frac{1}{2}\right)$ and a similar formula for $s \in\left(-\frac{1}{2}, 0\right)$. Multiplying both sides with $\xi(s)$, integrating over $s$, and using Fubini's theorem, we obtain

$$
\begin{aligned}
& \quad \int_{B_{1 / 2}^{2} \times\left(-\frac{1}{2}, \frac{1}{2}\right)} \tilde{\eta} \tilde{\xi}\left(\frac{1}{2}|\nabla u|^{2}-\left|\frac{\partial u}{\partial x^{3}}\right|^{2}\right) d x-\int_{B_{1 / 2}^{2} \times\{0\}} \tilde{\eta}\left(\frac{1}{2}|\nabla u|^{2}-\left|\frac{\partial u}{\partial x^{3}}\right|\right) d x^{\prime} \\
& =-\int_{B_{1 / 2}^{2} \times\left(-\frac{1}{2}, \frac{1}{2}\right)} \tilde{\sigma}\left(\tilde{\eta}\left\langle\frac{\partial u}{\partial x^{3}}, \Delta u\right\rangle+\nabla \tilde{\eta} \cdot\left\langle\frac{\partial u}{\partial x^{3}}, \nabla u\right\rangle\right) d x,
\end{aligned}
$$

where

$$
\tilde{\sigma}\left(x^{\prime}, x^{3}\right)= \begin{cases}\int_{x^{3}}^{1 / 2} \xi(s) d s & x^{3}>0 \\ -\int_{-1 / 2}^{x^{3}} \xi(s) d s & x^{3}<0\end{cases}
$$

The estimate in (i) follows with Young's inequality.
Part (ii) is an easy consequence of the energy identity (8).

For part (iii), it suffices to show that there exists a constant $C_{1}$ that depends only on $a$ and $b$, such that

$$
\begin{equation*}
\int_{B_{1 / 2} \times(-1,1)}|\nabla u|^{2} d z \leqslant C_{1}\left(\int_{B_{1}}|\nabla u(x, 0)|^{2} d x+\int_{P_{1}}\left|\frac{\partial u}{\partial t}\right|^{2} d z\right) . \tag{24}
\end{equation*}
$$

Once this is verified, also the inequality in (iii) follows directly from (8). To prove (24), choose a cut-off function $\omega \in C_{0}^{\infty}\left(B_{1}\right)$ with $\omega \equiv 1$ in $B_{1 / 2}$. As a consequence of (8), we have a constant $C_{2}=C_{2}(a, b, \omega)$ such that

$$
\left.\left.\left|\frac{d}{d t} \int_{B_{1}} \omega^{2}(x)\right| \nabla u(x, t)\right|^{2} d x\left|\leqslant \int_{B_{1}} \omega^{2}(x)\right| \nabla u(x, t)\right|^{2} d x+C_{2} \int_{B_{1}}\left|\frac{\partial u}{\partial t}\right|^{2} d x .
$$

Now (24) follows easily from this inequality.
Next we consider a sequence of maps $u_{k} \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$ such that

$$
\sup _{k \in \mathbb{N}}\left(E\left(u_{k}\right)+\int_{\Omega}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x\right)<\infty
$$

For such a sequence, a similar energy concentration phenomenon as described in Theorem 1.1 can be observed. This behaviour and the corresponding energy concentration measure have been studied in [19]. We repeat here the main results of that paper (in the case of a three-dimensional domain and the target manifold $\mathbb{S}^{2}$ ).

Theorem 3.1. Under the above conditions, there exists a subsequence $\left\{u_{k_{i}}\right\}$ with the following properties.
(i) There exists a map $u \in H^{1}\left(\Omega, \mathbb{S}^{2}\right)$ with

$$
\Delta u+|\nabla u|^{2} u \in L^{2}\left(\Omega, \mathbb{R}^{3}\right),
$$

such that $u_{k_{i}} \rightharpoonup u$ weakly in $H^{1}\left(\Omega, \mathbb{R}^{3}\right)$ and

$$
\Delta u_{k_{i}}+\left|\nabla u_{k_{i}}\right|^{2} u_{k_{i}} \rightharpoonup \Delta u+|\nabla u|^{2} u \quad \text { weakly in } L^{2}\left(\Omega, \mathbb{R}^{3}\right) .
$$

(ii) There exists an $\mathscr{H}^{1}$-integrable function $\theta: \Omega \rightarrow[0, \infty)$, such that for any $\eta \in C_{0}^{0}(\Omega)$,

$$
\lim _{i \rightarrow \infty} \int_{\Omega} \eta\left|\nabla u_{k_{i}}\right|^{2} d x=\int_{\Omega} \eta|\nabla u|^{2} d x+8 \pi \int_{\Omega} \eta \theta d \mathscr{H}^{1}
$$

and the set $\Sigma=\theta^{-1}((0, \infty))$ is closed relative to $\Omega$ and countably 1-rectifiable.
(iii) There exists an $\mathscr{H}^{1}$-measurable vector field $H: \Omega \rightarrow \mathbb{R}^{3}$ with $H(x) \perp T_{x} \Sigma$ for $\mathscr{H}^{1}$-almost every $x \in \Sigma$, such that

$$
\int_{\Omega}\left(\frac{1}{2}|\nabla u|^{2} \operatorname{div} \phi-\operatorname{tr}((\nabla u \otimes \nabla u) \nabla \phi)-\langle\Delta u, \nabla u\rangle \cdot \phi\right) d x+4 \pi \int_{\Omega}\left(\operatorname{div}_{\Sigma} \phi+H \cdot \phi\right) \theta d \mathscr{H}^{1}=0
$$

for every $\phi \in C_{0}^{1}\left(\Omega, \mathbb{R}^{3}\right)$.
(iv) The inequality

$$
\int_{\Omega}|u \wedge \Delta u|^{2} d x+4 \pi \int_{\Omega}|H|^{2} \theta d \mathscr{H}^{1} \leqslant \liminf _{k \rightarrow \infty} \int_{\Omega}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x
$$

holds.
If we have a sequence of solutions $u_{k} \in C^{\infty}\left(\Omega, \mathbb{S}^{2}\right)$ of (1) that satisfy the hypotheses of Theorem 1.1, then the energy identity (8) implies that for any precompact set $K \Subset \Omega$ and for every $t_{0} \in(0, T)$, we have

$$
\sup _{k \in \mathbb{N}}\left(\int_{K}\left|\nabla u_{k}\left(x, t_{0}\right)\right|^{2} d x+\int_{t_{0}}^{T} \int_{K}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x d t\right)<\infty .
$$

Thus, the above result can be applied to restrictions of $u_{k}$ to $K \times\left\{t_{0}\right\}$ for almost every $t_{0} \in[0, T)$, at least after the choice of a subsequence (which may depend on $t_{0}$ ). Parts (i)-(iii) of Theorem 1.1 then follow with relatively little extra work, except for the statement that $\theta$ is integer-valued. Most of the rest of the paper is therefore dedicated to this quantization property and part (iv) of Theorem 1.1.

## 4. Two-dimensional blow-up analysis

In this section we work in two-dimensional domains. The results will later be applied to cross-sections of a higherdimensional domain. We now write $B_{r}\left(x_{0}\right)$ for an open ball in $\mathbb{R}^{2}$ (or $B_{r}$ if $x_{0}=0$ ). We also use the notations $\hat{F}=\mathbb{R}^{2 \times 2} \times \mathbb{R}^{2 \times 2}$ and

$$
I=\left(\begin{array}{ll}
1 & 0 \\
0 & 1
\end{array}\right), \quad J=\left(\begin{array}{cc}
0 & -1 \\
1 & 0
\end{array}\right)
$$

Remember that a Dirac measure centred at a point $p$ is denoted by $\delta_{p}$.
We study a sequence of maps $u_{k} \in C^{\infty}\left(B_{1}, \mathbb{S}^{2}\right)$. We consider the corresponding measures $\mu_{k}$ on $B_{1} \times \hat{F}$ given by

$$
\begin{equation*}
\int_{B_{1} \times \hat{F}} \psi d \mu_{k}=\frac{1}{2} \int_{B_{1}} \psi\left(x, 2 \frac{\nabla u_{k} \otimes \nabla u_{k}}{\left|\nabla u_{k}\right|^{2}}, 2 \frac{\left(u_{k} \wedge \nabla u_{k}\right) \otimes \nabla u_{k}}{\left|\nabla u_{k}\right|^{2}}\right)\left|\nabla u_{k}\right|^{2} d x \tag{25}
\end{equation*}
$$

for $\psi \in C_{0}^{0}\left(B_{1} \times K\right)$; in other words, the equivalent of (18) for the maps $u_{k}$. The weight measures $\left\|\mu_{k}\right\|$ on $B_{1}$ are then defined similarly as in Section 2.

We examine the blow-up behaviour of this sequence under the assumption that there exists a number $\epsilon_{0}$ such that the following conditions are satisfied.
(I) There exists a number $\theta>0$ such that $\left\|\mu_{k}\right\| \rightarrow 4 \pi \theta \delta_{0}$.
(II) $\sup _{k \in \mathbb{N}}\left\|\nabla^{2} u_{k}\right\|_{L^{1}\left(B_{1}\right)}<\infty$.
(III) Suppose $x_{k} \in B_{1}$ and $r_{k}>0$ are such that $x_{k} \rightarrow 0$ and $r_{k} \rightarrow 0$. If

$$
0<\limsup _{k \rightarrow \infty} \int_{B_{r_{k}}\left(x_{k}\right)}\left|\nabla u_{k}\right|^{2} d x \leqslant 2 \epsilon_{0}
$$

then the sequence of rescaled maps $v_{k}(x)=u_{k}\left(r_{k} x+x_{k}\right)$ subconverges weakly in $H_{\text {loc }}^{1}\left(\mathbb{R}^{2}, \mathbb{R}^{3}\right)$. The limit is a non-constant harmonic map $v \in C^{\infty}\left(\mathbb{R}^{2}, \mathbb{S}^{2}\right)$ (i.e., it satisfies $v \wedge \Delta v=0$ ) with

$$
\int_{\mathbb{R}^{2}}|\nabla v|^{2} d x<\infty
$$

(IV) There exists a constant $C_{0}$ such that for any $\epsilon \in\left(0, \epsilon_{0}\right]$ and for any sequence of balls $B_{r_{k}}\left(x_{k}\right) \subset B_{1}$ with

$$
\limsup _{k \rightarrow \infty} \int_{B_{r_{k}}\left(x_{k}\right)}\left|\nabla u_{k}\right|^{2} d x \leqslant 2 \epsilon,
$$

the inequality

$$
\limsup _{k \rightarrow \infty} \sup _{x \in B_{r_{k} / 2}\left(x_{k}\right)}\left(r_{k}^{2}\left|\nabla u_{k}(x)\right|^{2}+r_{k}^{4}\left|\nabla^{2} u_{k}\right|^{2}\right) \leqslant C_{0} \epsilon
$$

holds.
We want to determine the possible limit measures of the sequence $\left\{\mu_{k}\right\}$ under these assumptions. We achieve this by a blow-up analysis similar to what has been done for harmonic maps by Jost [11] and for the harmonic map heat flow and Palais-Smale sequences for the Dirichlet energy by Qing [21], Ding and Tian [5], Qing and Tian [22], and Lin and Wang [17]. Typical for this method is that the limit measures are described in terms of so-called "harmonic bubbles". These bubbles are harmonic maps which are obtained from rescaled sequences as in (III). It turns out that
they contain all the information about the limit measure. Since harmonic maps of this type are well understood, we obtain a good description of any possible limit of $\left\{\mu_{k}\right\}$.

Proposition 4.1. Under the hypotheses (I)-(IV), the number $\theta$ is an integer. If $\mu_{k} \rightarrow \mu$ for a Radon measure $\mu$ on $B_{1} \times \hat{F}$, then there exists a number $\sigma \in[0,1]$, such that

$$
\mu=4 \pi \theta \delta_{0} \times \delta_{I} \times\left(\sigma \delta_{J}+(1-\sigma) \delta_{-J}\right)
$$

For the proof of this proposition, we use ideas from the papers mentioned earlier and also from Lin and Rivière [16]. One of the tools we need is the following lemma.

Lemma 4.1. There exists a constant $C$ with the following property. Suppose $r \in(0,1)$ and $f \in W^{1,1}\left(B_{1}\right)$ satisfy

$$
\begin{equation*}
\sup _{x \in B_{1} \backslash B_{r}}|x f(x)| \leqslant c . \tag{26}
\end{equation*}
$$

Then

$$
\begin{equation*}
\int_{B_{1} \backslash B_{r}} f^{2} d x \leqslant C c\left(\|f\|_{W^{1,1}\left(B_{1}\right)}+c\right) . \tag{27}
\end{equation*}
$$

Proof. It suffices to consider $r \in\left(0, \frac{1}{8}\right]$. Choose a cut-off function $\xi \in C_{0}^{1}\left(B_{1} \backslash B_{r}\right)$ with $0 \leqslant \xi \leqslant 1$ and $\xi \equiv 1$ in $B_{1 / 2} \backslash B_{2 r}$, such that $|\nabla \xi| \leqslant 4$ in $B_{1} \backslash B_{1 / 2}$ and $|\nabla \xi| \leqslant 2 / r$ in $B_{2 r} \backslash B_{r}$. Define $g=\xi f$. Then

$$
\|g\|_{W^{1,1}\left(\mathbb{R}^{2}\right)} \leqslant\|f\|_{W^{1,1}\left(B_{1}\right)}+16 \pi c .
$$

For every $s \geqslant c$, we have

$$
\left|\left\{x \in \mathbb{R}^{2}:|g(x)|>s\right\}\right| \leqslant \pi \frac{c^{2}}{s^{2}}
$$

by (26), thus we have the estimate

$$
\|g\|_{L^{(2, \infty)}\left(\mathbb{R}^{2}\right)} \leqslant \sqrt{\pi} c
$$

in the Lorentz space $L^{(2, \infty)}\left(\mathbb{R}^{2}\right)$. The Sobolev space $W^{1,1}\left(\mathbb{R}^{2}\right)$ is continuously embedded in $L^{(2,1)}\left(\mathbb{R}^{2}\right)$ (cf. Tartar [27]). Thus it follows that

$$
\int_{\mathbb{R}^{2}} g^{2} d x \leqslant C_{1}\|g\|_{L^{(2,1)}\left(\mathbb{R}^{2}\right)}\|g\|_{L^{(2, \infty)}\left(\mathbb{R}^{2}\right)} \leqslant C_{2} c\left(\|f\|_{W^{1,1}\left(B_{1}\right)}+c\right)
$$

for certain universal constants $C_{1}$ and $C_{2}$. To estimate the remaining part of the integral in (27), we use (26) again.
Proof of Proposition 4.1. We may assume that a limit $\mu$ of the sequence $\left\{\mu_{k}\right\}$ exists. Let $m \in \mathbb{N}$ be the minimal integer such that $\theta \leqslant m$. We prove the proposition by induction on $m$.

Suppose first that $m=1$. Then we choose a sequence of radii $r_{k} \rightarrow 0$ such that

$$
\frac{1}{2} \int_{B_{r_{k}}}\left|\nabla u_{k}\right|^{2} d x=\min \left\{\epsilon_{0}, 2 \pi \theta\right\} .
$$

By (III), there exists a subsequence of the sequence defined by $v_{k}(x)=u_{k}\left(r_{k} x\right)$ which converges weakly in $H_{\text {loc }}^{1}\left(\mathbb{R}^{2}, \mathbb{R}^{3}\right)$ to a non-constant harmonic map $v \in C^{\infty}\left(\mathbb{R}^{2}, \mathbb{S}^{2}\right)$ with finite energy. Such a harmonic map must be conformal and must satisfy

$$
\frac{1}{2} \int_{\mathbb{R}^{2}}|\nabla v|^{2} d x=4 \pi M
$$

for some $M \in \mathbb{N}$ (cf. Sacks and Uhlenbeck [23] and Section 11.5 in Eells and Lemaire [6]). It follows immediately that $M=\theta=1$. Moreover, the above convergence is strong in $H^{1}\left(K, \mathbb{R}^{3}\right)$ for every bounded set $K \subset \mathbb{R}^{2}$.

Fix $x_{0} \in \mathbb{R}^{2}$ and choose $r>0$ such that

$$
\int_{B_{r}\left(x_{0}\right)}|\nabla v|^{2} d x \leqslant \epsilon_{0} .
$$

Then it follows from (IV) that the above subsequence of $\left\{v_{k}\right\}$ converges to $v$ even in the topology of $C^{1}\left(B_{r / 2}\left(x_{0}\right), \mathbb{S}^{2}\right)$. Since $v$ is conformal, we have

$$
2 \frac{\nabla v \otimes \nabla v}{|\nabla v|^{2}}=I \quad \text { and } \quad 2 \frac{(v \wedge \nabla v) \otimes \nabla v}{|\nabla v|^{2}}= \pm J
$$

on $\mathbb{R}^{2}$. The claim of the proposition then follows immediately (with $\sigma=0$ or $\sigma=1$ ).
Next we assume that the statement of the proposition is true whenever $\theta$ does not exceed $m-1$. We want to show that it still holds for $\theta \leqslant m$.

We fix $\epsilon \in\left(0, \epsilon_{0}\right]$. We choose a sequence $r_{k} \rightarrow 0$ such that

$$
\frac{1}{2} \int_{B_{2 r_{k}} \backslash B_{r_{k}}}\left|\nabla u_{k}\right|^{2} d x=\epsilon
$$

but

$$
\frac{1}{2} \int_{B_{2 \rho} \backslash B_{\rho}}\left|\nabla u_{k}\right|^{2} d x \leqslant \epsilon
$$

for every $\rho \in\left(r_{k}, \frac{1}{2}\right]$. Then by (IV), there exists a constant $C_{1}$, such that

$$
\limsup _{k \rightarrow \infty} \sup _{x \in B_{1 / 4} \mid B_{4 r_{k}}}|x|\left|\nabla u_{k}(x)\right| \leqslant C_{1} \sqrt{\epsilon} .
$$

We apply Lemma 4.1 to $f=\left|\nabla u_{k}\right|$. Using also (II), we conclude that

$$
\begin{equation*}
\limsup _{k \rightarrow \infty} \int_{B_{1} \backslash B_{r_{k}}}\left|\nabla u_{k}\right|^{2} d x \leqslant C_{2} \sqrt{\epsilon} \tag{28}
\end{equation*}
$$

where $C_{2}$ is a constant that depends only $\epsilon_{0}, C_{0}$, and the supremum in (II).
Consider the rescaled maps

$$
v_{k}(x)=u_{k}\left(r_{k} x\right)
$$

and the corresponding measures $\tilde{\mu}_{k}$ which are defined similarly as in (25), but with $v_{k}$ instead of $u_{k}$. Using (III), we see that there exists a subsequence which converges weakly in $H_{\mathrm{loc}}^{1}\left(\mathbb{R}^{2}, \mathbb{R}^{3}\right)$ to a non-constant harmonic map $v \in C^{\infty}\left(\mathbb{R}^{2}, \mathbb{S}^{2}\right)$ with finite energy. We assume for simplicity that this convergence holds for the full sequence. We may also assume that $\tilde{\mu}_{k} \rightarrow \tilde{\mu}$ for a Radon measure $\tilde{\mu}$ on $\mathbb{R}^{2} \times \hat{F}$. If we can show that there exist $\hat{m} \in \mathbb{N}$ and $\sigma \in[0,1]$ such that

$$
\begin{equation*}
\int_{\mathbb{R}^{2} \times \hat{F}} \psi(A, B) d \tilde{\mu}(x, A, B)=4 \pi \hat{m}(\sigma \psi(I, J)+(1-\sigma) \psi(I,-J)) \tag{29}
\end{equation*}
$$

for every $\psi \in C_{0}^{0}(\hat{F})$, then the claim of the proposition follows, because we have (28) for an arbitrarily small number $\epsilon \in\left(0, \epsilon_{0}\right]$.

Define

$$
\Sigma_{0}=\left\{x \in \mathbb{R}^{2}:\|\tilde{\mu}\|(\{x\}) \geqslant \epsilon_{0} / 2\right\} .
$$

Because of (28), we have $\Sigma_{0} \subset \overline{B_{1}}$ whenever $\epsilon$ is chosen sufficiently small. Moreover, this is a finite set. At every point $x_{0} \in \Sigma_{0}$, it is easy to find a sequence $s_{k} \rightarrow 0$, such that the maps $w_{k}(x)=v_{k}\left(s_{k} x+x_{0}\right)$ satisfy

$$
\frac{1}{2} \mathcal{L}^{2}\left\llcorner\left|\nabla w_{k}\right|^{2} \rightarrow\|\tilde{\mu}\|\left(\left\{x_{0}\right\}\right) \delta_{0}\right.
$$

Such a sequence also satisfies the conditions (I)-(IV), but with the number

$$
\theta_{x_{0}}=\frac{\|\tilde{\mu}\|\left(\left\{x_{0}\right\}\right)}{4 \pi}
$$

in (I) instead of $\theta$, which satisfies $\theta_{x_{0}} \leqslant \theta-1$. By our assumptions, this means that we have

$$
\int_{\left\{x_{0}\right\} \times \hat{F}} \psi d \tilde{\mu}=4 \pi m_{x_{0}}\left(\sigma_{x_{0}} \psi(I, J)+\left(1-\sigma_{x_{0}}\right) \psi(I,-J)\right)
$$

for certain numbers $m_{x_{0}} \in \mathbb{N}$ and $\sigma_{x_{0}} \in[0,1]$. Similarly as in the first part of the proof we obtain

$$
\int_{\left(\mathbb{R}^{2} \backslash \Sigma_{0}\right) \times \hat{F}} \psi d \tilde{\mu}=4 \pi \tilde{m}(\tilde{\sigma} \psi(I, J)+(1-\tilde{\sigma}) \psi(I,-J))
$$

for some $\tilde{m} \in \mathbb{N}$ and $\tilde{\sigma} \in[0,1]$. Combining these identities, we obtain (29) for

$$
\hat{m}=\tilde{m}+\sum_{x_{0} \in \Sigma_{0}} m_{x_{0}}
$$

and

$$
\sigma=\frac{1}{\hat{m}}\left(\tilde{m} \tilde{\sigma}+\sum_{x_{0} \in \Sigma_{0}} m_{x_{0}} \sigma_{x_{0}}\right),
$$

and the proof is complete.

## 5. Three-dimensional blow-up analysis

We now examine a sequence of solutions $u_{k} \in C^{\infty}\left(P_{1}, \mathbb{S}^{2}\right)$ of the Landau-Lifshitz equation (1) that arises from rescaling a sequence as in Theorem 1.1 around a typical blow-up point. (We will see later what "typical" means here.)

We define the measures $\mu_{k}$ on $P_{1} \times F$ which belong to $u_{k}$ and are given by

$$
\int_{P_{1} \times F} \psi d \mu_{k}=\frac{1}{2} \int_{P_{1}} \psi\left(z, A_{u_{k}}(z), B_{u_{k}}(z)\right)\left|\nabla u_{k}(z)\right|^{2} d z \quad \text { for } \psi \in C_{0}^{0}\left(P_{1} \times F\right)
$$

Here $A_{u_{k}}$ and $B_{u_{k}}$ are the functions defined similarly as in (17). We assume that $\mu_{k} \rightarrow \mu$ for a Radon measure $\mu$ on $P_{1} \times F$. We also assume that this limit measure has a special structure: namely, that there exist a one-dimensional linear subspace $L \subset \mathbb{R}^{3}$ and a number $\theta>0$ such that

$$
\begin{equation*}
\|\mu\|=4 \pi \theta \mathscr{H}^{2}\left\llcorner\left((L \times \mathbb{R}) \cap P_{1}\right)\right. \tag{30}
\end{equation*}
$$

We then fix a unit vector $\tau \in L \cap \mathbb{S}^{2}$. Finally, we assume that

$$
\begin{equation*}
\lim _{k \rightarrow 0} \int_{P_{1}}\left|\frac{\partial u_{k}}{\partial t}\right|^{2} d z=0 \tag{31}
\end{equation*}
$$

Proposition 5.1. Under the above hypotheses, there exist $m \in \mathbb{N}$ and $\sigma \in[0,1]$, such that for every $\psi \in C^{0}(F)$,

$$
\int_{P_{1 / 4} \times F} \psi(A, B) d \mu(z, A, B)=\frac{\pi m}{4}\left(\sigma \psi\left(\Pi_{L}^{\perp}, \Lambda_{\tau}\right)+(1-\sigma) \psi\left(\Pi_{L}^{\perp},-\Lambda_{\tau}\right)\right)
$$

Proof. We consider the functional on $C^{0}(F)$ that assigns to $\psi$ the number

$$
\int_{P_{1 / 4} \times F} \psi(A, B) d \mu(z, A, B)
$$

Clearly this is represented by a Radon measure on $F$ with total mass $\pi \theta / 4$. Thus it suffices to show that
(i) $\theta \in \mathbb{N}$ and
(ii) for any $\psi \in C^{0}(F)$ with $\psi \geqslant 0$ and $\psi\left(\Pi_{L}^{\perp}, \Lambda_{\tau}\right)=\psi\left(\Pi_{L}^{\perp},-\Lambda_{\tau}\right)=0$, we have

$$
\int_{P_{1 / 4 \times F}} \psi(A, B) d \mu(z, A, B)=0 .
$$

For the proofs of both statements we use ideas of Lin [13] and of Lin and Rivière [16].
We first exploit the fact that the Landau-Lifshitz equation can be represented in the form (3) or (4). We note that

$$
\left|\nabla u_{k}\right|^{2} u_{k}=\sum_{\alpha=1}^{3} \frac{\partial u_{k}}{\partial x^{\alpha}} \wedge\left(u_{k} \wedge \frac{\partial u_{k}}{\partial x^{\alpha}}\right)
$$

and

$$
\sqrt{a^{2}+b^{2}}\left|\operatorname{div}\left(u_{k} \wedge \nabla u_{k}\right)\right|=\left|\frac{\partial u_{k}}{\partial t}\right| .
$$

With the compensated compactness method of Coifman, Lions, Meyer, and Semmes [4], we can prove the estimate

$$
\left\|\left|\nabla u_{k}(\cdot, t)\right|^{2} u_{k}(\cdot, t)\right\|_{\mathcal{H}^{1}\left(B_{3 / 4}\right)} \leqslant C_{1}\left\|\nabla u_{k}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}^{2}+C_{1}\left\|\nabla u_{k}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}\left\|\frac{\partial u_{k}}{\partial t}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}
$$

in the Hardy space $\mathcal{H}^{1}\left(B_{3 / 4}\right)$ for every $t \in(-1,1)$, where $C_{1}$ is a constant that depends only on $a$ and $b$. We also have

$$
\sqrt{a^{2}+b^{2}}\left|\Delta u_{k}+\left|\nabla u_{k}\right|^{2} u_{k}\right|=\left|\frac{\partial u_{k}}{\partial t}\right|
$$

thus

$$
\left\|\Delta u_{k}(\cdot, t)\right\|_{\mathcal{H}^{1}\left(B_{3 / 4}\right)} \leqslant C_{2}\left\|\nabla u_{k}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}^{2}+C_{2}\left(\left\|\nabla u_{k}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}+1\right)\left\|\frac{\partial u_{k}}{\partial t}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}
$$

for another constant $C_{2}$ that depends only on $a$ and $b$. Standard estimates for singular integrals involving Hardy spaces (see, e.g., Theorem 3 in Section III.3.1 of Stein [25]) now imply that there exists a constant $C_{3}=C_{3}(a, b)$, such that

$$
\left\|\nabla^{2} u_{k}(\cdot, t)\right\|_{L^{1}\left(B_{1 / 2}\right)} \leqslant C_{3}\left(\left\|\nabla u_{k}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}+1\right)\left(\left\|\nabla u_{k}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}+\left\|\frac{\partial u_{k}}{\partial t}(\cdot, t)\right\|_{L^{2}\left(B_{1}\right)}\right)
$$

In particular we have

$$
\begin{equation*}
\underset{k \rightarrow \infty}{\limsup }\left\|\nabla^{2} u_{k}\right\|_{L^{1}\left(P_{1 / 2}\right)}<\infty \tag{32}
\end{equation*}
$$

We now assume that $L=\{(0,0)\} \times \mathbb{R}$ for simplicity. Because of (30) and (31), the monotonicity formula of Lemma 3.2 (applied to $u_{k}(\cdot, t)$ for every $k \in \mathbb{N}$ and $t \in(-1,1)$ ) implies

$$
\lim _{k \rightarrow \infty} \int_{-1}^{1} \int_{B_{1} \backslash B_{r}} \frac{\left|x \cdot \nabla u_{k}\right|^{2}}{|x|^{3}} d x d t=0
$$

for every $r>0$. Since $\|\mu\|$ is supported on $\{(0,0)\} \times \mathbb{R}^{2}$, this means that

$$
\lim _{k \rightarrow \infty} \int_{P_{1}}\left|\frac{\partial u_{k}}{\partial x^{3}}\right|^{2} d z=0
$$

For $r>0$, let $P_{r}^{2}=(-r, r) \times\left(-r^{2}, r^{2}\right)$. On $P_{1 / 2}^{2}$, define the functions

$$
F_{k}(s, t)=\int_{B_{1 / 2}^{2} \times\{(s, t)\}}\left|\frac{\partial u_{k}}{\partial x^{3}}\right|^{2} d x^{\prime}
$$

and

$$
G_{k}(s, t)=\int_{B_{1 / 2}^{2} \times\{(s, t)\}}\left|\frac{\partial u_{k}}{\partial t}\right|^{2} d x^{\prime}
$$

Moreover, for $\left(s_{0}, t_{0}\right) \in P_{1 / 4}^{2}$, let

$$
F_{k}^{*}\left(s_{0}, t_{0}\right)=\sup _{0<r \leqslant 1 / 4}\left(r^{-3} \int_{t_{0}-r^{2}}^{t_{0}+r^{2}} \int_{s_{0}-r}^{s_{0}+r} F_{k} d s d t+r^{-1} \int_{s_{0}-r}^{s_{0}+r} F_{k}\left(s, t_{0}\right) d s\right)
$$

and

$$
G_{k}^{*}\left(s_{0}, t_{0}\right)=\sup _{0<r \leqslant 1 / 4}\left(r^{-1} \int_{t_{0}-r^{2}}^{t_{0}+r^{2}} \int_{s_{0}-r}^{s_{0}+r} G_{k} d s d t+r \int_{s_{0}-r}^{s_{0}+r} G_{k}\left(s, t_{0}\right) d s\right)
$$

We have

$$
\lim _{k \rightarrow \infty}\left\|F_{k}\right\|_{L^{1}\left(P_{1 / 2}^{2}\right)}=\lim _{k \rightarrow \infty}\left\|G_{k}\right\|_{L^{1}\left(P_{1 / 2}^{2}\right)}=0
$$

hence for any $c>0$,

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \mathcal{L}^{2}\left(\left\{(s, t) \in P_{1 / 4}^{2}: F_{k}^{*}(s, t) \geqslant c\right\}\right)=0 \tag{33}
\end{equation*}
$$

and

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \mathcal{L}^{2}\left(\left\{(s, t) \in P_{1 / 4}^{2}: G_{k}^{*}(s, t) \geqslant c\right\}\right)=0 \tag{34}
\end{equation*}
$$

Using also (32), we can find a sequence of points $\left(s_{k}, t_{k}\right) \in P_{1 / 4}^{2}$, such that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} F^{*}\left(s_{k}, t_{k}\right)=\lim _{k \rightarrow \infty} G^{*}\left(s_{k}, t_{k}\right)=0 \tag{35}
\end{equation*}
$$

and

$$
\begin{equation*}
\limsup _{k \rightarrow \infty}\left\|\nabla^{2} u_{k}\left(\cdot, s_{k}, t_{k}\right)\right\|_{L^{1}\left(B_{1 / 4}^{2}\right)}<\infty \tag{36}
\end{equation*}
$$

Define

$$
v_{k}\left(x^{\prime}, x^{3}, t\right)=u_{k}\left(\frac{x^{\prime}}{4}, \frac{x^{3}}{4}+s_{k}, \frac{t}{16}+t_{k}\right)
$$

These are again solutions of (1). Because of (35), there exists a sequence $\epsilon_{k} \rightarrow 0$ such that

$$
\begin{equation*}
\int_{B_{1}^{2} \times P_{r}^{2}}\left(r^{-3}\left|\frac{\partial v_{k}}{\partial x^{3}}\right|^{2}+r^{-1}\left|\frac{\partial v_{k}}{\partial t}\right|^{2}\right) d z+\int_{B_{1}^{2} \times(-r, r)}\left(r^{-1}\left|\frac{\partial v_{k}}{\partial x^{3}}(x, 0)\right|^{2}+r\left|\frac{\partial v_{k}}{\partial t}(x, 0)\right|^{2}\right) d x \leqslant \epsilon_{k} \tag{37}
\end{equation*}
$$

for every $k \in \mathbb{N}$ and every $r \in(0,1]$. Moreover, for any $\eta \in C_{0}^{0}\left(P_{1}\right)$, we have

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \int_{P_{1}} \eta\left|\nabla v_{k}\right|^{2} d z=8 \pi \theta \int_{-1}^{1} \int_{-1}^{1} \eta(0,0, s, t) d s d t \tag{38}
\end{equation*}
$$

by (30). Using part (ii) of Lemma 3.3, we see that

$$
\lim _{k \rightarrow \infty} \int_{B_{1}}\left|\nabla v_{k}(x, 0)\right|^{2} d x=16 \pi \theta
$$

Since we have (37), Lemma 3.2 then implies

$$
\begin{equation*}
\frac{1}{r} \limsup _{k \rightarrow \infty} \int_{B_{r}\left(x_{0}\right)}\left|\nabla v_{k}(x, 0)\right|^{2} d x \leqslant 32 \pi \theta \tag{39}
\end{equation*}
$$

uniformly in $x_{0} \in B_{1 / 2}$ and $r \in\left(0, \frac{1}{2}\right]$.
Choose two cut-off functions $\xi \in C_{0}^{\infty}\left(-\frac{1}{2}, \frac{1}{2}\right)$ and $\zeta \in C_{0}^{\infty}(-1,1)$ such that

$$
\int_{-1 / 2}^{1 / 2} \xi(s) d s=\int_{-1}^{1} \zeta(t) d t=1
$$

For $r \in\left(0, \frac{1}{2}\right]$, set

$$
\xi_{r}(s)=r^{-1} \xi(s / r), \quad \zeta_{r}(t)=r^{-2} \zeta\left(t / r^{2}\right) .
$$

Fix also a function $\eta \in C_{0}^{\infty}\left(B_{1 / 2}^{2}\right)$, and for $x_{0}^{\prime} \in B_{1 / 2}$ and $r \in\left(0, \frac{1}{2}\right]$ set

$$
\eta_{x_{0}^{\prime}, r}\left(x^{\prime}\right)=\eta\left(\frac{x^{\prime}-x_{0}^{\prime}}{r}\right) .
$$

Then (37), (39), and Lemma 3.3 imply

$$
\begin{equation*}
\left.\left|\int_{P_{1}} \eta_{x_{0}^{\prime}, r}\left(x^{\prime}\right) \xi_{r}\left(x^{3}\right) \zeta_{r}(t)\right| \nabla v_{k}\left(x^{\prime}, x^{3}, t\right)\right|^{2} d z-\int_{B_{1}^{2}} \eta_{x_{0}^{\prime}, r}\left(x^{\prime}\right)\left|\nabla v_{k}\left(x^{\prime}, 0,0\right)\right|^{2} d x^{\prime} \mid \rightarrow 0 \quad \text { as } k \rightarrow \infty \tag{40}
\end{equation*}
$$

uniformly in $x_{0}^{\prime}$ and $r$.
Define now

$$
w_{k}\left(x^{\prime}\right)=v_{k}\left(x^{\prime}, 0,0\right) .
$$

From (38) and (40), we obtain

$$
\lim _{k \rightarrow \infty} \int_{B_{1}^{2}} \eta\left|\nabla w_{k}\right|^{2} d x^{\prime}=8 \pi \theta \eta(0)
$$

for every $\eta \in C_{0}^{0}\left(B_{1}^{2}\right)$. That is, the sequence $\left\{w_{k}\right\}$ satisfies the hypothesis (I) from Section 4. Hypothesis (II) follows from (36). Combining (40) with Lemma 3.1, we find that (IV) is true. Using also (37) and the regularity results of Hélein [9], we easily show (III). Thus we can apply Proposition 4.1, and we immediately obtain (i).

To show (ii), we argue by contradiction. Suppose there exists a function $\psi \in C^{0}(F)$ with $\psi \geqslant 0$ and $\psi\left(\Pi_{L}^{\perp}, \Lambda_{\tau}\right)=$ $\psi\left(\Pi_{L}^{\perp},-\Lambda_{\tau}\right)=0$, such that

$$
\begin{equation*}
\int_{P_{1 / 4 \times F}} \psi(A, B) d \mu(z, A, B)>0 . \tag{41}
\end{equation*}
$$

For $(s, t) \in P_{1 / 4}^{2}$, define

$$
f_{k}(s, t)=\frac{1}{2} \int_{B_{1 / 4}^{2}} \psi\left(A_{u_{k}}\left(x^{\prime}, s, t\right), B_{u_{k}}\left(x^{\prime}, s, t\right)\right)\left|\nabla u_{k}\left(x^{\prime}, s, t\right)\right|^{2} d x^{\prime}
$$

Then we have

$$
\begin{equation*}
\liminf _{k \rightarrow \infty} \int_{-1 / 16}^{1 / 16} \int_{-1 / 4}^{1 / 4} f_{k}(s, t) d s d t>0 \tag{42}
\end{equation*}
$$

For $c>0$, define

$$
Z_{k}^{c}=\left\{(s, t) \in P_{1 / 4}^{2}: f_{k}(s, t) \geqslant c\right\} .
$$

We claim that there exists a number $c>0$ such that

$$
\begin{equation*}
\liminf _{k \rightarrow \infty} \mathcal{L}^{2}\left(Z_{k}^{c}\right)>0 \tag{43}
\end{equation*}
$$

To see this, suppose $c>0$ is such that

$$
\liminf _{k \rightarrow \infty} \mathcal{L}^{2}\left(Z_{k}^{c}\right)=0
$$

We may assume that $\mathcal{L}^{2}\left(Z_{k}^{c}\right) \leqslant 2^{-k}$ for every $k \in \mathbb{N}$; otherwise we replace our sequence by a subsequence. Define

$$
Y_{\ell}^{c}=\bigcup_{\ell=k}^{\infty} Z_{k}^{c} .
$$

Then $\mathcal{L}^{2}\left(Y_{\ell}^{c}\right) \leqslant 2^{1-\ell}$, and $f_{k}<c$ outside of $Y_{\ell}^{c}$ for any $k \geqslant \ell$. Thus

$$
\begin{aligned}
\liminf _{k \rightarrow \infty} \int_{-1 / 16}^{1 / 16} \int_{-1 / 4}^{1 / 4} f_{k}(s, t) d s d t & \leqslant \frac{c}{16}+\frac{1}{2}\|\psi\|_{C^{0}(F)} \liminf _{k \rightarrow \infty} \int_{B_{1 / 4}^{2} \times Y_{\ell}^{c}}\left|\nabla u_{k}\right|^{2} d z \\
& \leqslant \frac{c}{16}+2^{2-\ell} \pi \theta\|\psi\|_{C^{0}(F)}
\end{aligned}
$$

for any $\ell \in \mathbb{N}$. By (42), this is only possible for sufficiently large values of $c$.
Now because of (32)-(34) and (43), we can find a sequence of points $\left(s_{k}, t_{k}\right) \in P_{1 / 4}^{2}$ such that (35) and (36) hold true, and in addition,

$$
\begin{equation*}
\liminf _{k \rightarrow \infty} f_{k}\left(s_{k}, t_{k}\right)>0 . \tag{44}
\end{equation*}
$$

Similarly as in the first part of the proof, we see that a subsequence of the maps

$$
w_{k}\left(x^{\prime}\right)=u_{k}\left(\frac{x^{\prime}}{4}, s_{k}, t_{k}\right)
$$

satisfies the conditions (I)-(IV) in section 4. But then Proposition 4.1 contradicts (44). This concludes the proof.

## 6. Proof of Theorem 1.1

Suppose now that $u_{k} \in C^{\infty}\left(\Omega \times[0, T), \mathbb{S}^{2}\right)$ are solutions of the Landau-Lifshitz equation (1) that satisfy the hypotheses of Theorem 1.1. Then it follows from the energy identity (8) that for every precompact set $K \Subset \Omega$ and every $t_{0} \in(0, T)$, we have

$$
\begin{equation*}
\sup _{k \in \mathbb{N}}\left(\int_{K}\left|\nabla u_{k}\left(x, t_{0}\right)\right|^{2} d x+\int_{t_{0}}^{T} \int_{K}^{T}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x d t\right)<\infty . \tag{45}
\end{equation*}
$$

We consider the measures $\mu_{k}$ on $\Omega_{T} \times F$ given by

$$
\int_{\Omega_{T} \times F} \psi d \mu_{k}=\frac{1}{2} \int_{\Omega_{T}} \psi\left(z, A_{u_{k}}, B_{u_{k}}\right)\left|\nabla u_{k}\right|^{2} d z \quad \text { for } \psi \in C_{0}^{0}\left(\Omega_{T} \times F\right) .
$$

Because of (45), we may assume that there exist a function $u \in H_{\mathrm{loc}}^{1}\left(\Omega_{T}, \mathbb{S}^{2}\right)$ and a Radon measure $\mu$ on $\Omega_{T} \times F$ such that $u_{k} \rightharpoonup u$ weakly in $H_{\mathrm{loc}}^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right)$ and pointwise almost everywhere, and $\mu_{k} \rightarrow \mu$ (possibly after the choice of a subsequence). Using the representation (4) of the Landau-Lifshitz equation and passing to the limit, we see immediately that $u$ is a weak solution.

It also follows from (45) that there exist Radon measures $\mu^{t}$ on $\Omega \times F$ for almost every $t \in[0, T)$, such that

$$
\int_{\Omega_{T} \times F} \psi d \mu=\int_{0}^{T} \int_{\Omega \times F} \psi d \mu^{t} d t
$$

for every $\psi \in C_{0}^{0}\left(\Omega_{T} \times F\right)$. We define $m^{t}=\left\|\mu^{t}\right\|$. We also consider the measures

$$
m_{k}^{t}=\frac{1}{2} \mathcal{L}^{3}\left\llcorner\left|\nabla u_{k}(\cdot, t)\right|^{2} .\right.
$$

Fix a function $\xi \in C_{0}^{1}(\Omega)$ and define

$$
f_{k \xi}(t)=\int_{\Omega} \xi d m_{k}^{t}
$$

and

$$
f_{\xi}(t)=\int_{\Omega} \xi d m^{t}
$$

For any $\zeta \in C_{0}^{0}(0, T)$, we have

$$
\lim _{k \rightarrow \infty} \int_{0}^{T} f_{k \xi}(t) \zeta(t) d t=\int_{0}^{T} f_{\xi}(t) \zeta(t) d t
$$

By (8), we have

$$
\frac{d}{d t} f_{k \xi}(t)=-\int_{\Omega \times\{t\}}\left(a \xi\left|\frac{\partial u_{k}}{\partial t}\right|^{2}+a \nabla \xi \cdot\left\langle\nabla u_{k}, \Delta u_{k}\right\rangle-b \nabla \xi \cdot\left\langle\nabla u_{k}, u_{k} \wedge \Delta u_{k}\right\rangle\right) d x
$$

Using (45), we find that

$$
\sup _{k \in \mathbb{N}} \int_{t_{0}}^{T}\left|\frac{d}{d t} f_{k \xi}(t)\right| d t<\infty
$$

for every $t_{0}>0$. That is, the sequence $\left\{f_{k \xi}\right\}$ is bounded in $\mathrm{BV}\left(t_{0}, T\right)$. Hence we have $f_{k \xi}(t) \rightarrow f_{\xi}(t)$ for almost every $t \in[0, T)$. Since $C_{0}^{1}(\Omega)$ is dense in $C_{0}^{0}(\Omega)$, this means that

$$
\begin{equation*}
m_{k}^{t} \rightarrow m^{t} \tag{46}
\end{equation*}
$$

for almost every $t \in[0, T)$.
For almost every $t \in[0, T)$, we also have $u_{k}(x, t) \rightarrow u(x, t)$ at almost every $x \in \Omega$. By (45) again, the sequence $\left\{u_{k}(\cdot, t)\right\}$ is bounded in $H^{1}\left(K, \mathbb{R}^{3}\right)$ for every $K \Subset \Omega$. We conclude that

$$
\begin{equation*}
u_{k}(\cdot, t) \rightharpoonup u(\cdot, t) \quad \text { weakly in } H_{\mathrm{loc}}^{1}\left(\Omega, \mathbb{R}^{3}\right) \tag{47}
\end{equation*}
$$

for almost every $t \in[0, T)$. Using (45) and Fatou's lemma, we also find that

$$
\begin{equation*}
\liminf _{k \rightarrow \infty} \int_{\Omega}\left|\frac{\partial u_{k}}{\partial t}(x, t)\right|^{2} d x<\infty \tag{48}
\end{equation*}
$$

for almost every $t \in[0, T)$.
Fix $t \in[0, T)$, such that (46), (47), and (48) hold. Then Theorem 3.1 can be applied locally to a subsequence of $\left\{u_{k}(\cdot, t)\right\}$. Hence there exists a locally $\mathscr{H}^{1}$-integrable function $\theta_{t}: \Omega \rightarrow[0, \infty)$ such that $\Sigma_{t}=\theta_{t}^{-1}((0, \infty))$ is closed and countably 1 -rectifiable, and

$$
\begin{equation*}
m^{t}=\frac{1}{2} \mathcal{L}^{3}\left\llcorner|\nabla u(\cdot, t)|^{2}+4 \pi \mathscr{H}^{1}\left\llcorner\theta_{t} .\right.\right. \tag{49}
\end{equation*}
$$

Moreover, we have

$$
\int_{\Omega}|u(x, t) \wedge \Delta u(x, t)|^{2} d x \leqslant \liminf _{k \rightarrow \infty} \int_{\Omega}\left|u_{k}(x, t) \wedge \Delta u_{k}(x, t)\right|^{2} d x .
$$

Fatou's lemma then implies that $\Delta u+|\nabla u|^{2} u \in L_{\text {loc }}^{2}\left(\Omega_{T}, \mathbb{R}^{3}\right)$.
Consider the measures

$$
\lambda_{k}=\mathcal{L}^{4}\left\llcorner\left|\frac{\partial u_{k}}{\partial t}\right|^{2}\right.
$$

on $\Omega_{T}$. We may assume that there exists a Radon measure $\lambda$ on $\Omega_{T}$ such that $\lambda_{k} \rightarrow \lambda$. Define

$$
S=\left\{z_{0} \in \Omega_{T}: \liminf _{r \searrow 0} r^{-1} \lambda\left(P_{r}\left(z_{0}\right)\right)>0\right\} .
$$

We claim that the parabolic Hausdorff dimension of $S$ is at most 1 . To see this, define

$$
S_{i j}=\left\{z_{0} \in \Omega_{T}: r^{-1} \lambda\left(P_{r}\left(z_{0}\right)\right) \geqslant i^{-1} \text { for every } r \in\left(0, j^{-1}\right)\right\} .
$$

These are closed sets, and therefore a standard covering argument involving Vitali's covering lemma shows that $S_{i j}$ has locally finite $\mathscr{H}_{d}^{1}$-measure. Consequently, the set

$$
S=\bigcup_{i=1}^{\infty} \bigcup_{j=1}^{\infty} S_{i j}
$$

satisfies $\mathscr{H}_{d}^{\gamma}(S)=0$ for every $\gamma>1$. Hence for almost every $t_{0} \in[0, T)$, we have

$$
\begin{equation*}
S \cap\left(\Omega \times\left\{t_{0}\right\}\right)=\emptyset \tag{50}
\end{equation*}
$$

Fix a $t_{0}$ such that (49) and (50) hold. Choose a point $x_{0} \in \Sigma_{t_{0}}$ such that $\theta_{t_{0}}$ is approximately continuous at $x_{0}$ and the approximate tangent space $T_{x_{0}} \Sigma_{t_{0}}$ exists. (This is true $m^{t_{0}}$-almost everywhere on $\Sigma_{t_{0}}$.) Then there exists a sequence $r_{k} \searrow 0$ such that

$$
\lim _{k \rightarrow \infty} \frac{1}{r_{k}} \int_{P_{r_{k}}\left(x_{0}, t_{0}\right)}\left|\frac{\partial u_{k}}{\partial t}\right|^{2} d z=0
$$

and such that the rescaled maps

$$
v_{k}(x, t)=u_{k}\left(r_{k} x+x_{0}, r_{k}^{2} t+t_{0}\right)
$$

satisfy

$$
\lim _{k \rightarrow \infty} \int_{B_{1}} \xi(x)\left|\nabla v_{k}(x, 0)\right|^{2} d x=8 \pi \theta_{t_{0}}\left(x_{0}\right) \int_{T_{x_{0}} \Sigma_{t_{0}}} \xi d \mathscr{H}^{1}
$$

for every $\xi \in C_{0}^{1}\left(B_{2}\right)$. With the help of part (iii) of Lemma 3.3, we conclude that

$$
\lim _{k \rightarrow \infty} \int_{P_{1}} \eta\left|\nabla v_{k}\right|^{2} d z=8 \pi \theta_{t_{0}}\left(x_{0}\right) \int_{-1}^{1} \int_{T_{x_{0}} \Sigma_{t_{0}}} \eta d \mathscr{H}^{1} d t
$$

for every $\eta \in C_{0}^{0}\left(P_{1}\right)$. That is, a subsequence of $\left\{v_{k}\right\}$ satisfies the hypotheses of Proposition 5.1. Therefore, the measures $\tilde{\mu}_{k}$ on $P_{1} \times F$ that belong to $v_{k}$ subconverge to a measure $\tilde{\mu}$ on $P_{1} \times F$ which satisfies

$$
\int_{P_{1 / 4}} \psi(A, B) d \tilde{\mu}(z, A, B)=\frac{\pi \tilde{m}}{4}\left(\tilde{\sigma} \psi\left(\Pi_{\tilde{L}}^{\perp}, \Lambda_{\tilde{\tau}}\right)+(1-\tilde{\sigma}) \psi\left(\Pi_{\tilde{L}}^{\perp},-\Lambda_{\tilde{\tau}}\right)\right)
$$

for two numbers $\tilde{m} \in \mathbb{N}$ and $\tilde{\sigma} \in[0,1]$, where $\tilde{L}=T_{x_{0}} \Sigma_{t_{0}}$ and $\tilde{\tau} \in \tilde{L} \cap \mathbb{S}^{2}$. By the construction of $\tilde{\mu}$, this formula determines $\theta_{t_{0}}\left(x_{0}\right)$ as well as the fibre measure $\mu^{\left(x_{0}, t_{0}\right)}$ of $\mu$ at the point $\left(x_{0}, t_{0}\right)$. Namely, we have

$$
\theta_{t_{0}}\left(x_{0}\right)=\tilde{m}
$$

and

$$
\mu^{\left(x_{0}, t_{0}\right)}=\delta_{\Pi_{\tilde{L}}} \times\left(\tilde{\sigma} \delta_{\Lambda_{\tilde{\tau}}}+(1-\tilde{\sigma}) \delta_{-\Lambda_{\tilde{\tau}}}\right) .
$$

These arguments work for almost every $t_{0} \in[0, T)$ and $\mathscr{H}^{1}$-almost every $x_{0} \in \Sigma_{t_{0}}$. Hence there exists a function $\theta: \Omega_{T} \rightarrow \mathbb{N} \cup\{0\}$, such that $\theta(x, t)=\theta_{t}(x)$ for $\|\mu\|$-almost every point $(x, t) \in \Omega_{T}$, and thus

$$
\int_{\Omega_{T}} \eta d\|\mu\|=\frac{1}{2} \int_{\Omega_{T}} \eta|\nabla u|^{2} d x+4 \pi \int_{0}^{T} \int_{\Sigma_{t}} \eta \theta d \mathscr{H}^{1} d t
$$

for every $\eta \in C_{0}^{0}\left(\Omega_{T}\right)$. Moreover, there exist two functions $\sigma: \Omega_{T} \rightarrow[0,1]$ and $\tau: \Omega_{T} \rightarrow \mathbb{S}^{2}$ with $\tau(x, t) \in T_{x} \Sigma_{t}$ for almost every $t$ and $\mathscr{H}^{1}$-almost every $x \in \Sigma_{t}$, such that for every $\psi \in C_{0}^{0}\left(\Omega_{T} \times F\right)$, we have

$$
\begin{align*}
\int_{\Omega_{T} \times F} \psi d \mu= & \frac{1}{2} \int_{\Omega_{T}} \psi\left(z, A_{u}, B_{u}\right)|\nabla u|^{2} d z+4 \pi \int_{0}^{T} \int_{\Omega}\left[\sigma \psi\left(x, t, \Pi_{T_{x} \Sigma_{t}}, \Lambda_{\tau}\right)\right. \\
& \left.+(1-\sigma) \psi\left(x, t, \Pi_{T_{x} \Sigma_{t}},-\Lambda_{\tau}\right)\right] \theta d \mathscr{H}^{1} d t \tag{51}
\end{align*}
$$

So far we have proved parts (i) and (ii) of Theorem 1.1. For part (iii), it suffices to use Theorem 3.1 again. Only part (iv) remains to be proved.

Let $\epsilon_{0}$ be the constant from Lemma 3.1. Define the set

$$
\Sigma=\left\{z_{0} \in \Omega_{T}: \liminf _{r \searrow 0} r^{-3}\|\mu\|\left(P_{r}\left(z_{0}\right)\right) \geqslant \epsilon_{0}\right\}
$$

Using Lemma 3.1, we see that for every point $z_{0} \in \Omega_{T} \backslash \Sigma$, there exists a radius $r>0$ such that a subsequence of $\left\{u_{k}\right\}$ converges to $u$ in $C^{1}\left(P_{r}\left(z_{0}\right), S^{2}\right)$. We conclude that $\Sigma$ is closed and $\Sigma_{t} \times\{t\} \subset \Sigma$ for almost every $t \in[0, T)$. With a covering argument, we see that $\mathscr{H}_{d}^{3}(\Sigma \cap K)<\infty$ for every compact set $K \Subset \Omega_{T}$. Thus in particular $\mathcal{L}^{4}(\Sigma)=0$.

For $\delta>0$, define

$$
U_{\delta}=\left\{z \in \Omega_{T}: \operatorname{dist}(z, \Sigma)<\delta\right\},
$$

where dist is the distance function with respect to the parabolic metric $d$. Then we have

$$
\lim _{\delta \searrow 0} \int_{U_{\delta}} \eta|\nabla u|^{2} d z=0
$$

for every $\eta \in C_{0}^{0}\left(\Omega_{T}\right)$. Hence

$$
\int_{0}^{T} \int_{\Sigma_{t}} \eta \theta d \mathscr{H}^{1} d t=\lim _{\delta \rightarrow 0} \int_{U_{\delta}} \eta d\|\mu\| .
$$

We finally consider the vector fields $H_{k}: \Omega_{T} \rightarrow \mathbb{R}^{3}$ that satisfy $H_{k}(x) \perp$ ker $A_{u_{k}}(x)$ for $\left\|\mu_{k}\right\|$-almost every $x \in \Omega_{T}$ and

$$
\int_{\Omega_{T} \times F} \phi \cdot A H_{k} d \mu_{k}=-\delta \mu_{k}(\phi)
$$

for every $\phi \in C_{0}^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right)$. Also consider the function

$$
g(z, A, B)=|A|^{2}-2
$$

on $\Omega_{T} \times F$. Note that $g$ is non-negative by condition (c) in the definition of $F$ in Section 2. Now define

$$
V_{\delta k}=\left\{z \in U_{\delta}: g\left(z, A_{u_{k}}(z), B_{u_{k}}(z)\right) \leqslant 1\right\} .
$$

Then for every $z \in V_{\delta k}$, we have either $\nabla u_{k}=0$ or the rank of $\nabla u_{k}$ is 2 . By (22), the identity

$$
\int_{V_{\delta k}} \phi \cdot B H_{k} d \mu_{k}=-\int_{0}^{T} \int_{\Omega} \phi \cdot\left\langle u_{k} \wedge \Delta u_{k}, \nabla u_{k}\right\rangle d x d t
$$

holds for every $\phi \in C_{0}^{0}\left(\Omega_{T}, \mathbb{R}^{3}\right)$. Integrating both sides of (8), we thus obtain

$$
\begin{aligned}
0= & \int_{V_{\delta k}}\left(\frac{\partial \eta}{\partial t}-a \eta H_{k} \cdot A H_{k}+a \nabla \eta \cdot A H_{k}-b \nabla \eta \cdot B H_{k}\right) d \mu+\int_{\Omega \backslash V_{\delta k}}\left(\frac{1}{2}\left|\nabla u_{k}\right|^{2} \frac{\partial \eta}{\partial t}-a \eta\left|u_{k} \wedge \Delta u_{k}\right|^{2}\right. \\
& \left.-a \nabla \eta \cdot\left\langle\Delta u_{k}, \nabla u_{k}\right\rangle+b \nabla \eta \cdot\left\langle u_{k} \wedge \Delta u_{k}, \nabla u_{k}\right\rangle\right) d x d t
\end{aligned}
$$

for every $\eta \in C_{0}^{1}\left(\Omega_{T}\right)$.
Note that

$$
\int_{\Sigma \times F} g d \mu=0
$$

by (51). Hence

$$
\lim _{\delta \searrow 0} \lim _{k \rightarrow \infty} \int_{U_{\delta} \times F} g d \mu_{k}=0
$$

and we conclude that

$$
\begin{equation*}
\lim _{\delta \searrow 0} \lim _{k \rightarrow \infty}\left\|\mu_{k}\right\|\left(U_{\delta} \backslash V_{\delta k}\right)=0 . \tag{52}
\end{equation*}
$$

Let $\chi_{\delta k}$ be the characteristic function of $V_{\delta k}$, and define the functions

$$
\tilde{H}_{\delta k}(z, A, B)=\chi_{\delta k}(z) H_{k}(z)
$$

on $\Omega_{T} \times F$. Observe that

$$
\left|\tilde{H}_{\delta k}(z, A, B)\right|^{2} \leqslant C_{1}\left|u_{k}(z) \wedge \Delta u_{k}(z)\right|^{2}
$$

for a universal constant $C_{1}$ at every point $(z, A, B) \in \Omega_{T} \times F$ such that $\nabla u_{k}(z) \neq 0$. Thus for any $K \Subset \Omega$ and $t_{0} \in(0, T)$, we have

$$
\limsup _{k \rightarrow \infty} \int_{K \times\left(t_{0}, T\right)}\left|\tilde{H}_{\delta k}\right|^{2} d \mu_{k}<\infty
$$

for every $\delta>0$. By Proposition 2.1, we may assume that the measure-function pairs ( $\mu_{k}, \tilde{H}_{\delta k}$ ) converge weakly to a pair ( $\mu, \tilde{H}_{\delta}$ ). Clearly $\tilde{H}_{\delta^{\prime}}=\tilde{H}_{\delta}$ in $U_{\delta^{\prime}} \times F$ whenever $\delta^{\prime} \leqslant \delta$. Let now

$$
\tilde{H}(z, A, B)= \begin{cases}\tilde{H}_{1}(z, A, B) & \text { if } z \in \Sigma, \\ 0 & \text { else }\end{cases}
$$

This is the pointwise limit of $\tilde{H}_{\delta}$ as $\delta \searrow 0$. We claim that

$$
\begin{equation*}
\delta \mu(\phi)=\int_{0}^{T} \int_{\Omega} \phi \cdot\langle\Delta u, \nabla u\rangle d x d t-\int_{\Omega_{T} \times F} \phi \cdot A \tilde{H} d \mu \tag{53}
\end{equation*}
$$

for every $\phi \in C_{0}^{1}\left(\Omega_{T}, \mathbb{R}^{3}\right)$. To verify this, we calculate

$$
\delta \mu_{k}(\phi)=\int_{\Omega_{T} \backslash U_{\delta}} \phi \cdot\left\langle\Delta u_{k}, \nabla u_{k}\right\rangle d z+\int_{U_{\delta} \backslash V_{\delta k}} \phi \cdot\left\langle\Delta u_{k}, \nabla u_{k}\right\rangle d z-\int_{V_{\delta k}} \phi \cdot A \tilde{H}_{\delta k} d \mu
$$

We have

$$
\lim _{\delta \searrow 0} \lim _{k \rightarrow \infty} \int_{\Omega_{T} \backslash U_{\delta}} \phi \cdot\left\langle\Delta u_{k}, \nabla u_{k}\right\rangle d z=\int_{0}^{T} \int_{\Omega} \phi \cdot\langle\Delta u, \nabla u\rangle d x d t
$$

and

$$
\lim _{\delta \searrow 0} \lim _{k \rightarrow \infty} \int_{V_{\delta k}} \phi \cdot A \tilde{H}_{\delta k} d \mu=\lim _{\delta \searrow 0} \int_{U_{\delta}} \phi \cdot A \tilde{H}_{\delta} d \mu=\int_{\Sigma} \phi \cdot A \tilde{H} d \mu .
$$

Here we use Lebesgue's convergence theorem in the last step. Finally,

$$
\lim _{\delta \searrow 0} \lim _{k \rightarrow \infty}\left|\int_{U_{\delta} \backslash V_{\delta k}} \phi \cdot\left\langle\Delta u_{k}, \nabla u_{k}\right\rangle d z\right| \leqslant \lim _{\delta \searrow 0} \lim _{k \rightarrow \infty} \sup _{\Omega_{T}}|\phi|\left(\left\|\mu_{k}\right\|\left(U_{\delta} \backslash V_{\delta k}\right) \int_{\operatorname{supp} \phi}\left|u_{k} \wedge \Delta u_{k}\right|^{2} d z\right)^{1 / 2}=0
$$

by Hölder's inequality and (52) and (45). This shows (53). With a similar reasoning, we find that

$$
\int_{0}^{T} \int_{\Omega} \phi \cdot\langle u \wedge \Delta u, \nabla u\rangle d x d t-\int_{\Omega_{T} \times F} \phi \cdot B \tilde{H} d \mu=\lim _{k \rightarrow \infty} \int_{0}^{T} \int_{\Omega} \phi \cdot\left\langle u_{k} \wedge \Delta u_{k}, \nabla u_{k}\right\rangle d x d t
$$

for every $\phi \in C_{0}^{0}\left(\Omega_{T}, \mathbb{R}^{3}\right)$. With the help of Proposition 2.1, we also obtain the inequality

$$
\int_{0}^{T} \int_{\Omega} \eta|u \wedge \Delta u|^{2} d x d t+\int_{\Omega_{T} \times F} \eta|\tilde{H}|^{2} d \mu \leqslant \liminf _{k \rightarrow \infty} \int_{0}^{T} \int_{\Omega} \eta\left|u_{k} \wedge \Delta u_{k}\right|^{2} d x d t
$$

for every $\eta \in C_{0}^{1}\left(\Omega_{T}\right)$ with $\eta \geqslant 0$. Passing to the limit in (8), we therefore find

$$
\begin{aligned}
0 \leqslant & \int_{0}^{T} \int_{\Omega}\left(\frac{1}{2}|\nabla u|^{2} \frac{\partial \eta}{\partial t}-a \eta|u \wedge \Delta u|^{2}-a \nabla \eta \cdot\langle\Delta u, \nabla u\rangle+b \nabla \eta \cdot\langle u \wedge \Delta u, \nabla u\rangle\right) d x d t \\
& +\int_{\Omega_{T} \times F}\left(\frac{\partial \eta}{\partial t}-a \eta|\tilde{H}|^{2}+a \nabla \eta \cdot A \tilde{H}-b \nabla \eta \cdot B \tilde{H}\right) d \mu .
\end{aligned}
$$

Now we define

$$
H^{+}(x, t)=\sigma(x, t) \tilde{H}\left(x, t, \Pi_{T_{x} \Sigma_{t}}, \Lambda_{\tau(x, t)}\right)
$$

and

$$
H^{-}(x, t)=(1-\sigma(x, t)) \tilde{H}\left(x, t, \Pi_{T_{x} \Sigma_{t}},-\Lambda_{\tau(x, t)}\right) .
$$

Taking the representation (51) for $\mu$ and identity (53) into account, we obtain (15) as well as the other claims from part (iv) of Theorem 1.1.

Remark. We have mentioned earlier that we do not know whether the vector fields $H^{+}$and $\mathrm{H}^{-}$are parallel to $H$ almost everywhere. In other words, our arguments do not prove that $\tilde{H}$ is a function that depends only on the variable $z$. Note, however, that if we had strong convergence of the measure-function pairs $\left(\mu_{k}, \tilde{H}_{\delta k}\right)$ to ( $\mu, \tilde{H}_{\delta}$ ), then Proposition 2.1 , together with Jensen's inequality (applied to fibre measures of $\mu$ ) would imply that $\tilde{H}$ depends only on the first argument. Thus the statement in part (iv) of Theorem 1.1 could be improved in this case. Although strong convergence in this sense may be too much to expect in general, a further analysis of the higher order energies given by $u_{k} \wedge \Delta u_{k}$ might give a better understanding of the energy concentration.

## References

[1] W.K. Allard, On the first variation of a varifold, Ann. of Math. (2) 95 (1972) 417-491.
[2] L. Ambrosio, H.M. Soner, A measure-theoretic approach to higher codimension mean curvature flows, Ann. Scuola Norm. Sup. Pisa Cl. Sci. (4) 25 (1997) 27-49.
[3] K.A. Brakke, The Motion of a Surface by its Mean Curvature, Mathematical Notes, vol. 20, Princeton University Press, Princeton, NJ, 1978.
[4] R. Coifman, P.L. Lions, Y. Meyer, S. Semmes, Compensated compactness and Hardy spaces, J. Math. Pures Appl. 72 (1993) $247-286$.
[5] W. Ding, G. Tian, Energy identity for a class of approximate harmonic maps from surfaces, Comm. Anal. Geom. 3 (1995) 543-554.
[6] J. Eells, L. Lemaire, A report on harmonic maps, Bull. London Math. Soc. 10 (1978) 1-68.
[7] H. Federer, Geometric Measure Theory, Springer-Verlag, New York, 1969.
[8] M. Feldman, Partial regularity for harmonic maps of evolution into spheres, Comm. Partial Differential Equations 19 (1994) $761-790$.
[9] F. Hélein, Régularité des applications faiblement harmoniques entre une surface et une sphère, C. R. Acad. Sci. Paris Sér. I Math. 311 (1990) 519-524.
[10] J.E. Hutchinson, Second fundamental form for varifolds and the existence of surfaces minimising curvature, Indiana Univ. Math. J. 35 (1986) 45-71.
[11] J. Jost, Two-Dimensional Geometric Variational Problems, John Wiley \& Sons, Chichester, 1991.
[12] J. Li, G. Tian, The blow-up locus of heat flows for harmonic maps, Acta Math. Sin. (Engl. Ser.) 16 (2000) 29-62.
[13] F.-H. Lin, Gradient estimates and blow-up analysis for stationary harmonic maps, Ann. of Math. (2) 149 (1999) $785-829$.
[14] F.-H. Lin, Mapping problems, fundamental groups and defect measures, Acta Math. Sin. (Engl. Ser.) 15 (1999) 25-52.
[15] F.-H. Lin, Varifold type theory for Sobolev mappings, in: First International Congress of Chinese Mathematicians, Beijing, 1998, Amer. Math. Soc., Providence, 2001, pp. 423-430.
[16] F.-H. Lin, T. Rivière, Energy quantization for harmonic maps, Duke Math. J. 111 (2002) 177-193.
[17] F.-H. Lin, C. Wang, Energy identity of harmonic map flows from surfaces at finite singular time, Calc. Var. Partial Differential Equations 6 (1998) 369-380.
[18] F.-H. Lin, C. Wang, Harmonic and quasi-harmonic spheres. III. Rectifiability of the parabolic defect measure and generalized varifold flows, Ann. Inst. H. Poincaré Anal. Non Linéaire 19 (2002) 209-259.
[19] R. Moser, Energy concentration for almost harmonic maps and the Willmore functional, Math. Z. 251 (2005) 293-311.
[20] R. Moser, Partial Regularity for Harmonic Maps and Related Problems, World Scientific Publishing Co. Pte. Ltd, Singapore, 2005.
[21] J. Qing, On singularities of the heat flow for harmonic maps from surfaces into spheres, Comm. Anal. Geom. 3 (1995) $297-315$.
[22] J. Qing, G. Tian, Bubbling of the heat flows for harmonic maps from surfaces, Comm. Pure Appl. Math. 50 (1997) 295-310.
[23] J. Sacks, K. Uhlenbeck, The existence of minimal immersions of 2-spheres, Ann. of Math. (2) 113 (1981) 1-24.
[24] L. Simon, Lectures on Geometric Measure Theory, Australian National University Centre for Mathematical Analysis, Canberra, 1983.
[25] E.M. Stein, Harmonic Analysis: Real-Variable Methods, Orthogonality, and Oscillatory Integrals, Princeton University Press, Princeton, NJ, 1993.
[26] M. Struwe, On the evolution of harmonic maps in higher dimensions, J. Differential Geom. 28 (1988) 485-502.
[27] L. Tartar, Imbedding theorems of Sobolev spaces into Lorentz spaces, Boll. Unione Mat. Ital. Sez. B Artic. Ric. Mat. (8) 1 (1998) $479-500$.
[28] T.J. Willmore, Riemannian Geometry, Oxford Science Publications, The Clarendon Press, Oxford University Press, New York, 1993.


[^0]:    E-mail address: r.moser@maths.bath.ac.uk.

