# REGULARITY OF AREA MINIMIZING CURRENTS I: GRADIENT $L^{p}$ ESTIMATES 

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#### Abstract

In a series of papers, including the present one, we give a new, shorter proof of Almgren's partial regularity theorem for area minimizing currents in a Riemannian manifold, with a slight improvement on the regularity assumption for the latter. This note establishes a new a priori estimate on the excess measure of an area minimizing current, together with several statements concerning approximations with Lipschitz multiple valued graphs. Our new a priori estimate is an higher integrability type result, which has a counterpart in the theory of Dir-minimizing multiple valued functions and plays a key role in estimating the accuracy of the Lipschitz approximations.


## 0. Foreword: A new proof of Almgren's partial regularity

In the present work we continue the investigations started in [14, 17], which together with the forthcoming papers $[15,16]$ lead to a proof of the following theorem.

Theorem 0.1. Let $\Sigma \subset \mathbb{R}^{m+n}$ be a $C^{3, \varepsilon_{0}}$ submanifold for some $\varepsilon_{0}>0$ and $T$ an $m$ dimensional area minimizing integral current in $\Sigma$. Then, there is a closed set $\operatorname{Sing}(T)$ of Hausdorff dimension at most $m-2$ such that $T$ is a $C^{3, \varepsilon_{0}}$ embedded submanifold in $\Sigma \backslash(\operatorname{spt}(\partial T) \cup \operatorname{Sing}(T))$.

Theorem 0.1 was first proved by Almgren in his monumental work [3], assuming slightly better regularity on $\Sigma$, namely $\Sigma \in C^{5}$. The improvement itself is therefore not so significant, but our proof, besides being much shorter, introduces new ideas and establishes several new results, which we hope will provide useful tools for further investigations in the area. Indeed, although we still follow Almgren's program and use many of his groundbreaking discoveries, the main steps are achieved in a more efficient way thanks to new estimates and techniques. A striking example is the construction of the so-called center manifold, which is by far the most intricate part of Almgren's work and the least explored, in spite of its importance: in this respect, our construction in [15] is considerably simpler and shorter than [3, Chapter 4], and establishes better results.

Some of our improvements are more transparent, although not substantially simpler, when $\Sigma=\mathbb{R}^{m+n}$ and in a book in preparation [12] we will provide a complete and selfcontained account of Theorem 0.1 under such assumption. Moreover, building on our understanding of the various issues involved to the analysis of higher codimension singularities, we plan to tackle Chang's improvement [8], which shows that $\operatorname{Sing}(T)$ consists of isolated points when $m=2$. His arguments rely on a center manifold construction which does not match exactly the statements of [3] and it is not fully justified, but only briefly
sketched in the appendix of [8]. In [18], instead, we give a detailed, simple construction for such center manifold and a complete proof of this refined regularity result.

An alternative route to Chang's result for $J$-holomorphic currents in symplectic manifolds has been given recently in [27, 28]. The interest in the regularity theory for this class of area minimizing 2-dimensional currents has been generated by the seminal paper of Taubes [32] on the equivalence between Gromov and Seiberg-Witten invariants, where it plays an important role. Moreover, the papers [27, 28] have stimulated a lot of activity in the area, cf., for example, [6, 24, 25, 26]. In [6] Bellettini and Rivière proved that, when $T$ is a special Lagrangian cone in $\mathbb{R}^{6}, \operatorname{Sing}(T)$ consists of finitely many half-lines meeting at the origin. This is, to our knowledge, the only result of its type not covered by the Almgren-Chang works. We believe that the Bellettini-Rivière regularity theorem can be extended to general 3-dimensional area minimizing cones in any space dimension, combining the techniques developed in [14]-[18]. Most of the proofs in [6, 24, 25, 26, 27, 28, 32] take advantage of two specific assumptions, the underlying almost complex structure and the 2-dimensionality of the objects of study. Nonetheless these works have had a profound influence on our research.
0.1. A blow-up proof: a very brief overview. In the rest of this foreword we will give a rough outline of the proof of Theorem 0.1, highlighting the contents of this note and the way it merges with its companion papers [15, 16], while comparing them to [3]. Our discussion will be based on a well-known class of examples for which the statement of Theorem 0.1 is optimal, namely singular holomorphic curve of $\mathbb{C}^{2}$. As it was first observed by Federer (cf. [20, 5.4.19]), the integral currents induced by holomorphic subvarieties of $\mathbb{C}^{n}$ (with their natural orientation) are area minimizing.

We denote by $D_{Q}(T)$ the set of points in $\operatorname{spt}(T) \backslash \operatorname{spt}(\partial T)$ where the density of a current $T$ equals the natural number $Q \geq 1$. One first pioneering contribution by Almgren is an elementary, but very clever, generalization of Federer's reduction argument, which has been widely used in several contexts (see [29, Theorem 35.3] and [34]). This argument implies that, if $T$ is area minimizing, then $\operatorname{spt}(T) \backslash\left(\cup_{Q} D_{Q}(T) \cup \operatorname{spt}(\partial T)\right)$ has Hausdorff dimension at most $m-3$. Thus, to prove Theorem 0.1 it suffices to show that the Hausdorff dimension of $\operatorname{Sing}_{Q}(T):=\operatorname{Sing}(T) \cap D_{Q}(T)$ is at most $m-2$. Since the "classical" regularity theory ensures that $T$ is a $C^{1, \alpha}$ submanifold in the neighborhood of any point $x \in D_{1}(T)$, it is natural to argue by induction on $Q$.

Let us therefore consider the case $Q=2$ and a point $x \in D_{2}(T)$. By the monotonicity formula, in some neighborhood $U$ of $x,\|T\|$-almost all points have density 1 or 2 . If the points of density 1 are a set of $\|T\|$-measure zero, by the classical regularity theory $x$ is a regular point for $T$. So any $x \in \operatorname{Sing}_{2}(T)$ must be surrounded by many points of density 1 , as it is, for instance, for the complex curve $\left\{z^{2}=w^{3}\right\} \subset \mathbb{C}^{2}$ at $x=0$. On the other hand, in such an example 0 is an isolated singularity, whereas, if $T$ were to contradict Theorem 0.1, by standard measure theoretic arguments there would be a point $x \in \operatorname{Sing}_{2}(T)$ surrounded by many points of density 2 . From now on we argue by contradiction and assume that this happens for some area minimizing $T$ at the point $0 \in D_{2}(T)$. Moreover, by known facts in geometric measure theory, we can reduce the contradiction to the case that, for a suitable
sequence of radii $r_{k} \downarrow 0$, the homothetic rescalings of the current $T$ by a factor $1 / r_{k}$ (from now on denoted by $T_{k}$ ) converge to a double copy of an $m$-dimensional plane, while at the same time $D_{2}\left(T_{k}\right)$ remains rather large.

It was first recognized by De Giorgi that the convergence of $T_{k}$ to a single copy of a flat plane implies that $\operatorname{spt}\left(T_{k}\right)$ can be well approximated by the graph of Lipschitz functions which are "almost harmonic". However, the example $\left\{z^{2}=w^{3}\right\} \subset \mathbb{C}^{2}$ shows that this is not always the case if the limiting plane has higher multiplicity. Motivated by this fact, Almgren undertook in [3] the strikingly ambitious program of giving a rather complete existence and regularity theory for multiple valued functions minimizing a suitable generalization of the Dirichlet energy, called Dir-minimizers. The crowning achievement of this theory is that, except for a closed set of codimension at most 2, Dir-minimizers can be locally decomposed in classical (i.e. single-valued) non-intersecting harmonic sheets (possibly counted with multiplicity). Such "linear theory" is developed in [3, Chapter 2] and revisited in our paper [14]. Moreover, it is complemented by several technical statements linking the multiple valued graphs to the integral currents, a task which is accomplished in [3, Chapter 1] by Almgren and in [17] by us (we refer to the introduction to our previous two papers $[14,17]$ for more details).

The guiding idea in the contradiction argument is to approximate the currents $T_{k}$ with Lipschitz 2-valued functions and, after a suitable renormalization of their Dirichlet energy, show that they converge to a Dir-minimizer. If the limit inherits a large singular set from the currents $T_{k}$, then it contradicts the linear regularity theory. Obviously, this strategy requires suitable approximations of area minimizing currents with multiple valued graphs, accomplished by Almgren in [3, Chapter 3] and by us in the present paper. If one follows our approach, the convergence of these approximations to a Dir-minimizer can be concluded in a rather direct way. However, we cannot expect that such limit inherits the singular set of the current. For example, given the complex curve $\left\{(z, w):\left(z-w^{2}\right)^{2}=w^{5}\right\} \subset$ $\mathbb{C}^{2}$, any reasonable approximations of homothetic rescalings of this algebraic variety in a neighborhood of the origin converge to a double copy of the classical holomorphic graph $\left\{\left(w, w^{2}\right): w \in \mathbb{C}\right\}$, which has lost the singularity at the origin.

In order to perform the blow-up argument, we then need to "modulate lower order regularities out". This is accomplished by the construction of a center manifold (see [3, Chapter 4] and [15]): such an object is a regular $C^{3, \alpha}$ submanifold which is very close to the average of the sheets of the current at any scale where the latter is "very collapsed". The final blow-up argument is then carried over to a new sequence of 2 -valued approximations of $T_{k}$, performed on the normal bundles of the center manifolds (see [3, Chapter 5] and [16]). By a delicate unique continuation principle, based on a new monotonicity formula discovered by Almgren, a suitable normalization of the latter approximations does converge to a Dir-minimizer which would be forced to have a large singular set, reaching the desired contradiction. This final step builds upon very delicate computations, which thus require a lot of accuracy in the construction of the center manifold, that in turn needs very good estimates on the approximation results of this note. Thus, unlike the two works [14, 17], which can be considered separately, the present paper and [15, 16] are intimately interconnected.
0.2. Our contribution; or, what is new. In their overall structure, our five papers match bijectively the five chapters of [3]. Moreover, it is clear that the ultimate reason for the success of the program is the very same prodigious and celebrated discovery of Almgren: the monotonicity of the frequency function and its astonishing robustness, which enters twice in the plan: at the very beginning, in the linear regularity theory, and at the end, in the convergence of the final approximations (cf. [14, 16]).

So, what is new in our proof? Aside from finer details, which are explained in the introductions to each of our papers, there are some new contributions which come at a higher level. Our investigations started with the idea that the machinery developed in metric analysis and metric geometry in the last 30 years could reduce the complexity of several arguments in Almgren's program. This is, indeed, the case at many levels in the two papers $[14,17]$ and in this note. Approaching vast parts of Almgren's theory with these tools, we not only get shorter and more transparent proofs, but often also achieve stronger analytic estimates, which give a better starting point for the PDE parts of the program. Moreover, as it often happens when "abstract nonsense" simplifies preexisting mathematical theories, such machinery provides also a better insight to the material of [3], as it highlights the important points in the proofs therein.

However, this alone would not explain the shortness of our papers compared to [3, Chapters $3,4,5]$. The other important reason is that we also derive some fundamental, new "hard" estimates. A primary example is the present paper, where the main a priori estimate is a new higher integrability result, which comes from a Gehring-type argument and is inspired by a simple remark in the linear theory (the higher integrability of gradients of Dir-minimizers) which to our knowledge is not observed in Almgren's monograph. Similar instances are present in the papers [15, 16], where some new quantities and guiding principles are introduced (for instance, the "modified frequency" function in [16] and the "splitting-before-tilting" principle in [15], inspired by [26]), which probably lead to the improvement on the regularity assumptions of the ambient manifold $\Sigma$. In all these cases we provide more efficient tools compared to [3] and invoke more PDE theory at several levels, drawing connections with fairly classical concepts from other areas of analysis (such as maximal functions, Lipschitz truncations, elliptic systems, Sobolev capacity). Unfortunately we do not understand Almgren's arguments at a sufficiently deep level to draw a fine parallel between our papers $[15,16]$ and the last two chapters of his book, where the intricacy of the arguments in [3] is almost prohibitive. It remains the fact that our papers are much more accessible, and we hope that in the near future our work will be used to penetrate further in the richness and beauty of Almgren's monograph and to go beyond Theorem 0.1.

Acknowledgments. The first author is deeply indebted to Tristan Rivière who "infected him" with the beauty of the problem. Both authors are also warmly thankful to Giovanni Alberti, Bill Allard, Luigi Ambrosio and Bernd Kirchheim not only for several enlightening conversations, but also for their constant enthusiastic encouragement.

Several other colleagues and friends have contributed with important scientific conversations at some specific stage, for which they will be acknowledged specifically in the various
papers. In particular, for this first one we are grateful to Stefano Bianchini, Sergio Conti, Matteo Focardi, Jonas Hirsch and Luca Spolaor for very useful discussions and comments.

This work was carried over several years and the authors wish to thank many institutions where they spent very productive visits, namely: the University of Rome La Sapienza, the Scuola Normale Superiore and the University of Pisa, the Max Planck Institute for Mathematics in the Sciences and the University of Leipzig, the University of Zürich, the SISSA in Trieste, the University of Warwick and, most of all, Princeton University, which hosted the first author during his last sabbatical. We also acknowledge the support of the ERC grant RAM, ERC 306247.

We finally thank the anonymous referee for his/her careful reading and very valuable suggestions, which contributed to a substantial improvement of the initial manuscript.

## 1. Introduction

1.1. A priori gradient $L^{p}$ estimate. In order to state the main results, we start specifying some assumptions, which will hold throughout the paper. For the notation concerning submanifolds $\Sigma \subset \mathbb{R}^{m+n}$ we refer to [17, Section 1]. With $\mathbf{B}_{r}(p)$ and $B_{r}(x)$ we denote, respectively, the open ball with radius $r$ and center $p$ in $\mathbb{R}^{m+n}$ and the open ball with radius $r$ and center $x$ in $\mathbb{R}^{m} . \mathbf{C}_{r}(x)$ will always denote the cylinder $B_{r}(x) \times \mathbb{R}^{n}$ and the point $x$ will be omitted when it is the origin. In fact, by a slight abuse of notation, we will often treat the center $x$ as a point in $\mathbb{R}^{m+n}$, avoiding the correct, but more cumbersome, $(x, 0)$. Let $e_{i}$ be the unit vectors in the standard basis, $\pi_{0}$ the (oriented) plane $\mathbb{R}^{m} \times\{0\}$ and $\vec{\pi}_{0}$ the $m$-vector $e_{1} \wedge \ldots \wedge e_{m}$ orienting it. We denote by $\mathbf{p}$ and $\mathbf{p}^{\perp}$ the orthogonal projections onto, respectively, $\pi_{0}$ and its orthogonal complement $\pi_{0}^{\perp}$. In some cases we need orthogonal projections onto other planes $\pi$ and their orthogonal complements $\pi^{\perp}$, for which we use the notation $\mathbf{p}_{\pi}$ and $\mathbf{p}_{\pi}^{\perp}$. For what concerns integral currents we use the definitions and the notation of [29].

Assumption 1.1. $\Sigma \subset \mathbb{R}^{m+n}$ is a $C^{2}$ submanifold of dimension $m+\bar{n}=m+n-l$, which is the graph of an entire function $\Psi: \mathbb{R}^{m+\bar{n}} \rightarrow \mathbb{R}^{l}$ and satisfies the bounds

$$
\begin{equation*}
\|D \Psi\|_{0} \leq c_{0} \quad \text { and } \quad \mathbf{A}:=\left\|A_{\Sigma}\right\|_{0} \leq c_{0} \tag{1.1}
\end{equation*}
$$

where $c_{0}$ is a positive (small) dimensional constant. $T$ is an integral current of dimension $m$ with bounded support contained in $\Sigma$ and which, for some open cylinder $\mathbf{C}_{4 r}(x)$ (with $r \leq 1)$ and some positive integer $Q$, satisfies

$$
\begin{equation*}
\mathbf{p}_{\sharp} T=Q \llbracket B_{4 r}(x) \rrbracket \quad \text { and } \quad \partial T\left\llcorner\mathbf{C}_{4 r}(x)=0 .\right. \tag{1.2}
\end{equation*}
$$

If we say that $T$ is area minimizing we then mean that it is area-minimizing in $\Sigma \cap \mathbf{C}_{4 r}(x)$, namely that $\mathbf{M}(T) \leq \mathbf{M}(T+\partial S)$ for any integral $S$ with $\operatorname{spt}(S) \subset \Sigma \cap \mathbf{C}_{4 r}(x)$.

Definition 1.2 (Excess measure). For a current $T$ as in Assumption 1.1 we define the cylindrical excess $\mathbf{E}\left(T, \mathbf{C}_{4 r}(x)\right)$, the excess measure $\mathbf{e}_{T}$ and its density $\mathbf{d}_{T}$ :

$$
\begin{gathered}
\mathbf{E}\left(T, \mathbf{C}_{r}(x)\right):=\frac{\|T\|\left(\mathbf{C}_{r}(x)\right)}{\omega_{m} r^{m}}-Q, \\
\mathbf{e}_{T}(A):=\|T\|\left(A \times \mathbb{R}^{n}\right)-Q|A| \quad \text { for every Borel } A \subset B_{r}(x), \\
\mathbf{d}_{T}(y):=\underset{s \rightarrow 0}{\limsup } \frac{\mathbf{e}_{T}\left(B_{s}(y)\right)}{\omega_{m} s^{m}}=\underset{s \rightarrow 0}{\limsup } \mathbf{E}\left(T, \mathbf{C}_{s}(y)\right),
\end{gathered}
$$

where $\omega_{m}$ is the measure of the $m$-dimensional unit ball (the subscripts ${ }_{T}$ will be omitted if clear from the context).

Since $T$ has finite mass, the function $\mathbf{d}$ is naturally an $L^{1}$ function. However, we can show the following higher integrability estimate when $T$ is, in addition, area minimizing. We call it a gradient $L^{p}$ estimate because we will show that $\mathbf{d}$ coincides with the gradient of an appropriate Lipschitz function on a large region.

Theorem 1.3 (Gradient $L^{p}$ estimate). There exist constants $p_{1}>1$ and $C, \varepsilon_{10}>0$ (depending on $m, n, \bar{n}, Q$ ) with the following property. Let $T$ be as in Assumption 1.1 in the cylinder $\mathbf{C}_{4}$. If $T$ is area minimizing and $E=\mathbf{E}\left(T, \mathbf{C}_{4}\right)<\varepsilon_{10}$, then

$$
\begin{equation*}
\int_{\{\mathbf{d} \leq 1\} \cap B_{2}} \mathbf{d}^{p_{1}} \leq C E^{p_{1}-1}\left(E+\mathbf{A}^{2}\right) . \tag{1.3}
\end{equation*}
$$

In the case $Q=1$ or $\bar{n}=1$, it follows from the classical regularity theory (essentially due to De Giorgi, cf. [9]) that $T$ is a $C^{1, \alpha}$ submanifold in $\mathbf{C}_{2}$. However, when $\min \{Q, \bar{n}\} \geq 2$, $T$ is not necessarily regular and Theorem 1.3 gives in fact an a priori regularity estimate: in this case (1.3) cannot be improved (except for optimizing the constants $p_{1}, C$ and $\varepsilon_{10}$ ). Indeed, for $Q=m=2, \Sigma=\mathbb{R}^{4}$ and $p_{1}=2$, (1.3) is false no matter how large $\varepsilon_{10}^{-1}$ and $C$ are chosen (cf. [10, Section 6.2]).

In order to prove Theorem 1.3 we develop the following tools:
(a) a general scheme to approximate integer rectifiable currents with multiple valued functions, relying heavily on the "metric analysis" of [14] and on a modified "Jerrard-Soner" BV estimate for the slicing of currents (cf. Proposition 2.2);
(b) a simple and robust harmonic approximation of area minimizing currents with multiple valued functions (cf. Theorem 4.2);
(c) the higher integrability of the gradient of Dir-minimizing multiple valued functions (cf. Theorem 5.1 - see also [30] for a different proof and related results).
In turn, Theorem 1.3 will be combined with (a) to achieve a very accurate approximation result for area minimizing current, stated in Theorem 1.4. This theorem and some corollaries of our analysis play a fundamental role in the papers $[15,16]$ and, as explained in the Foreword, have a counterpart in [3, Chapter 3]. However, our derivation of Theorem 1.4 differs substantially from Almgren's and when we use some of his ideas, as it is for the existence of the almost projection $\boldsymbol{\rho}^{\star}$ of Section 7, we give independent arguments for the main steps of the proof.
1.2. Strong approximation of area minimizing currents. Concerning multiple valued functions we will follow the notation and terminology of [14, 17]. In particular, a $Q$-valued function is a map $f$ (usually defined over a measurable subset $\Omega$ of $\mathbb{R}^{m}$ ) taking values in the space $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ of unordered $Q$-tuples of points in $\mathbb{R}^{n}$, denoted by $\sum_{i} \llbracket P_{i} \rrbracket$. $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ can be equipped with a natural metric $\mathcal{G}$ (cf. [14, Definition 0.2$]$ ) and for $f$ measurable there exist measurable functions $f_{i}: \Omega \rightarrow \mathbb{R}^{n}$ such that $f(x)=\sum_{i} \llbracket f_{i}(x) \rrbracket \forall x \in \Omega$ (cf. [14, Proposition 0.4]). The functions $f_{i}$ are not uniquely determined, but in using this notation we assume to have fixed some suitable $f_{i}$ 's. Moreover, if $f$ is Lipschitz, resp. $f \in W^{1,2}\left(\Omega, \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)\right)$ (cf. [14, Definition 0.5$\left.]\right)$ and $\Omega$ is open, then there exist measurable functions $D f_{i} \in L^{\infty}$, resp. $L^{2}$, such that $\sum_{i} \llbracket D f_{i}(x) \rrbracket$ is the approximate differential of $f$ (cf. [14, Definition 2.6]) at a.e. $x$. In fact in this case the $f_{i}$ 's and $D f_{i}$ 's can be chosen so that the first are approximately differentiable a.e. and the second are their approximate differentials in the classical sense (cf. [17, Lemma 1.1]). The Dirichlet energy of $f$ is then $\operatorname{Dir}(f, \Omega):=\int_{\Omega}|D f|^{2}$, where $|D f|^{2}:=\sum_{i}\left|D f_{i}\right|^{2}$. Following [17, Definition 1.10], we denote by $\mathbf{G}_{f}$ the integer rectifiable current, in $\mathbb{R}^{m+n}$, naturally associated to the graph of a Lipschitz $Q$-valued map $f: \mathbb{R}^{m} \supset A \rightarrow \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$. Moreover, we will use the notation $\operatorname{osc}(f)$ for the quantity $\inf _{p} \sup _{x} \mathcal{G}(f(x), Q \llbracket p \rrbracket)$.

Theorem 1.4 (Almgren's strong approximation). There exist constants $C, \gamma_{1}, \varepsilon_{1}>0$ (depending on $m, n, \bar{n}, Q$ ) with the following property. Assume that $T$ is area minimizing, satisfies Assumption 1.1 in the cylinder $\mathbf{C}_{4 r}(x)$ and $E=\mathbf{E}\left(T, \mathbf{C}_{4 r}(x)\right)<\varepsilon_{1}$. Then, there is a map $f: B_{r}(x) \rightarrow \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$, with $\operatorname{spt}(f(x)) \subset \Sigma$ for every $x$, and a closed set $K \subset B_{r}(x)$ such that

$$
\begin{gather*}
\operatorname{Lip}(f) \leq C E^{\gamma_{1}},  \tag{1.4}\\
\mathbf{G}_{f}\left\llcorner\left(K \times \mathbb{R}^{n}\right)=T\left\llcorner\left(K \times \mathbb{R}^{n}\right) \quad \text { and } \quad\left|B_{r}(x) \backslash K\right| \leq C E^{\gamma_{1}}\left(E+r^{2} \mathbf{A}^{2}\right) r^{m},\right.\right.  \tag{1.5}\\
\left.\left.\left|\|T\|\left(\mathbf{C}_{\sigma r}(x)\right)-Q \omega_{m}(\sigma r)^{m}-\frac{1}{2} \int_{B_{\sigma r}(x)}\right| D f\right|^{2} \right\rvert\, \leq C E^{\gamma_{1}}\left(E+r^{2} \mathbf{A}^{2}\right) r^{m} \quad \forall 0<\sigma \leq 1 . \tag{1.6}
\end{gather*}
$$

$$
\begin{align*}
& \text { If in addition } \mathbf{h}\left(T, \mathbf{C}_{4 r}(x), \pi_{0}\right):=\sup \left\{\left|\mathbf{p}^{\perp}(x)-\mathbf{p}^{\perp}(y)\right|: x, y \in \operatorname{spt}(T) \cap \mathbf{C}_{4 r}(x)\right\} \leq r \text {, then } \\
& \qquad \operatorname{osc}(f) \leq C \mathbf{h}\left(T, \mathbf{C}_{4 r}(x), \pi_{0}\right)+C\left(E^{1 / 2}+r \mathbf{A}\right) r \tag{1.7}
\end{align*}
$$

The gain of a small power $E^{\gamma_{1}}$ in the three estimates (1.4)-(1.6) plays a crucial role in the papers $[15,16]$. When $Q=1$ and $\Sigma=\mathbb{R}^{m+1}$, this approximation theorem was first proved by De Giorgi in [9]. In the generality above it appears in the big regularity paper for the first time (cf. [3, Sections 3.28-3.30]). Its proof is an elementary consequence of Theorem 6.1 and the Lipschitz approximation algorithm mentioned above. In turn Theorem 6.1 will be derived from Theorem 1.3 using a suitable competitor argument. In the case $Q=1$, the competitor is the convolution of (a first) Lipschitz approximation with a smooth kernel, a classical argument which in fact appears already in De Giorgi's seminal paper [9], although in a slightly different form (cf. [13, Appendix]).

Here we need a similar approach in the framework of multiple valued functions. However, since $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ is highly nonlinear, it is not possible to regularize directly by convolution.

We exploit at this point a key idea of Almgren, embedding $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ in an Euclidean space and using some suitable "almost projections" $\boldsymbol{\rho}_{\delta}^{\star}$. Our proof of the existence of these almost projections is however different from the one given by Almgren in [3, Theorem 1.3] and, indeed, gives better bounds in terms of the relevant parameters (see Proposition 6.2).
1.3. Harmonic approximation. A second ingredient which in $[15,16]$ will play a key role is the harmonic approximation of Theorem 1.6 below (already mentioned in (b) above). In order to state it we need to set some notation about the ambient manifold $\Sigma$.

Remark 1.5 (Estimates on $\Psi$ in good Cartesian coordinates). Assume that $T$ is as in Assumption 1.1 in the cylinder $\mathbf{C}_{4 r}(x)$. If $E:=\mathbf{E}\left(T, \mathbf{C}_{4 r}(x)\right)$ is smaller than a geometric constant, we can assume, without loss of generality, that the function $\Psi: \mathbb{R}^{m+\bar{n}} \rightarrow \mathbb{R}^{l}$ parameterizing $\Sigma$ satisfies $\Psi(x)=0,\|D \Psi\|_{0} \leq C E^{1 / 2}+C \mathbf{A} r$ and $\left\|D^{2} \Psi\right\|_{0} \leq C \mathbf{A}$. Indeed observe that

$$
E=\mathbf{E}\left(T, \mathbf{C}_{4 r}(x)\right)=\frac{1}{2 \omega_{m}(4 r)^{m}} \int_{\mathbf{C}_{4 r}(x)}\left|\vec{T}(y)-\vec{\pi}_{0}\right|^{2} d\|T\|(y)
$$

Thus, we can fix a point $p \in \operatorname{spt}(T) \cap \mathbf{C}_{4 r}(x)$ such that $\left|\vec{T}(p)-\vec{\pi}_{0}\right| \leq C E^{1 / 2}$. Then, we can find an associated rotation $O \in O(m+n)$ such that $O_{\sharp} \vec{T}(p)=\vec{\pi}_{0}$ and $|O| \leq C E^{1 / 2}$. It follows that $\pi:=O\left(T_{p} \Sigma\right)$ is a $(m+\bar{n})$-dimensional plane such that $\pi_{0} \subset \pi$ and $\left\|\pi-T_{p} \Sigma\right\| \leq$ $C E^{1 / 2}$. We choose new coordinates so that $\pi_{0}$ remains equal to $\mathbb{R}^{m} \times\{0\}$ but $\mathbb{R}^{m+\bar{n}} \times\{0\}$ equals $\pi$. Since the excess $E$ is assumed to be sufficiently small, we can write $\Sigma$ as the graph of a function $\Psi: \pi \rightarrow \pi^{\perp}$. If $(z, \Psi(z))=p$, then $|D \Psi(z)| \leq C\left\|T_{p} \Sigma-\mathbb{R}^{m+\bar{n}} \times\{0\}\right\| \leq C E^{1 / 2}$. However, $\left\|D^{2} \Psi\right\|_{0} \leq C \mathbf{A}$ and so $\|D \Psi\|_{0} \leq C E^{1 / 2}+C \mathbf{A} r$. Moreover, $\Psi(x)=0$ is achieved translating the system of reference by a vector orthogonal to $\mathbb{R}^{m+\bar{n}} \times\{0\}$ and, hence, belonging to $\{0\} \times \mathbb{R}^{n}$.

From now on, we will often consider $Q$-valued maps $y \mapsto w(y) \in \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)=\mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}} \times \mathbb{R}^{l}\right)$ which take the form $w(y)=\sum_{i} \llbracket\left(u_{i}(y), \Psi\left(y, u_{i}(y)\right) \rrbracket\right.$, where $u=\sum_{i} \llbracket u_{i} \rrbracket$ is evidently a map taking values in $\mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)$. For $w$ we will then use the short-hand notation $w=(u, \Psi(y, u))$. We also recall the notation for the average map $\eta: \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}^{n}$ defined by

$$
\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right) \ni T=\sum_{i=1}^{Q} \llbracket P_{i} \rrbracket \mapsto \eta(T):=\frac{1}{Q} \sum_{i=1}^{Q} P_{i} \in \mathbb{R}^{n}
$$

Theorem 1.6 (Harmonic approximation). Let $\gamma_{1}$ be the constant of Theorem 1.4. Then, for every $\bar{\eta}, \bar{\delta}>0$, there is a positive constant $\bar{\varepsilon}_{1}$ with the following property. Assume that $T$ is as in Theorem 1.4, $E:=\mathbf{E}\left(T, \mathbf{C}_{4 r}(x)\right)<\bar{\varepsilon}_{1}$ and $r \mathbf{A} \leq E^{1 / 4+\bar{\delta}}$. If $f$ is the map in Theorem 1.4 and we fix Cartesian coordinates as in Remark 1.5, then there exists a Dir-minimizing function $u: B_{r}(x) \rightarrow \mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)$ such that $w:=(u, \Psi(y, u))$ satisfies

$$
\begin{equation*}
r^{-2} \int_{B_{r}(x)} \mathcal{G}(f, w)^{2}+\int_{B_{r}(x)}(|D f|-|D w|)^{2}+\int_{B_{r}(x)}|D(\boldsymbol{\eta} \circ f)-D(\boldsymbol{\eta} \circ w)|^{2} \leq \bar{\eta} E r^{m} \tag{1.8}
\end{equation*}
$$

This theorem is the multi-valued analog of De Giorgi's harmonic approximation (cf. [9]). We prove it via a compactness argument which, although very close in spirit to De Giorgi's
original one, is to our knowledge new (even when $n=\bar{n}=1$ ). Indeed, it uses neither the monotonicity formula nor a regularization by convolution of the Lipschitz approximation, and we expect it to be useful in different contexts.
1.4. Persistence of $Q$-points. A major ingredient in [16] is the persistence of points of maximal multiplicity in the approximation of Theorem 1.4, when interpreted in a suitable "limiting sense". If the current $T$ has a point of density $Q, f$ must satisfy the following integral bound (even though $f$ might have no values of multiplicity $Q$ ).

Theorem 1.7 (Persistence of $Q$-points). For every $\hat{\delta}, C^{\star}>0$, there is $\left.\bar{s} \in\right] 0, \frac{1}{2}[$ such that, for every $s<\bar{s}$, there exists $\hat{\varepsilon}\left(s, C^{*}, \hat{\delta}\right)>0$ with the following property. If $T$ is as in Theorem 1.4, $E:=\mathbf{E}\left(T, \mathbf{C}_{4 r}(x)\right)<\hat{\varepsilon}, r^{2} \mathbf{A}^{2} \leq C^{\star} E$ and $\Theta(T,(p, q))=Q$ at some $(p, q) \in \mathbf{C}_{r / 2}(x)$, then the approximation $f$ of Theorem 1.4 satisfies

$$
\begin{equation*}
\int_{B_{s r}(p)} \mathcal{G}(f, Q \llbracket \boldsymbol{\eta} \circ f \rrbracket)^{2} \leq \hat{\delta} s^{m} r^{2+m} E \tag{1.9}
\end{equation*}
$$

1.5. A remark on notation. Finally we remark that we follow closely the notation of $[14,17]$, except for a subtle point. We denote by $\boldsymbol{\xi}$ the map in [14, Corollary 2.2], which there was denoted by $\boldsymbol{\xi}_{B W}$, since the symbol $\boldsymbol{\xi}$ was in fact used for the "precursor map" of [14, Theorem 2.1]. So, here $\boldsymbol{\xi}: \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}^{N(Q, n)}$ is an injective function satisfying the following three properties:
(i) $\operatorname{Lip}(\boldsymbol{\xi}) \leq 1$;
(ii) $\operatorname{Lip}\left(\left.\boldsymbol{\xi}^{-1}\right|_{\mathcal{Q}}\right) \leq C(n, Q)$, where $\mathcal{Q}=\boldsymbol{\xi}\left(\mathcal{A}_{Q}\right)$;
(iii) $|D f|=|D(\boldsymbol{\xi} \circ f)|$ almost everywhere for every $f \in W^{1,2}\left(\Omega, \mathcal{A}_{Q}\right)$.

This "improved" $\boldsymbol{\xi}$ was suggested by Brian White and appears for the first time in [8]. The conclusion (iii) above is actually not explicitly stated in [14], but it follows easily: indeed [14, Corollary 2.2] implies the identity $|D f|=|D(\boldsymbol{\xi} \circ f)|$ at every point of differentiability of a Lipschitz map and, hence, almost everywhere. The case of a general $f \in W^{1,2}\left(\Omega, \mathcal{A}_{Q}\right)$ can then be concluded from [14, Proposition 2.5].

We will use the notation $C$ and $c$ for generic positive dimensional constants, which may possibly change from line to line: we will always understand that these constants depends only on the dimensional parameters $m, \bar{n}, n, Q, c_{0}$ of Assumption 1.1.

## 2. Lipschitz approximation

To begin with, we develop a robust algorithm to approximate currents $T$ as in Assumption 1.1 with graphs of multiple valued functions. Following the work of Ambrosio and Kirchheim [5], we view the slice map $x \mapsto\langle T, \mathbf{p}, x\rangle$ as a function taking values in the space $\boldsymbol{I}_{0}\left(\mathbb{R}^{n}\right)$ of 0-dimensional integral currents. A key estimate of Jerrard and Soner (cf. [5, 23]) implies that this map has bounded variation in the metric sense introduced by Ambrosio in [4]. On the other hand, following [14], $Q$-valued functions can be viewed as Sobolev maps taking values into (a subset of) $\boldsymbol{I}_{0}\left(\mathbb{R}^{n}\right)$. Thus, finding Lipschitz multiple valued approximations of $T$ can be seen as a particular case of the more general task of finding Lipschitz approximations of BV maps with a fairly general target space.

Definition 2.1 (Maximal function of the excess measure). Given a current $T$ as in Assumption 1.1 we introduce the "non-centered" maximal function of $\mathbf{e}_{T}$ :

$$
\mathbf{m e}_{T}(y):=\sup _{y \in B_{s}(w) \subset B_{4 r}(x)} \frac{\mathbf{e}_{T}\left(B_{s}(w)\right)}{\omega_{m} s^{m}}=\sup _{y \in B_{s}(w) \subset B_{4 r}(x)} \mathbf{E}\left(T, \mathbf{C}_{s}(w)\right) .
$$

We can now state the main result of the section, which provides the first Lipschitz approximation for rectifiable currents.

Proposition 2.2 (Lipschitz approximation). There exists a constant $C>0$ with the following property. Let $T$ and $\Psi$ be as in Assumption 1.1 in the cylinder $\mathbf{C}_{4 s}(x)$. Set $E=\mathbf{E}\left(T, \mathbf{C}_{4 s}(x)\right)$, let $0<\delta_{11}<1$ be such that $16^{m} E<\delta_{11}$, and define

$$
K:=\left\{\mathbf{m e}_{T}<\delta_{11}\right\} \cap B_{3 s}(x) .
$$

Then, there is $u \in \operatorname{Lip}\left(B_{3 s}(x), \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)\right)$ such that $\operatorname{spt}(u(y)) \subset \Sigma$ for every $y \in B_{3 s}(x)$ and

$$
\begin{gather*}
\operatorname{Lip}(u) \leq C\left(\delta_{11}^{1 / 2}+\|D \Psi\|_{0}\right), \quad \operatorname{osc}(u) \leq C \mathbf{h}\left(T, \mathbf{C}_{4 s}(x), \pi_{0}\right)+C s\|D \Psi\|_{0}, \\
\mathbf{G}_{u}\left\llcorner\left(K \times \mathbb{R}^{n}\right)=T\left\llcorner\left(K \times \mathbb{R}^{n}\right),\right.\right. \\
\left|B_{r}(x) \backslash K\right| \leq \frac{10^{m}}{\delta_{11}} \mathbf{e}_{T}\left(\left\{\mathbf{m e}_{T}>2^{-m} \delta_{11}\right\} \cap B_{r+r_{0} s}(x)\right) \quad \forall r \leq 3 s \tag{2.1}
\end{gather*}
$$

where $r_{0}=16 \sqrt[m]{E / \delta_{11}}<1$.
The proof of the proposition is based on a BV estimate which differs from the ones of [5, 23]. Note that we do not assume that $T$ is area minimizing. Indeed, even the assumption (1.2) could be relaxed, but we do not pursue this issue here.
2.1. The modified Jerrard-Soner estimate. Recall that each element $S \in \boldsymbol{I}_{0}\left(\mathbb{R}^{m+n}\right)$ is simply a finite sum of Dirac deltas, $S=\sum_{i=1}^{h} w_{i} \delta_{z_{i}}$, where $h \in \mathbb{N}, w_{i} \in\{-1,1\}$ and the $z_{i}$ 's are (not necessarily distinct) points in $\mathbb{R}^{m+n}$. Let $T$ be a current as in Assumption 1.1 in the cylinder $\mathbf{C}_{4}$. The slicing map $x \mapsto\langle T, \mathbf{p}, x\rangle$ takes values in $\boldsymbol{I}_{0}\left(\mathbb{R}^{m+n}\right)$ and is characterized by (cf. [29, Section 28]):

$$
\begin{equation*}
\int_{B_{4}}\langle T, \mathbf{p}, x\rangle(\varphi) d x=T(\varphi d x) \quad \text { for every } \varphi \in C_{c}^{\infty}\left(\mathbf{C}_{4}\right) \tag{2.2}
\end{equation*}
$$

Moreover $\operatorname{spt}(\langle T, \mathbf{p}, x\rangle) \subseteq \mathbf{p}^{-1}(\{x\})$ and therefore $\langle T, \mathbf{p}, x\rangle=\sum_{i} w_{i} \delta_{\left(x, y_{i}\right)}$. The assumption (1.2) guarantees that $\sum_{i} w_{i}=Q$ for almost every $x$. In order to state our BV estimate, we consider the push-forwards of $\langle T, \mathbf{p}, x\rangle$ into the vertical directions:

$$
\begin{equation*}
T_{x}:=\mathbf{p}_{\sharp}^{\perp}(\langle T, \mathbf{p}, x\rangle) \in \boldsymbol{I}_{0}\left(\mathbb{R}^{n}\right) . \tag{2.3}
\end{equation*}
$$

It follows from (2.2) that the currents $T_{x}$ are characterized through the identity:

$$
\begin{equation*}
\int_{B_{4}} T_{x}(\psi) \varphi(x) d x=T(\varphi(x) \psi(y) d x) \quad \text { for every } \varphi \in C_{c}^{\infty}\left(B_{4}\right), \psi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right) \tag{2.4}
\end{equation*}
$$

Proposition 2.3 (BV estimate). Assume $T$ satisfies Assumption 1.1 in $\mathbf{C}_{4}$ (i.e. $r=1$ and $x=0$ in Assumption 1.1). For every $\psi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, set $\Phi_{\psi}(x):=T_{x}(\psi)$. If $\|D \psi\|_{\infty} \leq 1$, then $\Phi_{\psi} \in B V\left(B_{4}\right)$ and satisfies

$$
\begin{equation*}
\left(\left|D \Phi_{\psi}\right|(A)\right)^{2} \leq 2 m^{2} \mathbf{e}_{T}(A)\|T\|\left(A \times \mathbb{R}^{n}\right) \quad \text { for every Borel set } A \subseteq B_{4} \tag{2.5}
\end{equation*}
$$

Note that in the usual Jerrard-Soner estimate the RHS of (2.5) would be $\left(\|T\|\left(A \times \mathbb{R}^{n}\right)\right)^{2}$. Proof. It is enough to prove (2.5) for every open set $A \subseteq B_{4}$. To this aim, recall that:

$$
\begin{equation*}
\left|D \Phi_{\psi}\right|(A)=\sup \left\{\int_{A} \Phi_{\psi}(x) \operatorname{div} \varphi(x) d x: \varphi \in C_{c}^{\infty}\left(A, \mathbb{R}^{m}\right),\|\varphi\|_{\infty} \leq 1\right\} \tag{2.6}
\end{equation*}
$$

For any smooth vector field $\varphi$ we have $(\operatorname{div} \varphi(x)) d x=d \Xi$, where

$$
\Xi=\sum_{j} \varphi_{j} d \hat{x}^{j} \quad \text { and } \quad d \hat{x}^{j}=(-1)^{j-1} d x^{1} \wedge \cdots \wedge d x^{j-1} \wedge d x^{j+1} \wedge \cdots \wedge d x^{m}
$$

From (2.4) and the assumption $\partial T\left\llcorner\mathbf{C}_{4}=0\right.$ in (1.2), we conclude that

$$
\begin{align*}
\int_{A} \Phi_{\psi}(x) \operatorname{div} \varphi(x) d x & =\int_{B_{4}} T_{x}(\psi) \operatorname{div} \varphi(x) d x=T(\psi \operatorname{div} \varphi d x) \\
& =T(\psi d \Xi)=T(d(\psi \Xi))-T(d \psi \wedge \Xi)=-T(d \psi \wedge \Xi) \tag{2.7}
\end{align*}
$$

Observe that the $m$-form $d \psi \wedge \Xi$ has no $d x$ component, since

$$
\begin{equation*}
d \psi \wedge \Xi=\sum_{j=1}^{m} \sum_{i=1}^{n} \frac{\partial \psi}{d y^{i}}(y) \varphi_{j}(x) d y^{i} \wedge d \hat{x}^{j} \tag{2.8}
\end{equation*}
$$

Write $\vec{T}=\left\langle\vec{T}, \vec{\pi}_{0}\right\rangle \vec{\pi}_{0}+\vec{S}$. Then,

$$
(T(d \psi \wedge \Xi))^{2}=\left(\int\langle\vec{S}, d \psi \wedge \Xi\rangle d\|T\|\right)^{2} \leq\||d \psi \wedge \Xi|\|_{\infty}^{2}\|T\|\left(A \times \mathbb{R}^{n}\right) \int_{A \times \mathbb{R}^{n}}|\vec{S}|^{2} d\|T\|
$$

( $|\cdot|$ denotes the norms on $\Lambda_{m}$ and $\Lambda^{m}$ induced by the natural inner products $\left.\langle\rangle,\right)$. Since $|\vec{S}|^{2}=1-\left\langle\vec{T}, \vec{\pi}_{0}\right\rangle^{2} \leq 2-2\left\langle\vec{T}, \vec{\pi}_{0}\right\rangle$, we have

$$
\int_{A \times \mathbb{R}^{n}}|\vec{S}|^{2} d\|T\| \leq 2 \int_{A \times \mathbb{R}^{n}}\left(1-\left\langle\vec{T}, \vec{\pi}_{0}\right\rangle\right) d\|T\|=2 \mathbf{e}_{T}(A)
$$

Moreover, by (2.8), $\|\mid d \psi \wedge \Xi\|_{\infty} \leq m\|D \psi\|_{\infty}\|\varphi\|_{\infty} \leq m$. Summarizing, we get

$$
\int_{A} \Phi_{\psi}(x) \operatorname{div} \varphi(x) d x \leq\left(2 m^{2} \mathbf{e}_{T}(A)\|T\|\left(A \times \mathbb{R}^{n}\right)\right)^{1 / 2} \quad \forall \varphi \in C_{c}^{\infty}\left(A, \mathbb{R}^{m}\right),\|\varphi\|_{\infty} \leq 1
$$

Taking the supremum over $\varphi$ 's we conclude (2.5) through (2.6).
2.2. Proof of Proposition 2.2. Since the statement is invariant under translations and dilations, without loss of generality we assume $x=0$ and $s=1$. Consider the slices $T_{x}:=\mathbf{p}_{\sharp}^{\perp}\langle T, \mathbf{p}, x\rangle \in \boldsymbol{I}_{0}\left(\mathbb{R}^{n}\right)$ and recall that $\|T\|\left(A \times \mathbb{R}^{n}\right) \geq \int_{A} \mathbf{M}\left(T_{x}\right) d x$ for every open set $A$ (cf. [29, Lemma 28.5]). Therefore,

$$
\mathbf{M}\left(T_{x}\right) \leq \lim _{r \rightarrow 0} \frac{\|T\|\left(\mathbf{C}_{r}(x)\right)}{\omega_{m} r^{m}} \leq \mathbf{m e}_{T}(x)+Q \quad \text { for almost every } x
$$

Since $\delta_{11}<1$, we infer $\mathbf{M}\left(T_{x}\right)<Q+1$ for a.e. $x \in K$. There are, then, $Q$ functions $g_{i}: K \rightarrow \mathbb{R}^{n}$ such that $T_{x}=\sum_{i=1}^{Q} \delta_{g_{i}(x)}$ for a.e. $x \in K$. Define $g: K \mapsto \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ as $g:=\sum_{i} \llbracket g_{i} \rrbracket$ and fix $\psi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$. Proposition 2.3 gives

$$
\left(\left|D \Phi_{\psi}\right|\left(B_{r}(y)\right)\right)^{2} \leq 2 m^{2} \mathbf{e}_{T}\left(B_{r}(y)\right)\|T\|\left(\mathbf{C}_{r}(y)\right)=2 m^{2} \mathbf{e}_{T}\left(B_{r}(y)\right)\left(Q\left|B_{r}(y)\right|+\mathbf{e}_{T}\left(B_{r}(y)\right)\right)
$$

Hence, if we define the maximal function

$$
\mathbf{m}\left|D \Phi_{\psi}\right|(x):=\sup _{x \in B_{r}(y) \subset B_{4 r}} \frac{\left|D \Phi_{\psi}\right|\left(B_{r}(y)\right)}{\left|B_{r}(y)\right|}
$$

we conclude that

$$
\left(\mathbf{m}\left|D \Phi_{\psi}\right|(x)\right)^{2} \leq 2 m \mathbf{m e}_{T}(x)^{2}+2 m Q \mathbf{m e}_{T}(x) \leq C \delta_{11} \quad \text { for every } x \in K
$$

Therefore, the theory of $B V$ functions gives a dimensional constant $C$ such that

$$
\begin{equation*}
\left|\Phi_{\psi}(x)-\Phi_{\psi}(y)\right| \leq C \delta_{11}^{1 / 2}|x-y| \quad \forall x, y \in K \text { Lebesgue points of } \Phi_{\psi}, \tag{2.9}
\end{equation*}
$$

(see for instance [19, Section 6.6.2]: although in that reference the authors use the centered maximal function, the proof works obviously also in our context). Consider next the Wasserstein distance of exponent 1 on 0-dimensional integral currents $S_{1}, S_{2}$ :

$$
\begin{equation*}
W_{1}\left(S_{1}, S_{2}\right):=\sup \left\{\left\langle S_{1}-S_{2}, \psi\right\rangle: \psi \in C^{1}\left(\mathbb{R}^{n}\right),\|D \psi\|_{\infty} \leq 1\right\} \tag{2.10}
\end{equation*}
$$

Obviously, when $S_{1}=\sum_{i} \llbracket S_{1 i} \rrbracket, S_{2}=\sum_{i} \llbracket S_{2 i} \rrbracket \in \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$, the supremum in (2.10) can be taken over a suitable countable subset of $\psi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, chosen independently of the $S_{i}$ 's. Moreover, we have that

$$
\begin{equation*}
W_{1}\left(S_{1}, S_{2}\right)=\min _{\sigma \in \mathscr{P}_{Q}} \sum_{i}\left|S_{1 i}-S_{2 \sigma(i)}\right| \geq \min _{\sigma \in \mathscr{P}_{Q}}\left(\sum_{i}\left|S_{1 i}-S_{2 \sigma(i)}\right|^{2}\right)^{1 / 2}=\mathcal{G}\left(S_{1}, S_{2}\right) \tag{2.11}
\end{equation*}
$$

So $\mathcal{G}(g(x), g(y)) \leq C \delta_{11}^{1 / 2}|x-y|$ for a.e. $x, y \in K$. (The first equality in (2.11) is wellknown, but not easy to find in the literature. It can be derived by suitably modifying the arguments of [20, 4.1.12]. Another quick derivation is the following. Consider the set $\Pi$ of probability measures $\pi$ on $\mathbb{R}^{n} \times \mathbb{R}^{n}$ of the form $\sum_{i, j} c_{i j} \delta_{\left(S_{1 i}, S_{2 j}\right)}$, where the matrix of coefficients $c_{i j}$ consists of nonnegative entries with $\sum_{k} c_{k j}=1$ and $\sum_{k} c_{i k}=1$ for every $i$ and $j$, i.e. it is a doubly stochastic matrix. It then follows from the Kantorovich duality, see for instance [33, Theorem 1.14], that $W_{1}\left(S_{1}, S_{2}\right)=\min _{\pi \in \Pi} \int|x-y| d \pi(x, y)$. Observe however that $\int|x-y| d \pi(x, y)$ is a linear function of the coefficients $c_{i j}$ : the space of such matrices, also called Birkhoff polytope, is a compact convex set and so
the minimum is attained on the subset of extremal points. By the classical Birkhoffvon Neumann Theorem this set consists of the permutations matrices (see [7]) and so $\left.\min _{\pi \in \Pi} \int|x-y| d \pi=\min _{\sigma \in \mathscr{P}_{Q}} \sum_{i}\left|S_{1 i}-S_{2 \sigma(i)}\right|.\right)$

Next, write $g(x)=\sum_{i} \llbracket\left(h_{i}(x), \Psi\left(x, h_{i}(x)\right)\right) \rrbracket$. Obviously $x \mapsto h(x):=\sum_{i} \llbracket h_{i}(x) \rrbracket \in$ $\mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)$ is a Lipschitz map on $K$ with Lipschitz constant $\leq C \delta_{11}^{1 / 2}$. Recalling [14, Theorem 1.7], we can extend it to a map $\bar{u} \in \operatorname{Lip}\left(B_{3}, \mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)\right)$ satisfying $\operatorname{Lip}(\bar{u}) \leq C \delta_{11}^{1 / 2}$ and $\operatorname{osc}(\bar{u}) \leq C \operatorname{osc}(h)$. Set finally $u(x)=\sum_{i} \llbracket\left(\bar{u}_{i}(x), \Psi\left(x, \bar{u}_{i}(x)\right)\right) \rrbracket$. We start showing the Lipschitz bound. Fix $x_{1}, x_{2} \in B_{3}$ and assume, without loss of generality, that $\mathcal{G}\left(\bar{u}\left(x_{1}\right), \bar{u}\left(x_{2}\right)\right)^{2}=\sum_{i}\left|\bar{u}_{i}\left(x_{1}\right)-\bar{u}_{i}\left(x_{2}\right)\right|^{2}$. Then

$$
\begin{aligned}
\mathcal{G}\left(u\left(x_{1}\right), u\left(x_{2}\right)\right)^{2} & \leq \sum_{i}\left|\left(\bar{u}_{i}\left(x_{1}\right), \Psi\left(x_{1}, \bar{u}_{i}\left(x_{1}\right)\right)\right)-\left(\bar{u}_{i}\left(x_{2}\right), \Psi\left(x_{2}, \bar{u}_{i}\left(x_{2}\right)\right)\right)\right|^{2} \\
& \leq 2 \sum_{i}\left(\left(1+\left\|D_{y} \Psi\right\|_{0}^{2}\right)\left|\bar{u}_{i}\left(x_{1}\right)-\bar{u}_{i}\left(x_{2}\right)\right|^{2}+\left\|D_{x} \Psi\right\|_{0}^{2}\left|x_{1}-x_{2}\right|^{2}\right) \\
& \leq 2\left(1+\|D \Psi\|_{0}^{2}\right) \mathcal{G}\left(\bar{u}\left(x_{1}\right), \bar{u}\left(x_{2}\right)\right)^{2}+2\|D \Psi\|_{0}^{2}\left|x_{1}-x_{2}\right|^{2} \\
& \leq C\left(\delta_{11}+\|D \Psi\|_{0}^{2}\right)\left|x_{1}-x_{2}\right|^{2} .
\end{aligned}
$$

As for the $L^{\infty}$ bound, let $\eta>0$ be arbitrary and $p \in \mathbb{R}^{\bar{n}}$ be such that $\operatorname{osc}(\bar{u}) \leq$ $\sup _{x \in B_{3}} \mathcal{G}(\bar{u}(x), Q \llbracket p \rrbracket)+\eta$. Proceeding as above

$$
\begin{aligned}
\operatorname{osc}(u)^{2} & \leq \sup _{x \in B_{3}} \mathcal{G}(u(x), Q \llbracket(p, \Psi(0, p)) \rrbracket)^{2} \\
& \leq 2 \sup _{x \in B_{3}}\left(\left(1+\|D \Psi\|_{0}^{2}\right) \mathcal{G}(\bar{u}(x), Q \llbracket p \rrbracket)^{2}+\|D \Psi\|_{0}^{2}|x|^{2}\right) \\
& \leq 4\left(1+\|D \Psi\|_{0}^{2}\right)\left(\operatorname{osc}(\bar{u})^{2}+\eta^{2}\right)+18\|D \Psi\|_{0}^{2} .
\end{aligned}
$$

Since $\operatorname{osc}(h) \leq \mathbf{h}\left(T, \mathbf{C}_{4}, \pi_{0}\right)$, the estimate on $\operatorname{osc}(u)$ follows letting $\eta \downarrow 0$.
The identity $\mathbf{G}_{u}\left\llcorner\left(K \times \mathbb{R}^{n}\right)=T\left\llcorner\left(K \times \mathbb{R}^{n}\right)\right.\right.$ is a consequence of $u(x)=T_{x}$ for a.e. $x \in K$. Indeed, recall that both $T$ and $\mathbf{G}_{u}$ are rectifiable and observe that $\left\langle\vec{T}, \vec{\pi}_{0}\right\rangle \neq 0\|T\|$-a.e. on $K \times \mathbb{R}^{n}$, because $\mathbf{m e}_{T}<\infty$ on $K$. Similarly, $\left\langle\overrightarrow{\mathbf{G}}_{u}, \vec{\pi}_{0}\right\rangle \neq 0\left\|\mathbf{G}_{u}\right\|$-a.e. on $K \times \mathbb{R}^{n}$, by [17, Proposition 1.4]. Thus, $\left(\mathbf{G}_{u}-T\right)\left\llcorner K \times \mathbb{R}^{n}=0\right.$ if and only if $\left(\mathbf{G}_{u}-T\right)\left\llcorner d x \mathbf{1}_{K \times \mathbb{R}^{n}}=0\right.$. The latter identity follows from the slicing formula and the property $\langle T, \mathbf{p}, x\rangle=\left\langle\mathbf{G}_{u}, \mathbf{p}, x\right\rangle=$ $\sum_{i} \delta_{\left(x, u_{i}(x)\right)}$, valid for a.e. $x \in K$.

Finally, for each $x \in B_{r} \backslash K$ choose a ball $x \in B^{x}=B_{r(x)}(y(x)) \subset B_{4}$ such that $\mathbf{e}_{T}\left(B^{x}\right) \geq$ $2^{-m} \delta_{11} \omega_{m} r(x)^{m}$. By the $5 r$-Covering theorem, we choose balls $\hat{B}^{i}=B_{5 r\left(x_{i}\right)}\left(y\left(x_{i}\right)\right)$ which cover $B_{r} \backslash K$ and such that the balls $B^{x_{i}}$ are pairwise disjoint. We then conclude

$$
\begin{equation*}
\left|B_{r} \backslash K\right| \leq 10^{m} \delta_{11}^{-1} \mathbf{e}_{T}\left(\bigcup_{i} B^{x_{i}}\right) \tag{2.12}
\end{equation*}
$$

Fix $y \in B^{x_{i}}$. Since $B^{x_{i}} \subset B_{4}$, we have $2^{-m} \delta_{11} \omega_{m} r\left(x_{i}\right)^{m} \leq \mathbf{e}_{T}\left(B^{x_{i}}\right) \leq \mathbf{e}_{T}\left(B_{4}\right)=4^{m} \omega_{m} E$, which implies $2 r\left(x_{i}\right) \leq r_{0}<1$. Thus, $y \in B_{r+r_{0}} \subset B_{4}$. By definition of me $\mathbf{m e}_{T}$ we obviously have $\mathbf{m e}_{T}(y) \geq 2^{-m} \delta_{11}$. So $\cup_{i} B^{x_{i}} \subset B_{r+r_{0}} \cap\left\{\mathbf{m e}_{T}>2^{-m} \delta_{11}\right\}$ and (2.12) implies (2.1).

## 3. Patching multiple valued graphs

In this section we prove some complementary results to the theory of multiple valued functions as exposed in $[14,17]$. In particular, we show here a concentration compactness principle for $Q$-valued functions, and give an algorithm to construct suitable competitors for the Dirichlet energy, which will be also used in [16]. We first introduce some terminology.

Definition 3.1 (Translating sheets). Let $\Omega \subset \mathbb{R}^{m}$ be a bounded open set. A sequence of maps $\left\{h_{k}\right\}_{i \in \mathbb{N}} \subset W^{1,2}\left(\Omega, \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)\right)$ is called a sequence of translating sheets if there are:
(a) integers $J \geq 1$ and $Q_{1}, \ldots, Q_{J} \geq 1$ satisfying $\sum_{j=1}^{J} Q_{j}=Q$,
(b) vectors $y_{k}^{j} \in \mathbb{R}^{n}$ (for $j \in\{1, \ldots, J\}$ and $k \in \mathbb{N}$ ) with

$$
\begin{equation*}
\lim _{k}\left|y_{k}^{j}-y_{k}^{i}\right|=+\infty \quad \forall i \neq j \tag{3.1}
\end{equation*}
$$

(c) and maps $\zeta^{j} \in W^{1,2}\left(\Omega, \mathcal{A}_{Q_{j}}\right)$ for $j \in\{1, \ldots, J\}$,
such that $h_{k}=\sum_{j=1}^{J} \llbracket \tau_{y_{k}^{j}} \circ \zeta^{j} \rrbracket$, where for any generic $y \in \mathbb{R}^{n}$ we denote by $\tau_{y}: \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right) \rightarrow$ $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ the translation map (cp. [14, Section 3.3.3])

$$
\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right) \ni T=\sum_{i} \llbracket P_{i} \rrbracket \mapsto \tau_{y}(T):=\sum_{i} \llbracket P_{i}-y \rrbracket \in \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)
$$

Remark 3.2. Assume that $h_{k}, Q_{j}, y_{k}^{j}$ and $\zeta^{k}$ satisfy all the requirements of Definition 3.1 except for (3.1). Up to subsequences and relabellings, assume that $y_{k}^{1}-y_{k}^{2}$ converges to a vector $2 \bar{y}$. We can replace

- the integers $Q_{1}$ and $Q_{2}$ with $Q^{\prime}=Q_{1}+Q_{2}$;
- the vectors $y_{k}^{1}$ and $y_{2}^{k}$ with $y_{k}^{\prime}=\left(y_{k}^{1}+y_{k}^{2}\right) / 2$;
- the maps $\zeta^{1}$ and $\zeta^{2}$ with $\zeta^{\prime}:=\llbracket \tau_{\bar{y}} \circ \zeta^{1} \rrbracket+\llbracket \tau_{-\bar{y}} \circ \zeta^{2} \rrbracket$.

The new collections $Q^{\prime}, Q_{3}, \ldots, Q_{J}, y_{k}^{\prime}, y_{k}^{3}, \ldots, y_{k}^{J}$ and $\zeta^{\prime}, \zeta^{3}, \ldots, \zeta^{J}$, and the function $h_{k}^{\prime}:=$ $\llbracket \zeta^{\prime} \rrbracket+\sum_{j=3}^{J} \llbracket \zeta^{j} \rrbracket$, satisfy again all the requirements of Definition 3.1 except, possibly, for (3.1). Moreover, $\left\|\mathcal{G}\left(h_{k}^{\prime}, h_{k}\right)\right\|_{L^{2}} \rightarrow 0$ and $\left|D h_{k}^{\prime}\right|=\left|D h_{k}\right|$. Obviously, we can iterate this procedure only a finite number of times, obtaining a subsequence of translating sheets $\hat{h}_{k}$ asymptotic to $h_{k}$ in the $L^{2}$ distance with $\left|D \hat{h}_{k}\right|=\left|D h_{k}\right|$.
3.1. Concentration compactness. Translating sheets give a useful device to recover a suitable "compactness statement" for sequences of maps with equi-bounded energy.

Proposition 3.3 (Concentration compactness). Let $\Omega \subset \mathbb{R}^{m}$ be a Lipschitz bounded open set and $\left(g_{k}\right)_{k \in \mathbb{N}} \subset W^{1,2}\left(\Omega, \mathcal{A}_{Q}\right)$ a sequence of functions with $\sup _{k} \int_{\Omega}\left|D g_{k}\right|^{2}<\infty$. Then, there exist a subsequence (not relabeled) and a sequence of translating sheets $h_{k}$ such that $\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}} \rightarrow 0$ and the following inequalities hold for every open $\Omega^{\prime} \subset \Omega$ and any
sequence of measurable sets $J_{k}$ with $\left|J_{k}\right| \rightarrow 0$ :

$$
\begin{gather*}
\liminf _{k \rightarrow+\infty}\left(\int_{\Omega^{\prime} \backslash J_{k}}\left|D g_{k}\right|^{2}-\int_{\Omega^{\prime}}\left|D h_{k}\right|^{2}\right) \geq 0  \tag{3.2}\\
\limsup _{k \rightarrow+\infty} \int_{\Omega}\left(\left|D g_{k}\right|-\left|D h_{k}\right|\right)^{2} \leq \limsup _{k} \int_{\Omega}\left(\left|D g_{k}\right|^{2}-\left|D h_{k}\right|^{2}\right) . \tag{3.3}
\end{gather*}
$$

Proof. We start proving, by induction on $Q$, the existence of translating sheets $\left\{h_{k}\right\}$ (and a subsequence) with $\left\|\mathcal{G}\left(h_{k}, g_{k}\right)\right\|_{L^{2}} \rightarrow 0$ and satisfying the following additional property. If $J, Q_{j}, y_{k}^{j}$ and $\zeta^{j}$ are as in Definition 3.1, then there are $Q_{j}$ valued functions $w_{k}^{j}$ such that, after setting $f_{k}=\sum_{j} \llbracket w_{k}^{j} \rrbracket$, we have

$$
\begin{equation*}
\left\|\mathcal{G}\left(f_{k}, g_{k}\right)\right\|_{L^{2}}+\left|\left\{g_{k} \neq f_{k}\right\}\right| \rightarrow 0, \quad\left\|\mathcal{G}\left(\tau_{-y_{k}^{j}} \circ w_{k}^{j}, \zeta^{j}\right)\right\|_{L^{2}} \rightarrow 0 \quad \text { and } \quad\left|D f_{k}\right| \leq\left|D g_{k}\right| . \tag{3.4}
\end{equation*}
$$

If $Q=1$ the claim with $f_{k}=g_{k}$ is an easy corollary of the Poincaré inequality and the compact embedding $W^{1,2} \hookrightarrow L^{2}$. Assuming that the claim holds for any $Q^{*}<Q$, we prove it for $Q$. By the generalized Poincaré inequality [14, Proposition 2.12], there exist points $\bar{g}_{k} \in \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ and a real number $M$ such that

$$
\int_{\Omega} \mathcal{G}\left(g_{k}, \bar{g}_{k}\right)^{2} \leq C \int_{\Omega}\left|D g_{k}\right|^{2} \leq M<\infty \quad \forall k \in \mathbb{N}
$$

Recall the separation $s(T)$ and the diameter $d(T)$ of a point $T=\sum_{i} \llbracket P_{i} \rrbracket$ introduced in [14, Definition 3.4]: $s(T):=\min \left\{\left|P_{i}-P_{j}\right|: P_{i} \neq P_{j}\right\}$ and $d(T):=\max \left\{\left|P_{i}-P_{j}\right|\right\}$. We distinguish between to cases.

Case 1: $\lim \inf _{k} d\left(\bar{g}_{k}\right)<\infty$. After passing to a subsequence, we find $y_{k} \in \mathbb{R}^{n}$ such that the functions $\tau_{y_{k}} \circ g_{k}$ are equi-bounded in the $W^{1,2}$-metric. By the Sobolev embedding [14, Proposition 2.11], there exists a $Q$-valued map $\zeta \in W^{1,2}$ such that $\tau_{y_{k}} \circ g_{k} \rightarrow \zeta$ in $L^{2}(\Omega)$.
Case 2: $\lim _{k} d\left(\bar{g}_{k}\right)=+\infty$. By [14, Lemma 3.8] there are points $S_{k} \in \mathcal{A}_{Q}$ such that

$$
\beta d\left(\bar{g}_{k}\right) \leq s\left(S_{k}\right)<+\infty \quad \text { and } \quad \mathcal{G}\left(S_{k}, \bar{g}_{k}\right) \leq s\left(S_{k}\right) / 32
$$

where $\beta$ is a dimensional constant. Write $S_{k}=\sum_{i=1}^{J} \kappa_{i} \llbracket P_{k}^{i} \rrbracket$, with $P_{k}^{i} \neq P_{k}^{j}$ for $i \neq j$. Both $J$ and $\kappa_{i}$ may depend on $k$ but they have a finite range: therefore, after extracting a subsequence, we can assume that they do not depend on $k$. Set next $r_{k}=\frac{s\left(S_{k}\right)}{16}$ and let $\vartheta_{k}$ be the retraction of $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ into $\overline{B_{r_{k}}\left(S_{k}\right)}$ provided by [14, Lemma 3.7]. Clearly, the functions $\hat{f}_{k}=\vartheta_{k} \circ g_{k}$ satisfy $\left|D \hat{f}_{k}\right| \leq\left|D g_{k}\right|$ and there are $\kappa_{i}$-valued functions $z_{k}^{i}$ such that

$$
\hat{f}_{k}=\sum_{i=1}^{J} \llbracket z_{k}^{i} \rrbracket, \quad \text { with } \quad\left\|\mathcal{G}\left(z_{k}^{i}, \kappa_{i} \llbracket P_{k}^{i} \rrbracket\right)\right\|_{\infty} \leq r_{k} .
$$

Since $\kappa_{i}<Q$, we apply the inductive hypothesis to each sequence $\left(z_{k}^{i}\right)_{k}$ and, using Remark 3.2 reach a subsequence (not relabeled) of $\hat{f}_{k}$, a sequence of translating sheets $h_{k}$ and corresponding functions $f_{k}$ which satisfy (3.4) with $\hat{f}_{k}$ replacing $g_{k}$.

We next claim that (3.4) holds even for $g_{k}$, i.e. that $\lim _{k}\left(\left\|\mathcal{G}\left(f_{k}, g_{k}\right)\right\|_{L^{2}}+\left|\left\{f_{k} \neq g_{k}\right\}\right|\right)=$ 0 . To this aim, recall first that

$$
\left\{g_{k} \neq \hat{f}_{k}\right\}=\left\{\mathcal{G}\left(g_{k}, S_{k}\right)>r_{k}\right\} \subseteq\left\{\mathcal{G}\left(g_{k}, \bar{g}_{k}\right)>r_{k} / 2\right\} .
$$

Thus,

$$
\begin{equation*}
\left|\left\{g_{k} \neq \hat{f}_{k}\right\}\right| \leq\left|\left\{\mathcal{G}\left(g_{k}, \bar{g}_{k}\right)>r_{k} / 2\right\}\right| \leq \frac{C}{r_{k}^{2}} \int_{\left\{\mathcal{G}\left(g_{k}, \bar{g}_{k}\right)>\frac{r_{k}}{2}\right\}} \mathcal{G}\left(g_{k}, \bar{g}_{k}\right)^{2} \leq \frac{C M}{\left(d\left(\bar{g}_{k}\right)\right)^{2}} . \tag{3.5}
\end{equation*}
$$

Since $d\left(\bar{g}_{k}\right) \rightarrow+\infty$ and (3.4) holds with $\hat{f}_{k}$ replacing $g_{k}$, we conclude $\left|\left\{f_{k} \neq g_{k}\right\}\right| \rightarrow 0$. Next, since $\vartheta_{k}\left(\bar{g}_{k}\right)=\bar{g}_{k}$ and $\operatorname{Lip}\left(\vartheta_{k}\right)=1$, we have $\mathcal{G}\left(\hat{f}_{k}, \bar{g}_{k}\right) \leq \mathcal{G}\left(g_{k}, \bar{g}_{k}\right)$. Therefore, by the Sobolev embedding and the Poincaré inequality, for any $p \in] 2,2^{*}[$, we infer

$$
\begin{aligned}
& \int_{\Omega} \mathcal{G}\left(\hat{f}_{k}, g_{k}\right)^{2}=\int_{\left\{g_{k} \neq \hat{f}_{k}\right\}} \mathcal{G}\left(\hat{f}_{k}, g_{k}\right)^{2} \leq 2 \int_{\left\{\hat{f}_{k} \neq g_{k}\right\}} \mathcal{G}\left(\hat{f}_{k}, \bar{g}_{k}\right)^{2}+2 \int_{\left\{\hat{f}_{k} \neq g_{k}\right\}} \mathcal{G}\left(\bar{g}_{k}, g_{k}\right)^{2} \\
& \leq 4 \int_{\left\{\hat{f}_{k} \neq g_{k}\right\}} \mathcal{G}\left(\bar{g}_{k}, g_{k}\right)^{2} \leq C\left\|\mathcal{G}\left(g_{k}, \bar{g}_{k}\right)\right\|_{L^{p}}^{2}\left|\left\{\hat{f}_{k} \neq g_{k}\right\}\right|^{1-\frac{2}{p}} \stackrel{(3.5)}{\leq} \frac{C M^{1-2 / p}}{d\left(\bar{g}_{k}\right)^{2-4 / p}} \int_{\Omega}\left|D g_{k}\right|^{2} .
\end{aligned}
$$

Since $d\left(\bar{g}_{k}\right)$ diverges, this shows $\left\|\mathcal{G}\left(\hat{f}_{k}, g_{k}\right)\right\|_{L^{2}} \rightarrow 0$ and by inductive hypothesis that $\left\|\mathcal{G}\left(f_{k}, g_{k}\right)\right\|_{L^{2}} \rightarrow 0$.
We now show that (3.2) and (3.3) are consequences of (3.4). For each $j$ we consider the corresponding embedding $\boldsymbol{\xi}_{j}: \mathcal{A}_{Q_{j}}\left(\mathbb{R}^{n}\right) \rightarrow \mathbb{R}^{N\left(Q_{j}, n\right)}$ and, by a slight abuse of notation, we drop the $j$ subscript. Then, we conclude that $\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j} \rightarrow \boldsymbol{\xi} \circ \zeta^{j}$ in $L^{2}$ and $\left\|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)\right\|_{L^{2}}$ is a bounded sequence, from which

$$
\begin{equation*}
D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right) \rightharpoonup D\left(\boldsymbol{\xi} \circ \zeta^{j}\right) \quad \text { in } L^{2}(\Omega) . \tag{3.6}
\end{equation*}
$$

If $J_{k}$ is a sequence of measurable sets with $\left|J_{k}\right| \downarrow 0$, then $\mathbf{1}_{\Omega^{\prime} \backslash J_{k}} \rightarrow \mathbf{1}_{\Omega^{\prime}}$ in $L^{2}(\Omega)$ and it follows from (3.6) that

$$
D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right) \mathbf{1}_{\Omega^{\prime} \backslash J_{k}} \rightharpoonup D\left(\boldsymbol{\xi} \circ \zeta^{j}\right) \mathbf{1}_{\Omega^{\prime}} \quad \text { in } L^{2}(\Omega),
$$

and, hence,

$$
\operatorname{Dir}\left(\zeta^{j}, \Omega^{\prime}\right)=\int_{\Omega^{\prime}}\left|D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right|^{2} \leq \liminf _{k} \int_{\Omega^{\prime} \backslash J_{k}}\left|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)\right|^{2}=\underset{k}{\lim \inf } \int_{\Omega^{\prime} \backslash J_{k}}\left|D w_{k}^{j}\right|^{2} .
$$

Summing over $j$, we obtain (3.2). As for (3.3), set $J_{k}:=\left\{g_{k} \neq f_{k}\right\}$. Thus,

$$
\begin{align*}
& \int_{\Omega \backslash J_{k}}\left(\left|D g_{k}\right|-\left|D h_{k}\right|\right)^{2} \leq \sum_{j} \int_{\Omega \backslash J_{k}}\left(\left|D w_{k}^{j}\right|-\left|D \zeta^{j}\right|\right)^{2} \\
= & \sum_{j} \int_{\Omega \backslash J_{k}}\left(\left|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)\right|-\left|D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right|\right)^{2} \leq \sum_{j} \int_{\Omega \backslash J_{k}}\left|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)-D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right|^{2} \\
= & \sum_{j} \int_{\Omega \backslash J_{k}}\left(\left|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)\right|^{2}+\left|D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right|^{2}-2 D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right) \cdot D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right) . \tag{3.7}
\end{align*}
$$

Therefore, by (3.6) (and taking into account that $\left|J_{k}\right| \rightarrow 0$ ) one gets

$$
\begin{align*}
& \limsup _{k \rightarrow+\infty} \int_{\Omega \backslash J_{k}}\left(\left|D g_{k}\right|-\left|D h_{k}\right|\right)^{2} \\
\leq & \lim _{k \rightarrow+\infty} \sum_{j} \int_{\Omega \backslash J_{k}}\left(\left|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)\right|^{2}+\left|D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right|^{2}-2 D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right) \cdot D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right) \\
= & \limsup _{k \rightarrow+\infty} \int_{\Omega \backslash J_{k}} \sum_{j}\left|D\left(\boldsymbol{\xi} \circ \tau_{-y_{k}^{j}} \circ w_{k}^{j}\right)\right|^{2}-\int_{\Omega} \sum_{j}\left|D\left(\boldsymbol{\xi} \circ \zeta^{j}\right)\right|^{2} \\
= & \limsup _{k \rightarrow+\infty} \int_{\Omega \backslash J_{k}}\left|D g_{k}\right|^{2}-\int_{\Omega}\left|D h_{k}\right|^{2} . \tag{3.8}
\end{align*}
$$

On the other hand, since $\left|J_{k}\right| \rightarrow 0$ we conclude

$$
\limsup _{k \rightarrow \infty} \int_{J_{k}}\left(\left|D g_{k}\right|-\left|D h_{k}\right|\right)^{2}=\limsup _{k \rightarrow \infty} \int_{J_{k}}\left|D g_{k}\right|^{2}
$$

Observe that, after passing to a subsequence, we can actually assume that all limsups are in fact limits. Summing (3.8) and the last equation we then conclude (3.3).
3.2. Dirichlet competitors. We consider next a standard procedure to construct competitors for the Dirichlet energy of a sequence of functions with equi-bounded energy.

Proposition 3.4 (Construction of a competitor). Consider two radii $1 \leq r_{0}<r_{1}<4$ and maps $g_{k}, h_{k} \in W^{1,2}\left(B_{r_{1}}, \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)\right)$ such that $\left\{h_{k}\right\}_{k}$ is a sequence of translating sheets,

$$
\sup _{k} \operatorname{Dir}\left(g_{k}, B_{r_{1}}\right)<+\infty \quad \text { and } \quad\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}\left(B_{r_{1}} \backslash B_{r_{0}}\right)} \rightarrow 0
$$

For every $\eta>0$, there exist $r \in] r_{0}, r_{1}\left[\right.$, a subsequence of $\left\{g_{k}\right\}_{k}$ (not relabeled) and functions $H_{k} \in W^{1,2}\left(B_{r_{1}}, \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)\right)$ such that $\left.H_{k}\right|_{B_{r_{1}} \backslash B_{r}}=\left.g_{k}\right|_{B_{r_{1}} \backslash B_{r}}$ and $\operatorname{Dir}\left(H_{k}, B_{r_{1}}\right) \leq \operatorname{Dir}\left(h_{k}, B_{r_{1}}\right)+$ $\eta$. In addition, there is a dimensional constant $C$ and a constant $C^{*}$ (depending on $\eta$ and the two sequences, but not on $k$ ) such that

$$
\begin{gather*}
\operatorname{Lip}\left(H_{k}\right) \leq C^{*}\left(\operatorname{Lip}\left(g_{k}\right)+1\right),  \tag{3.9}\\
\left\|\mathcal{G}\left(H_{k}, h_{k}\right)\right\|_{L^{2}\left(B_{r}\right)} \leq C \operatorname{Dir}\left(g_{k}, B_{r}\right)+C \operatorname{Dir}\left(H_{k}, B_{r}\right),  \tag{3.10}\\
\left\|\boldsymbol{\eta} \circ H_{k}\right\|_{L^{1}\left(B_{r_{1}}\right)} \leq C^{*}\left\|\boldsymbol{\eta} \circ g_{k}\right\|_{L^{1}\left(B_{r_{1}}\right)}+C\left\|\boldsymbol{\eta} \circ h_{k}\right\|_{L^{1}\left(B_{r_{1}}\right)} . \tag{3.11}
\end{gather*}
$$

In order to prove the proposition, we need to recall the following two lemmas, which are slight variants of [14, Proposition 4.4] and [14, Lemma 2.15].

Lemma 3.5 (Lipschitz approximation). Let $f \in W^{1,2}\left(B_{r}, \mathcal{A}_{Q}\right)$. Then, for every $\varepsilon>0$, there exists $f_{\varepsilon} \in \operatorname{Lip}\left(B_{r}, \mathcal{A}_{Q}\right)$ such that

$$
\begin{equation*}
\int_{B_{r}} \mathcal{G}\left(f, f_{\varepsilon}\right)^{2}+\int_{B_{r}}\left(|D f|-\left|D f_{\varepsilon}\right|\right)^{2}+\int_{B_{r}}\left(|D(\boldsymbol{\eta} \circ f)|-\left|D\left(\boldsymbol{\eta} \circ f_{\varepsilon}\right)\right|\right)^{2} \leq \varepsilon \tag{3.12}
\end{equation*}
$$

If $\left.f\right|_{\partial B_{r}} \in W^{1,2}\left(\partial B_{r}, \mathcal{A}_{Q}\right)$, then $f_{\varepsilon}$ can be chosen to satisfy also

$$
\begin{equation*}
\int_{\partial B_{r}} \mathcal{G}\left(f, f_{\varepsilon}\right)^{2}+\int_{\partial B_{r}}\left(|D f|-\left|D f_{\varepsilon}\right|\right)^{2} \leq \varepsilon \tag{3.13}
\end{equation*}
$$

Proof. By an obvious scaling argument we can assume $r=1$. We start noticing that (3.12) is a corollary of [14, Proposition 4.4]. On the other hand, if $\left.f\right|_{\partial B_{1}} \in W^{1,2}\left(\partial B_{1}\right)$, we extend the map to $B_{2}$ by setting $f(x)=f\left(\frac{x}{|x|}\right)$ if $|x| \geq 1$. We then can apply [14, Proposition 2.5] to find a sequence of Lipschitz maps $f_{k}$ such that $f_{k} \rightarrow f$ strongly in $W^{1,2}\left(B_{2}\right)$. Given $\delta>0$, define the maps $f^{\delta}(x)=f((1+\delta) x)$ and $f_{k}^{\delta}(x)=f_{k}((1+\delta) x)$. Obviously, $f_{k}^{\delta} \rightarrow f^{\delta}$ strongly in $W^{1,2}\left(B_{1}\right)$ and $f^{\delta} \rightarrow f$ strongly in $W^{1,2}\left(B_{1}\right)$ as $\delta \downarrow 0$. By a standard Fubini argument, for each $j$ we can find a $\delta_{j}<\frac{1}{j}$ and a subsequence $\left\{f_{k, j}\right\}_{k}$ such that $\left.\left.f_{k, j}\right|_{\partial B_{1+\delta_{j}}} \rightarrow f\right|_{\partial B_{1+\delta_{j}}}$ (i.e. $\left.\left.\left.f_{k, j}^{\delta_{j}}\right|_{\partial B_{1}} \rightarrow f^{\delta_{j}}\right|_{\partial B_{1}}=\left.f\right|_{\partial B_{1}}\right)$ strongly in $W^{1,2}\left(\partial B_{1+\delta_{j}}\right)$ as $k \uparrow \infty$. By standard diagonal argument we can arrange the subsequences so that $\left\{f_{k, j}\right\} \supset\left\{f_{k, j+1}\right\}$. Thus, a suitable diagonal sequence $\bar{f}_{j}:=f_{k(j), j}^{\delta_{j}}$ has the property that $\bar{f}_{j} \rightarrow f$ in $W^{1,2}\left(B_{1}\right)$ and $\left.\left.\bar{f}_{j}\right|_{\partial B_{1}} \rightarrow f\right|_{\partial B_{1}}$ in $W^{1,2}\left(\partial B_{1}\right)$.
Lemma 3.6 (Interpolation). There exists a constant $C_{0}=C_{0}(m, n, Q)>0$ with the following property. Assume $r \in] 1,3\left[, f \in W^{1,2}\left(B_{r}, \mathcal{A}_{Q}\right)\right.$ and $g \in W^{1,2}\left(\partial B_{r}, \mathcal{A}_{Q}\right)$ are given maps such that $\left.f\right|_{\partial B_{r}} \in W^{1,2}\left(\partial B_{r}, \mathcal{A}_{Q}\right)$. Then, for every $\left.\varepsilon \in\right] 0, r[$ there exists a function $h \in W^{1,2}\left(B_{r}, \mathcal{A}_{Q}\right)$ such that $\left.h\right|_{\partial B_{r}}=g$ and

$$
\begin{align*}
\int_{B_{r}}|D h|^{2} \leq & \int_{B_{r}}|D f|^{2}+\varepsilon \int_{\partial B_{r}}\left(\left|D_{\tau} f\right|^{2}+\left|D_{\tau} g\right|^{2}\right)+\frac{C_{0}}{\varepsilon} \int_{\partial B_{r}} \mathcal{G}(f, g)^{2}  \tag{3.14}\\
& \operatorname{Lip}(h) \leq C_{0}\left\{\operatorname{Lip}(f)+\operatorname{Lip}(g)+\varepsilon^{-1} \sup _{\partial B_{r}} \mathcal{G}(f, g)\right\}  \tag{3.15}\\
& \int_{B_{r}}|\boldsymbol{\eta} \circ h| \leq C_{0} \int_{\partial B_{r}}|\boldsymbol{\eta} \circ g|+C_{0} \int_{B_{r}}|\boldsymbol{\eta} \circ f| \tag{3.16}
\end{align*}
$$

(here $D_{\tau}$ denotes the tangential derivative).
Proof. The first conclusion is an obvious corollary of [14, Lemma 2.15]. It is then straightforward to see that the map constructed in the proof of [14, Lemma 2.15] satisfies also (3.15). As for the final claim, let $\bar{g}:=\sum \llbracket g_{i}-\boldsymbol{\eta} \circ g \rrbracket, \bar{f}:=\sum \llbracket f_{i}-\boldsymbol{\eta} \circ f \rrbracket$ and consider the interpolation map $\bar{h}$ between $\bar{f}$ and $\bar{g}$ given by [14, Lemma 2.15]. Set $\hat{h}=\sum_{i} \llbracket \bar{h}_{i}-\boldsymbol{\eta} \circ \bar{h} \rrbracket$ and observe that $\operatorname{Lip}(\hat{h}) \leq \operatorname{Lip}(\bar{h})$ and $\operatorname{Dir}(\hat{h}) \leq \operatorname{Dir}(\bar{h})$. We apply again [14, Lemma 2.15] in the case $Q=1$ to $\boldsymbol{\eta} \circ f$ and $\boldsymbol{\eta} \circ g$, and get the interpolation $u$. It is then easy to check that the map $h:=\sum_{i} \llbracket \hat{h}_{i}+u \rrbracket$ has all the desired properties.
Proof of Proposition 3.4. Set for simplicity $A_{k}:=\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}\left(B_{r_{1}} \backslash B_{r_{0}}\right)}$ and $B_{k}:=\| \boldsymbol{\eta} \circ$ $g_{k} \|_{L^{1}\left(B_{r_{1}}\right)}$. If $A_{k} \equiv 0$, then there is nothing to prove and so we can assume that, for a subsequence, not relabeled, $A_{k}>0$. Assuming that for yet another subsequence (not relabeled) $B_{k}>0$, we consider the function

$$
\begin{equation*}
\psi_{k}(r):=\int_{\partial B_{r}}\left(\left|D g_{k}\right|^{2}+\left|D h_{k}\right|^{2}\right)+A_{k}^{-2} \int_{\partial B_{r}} \mathcal{G}\left(g_{k}, h_{k}\right)^{2}+B_{k}^{-1} \int_{\partial B_{r}}\left|\boldsymbol{\eta} \circ g_{k}\right| . \tag{3.17}
\end{equation*}
$$

By assumption $\lim \inf _{k} \int_{r_{0}}^{r_{1}} \psi_{k}(r) d r<\infty$. So, by Fatou's Lemma, there is $\left.r \in\right] r_{0}, r_{1}$ [ and a subsequence, not relabeled, such that $\lim _{k} \psi_{k}(r)<\infty$. Thus, for some $M>0$ we have

$$
\begin{gather*}
\int_{\partial B_{r}} \mathcal{G}\left(g_{k}, h_{k}\right)^{2} \rightarrow 0  \tag{3.18}\\
\operatorname{Dir}\left(h_{k}, \partial B_{r}\right)+\operatorname{Dir}\left(g_{k}, \partial B_{r}\right) \leq M  \tag{3.19}\\
\int_{\partial B_{r}}\left|\boldsymbol{\eta} \circ g_{k}\right| \leq M\left\|\boldsymbol{\eta} \circ g_{k}\right\|_{L^{1}\left(B_{r_{1}}\right)} . \tag{3.20}
\end{gather*}
$$

In case $B_{k}=0$ for all $k$ large enough, we define $\psi_{k}$ dropping the last summand in (3.17) and reach the same conclusion.

Let $\zeta^{j}$ be the blocks of the translating sheets $h_{k}$ as in Definition 3.1. We apply Lemma 3.5 to each $\zeta^{j}$ and find Lipschitz functions $\zeta_{\eta}^{j}$ satisfying the conclusion of the lemma with $\bar{\varepsilon}_{1}=$ $\bar{\varepsilon}_{1}(\eta, M)>0$ (which will be chosen later). We also choose a standard radial convolution kernel $\varphi$ in $\mathbb{R}^{m}$ and a small parameter $\bar{\rho}$ (also to be chosen later). Then, set

$$
h_{k, \eta}:=\sum_{j=1}^{J} \llbracket \tau_{y_{k}^{j}} \circ \zeta_{\eta}^{j} \rrbracket \quad \text { and } \quad \bar{h}_{k, \eta}:=\sum_{i=1}^{Q} \llbracket\left(h_{k, \eta}\right)_{i}-\boldsymbol{\eta} \circ h_{k, \eta}+\left(\boldsymbol{\eta} \circ h_{k}\right) * \varphi_{\bar{\rho}} \rrbracket \text {, }
$$

and choose $\bar{\rho}$ so small that

$$
\begin{gather*}
Q^{2}\left\|\boldsymbol{\eta} \circ h_{k}-\left(\boldsymbol{\eta} \circ h_{k}\right) * \varphi_{\bar{\rho}}\right\|_{L^{2}}^{2} \leq \bar{\varepsilon}_{1},  \tag{3.21}\\
\int_{B_{r}}\left(\left|D\left(\boldsymbol{\eta} \circ h_{k}\right)\right|^{2}-\left|D\left(\boldsymbol{\eta} \circ h_{k} * \varphi_{\bar{\rho}}\right)\right|^{2}\right) \leq \bar{\varepsilon}_{1} . \tag{3.22}
\end{gather*}
$$

Note that this is possible because, from the fact that $h_{k}$ is a sequence of translating sheets, it follows that $\eta \circ h_{k}(x)=F(x)+p_{k}$ for some $F \in W^{1,2}$ and a sequence of vectors $p_{k} \in \mathbb{R}^{n}$. Therefore $\left(\eta \circ h_{k}\right) * \varphi_{\bar{\rho}}=F * \varphi_{\bar{\rho}}+p_{k}$ and $D\left(\eta \circ h_{k}\right) * \varphi_{\bar{\rho}}=D F * \varphi_{\bar{\rho}}$, and (3.21) and (3.22) follows if $\bar{\rho}$ is sufficiently small by the usual convolution estimates. In particular by very rough estimates,

$$
\begin{gather*}
\left\|\mathcal{G}\left(g_{k}, \bar{h}_{k, \eta}\right)\right\|_{L^{2}} \stackrel{(3.21)}{\leq}\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}}+2\left\|\mathcal{G}\left(h_{k}, h_{k, \eta}\right)\right\|_{L^{2}}+\bar{\varepsilon}_{1} \leq o(1)+3 \bar{\varepsilon}_{1}  \tag{3.23}\\
\operatorname{Dir}\left(\bar{h}_{k, \eta}, \partial B_{r}\right) \leq 2 M+2 \bar{\varepsilon}_{1} \tag{3.24}
\end{gather*}
$$

and

$$
\begin{aligned}
\operatorname{Dir}\left(\bar{h}_{k, \eta}, B_{r}\right)= & \sum_{i} \int_{B_{r}}\left|D\left(h_{k, \eta}\right)_{i}-D\left(\boldsymbol{\eta} \circ h_{k, \eta}\right)+D\left(\boldsymbol{\eta} \circ h_{k} * \varphi_{\bar{\rho}}\right)\right|^{2} \\
= & \int_{B_{r}}\left(\left|D h_{k, \eta}\right|^{2}-Q\left|D\left(\boldsymbol{\eta} \circ h_{k, \eta}\right)\right|^{2}+Q\left|D\left(\boldsymbol{\eta} \circ h_{k} * \varphi_{\bar{\rho}}\right)\right|^{2}\right) \\
= & \operatorname{Dir}\left(h_{k, \eta}, B_{r}\right)+Q \int_{B_{r}}\left(\left|D\left(\boldsymbol{\eta} \circ h_{k}\right)\right|^{2}-\left|D\left(\boldsymbol{\eta} \circ h_{k, \eta}\right)\right|^{2}\right) \\
& +Q \int_{B_{r}}\left(\left|D\left(\boldsymbol{\eta} \circ h_{k} * \varphi_{\bar{\rho}}\right)\right|^{2}-\left|D\left(\boldsymbol{\eta} \circ h_{k}\right)\right|^{2}\right)
\end{aligned}
$$

$$
\begin{equation*}
\stackrel{(3.12),(3.22)}{\leq} \operatorname{Dir}\left(h_{k, \eta}, B_{r}\right)+2 Q \bar{\varepsilon}_{1} \tag{3.25}
\end{equation*}
$$

We can then apply Lemma 3.6 to $\bar{h}_{k, \eta}$ and $g_{k}$ with $\bar{\varepsilon}_{2}=\bar{\varepsilon}_{2}(\eta, M)>0$, and get (up to subsequences) maps $H_{k}$ satisfying $\left.H_{k}\right|_{\partial B_{r}}=\left.g_{k}\right|_{\partial B_{r}}$ and

$$
\begin{aligned}
\operatorname{Dir}\left(H_{k}, B_{r}\right) & \leq \operatorname{Dir}\left(\bar{h}_{k, \eta}, B_{r}\right)+\bar{\varepsilon}_{2} \operatorname{Dir}\left(\bar{h}_{k, \eta}, \partial B_{r}\right)+\bar{\varepsilon}_{2} \operatorname{Dir}\left(g_{k}, \partial B_{r}\right)+\frac{C_{0}}{\bar{\varepsilon}_{2}} \int_{\partial B_{r}} \mathcal{G}\left(\bar{h}_{k, \eta}, g_{k}\right)^{2} \\
& \leq \operatorname{Dir}\left(h_{k}, B_{r}\right)+Q \bar{\varepsilon}_{1}+3 \bar{\varepsilon}_{2}\left(M+\bar{\varepsilon}_{1}\right)+3 C_{0} \bar{\varepsilon}_{2}^{-1} \bar{\varepsilon}_{1}
\end{aligned}
$$

where in the last line we have used (3.18), (3.19) and (3.23) - (3.25). An appropriate choice of the parameters $\varepsilon_{1}$ and $\varepsilon_{2}$ gives the desired bound $\operatorname{Dir}\left(H_{k}, B_{r}\right) \leq \operatorname{Dir}\left(h_{k}, B_{r}\right)+\eta$.

Observe next that, by construction, $\lim \sup _{k} \operatorname{Lip}\left(\bar{h}_{k, \eta}\right) \leq C^{*}$, for some constant which depends on $\eta$ and the two sequences, but not on $k$. Moreover,

$$
\left\|\mathcal{G}\left(\bar{h}_{k, \eta}, g_{k}\right)\right\|_{L^{\infty}\left(\partial B_{r}\right)} \leq\left\|\mathcal{G}\left(\bar{h}_{k, \eta}, g_{k}\right)\right\|_{L^{2}\left(\partial B_{r}\right)}+C \operatorname{Lip}\left(g_{k}\right)+C \operatorname{Lip}\left(\bar{h}_{k, \eta}\right) .
$$

Thus (3.9) follows from (3.15).
Finally, (3.10) follows from the Poincaré inequality applied to $\mathcal{G}\left(H_{k}, g_{k}\right)$ (which vanishes identically on $\partial B_{r}$ ), and (3.11) follows from (3.16), because of (3.20) and $\left\|\boldsymbol{\eta} \circ \bar{h}_{k, \eta}\right\|_{L^{1}\left(B_{r}\right)}=$ $\left\|\left(\boldsymbol{\eta} \circ h_{k}\right) * \varphi_{\bar{\rho}}\right\|_{L^{1}\left(B_{r}\right)} \leq\left\|\boldsymbol{\eta} \circ h_{k}\right\|_{L^{1}\left(B_{r_{1}}\right)}$ if $\bar{\rho}$ is also chosen small enough such that $r+\bar{\rho}<r_{1}$.

## 4. Harmonic approximation

In what follows we will always apply Proposition 2.2 with $\delta_{11}=E^{2 \beta}$ and under a certain scaling of A.

Definition 4.1 ( $E^{\beta}$-Lipschitz approximation). Let $\beta \in\left(0, \frac{1}{2 m}\right), T$ be as in Proposition 2.2 such that $32 E^{(1-2 \beta) / m}<1$ and $s \mathbf{A} \leq E^{1 / 4+\delta}$ for some $\delta>0$. If the coordinates are fixed as in Remark 1.5, the map $u$ given by Proposition 2.2 for $\delta_{11}=E^{2 \beta}$ is then called the $E^{\beta}$-Lipschitz approximation of $T$ in $\mathbf{C}_{3 s}(x)$ and will be denoted by $f$.

In this section we prove that, if $T$ is also area minimizing, the corresponding $E^{\beta}{ }^{\beta}$ Lipschitz approximation is close to a Dir-minimizing function $w$. This comes with an $o(E)$-improvement of the estimates in Proposition 2.2.

Theorem 4.2 (First harmonic approximation). For every $\eta_{1}, \delta>0$ and every $\beta \in\left(0, \frac{1}{2 m}\right)$, there exist constants $\varepsilon_{12}, C_{12}>0$ with the following property. Let $T$ be as in Assumption 1.1 in $\mathbf{C}_{4 s}(x)$ and assume it is area minimizing. If $E=\mathbf{E}\left(T, \mathbf{C}_{4 s}(x)\right) \leq \varepsilon_{12}$ and $s \mathbf{A} \leq E^{1 / 4+\delta}$, then the $E^{\beta}$-Lipschitz approximation $f$ in $\mathbf{C}_{3 s}(x)$ satisfies

$$
\begin{equation*}
\int_{B_{2 s}(x) \backslash K}|D f|^{2} \leq \eta_{1} E \omega_{m}(4 s)^{m}=\eta_{1} \mathbf{e}_{T}\left(B_{4 s}(x)\right) . \tag{4.1}
\end{equation*}
$$

Moreover, if we consider the coordinates of Remark 1.5, there exists a Dir-minimizing function $u: B_{2 s}(x) \rightarrow \mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)$ such that the map $B_{2 s}(x) \ni y \mapsto w=(u, \Psi(y, u))$ satisfies

$$
\begin{gather*}
s^{-2} \int_{B_{2 s}(x)} \mathcal{G}(f, w)^{2}+\int_{B_{2 s}(x)}(|D f|-|D w|)^{2} \leq \eta_{1} E \omega_{m}(4 s)^{m}=\eta_{1} \mathbf{e}_{T}\left(B_{4 s}(x)\right)  \tag{4.2}\\
\int_{B_{2 s}(x)}|D(\boldsymbol{\eta} \circ f)-D(\boldsymbol{\eta} \circ w)|^{2} \leq \eta_{1} E \omega_{m}(4 s)^{m}=\eta_{1} \mathbf{e}_{T}\left(B_{4 s}(x)\right) \tag{4.3}
\end{gather*}
$$

Remark 4.3 (Isoperimetric inequality). If $S \subset \mathbb{R}^{m+n}$ is an integral current of dimension $m-1$ with $\partial S=0$, then there is an $m$-dimensional integral current $R \subset \mathbb{R}^{m+n}$ such that $\partial R=S$ and $\mathbf{M}(R) \leq C \mathbf{M}(S)^{m /(m-1)}$, where the constant $C$ is only dimensional (see [29, Theorem 30.1]). It is also well-known that, when $\operatorname{spt}(S) \subset \Sigma$ and $\Sigma$ is as in Assumption 1.1 the same inequality holds for some $\bar{R}$ with $\operatorname{spt}(\bar{R}) \subset \Sigma$ and $\partial \bar{R}=S$, with a dimensional constant $C$ which depends additionally on the constant $c_{0}$. This can be easily seen as follows: let $\mathbf{q}: \mathbb{R}^{m+n} \rightarrow \mathbb{R}^{m+\bar{n}}$ be the orthogonal projection and $\Lambda: \mathbb{R}^{m+n} \rightarrow \Sigma$ be the $\operatorname{map} \Lambda(p)=(\mathbf{q}(p), \Psi(\mathbf{q}(p))) . \Lambda$ is a global Lipschitz retraction of $\mathbb{R}^{m+n}$ onto $\Sigma$ which is the identity on $\Sigma$ : thus we can simply set $\bar{R}=\Lambda_{\sharp} R$.

Remark 4.4 (Taylor expansion of the mass). There are dimensional constants $c, C>0$ such that the following holds. Let $V \subset \mathbb{R}^{m}$ be a bounded measurable set and let $u: V \rightarrow$ $\mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ be a Lipschitz function with $\operatorname{Lip}(u) \leq c$. Denote by $\mathbf{G}_{u}$ the integer rectifiable current associated to the graph of $u$ as in [17, Definition 1.10]. Then, the following Taylor expansion of the mass of $\mathbf{G}_{u}$ holds:

$$
\mathbf{M}\left(\mathbf{G}_{u}\right)=Q|V|+\int_{V} \frac{|D u|^{2}}{2}+\int_{V} \sum_{i} R\left(D u_{i}\right)
$$

where $R: \mathbb{R}^{n \times m} \rightarrow R$ is a $C^{1}$ function satisfying $|R(D)|=|D|^{3} L(D)$ for some positive function $L$ such that $L(0)=0$ and $\operatorname{Lip}(L) \leq C$. This Taylor expansion is proven in [17, Corollary 3.3] (although the corollary is stated for $V$ open, the proof works obviously when $V$ is merely measurable).

Remark 4.5. There exists a dimensional constant $c>0$ such that, if $E \leq c$, then the $E^{\beta}$-Lipschitz approximation satisfies the following estimates:

$$
\begin{gather*}
\operatorname{Lip}(f) \leq C E^{\beta}  \tag{4.4}\\
\int_{B_{3 s}(x)}|D f|^{2} \leq C E s^{m} \tag{4.5}
\end{gather*}
$$

Indeed (4.4) follows from Proposition 2.2, Remark 1.5 and $\|D \Psi\|_{0} \leq C\left(E^{1 / 2}+\mathbf{A}\right) \leq C E^{\beta}$ by the choice of $\beta$ and the scaling of $\mathbf{A}$. While (4.5) follows from Remark 4.4 since for $E$ sufficiently small

$$
\int_{B_{3 s}(x)} \sum_{i} R\left(D f_{i}\right) \leq C E^{2 \beta} \int_{B_{3 s}(x)}|D f|^{2}<\frac{1}{4} \int_{B_{3 s}(x)}|D f|^{2}
$$

and therefore

$$
\begin{aligned}
\int_{B_{3 s}(x)}|D f|^{2} & \leq C\left(\mathbf{M}\left(\mathbf{G}_{f}\left\llcorner\mathbf{C}_{3 s}(x)\right)-Q \omega_{m}(3 s)^{m}\right)\right. \\
& \leq C\left(\mathbf{M}\left(T\left\llcorner\mathbf{C}_{3 s}(x)\right)-Q \omega_{m}(3 s)^{m}\right)+C \mathbf{M}\left(\mathbf{G}_{f}\left\llcorner\left(B_{3 s}(x) \backslash K\right) \times \mathbb{R}^{n}\right)\right.\right. \\
& \leq C E s^{m}+C E^{2 \beta}\left|B_{3 s}(x) \backslash K\right| \leq C E s^{m} .
\end{aligned}
$$

(4.5) is therefore a rather simple corollary of the "maximal function truncation" argument employed in Proposition 2.2. Other approximation schemes give instead worse bounds for the Lipschitz constant of the approximating map, cf. for instance [29, Theorem 5.1.1].

Proof of Theorem 4.2. By rescaling and translating, it is not restrictive to assume that $x=0$ and $s=1$. Thus, by Remark 1.5 we can assume $\Psi(0)=0,\|D \Psi\|_{0} \leq C\left(E^{1 / 2}+\mathbf{A}\right)$ and $\left\|D^{2} \Psi\right\|_{0} \leq \mathbf{A}$. The proof of (4.1) is by contradiction. Assume there exist a constant $c_{1}>0$, a sequence of currents $\left(T_{k}\right)_{k \in \mathbb{N}}$ satisfying Assumption 1.1 and area minimizing, ambient manifolds $\Sigma_{k}$ (parametrized by $\Psi_{k}$, with second fundamental forms bounded by $\mathbf{A}_{k}$ ) and corresponding $E_{k}^{\beta}$-Lipschitz approximations $\left(f_{k}\right)_{k \in \mathbb{N}}$ such that

$$
\begin{equation*}
E_{k}:=\mathbf{E}\left(T_{k}, \mathbf{C}_{4}\right) \rightarrow 0, \quad \mathbf{A}_{k} \leq E_{k}^{1 / 4+\delta} \quad \text { and } \quad \int_{B_{2} \backslash K_{k}}\left|D f_{k}\right|^{2} \geq c_{1} E_{k} \tag{4.6}
\end{equation*}
$$

where $K_{k}:=\left\{x \in B_{3}: \mathbf{m e}_{T_{k}}(x)<E_{k}^{2 \beta}\right\}$. Set $\Gamma_{k}:=\left\{x \in B_{4}: \mathbf{m e}_{T_{k}}(x) \leq 2^{-m} E_{k}^{2 \beta}\right\}$ and observe that $\Gamma_{k} \cap B_{3} \subset K_{k}$. From Proposition 2.2, it follows that

$$
\begin{gather*}
\operatorname{Lip}\left(f_{k}\right) \leq C E_{k}^{\beta}  \tag{4.7}\\
\left|B_{r} \backslash K_{k}\right| \leq C E_{k}^{-2 \beta} \mathbf{e}_{T}\left(B_{r+r_{0}(k)} \backslash \Gamma_{k}\right) \quad \text { for every } r \leq 3 \tag{4.8}
\end{gather*}
$$

where $r_{0}(k)=16 E_{k}^{(1-2 \beta) / m}<\frac{1}{2}$. We also assume

$$
\begin{equation*}
\Psi_{k}(0)=0 \quad \text { and } \quad\left\|D \Psi_{k}\right\|_{0}+\left\|D^{2} \Psi_{k}\right\|_{0} \leq C E_{k}^{1 / 4+\delta} \tag{4.9}
\end{equation*}
$$

Then, (4.6), (4.7) and (4.8) give

$$
c_{1} E_{k} \leq \int_{B_{2} \backslash K_{k}}\left|D f_{k}\right|^{2} \leq C \mathbf{e}_{T_{k}}\left(B_{s} \backslash \Gamma_{k}\right) \quad \forall s \in\left[\frac{5}{2}, 3\right]
$$

Setting $c_{2}:=c_{1} /(2 C)$, we have $2 c_{2} E_{k} \leq \mathbf{e}_{T_{k}}\left(B_{s} \backslash \Gamma_{k}\right)=\mathbf{e}_{T_{k}}\left(B_{s}\right)-\mathbf{e}_{T_{k}}\left(B_{s} \cap \Gamma_{k}\right)$, thus leading to

$$
\begin{equation*}
\mathbf{e}_{T_{k}}\left(\Gamma_{k} \cap B_{s}\right) \leq \mathbf{e}_{T_{k}}\left(B_{s}\right)-2 c_{2} E_{k} \tag{4.10}
\end{equation*}
$$

Next observe that $\omega_{m} 4^{m} E_{k}=\mathbf{e}_{T_{k}}\left(B_{4}\right) \geq \mathbf{e}_{T_{k}}\left(B_{s}\right)$. Therefore, by the Taylor expansion in Remark 4.4, (4.10) and $E_{k} \downarrow 0$, it follows that, for every $s \in[5 / 2,3]$,

$$
\begin{align*}
\int_{\Gamma_{k} \cap B_{s}} \frac{\left|D f_{k}\right|^{2}}{2} & \leq\left(1+C E_{k}^{2 \beta}\right) \mathbf{e}_{T_{k}}\left(\Gamma_{k} \cap B_{s}\right) \\
& \leq\left(1+C E_{k}^{2 \beta}\right)\left(\mathbf{e}_{T_{k}}\left(B_{s}\right)-2 c_{2} E_{k}\right) \leq \mathbf{e}_{T_{k}}\left(B_{s}\right)-c_{2} E_{k} \tag{4.11}
\end{align*}
$$

Our aim is to show that (4.11) contradicts the minimizing property of $T_{k}$. To construct a competitor we write $f_{k}(x)=\sum_{i} \llbracket f_{k}^{i}(x) \rrbracket \in \mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}} \times \mathbb{R}^{l}\right)$, and denote by $\left(f_{k}^{i}\right)^{\prime}(x)$ the first
$\bar{n}$ components of the points $f_{k}^{i}(x)$. This induces a map $f_{k}^{\prime}:=\sum_{i} \llbracket\left(f_{k}^{i}\right)^{\prime} \rrbracket$ taking values into $\mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)$. Observe that, since $f_{k}^{i}(x)$ are indeed points of the manifold $\Sigma_{k}$

$$
\begin{equation*}
f_{k}=\sum_{i} \llbracket\left(\left(f_{k}^{i}\right)^{\prime}(x), \Psi_{k}\left(x,\left(f_{k}^{i}\right)^{\prime}(x)\right)\right) \rrbracket . \tag{4.12}
\end{equation*}
$$

We consider $g_{k}:=E_{k}^{-1 / 2} f_{k}^{\prime}$. Since by Remark $4.5 \sup _{k} \operatorname{Dir}\left(g_{k}, B_{3}\right)<\infty$ and $\left|B_{3} \backslash \Gamma_{k}\right| \rightarrow 0$, by Proposition 3.3 we can find a subsequence (not relabelled) of translating sheets $h_{k}$ satisfying (3.2) - (3.3) and $\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}\left(B_{3}\right)} \rightarrow 0$. In particular, we are in the position to apply Proposition 3.4 to $g_{k}$ and $h_{k}$, with $r_{0}=\frac{5}{2}, r_{1}=3$ and $\eta=\frac{c_{2}}{4}$, and find $r \in\left(\frac{5}{2}, 3\right)$ and competitor functions $H_{k}$ satisfying $\left.H_{k}\right|_{B_{3} \backslash B_{r}}=\left.g_{k}\right|_{B_{3} \backslash B_{r}}$,

$$
\begin{gather*}
\operatorname{Dir}\left(H_{k}, B_{r}\right) \leq \operatorname{Dir}\left(h_{k}, B_{r}\right)+\frac{c_{2}}{4}  \tag{4.13}\\
\operatorname{Lip}\left(H_{k}\right) \leq C^{*} E_{k}^{\beta-1 / 2}  \tag{4.14}\\
\left\|\mathcal{G}\left(H_{k}, g_{k}\right)\right\|_{L^{2}\left(B_{r}\right)} \leq C^{*} \operatorname{Dir}\left(g_{k}, B_{r}\right)+C \operatorname{Dir}\left(H_{k}, B_{r}\right) \leq M<\infty \tag{4.15}
\end{gather*}
$$

Moreover, Proposition 3.3 implies that, for $k$ is large enough,

$$
\begin{equation*}
\operatorname{Dir}\left(h_{k}, B_{r}\right) \leq \operatorname{Dir}\left(g_{k}, B_{r} \cap \Gamma_{k}\right)+\frac{c_{2}}{4} \stackrel{(4.11)}{\leq} \frac{\mathbf{e}_{T_{k}}\left(B_{r}\right)}{E_{k}}-\frac{3 c_{2}}{4} E_{k} . \tag{4.16}
\end{equation*}
$$

Note that (4.14) follows from (3.9) observing that $E_{k}^{\beta-1 / 2} \uparrow \infty$ : thus $C^{*}$ depends on $c_{2}$ and the two chosen sequences, but not on $k$. From now on, although this and similar constants are not dimensional, we will keep denoting them by $C$, with the understanding that they do not depend on $k$. Note that, from (4.7) and (4.8), one gets

$$
\begin{align*}
\left\|T_{k}-\mathbf{G}_{f_{k}}\right\|\left(\mathbf{C}_{3}\right) & \leq\left\|T_{k}\right\|\left(\left(B_{3} \backslash K_{k}\right) \times \mathbb{R}^{n}\right)+\left\|\mathbf{G}_{f_{k}}\right\|\left(\left(B_{3} \backslash K_{k}\right) \times \mathbb{R}^{n}\right) \\
& \leq Q\left|B_{3} \backslash K_{k}\right|+E_{k}+Q\left|B_{3} \backslash K_{k}\right|+C\left|B_{3} \backslash K_{k}\right| \operatorname{Lip}\left(f_{k}\right) \\
& \leq E_{k}+C E_{k}^{1-2 \beta} \leq C E_{k}^{1-2 \beta} \tag{4.17}
\end{align*}
$$

Let $(z, y)$ be coordinates on $\mathbb{R}^{m} \times \mathbb{R}^{n}$ and consider the function $\varphi(z, y)=|z|$ and the slice $\left\langle T_{k}-\mathbf{G}_{f_{k}}, \varphi, r\right\rangle$. Observe that, by the coarea formula and Fatou's Lemma,

$$
\int_{r}^{3} \liminf _{k} E_{k}^{2 \beta-1} \mathbf{M}\left(\left\langle T_{k}-\mathbf{G}_{f_{k}}, \varphi, s\right\rangle\right) d s \leq \liminf _{k} E_{k}^{2 \beta-1}\left\|T_{k}-\mathbf{G}_{f_{k}}\right\|\left(\mathbf{C}_{3}\right) \leq C
$$

Therefore, for some $\bar{r} \in(r, 3)$ and a subsequence, not relabeled, $\mathbf{M}\left(\left\langle T_{k}-\mathbf{G}_{f_{k}}, \varphi, \bar{r}\right\rangle\right) \leq$ $C E_{k}^{1-2 \beta}$.

Let now $v_{k}:=\left.E_{k}^{1 / 2} H_{k}\right|_{B_{\bar{r}}}, u_{k}:=\left(v_{k}, \Psi_{k}\left(x, v_{k}\right)\right)$ and consider the current $Z_{k}:=\mathbf{G}_{u_{k}}\left\llcorner\mathbf{C}_{\bar{r}}\right.$. Since $\left.u_{k}\right|_{\partial B_{\bar{r}}}=\left.f_{k}\right|_{\partial B_{\bar{r}}}$, one gets $\partial Z_{k}=\left\langle\mathbf{G}_{f_{k}}, \varphi, \bar{r}\right\rangle$ and, hence, $\mathbf{M}\left(\partial\left(T_{k}\left\llcorner\mathbf{C}_{\bar{r}}-Z_{k}\right)\right) \leq\right.$ $C E_{k}^{1-2 \beta}$. We define

$$
\begin{equation*}
S_{k}=T_{k}\left\llcorner\left(\mathbf{C}_{4} \backslash \mathbf{C}_{\bar{r}}\right)+Z_{k}+R_{k}\right. \tag{4.18}
\end{equation*}
$$

where (cp. Remark 4.3) $R_{k}$ is an integral current supported in $\Sigma_{k}$ such that

$$
\partial R_{k}=\partial\left(T_{k}\left\llcorner\mathbf{C}_{\bar{r}}-Z_{k}\right) \quad \text { and } \quad \mathbf{M}\left(R_{k}\right) \leq C E_{k}^{\frac{(1-2 \beta) m}{m-1}}\right.
$$

$S_{k}$ is supported in $\Sigma_{k}$ and $\partial S_{k}=\partial\left(T_{k}\left\llcorner\mathbf{C}_{4}\right)\right.$. We now show that, since $\beta<\frac{1}{2 m}$, for $k$ large enough, the mass of $S_{k}$ is smaller than that of $T_{k}$. To this aim we write

$$
\begin{aligned}
& \operatorname{Dir}\left(u_{k}, B_{\bar{r}}\right)-\operatorname{Dir}\left(f_{k}, B_{\bar{r}} \cap \Gamma_{k}\right)=\underbrace{\int_{B_{\bar{r}}}\left|D v_{k}\right|^{2}-\int_{B_{\bar{r}} \cap \Gamma_{k}}\left|D f_{k}^{\prime}\right|^{2}}_{I_{1}} \\
& \quad+\underbrace{\int_{B_{\bar{r}}}\left|D\left(\Psi_{k}\left(x, v_{k}\right)\right)\right|^{2}-\int_{B_{\bar{r}}}\left|D\left(\Psi_{k}\left(x, f_{k}^{\prime}\right)\right)\right|^{2}}_{I_{2}}+\underbrace{\int_{B_{\bar{r} \backslash \Gamma_{k}}\left|D\left(\Psi_{k}\left(x, f_{k}^{\prime}\right)\right)\right|^{2}}}_{I_{3}} .
\end{aligned}
$$

The first term is estimated by (4.13) and (3.2): recalling that $v_{k}=E_{k}^{1 / 2} H_{k}$ and $f_{k}^{\prime}=E_{k}^{1 / 2} g_{k}$ (but also that the two functions coincide on $B_{\bar{r}} \backslash B_{r}$ ) we achieve $I_{1} \leq \frac{c_{2}}{2} E_{k}$ for $k$ large enough. For what concerns the second, we proceed as follows. First we write

$$
I_{2}=\sum_{i} \int_{B_{\bar{r}}}\left(D \left(\Psi_{k}\left(x, u_{k}(x)\right)_{i}-D\left(\Psi_{k}\left(x, f_{k}^{\prime}(x)\right)_{i}\right):\left(D \left(\Psi_{k}\left(x, u_{k}(x)\right)_{i}+D\left(\Psi_{k}\left(x, f_{k}^{\prime}(x)\right)_{i}\right)\right.\right.\right.\right.
$$

Next, recalling the chain rule [14, Proposition 1.12], we get

$$
\begin{aligned}
\mid D\left(\Psi_{k}\left(x, u_{k}(x)\right)_{i}+D\left(\Psi_{k}\left(x, f_{k}^{\prime}(x)\right)_{i} \mid\right.\right. & \leq C\left\|D_{x} \Psi_{k}\right\|_{0}+C\left\|D_{u} \Psi_{k}\right\|_{0}\left(\operatorname{Lip}\left(u_{k}\right)+\operatorname{Lip}\left(f_{k}^{\prime}\right)\right) \\
& \stackrel{(4.9)}{\leq} C E_{k}^{1 / 4+\delta} .
\end{aligned}
$$

Using the letter inequality, the chain rule and (4.9), once again we achieve

$$
\begin{align*}
& I_{2} \leq C E_{k}^{1 / 4+\delta} \int_{B_{\bar{r}}}( \sum_{i}\left|D_{x} \Psi_{k}\left(x, v_{k}^{i}(x)\right)-D_{x} \Psi_{k}\left(x,\left(f_{k}^{i}\right)^{\prime}(x)\right)\right| \\
&\left.\quad+\left\|D_{u} \Psi_{k}\right\|_{0}\left(\left|D v_{k}\right|+\left|D f_{k}^{\prime}\right|\right)\right) \\
& \leq C E_{k}^{1 / 4+\delta}\left\|D^{2} \Psi_{k}\right\|_{0} \int_{B_{\bar{r}}} \mathcal{G}\left(v_{k}, f_{k}^{\prime}\right)+C E_{k}^{1 / 2+2 \delta} \int_{B_{\bar{r}}}\left(\left|D v_{k}\right|+\left|D f_{k}^{\prime}\right|\right) \\
& \leq C E_{k}^{1 / 2+2 \delta} E_{k}^{1 / 2}+C E_{k}^{1+2 \delta} \leq C E_{k}^{1+2 \delta} \tag{4.19}
\end{align*}
$$

Finally, $I_{3} \leq C\left\|D \Psi_{k}\right\|_{\infty}^{2}\left|B_{3} \backslash \Gamma_{k}\right| \leq C E_{k}^{1+\beta}$. Thus, for $k$ large enough we achieve $I_{2}+I_{3} \leq$ $\frac{c_{2}}{4} E_{k}$, thereby reaching $\operatorname{Dir}\left(u_{k}, B_{\bar{r}}\right)-\operatorname{Dir}\left(f_{k}, B_{\bar{r}} \cap \Gamma_{k}\right) \leq \frac{3 c_{2}}{4} E_{k}$. Hence,

$$
\begin{align*}
& \mathbf{M}\left(S_{k}\right)-\mathbf{M}\left(T_{k}\right) \leq \mathbf{M}\left(Z_{k}\right)+C \mathbf{M}\left(R_{k}\right)-\mathbf{M}\left(T_{k}\left\llcorner\mathbf{C}_{\bar{r}}\right)\right. \\
& \leq Q\left|B_{\bar{r}}\right|+\int_{B_{\bar{r}}} \frac{\left|D u_{k}\right|^{2}}{2}+C E_{k}^{1+2 \beta}+C E_{k}^{\frac{(1-2 \beta) m}{m-1}}-Q\left|B_{\bar{r}}\right|-\mathbf{e}_{T_{k}}\left(B_{\bar{r}}\right) \\
& \leq \int_{B_{\bar{r}} \cap \Gamma_{k}} \frac{\left|D f_{k}\right|^{2}}{2}+\frac{3}{4} c_{2} E_{k}+C E_{k}^{1+2 \beta}+C E_{k}^{\frac{(1-2 \beta) m}{m-1}}-\mathbf{e}_{T_{k}}\left(B_{\bar{r}}\right) \\
&(4.11)  \tag{4.20}\\
& c_{2} E_{k} \\
& 4
\end{align*} C E_{k}^{1+\beta}+C E_{k}^{\frac{(1-2 \beta) m}{m-1}}<0, ~ 又 土 \text {. }
$$

as soon as $E_{k}$ is small enough. This gives the desired contradiction and proves (4.1).

For what concerns (4.2) and (4.3), we argue similarly. Without loss of generality we assume $x=0$ and $s=1$. Hence, we let $\left(T_{k}\right)_{k},\left(\Sigma_{k}\right)_{k}$ and $\left(\Psi_{k}\right)_{k}$ be sequences with vanishing $E_{k}:=\mathbf{E}\left(T_{k}, \mathbf{C}_{4}\right)$ and satisfying (4.9), but contradicting (4.2) or (4.3). So, being $f_{k}$ the $E_{k}^{\beta}$-Lipschitz approximations, we know that, for any sequence of Dir-minimizing functions $\bar{u}_{k}$ which we might choose, when we set $w_{k}=\left(\bar{u}_{k}, \Psi_{k}\left(x, \bar{u}_{k}\right)\right)$ we will have

$$
\begin{equation*}
\liminf _{k} \underbrace{E_{k}^{-1} \int_{B_{2}}\left(\mathcal{G}\left(f_{k}, w_{k}\right)^{2}+\left(\left|D f_{k}\right|-\left|D w_{k}\right|\right)^{2}+\left|D\left(\boldsymbol{\eta} \circ f_{k}-\boldsymbol{\eta} \circ w_{k}\right)\right|^{2}\right)}_{=: I(k)}>0 \tag{4.21}
\end{equation*}
$$

As in the previous argument we introduce the maps $f_{k}^{\prime}$ satisfying (4.12), the normalized functions $g_{k}=E_{k}^{-1 / 2} f_{k}^{\prime}$ and, after extraction of a subsequence, the translating sheets $h_{k}$ satisfying (3.2) - (3.3) and $\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}\left(B_{3}\right)} \rightarrow 0$. We next claim that
(i) $\lim _{k} \int_{B_{2}}\left|D g_{k}\right|^{2}=\int_{B_{2}}\left|D h_{k_{0}}\right|^{2}$, for any $k_{0}$ (recall that $\int_{B_{2}}\left|D h_{k}\right|^{2}$ is constant);
(ii) $h_{k}$ is Dir-minimizing in $B_{2}$.

If (i) is false, then there is a positive constant $c_{2}$ such that, for any $r \in[5 / 2,3]$,

$$
\begin{equation*}
\int_{B_{r}} \frac{\left|D h_{k}\right|^{2}}{2} \leq \int_{B_{r}} \frac{\left|D g_{k}\right|^{2}}{2}-c_{2} \leq \frac{\mathbf{e}_{T_{k}}\left(B_{r}\right)}{E_{k}}-\frac{c_{2}}{2} \tag{4.22}
\end{equation*}
$$

provided $k$ large enough (where the last inequality is again an effect of the Taylor expansion of Remark 4.4). We next define the competitor currents $S_{k}$ as in the argument leading to (4.20): this latter inequality is reached thanks to (4.22), which substitutes (4.11) and (4.16). On the other hand (4.20) contradicts the minimizing property of $T_{k}$. If (ii) is false, then $h_{k}$ is not Dir-minimizing in $B_{2}$. This implies that one of the $\zeta^{j}$ in the translating sheets $h_{k}$ is not Dir-minimizing in $B_{2}$. Indeed, in the opposite case, by [14, Theorem 3.9], $\left\|\mathcal{G}\left(\zeta^{j}, Q \llbracket 0 \rrbracket\right)\right\|_{C^{0}\left(B_{2}\right)}<\infty$ and, since $h_{k}=\sum_{i} \llbracket \tau_{y_{k}^{i}} \circ \zeta^{i} \rrbracket$ and $\left|y_{k}^{i}-y_{k}^{j}\right| \rightarrow \infty$ for $i \neq j$, by the maximum principle of [14, Proposition 3.5], $h_{k}$ would be Dir-minimizing. Thus, for some $\zeta^{j}$ we can find a competitor $\hat{\zeta}^{j}$ with less energy in the ball $B_{2}$. So the functions $F_{k}=\sum_{j} \llbracket \tau_{y_{k}^{j}} \circ \hat{\zeta}^{j} \rrbracket$ satisfy, for any $r \in[5 / 2,3]$,

$$
\begin{equation*}
\int_{B_{r}} \frac{\left|D F_{k}\right|^{2}}{2} \leq \int_{B_{r}} \frac{\left|D h_{k}\right|^{2}}{2}-c_{2}=\lim _{k} \int_{B_{r}} \frac{\left|D g_{k}\right|^{2}}{2}-c_{2} \leq \frac{\mathbf{e}_{T}\left(B_{r}\right)}{E_{k}}-\frac{c_{2}}{2} \tag{4.23}
\end{equation*}
$$

provided $k$ is large enough (here $c_{2}>0$ is some constant independent of $r$ and $k$ ). On the other hand $F_{k}=h_{k}$ on $B_{3} \backslash B_{5 / 2}$ and therefore $\left\|\mathcal{G}\left(F_{k}, g_{k}\right)\right\|_{L^{2}\left(B_{3} \backslash B_{5 / 2}\right)} \rightarrow 0$. We then construct the competitor current $S_{k}$ of (4.18): this time we use, however, the map $F_{k}$ in place of $h_{k}$ to construct $H_{k}$ via Proposition 3.4 and we reach the contradiction (4.20) using (4.23) in place of (4.11) and (4.16).

We next set $\bar{u}_{k}:=E_{k}^{1 / 2} h_{k}$ and we aim at showing that, for $w_{k}=\left(\bar{u}_{k}, \Psi_{k}\left(x, \bar{u}_{k}\right)\right), I(k) \rightarrow 0$, a contradiction to (4.21). Observe first that, by $\left\|\mathcal{G}\left(g_{k}, h_{k}\right)\right\|_{L^{2}} \rightarrow 0$, we have $D\left(\boldsymbol{\xi} \circ g_{k}\right)-$ $D\left(\boldsymbol{\xi} \circ h_{k}\right) \rightharpoonup 0$ in $L^{2}$ (recall the definition of $\boldsymbol{\xi}$ in Section 1.5). On the other hand, recall that $D\left(\boldsymbol{\xi} \circ h_{k}\right)$ is actually a single function, independent of $k$, because $h_{k}$ is a sequence of translating sheets. So, (i) and the identities $\left|D\left(\boldsymbol{\xi} \circ g_{k}\right)\right|=\left|D g_{k}\right|,\left|D\left(\boldsymbol{\xi} \circ h_{k}\right)\right|=\left|D h_{k}\right|$ imply
that $D\left(\boldsymbol{\xi} \circ g_{k}\right)-D\left(\boldsymbol{\xi} \circ h_{k}\right)$ converge strongly to 0 in $L^{2}$. If we next set $\hat{h}_{k}=\sum_{i} \llbracket h_{k}^{i}-\boldsymbol{\eta} \circ h_{k} \rrbracket$ and $\hat{g}_{k}=\sum_{i} \llbracket g_{k}^{i}-\boldsymbol{\eta} \circ g_{k} \rrbracket$, we obviously have $\left\|\mathcal{G}\left(\hat{h}_{k}, \hat{g}_{k}\right)\right\|_{L^{2}}+\left\|\boldsymbol{\eta} \circ h_{k}-\boldsymbol{\eta} \circ g_{k}\right\|_{L^{2}} \rightarrow 0$. Recall however that the Dirichlet energy enjoys the splitting

$$
\operatorname{Dir}\left(g_{k}\right)=Q \int\left|D\left(\boldsymbol{\eta} \circ g_{k}\right)\right|^{2}+\operatorname{Dir}\left(\hat{g}_{k}\right) \quad \operatorname{Dir}\left(h_{k}\right)=Q \int\left|D\left(\boldsymbol{\eta} \circ h_{k}\right)\right|^{2}+\operatorname{Dir}\left(\hat{h}_{k}\right)
$$

So (i) implies that the Dirichlet energies of $\boldsymbol{\eta} \circ g_{k}$ and $\hat{g}_{k}$ converge, respectively, to those of $\boldsymbol{\eta} \circ h_{k}$ and $\hat{h}_{k}$ (which, we recall again, are independent of $k$ because the $h_{k}$ 's are translating sheets). We thus infer that $D\left(\boldsymbol{\eta} \circ h_{k}\right)-D\left(\boldsymbol{\eta} \circ g_{k}\right)$ converges to 0 strongly in $L^{2}$.

Coming back to $w_{k}$ we observe that

$$
\begin{equation*}
E_{k}^{-1} \int_{B_{2}} \mathcal{G}\left(w_{k}, f_{k}\right)^{2} \leq\left(2+\operatorname{Lip}(D \Psi)^{2}\right) E_{k}^{-1} \int_{B_{2}} \mathcal{G}\left(\bar{u}_{k}, f_{k}^{\prime}\right)^{2}=C \int_{B_{2}} \mathcal{G}\left(h_{k}, g_{k}\right)^{2} \rightarrow 0 \tag{4.24}
\end{equation*}
$$

So,

$$
\begin{align*}
\underset{k}{\limsup } I(k) \leq & \left.\underset{k}{\lim \sup } \int_{B_{2}}\left(\left|D g_{k}\right|-\left|D h_{k}\right|\right)^{2}+\left|D\left(\boldsymbol{\eta} \circ g_{k}-\boldsymbol{\eta} \circ h_{k}\right)\right|^{2}\right) \\
& +C(Q) \limsup _{k} E_{k}^{-1} \int_{B_{2}} \mathcal{G}\left(D\left(\Psi\left(x, f_{k}^{\prime}\right)\right), D\left(\Psi\left(x, \bar{u}_{k}\right)\right)\right)^{2} \\
\leq & C \underset{k}{\limsup } E_{k}^{-1} \int_{B_{2}} \mathcal{G}\left(D\left(\Psi\left(x, f_{k}^{\prime}\right)\right), D\left(\Psi\left(x, \bar{u}_{k}\right)\right)\right)^{2}=\underset{k}{\lim \sup ^{2}} E_{k}^{-1} J(k) . \tag{4.25}
\end{align*}
$$

Recalling the chain rule of [14, Proposition 1.12], we have

$$
\begin{aligned}
& D\left(\Psi\left(x, f_{k}^{\prime}\right)\right)(x)=\sum_{i} \llbracket D_{x} \Psi\left(x,\left(f_{k}^{i}\right)^{\prime}(x)\right)+D_{v} \Psi\left(x,\left(f_{k}^{i}\right)^{\prime}(x)\right) \cdot D\left(f_{k}^{i}\right)^{\prime}(x) \rrbracket \\
& D\left(\Psi\left(x, \bar{u}_{k}\right)\right)(x)=\sum_{i} \llbracket D_{x} \Psi\left(x, \bar{u}_{k}^{i}(x)\right)+D_{v} \Psi\left(x, \bar{u}_{k}^{i}(x)\right) \cdot D \bar{u}_{k}^{i}(x) \rrbracket
\end{aligned}
$$

So we can estimate

$$
J(k) \leq C \operatorname{Lip}\left(D_{x} \Psi\right)^{2} \int_{B_{2}} \mathcal{G}\left(f_{k}^{\prime}, \bar{u}_{k}\right)^{2}+C\|D \Psi\|_{0}^{2} \int_{B_{2}}\left(\left|D f_{k}^{\prime}\right|^{2}+\left|D \bar{u}_{k}\right|^{2}\right)^{(4.9)} \leq C E_{k}^{3 / 2+2 \delta}
$$

We therefore conclude that $E_{k}^{-1} J(k) \rightarrow 0$ and thus $I(k) \rightarrow 0$, which contradicts (4.21).

## 5. Gradient $L^{p}$ estimate

In this section we prove Theorem 1.3. The result is a consequence of an higher integrability estimate for the gradient of Dir-minimizing functions, the $o(E)$-improved estimate for the excess measure given in Proposition 5.4 and a very careful "covering and stopping radius" argument (cf. [31] for an exposition in a more elementary context).
5.1. Higher integrability of the gradient of Dir-minimizers. Most of the energy of a Dir-minimizer lies where the gradient is relatively small. We prove indeed the following a priori estimate (cf. [30] for a different proof and some improvements).

Theorem 5.1 (Higher integrability of Dir-minimizers). There exists $p_{10}>2$ such that, for every $\Omega^{\prime} \subset \subset \Omega \subset \mathbb{R}^{m}$ open domains, there is a constant $C>0$ such that

$$
\begin{equation*}
\|D u\|_{L^{p_{10}\left(\Omega^{\prime}\right)}} \leq C\|D u\|_{L^{2}(\Omega)} \quad \text { for every } \operatorname{Dir-minimizing~} u \in W^{1,2}\left(\Omega, \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)\right) \tag{5.1}
\end{equation*}
$$

Proof. The statement is a corollary of Proposition 5.2 below and a Gehring type lemma, cf. [21, Proposition 5.1].

Proposition 5.2. Let $\frac{2(m-1)}{m}<p_{11}<2$. Then, there exists $C=C\left(m, n, Q, p_{11}\right)$ such that, for every $u: \Omega \rightarrow \mathcal{A}_{Q}$ Dir-minimizing, the following holds

$$
\left(f_{B_{r}(x)}|D u|^{2}\right)^{1 / 2} \leq C\left(f_{B_{2 r}(x)}|D u|^{p_{11}}\right)^{1 / p_{11}} \quad \forall x \in \Omega, \forall r<\min \{1, \operatorname{dist}(x, \partial \Omega) / 2\} .
$$

Proof. Since the estimate is invariant under translations and rescalings, it is enough to prove it for $x=0$ and $r=1$. We assume, therefore $\Omega=B_{2}$. Let $u: \Omega \rightarrow \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right)$ be Dir-minimizing and let $F=\boldsymbol{\xi} \circ u: \Omega \rightarrow \mathcal{Q} \subset \mathbb{R}^{N}$. Denote by $\bar{F} \in \mathbb{R}^{N}$ the average of $F$ on $B_{2}$. By Fubini's theorem and the Poincaré inequality, there exists $s \in[1,2]$ such that

$$
\int_{\partial B_{s}}\left(|F-\bar{F}|^{p_{11}}+|D F|^{p_{11}}\right) \leq C \int_{B_{2}}\left(|F-\bar{F}|^{p_{11}}+|D F|^{p_{11}}\right) \leq C\|D F\|_{L^{p_{11}\left(B_{2}\right)}}^{p_{p_{1}}}
$$

Consider $\left.F\right|_{\partial B_{s}}$. Since $\frac{1}{2}>\frac{1}{p_{11}}-\frac{1}{2(m-1)}$, we can use the embedding $W^{1, p_{11}}\left(\partial B_{s}\right) \hookrightarrow$ $H^{1 / 2}\left(\partial B_{s}\right)$ (see, for example, [1]). Hence, we infer that

$$
\begin{equation*}
\|F-\bar{F}\|_{H^{1 / 2}\left(\partial B_{s}\right)} \leq C\|D F\|_{L^{p_{11}\left(B_{2}\right)}} \tag{5.2}
\end{equation*}
$$

Let $\hat{F}$ be the harmonic extension of $\left.F\right|_{\partial B_{s}}$ in $B_{s}$. It is well known (one could, for example, use the result in [1] on the half-space together with a partition of unity) that

$$
\begin{equation*}
\|D \hat{F}\|_{L^{2}\left(B_{s}\right)} \leq C(m) \min _{p \in \mathbb{R}^{N}}\|\hat{F}-p\|_{H^{1 / 2}\left(\partial B_{s}\right)} \stackrel{(5.2)}{\leq} C\|D F\|_{L^{p_{11}\left(B_{2}\right)}} \tag{5.3}
\end{equation*}
$$

Consider the map $\boldsymbol{\rho}$ of [14, Theorem 2.1]. Since $\left.\boldsymbol{\rho} \circ \hat{F}\right|_{\partial B_{s}}=\left.u\right|_{\partial B_{s}}$ and $\boldsymbol{\rho} \circ \hat{F}$ takes values in $\mathcal{Q}$, by the minimizing property of $u$ and the Lipschitz continuity of $\boldsymbol{\xi}, \boldsymbol{\xi}^{-1}$ and $\boldsymbol{\rho}$, we conclude:

$$
\left(\int_{B_{1}}|D u|^{2}\right)^{1 / 2} \leq C\left(\int_{B_{s}}|D \hat{F}|^{2}\right)^{1 / 2} \leq C\left(\int_{B_{2}}|D F|^{p_{11}}\right)^{1 / p_{11}}=C\left(\int_{B_{2}}|D u|^{p_{11}}\right)^{1 / p_{11}}
$$

Remark 5.3. Proposition 5.2 can be proved in several different ways, which are based on more common test function arguments: cf. the intrinsic proof (i.e. which does not use the biLipschitz embedding $\boldsymbol{\xi}$ ) in [30] or the usual Caccioppoli's inequality for quasi minima [22, Theorem 6.5].
5.2. Improved excess estimate. The higher integrability of the Dir-minimizing functions and the harmonic approximation lead to the following estimate, which we call "weak" since we will improve it in Theorem 6.1.

Proposition 5.4 (Weak excess estimate). For every $\eta_{10}>0$, there exists $\varepsilon_{13}>0$ with the following property. Let $T$ be area minimizing and assume it satisfies Assumption 1.1 in $\mathbf{C}_{4 s}(x)$. If $E=\mathbf{E}\left(T, \mathbf{C}_{4 s}(x)\right) \leq \varepsilon_{13}$, then

$$
\begin{equation*}
\mathbf{e}_{T}(A) \leq \eta_{10} E s^{m}+C \mathbf{A}^{2} s^{m+2} \tag{5.4}
\end{equation*}
$$

for every $A \subset B_{s}(x)$ Borel with $|A| \leq \varepsilon_{13}\left|B_{s}(x)\right|$.
Proof. Without loss of generality, we can assume $s=1$ and $x=0$. We distinguish the two regimes: $E \leq \mathbf{A}^{2}$ and $\mathbf{A}^{2} \leq E$. In the former, clearly $\mathbf{e}_{T}(A) \leq C E \leq C \mathbf{A}^{2}$. In the latter, we let $f$ be the $E^{\frac{1}{4 m}}$-Lipschitz approximation of $T$ in $\mathbf{C}_{3}$ and, arguing as for the proof of Theorem 4.2, we find a radius $r \in(1,2)$ and a current $R$ such that

$$
\partial R=\left\langle T-\mathbf{G}_{f}, \varphi, r\right\rangle \quad \text { and } \quad \mathbf{M}(R) \leq C E^{\left(1-\frac{1}{2 m}\right) \frac{m}{m-1}}
$$

Therefore, by the Taylor expansion in Remark 4.4, we have:

$$
\begin{equation*}
\|T\|\left(\mathbf{C}_{r}\right) \leq \mathbf{M}\left(\mathbf{G}_{f}\left\llcorner\mathbf{C}_{r}+R\right) \leq\left\|\mathbf{G}_{f}\right\|\left(\mathbf{C}_{r}\right)+C E^{\frac{2 m-1}{2 m-2}} \leq Q\left|B_{r}\right|+\int_{B_{r}} \frac{|D f|^{2}}{2}+C E^{1+\gamma}\right. \tag{5.5}
\end{equation*}
$$

where $\gamma=\frac{1}{2 m}$. On the other hand, using again the Taylor expansion for the part of the current which coincides with the graph of $f$, we deduce as well that

$$
\begin{equation*}
\|T\|\left(\left(B_{r} \cap K\right) \times \mathbb{R}^{n}\right) \geq Q\left|B_{r} \cap K\right|+\int_{B_{r} \cap K} \frac{|D f|^{2}}{2}-C E^{1+\gamma} \tag{5.6}
\end{equation*}
$$

Subtracting (5.6) from (5.5), we deduce

$$
\begin{equation*}
\mathbf{e}_{T}\left(B_{r} \backslash K\right) \leq \int_{B_{r} \backslash K} \frac{|D f|^{2}}{2}+C E^{1+\gamma} \tag{5.7}
\end{equation*}
$$

If $\varepsilon_{13}$ is chosen small enough, we infer from (5.7) and (4.1) in Theorem 4.2 that

$$
\begin{equation*}
\mathbf{e}_{T}\left(B_{r} \backslash K\right) \leq \eta E+C E^{1+\gamma} \tag{5.8}
\end{equation*}
$$

for a suitable $\eta>0$ to be chosen. Let now $A \subset B_{1}$ be such that $|A| \leq \varepsilon_{13} \omega_{m}$. Combining (5.8) with the Taylor expansion, we have

$$
\begin{equation*}
\mathbf{e}_{T}(A) \leq \mathbf{e}_{T}(A \backslash K)+\int_{A} \frac{|D f|^{2}}{2}+C E^{1+\gamma} \leq \int_{A} \frac{|D f|^{2}}{2}+\eta E+C E^{1+\gamma} \tag{5.9}
\end{equation*}
$$

If $\varepsilon_{13}$ is small enough, we can again apply Theorem 4.2. Using the coordinates of Remark 1.5, there is a Dir-minimizing $u$ such that $|D f|$ is close in $L^{2}$ (with an error $\eta E$ ) to $|D w|$ with $w=(u, \Psi(x, u))$ and by Remark 4.5 $\operatorname{Dir}(u) \leq C E$. On the other hand $|D w(x)| \leq$
$\left(1+\|D \Psi\|_{0}\right)|D u|+\|D \Psi\|_{0}$. Since $\|D \Psi\|_{0} \leq C E^{1 / 2}$, by Theorem 5.1 $\||D w|\|_{L^{p_{10}\left(B_{1}\right)}} \leq C E^{1 / 2}$. Therefore,

$$
\begin{equation*}
\mathbf{e}_{T}(A) \stackrel{(4.2)}{\leq} \int_{A}|D w|^{2}+3 \eta E+C E^{1+\gamma} \leq C\left(|A|^{1-2 / p_{10}}+\eta\right) E+C E^{1+\gamma} \tag{5.10}
\end{equation*}
$$

Hence, if $\varepsilon_{13}$ and $\eta$ are suitably chosen, (5.4) follows from (5.10).
5.3. Proof of Theorem 1.3. We assume without loss of generality that $E>0$ and divide the proof into two steps.

Step 1. There exist constants $\gamma \geq 2^{m}$ and $\varrho>0$ such that, for every $c \in\left[1,(\gamma E)^{-1}\right]$ and $s \in[2,4]$ with $\bar{s}=s+4 c^{-1 / m} \leq 4$, we have

$$
\begin{equation*}
\int_{\{\gamma c E \leq \mathbf{d} \leq 1\} \cap B_{s}} \mathbf{d} \leq \gamma^{-\varrho} \int_{\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\} \cap B_{\bar{s}}} \mathbf{d}+C c^{-2 / m} \mathbf{A}^{2} \tag{5.11}
\end{equation*}
$$

In order to prove it, let $N_{B}$ be the constant in Besicovich's covering theorem [19, Section 1.5.2] and choose $N \in \mathbb{N}$ so large that $N_{B}<2^{N-1}$. Let $\varepsilon_{13}$ be as in Proposition 5.4 when we choose $\eta_{10}=2^{-2 m-N}$, and set

$$
\gamma=\max \left\{2^{m}, \varepsilon_{13}^{-1}\right\} \quad \text { and } \quad \varrho=\min \left\{-\log _{\gamma}\left(N_{B} / 2^{N-1}\right), \frac{1}{2 m}\right\}
$$

Let $c$ and $s$ be any real numbers as above. For almost every $x \in\{\gamma c E \leq \mathbf{d} \leq 1\} \cap B_{s}$, there exists $r_{x}$ such that

$$
\begin{equation*}
\left.E\left(T, \mathbf{C}_{4 r_{x}}(x)\right) \leq c E \quad \text { and } \quad E\left(T, \mathbf{C}_{t}(x)\right) \geq c E \quad \forall t \in\right] 0,4 r_{x}[ \tag{5.12}
\end{equation*}
$$

Indeed, since $\mathbf{d}(x)=\lim _{r \rightarrow 0} E\left(T, \mathbf{C}_{r}(x)\right) \geq \gamma c E \geq 2^{m} c E$ and

$$
E\left(T, \mathbf{C}_{t}(x)\right)=\frac{\mathbf{e}_{T}\left(B_{t}(x)\right)}{\omega_{m} t^{m}} \leq \frac{4^{m} E}{t^{m}} \leq c E \quad \text { for } t \geq \frac{4}{\sqrt[m]{c}}
$$

we just choose $4 r_{x}=\min \left\{t \leq 4 / \sqrt[m]{c}: E\left(T, \mathbf{C}_{t}(x)\right) \leq c E\right\}$. Note also that $r_{x} \leq 1 / \sqrt[m]{c}$. Consider the current $T$ in $\mathbf{C}_{4 r_{x}}(x)$. Setting $A=\{\gamma c E \leq \mathbf{d}\} \cap B_{4 r_{x}}(x)$, we have that

$$
\mathbf{E}\left(T, \mathbf{C}_{4 r_{x}}(x)\right) \leq c E \leq \frac{E}{\gamma E} \leq \varepsilon_{13} \quad \text { and } \quad|A| \leq \frac{c E\left|B_{4 r_{x}}(x)\right|}{\gamma c E} \leq \varepsilon_{13}\left|B_{4 r_{x}}(x)\right|
$$

Hence, we can apply Proposition 5.4 to $T\left\llcorner\mathbf{C}_{4 r_{x}}(x)\right.$ to get

$$
\begin{align*}
& \int_{B_{r_{x}}(x) \cap\{\gamma c E \leq \mathbf{d} \leq 1\}} \mathbf{d} \leq \int_{A} \mathbf{d} \leq \mathbf{e}_{T}(A) \leq 2^{-2 m-N} \mathbf{e}_{T}\left(B_{4 r_{x}}(x)\right)+C r_{x}^{m+2} \mathbf{A}^{2} \\
\leq & 2^{-2 m-N}\left(4 r_{x}\right)^{m} \omega_{m} \mathbf{E}\left(T, \mathbf{C}_{4 r_{x}}(x)\right)+C r_{x}^{m+2} \mathbf{A}^{2} \stackrel{(5.12)}{\leq} 2^{-N} \mathbf{e}_{T}\left(B_{r_{x}}(x)\right)+C r_{x}^{m+2} \mathbf{A}^{2} . \tag{5.13}
\end{align*}
$$

Thus,

$$
\begin{align*}
\mathbf{e}_{T}\left(B_{r_{x}}(x)\right) & =\int_{B_{r_{x}}(x) \cap\{\mathbf{d}>1\}} \mathbf{d}+\int_{B_{r_{x}}(x) \cap\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\}} \mathbf{d}+\int_{B_{r_{x}}(x) \cap\left\{\mathbf{d}<\frac{c E}{\gamma}\right\}} \mathbf{d} \\
& \leq \int_{A} \mathbf{d}+\int_{B_{r_{x}}(x) \cap\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\}} \mathbf{d}+\frac{c E}{\gamma} \omega_{m} r_{x}^{m} \\
& \stackrel{(5.12),(5.13)}{\leq}\left(2^{-N}+\gamma^{-1}\right) \mathbf{e}_{T}\left(B_{r_{x}}(x)\right)+C r_{x}^{m+2} \mathbf{A}^{2}+\int_{B_{r_{x}}(x) \cap\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\}} \mathbf{d} . \tag{5.14}
\end{align*}
$$

Therefore, recalling that $\gamma \geq 2^{m} \geq 4$, from (5.13) and (5.14) we infer:

$$
\begin{aligned}
\int_{B_{r_{x}}(x) \cap\{\gamma c E \leq \mathbf{d} \leq 1\}} \mathbf{d} & \leq \frac{2^{-N}}{1-2^{-N}-\gamma^{-1}} \int_{B_{r_{x}}(x) \cap\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\}} \mathbf{d}+C r_{x}^{m+2} \mathbf{A}^{2} \\
& \leq 2^{-N+1} \int_{B_{r_{x}}(x) \cap\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\}} \mathbf{d}+C r_{x}^{m+2} \mathbf{A}^{2}
\end{aligned}
$$

By Besicovich's covering theorem, we choose $N_{B}$ families of disjoint balls $\bar{B}_{r_{x}}(x)$ whose union covers $\{\gamma c E \leq \mathbf{d} \leq 1\} \cap B_{s}$ and, since as already noticed $r_{x} \leq 1 / \sqrt[m]{c}$ for every $x$, we conclude:

$$
\int_{\{\gamma c E \leq \mathbf{d} \leq 1\} \cap B_{s}} \mathbf{d} \leq N_{B} 2^{-N+1} \int_{\left\{\frac{c E}{\gamma} \leq \mathbf{d} \leq 1\right\} \cap B_{s+\frac{2}{m} /}^{\sqrt[m]{c}}} \mathbf{d}+C c^{-\frac{2}{m}} \mathbf{A}^{2},
$$

which, for the above defined $\varrho$, implies (5.11).
Step 2. We iterate (5.11) in order to conclude (1.3). Denote by $L$ the largest integer smaller than $2^{-1}\left(\left(\log _{\gamma} E^{-1}\right)-1\right), s_{L}=2$ and recursively $s_{k}=s_{k+1}+2 \gamma^{-\frac{2 k}{m}}$ for $k \in$ $\{L-1, \ldots, 1\}$. Notice that, since $\gamma \geq 2^{m}, s_{k}<4$ for every $k$. Thus, we can apply (5.11) with $c=\gamma^{2 k}, s=s_{k}$ and $\bar{s}=s_{k-1}$ to conclude

$$
\int_{\left\{\gamma^{2 k+1} E \leq \mathbf{d} \leq 1\right\} \cap B_{s_{k}}} \mathbf{d} \leq \gamma^{-\varrho} \int_{\left\{\gamma^{2 k-1} E \leq \mathbf{d} \leq 1\right\} \cap B_{s_{k-1}}} \mathbf{d}+C \gamma^{-\frac{4 k}{m}} \mathbf{A}^{2} \quad \forall k \in\{2, \ldots, L\}
$$

In particular, iterating this estimate we get

$$
\begin{equation*}
\left.\left.\int_{\left\{\gamma^{2 k+1}\right.} E \leq \mathbf{d} \leq 1\right\} \cap B_{2}\right) ~ d \leq \gamma^{-(k-1) \varrho} \int_{\{\gamma E \leq \mathbf{d} \leq 1\} \cap B_{s_{1}}} \mathbf{d}+C \mathbf{A}^{2} \sum_{\ell=0}^{k-2} \gamma^{-\left(\frac{4(k-\ell)}{m}+\ell \varrho\right)} . \tag{5.15}
\end{equation*}
$$

Set $A_{0}=\{\mathbf{d}<\gamma E\}, A_{k}=\left\{\gamma^{2 k-1} E \leq \mathbf{d}<\gamma^{2 k+1} E\right\}$ for $k=1, \ldots, L$, and $A_{L+1}=$ $\left\{\gamma^{2 L+1} E \leq \mathbf{d} \leq 1\right\}$. Since $\cup A_{k}=\{\mathbf{d} \leq 1\}$, for $p_{1}<1+\frac{\varrho}{2} \leq 1+\frac{1}{m}$, we conclude:

$$
\begin{aligned}
& \int_{B_{2} \cap\{\mathbf{d} \leq 1\}} \mathbf{d}^{p_{1}}=\sum_{k=0}^{L+1} \int_{A_{k} \cap B_{2}} \mathbf{d}^{p_{1}} \leq \sum_{k} \gamma^{(2 k+1)\left(p_{1}-1\right)} E^{p_{1}-1} \int_{A_{k} \cap B_{2}} \mathbf{d} \\
& \stackrel{(5.15)}{\leq} C \sum_{k} \gamma^{k\left(2\left(p_{1}-1\right)-\varrho\right)} E^{p_{1}}+C \sum_{k} \sum_{\ell=0}^{k-2} \gamma^{k\left(2\left(p_{1}-1\right)-\frac{4}{m}\right)+\ell\left(\frac{4}{m}-\varrho\right)} E^{p_{1}-1} \mathbf{A}^{2} \\
& \leq C E^{p_{1}}+C \sum_{k} \gamma^{k\left(2\left(p_{1}-1\right)-\varrho\right)} E^{p_{1}-1} \mathbf{A}^{2} .
\end{aligned}
$$

## 6. Almgren's approximation theorem

In this section we show how Theorem 1.3 gives a simple proof of the approximation result in Theorem 1.4. The key point is the following theorem.

Theorem 6.1 (Almgren's strong excess estimate). There are constants $\varepsilon_{11}, \gamma_{11}, C>0$ (depending on $m, n, \bar{n}, Q$ ) with the following property. Assume $T$ satisfies Assumption 1.1 in $\mathbf{C}_{4}$ and is area minimizing. If $E=\mathbf{E}\left(T, \mathbf{C}_{4}\right)<\varepsilon_{11}$, then

$$
\begin{equation*}
\mathbf{e}_{T}(A) \leq C\left(E^{\gamma_{11}}+|A|^{\gamma_{11}}\right)\left(E+\mathbf{A}^{2}\right) \quad \text { for every Borel } A \subset B_{\frac{9}{8}} . \tag{6.1}
\end{equation*}
$$

This estimate complements (1.3) enabling to control the excess in the region where $\mathbf{d}>1$. We call it strong Almgren's estimate because a similar formula can be found in the big regularity paper (cf. [3, Sections 3.24-3.26 \& 3.30(8)]) and is a strengthened version of Proposition 5.4. To achieve (6.1) we construct a suitable competitor to estimate the size of the set $K$ where the graph of the $E^{\beta}$-Lipschitz approximation $f$ differs from $T$. Following Almgren, we embed $\mathcal{A}_{Q}$ in a large Euclidean space, via a biLipschitz embedding $\boldsymbol{\xi}$. We then regularize $\boldsymbol{\xi} \circ f$ by convolution and project it back onto $\mathcal{Q}=\boldsymbol{\xi}\left(\mathcal{A}_{Q}\right)$. To avoid loss of energy we need a rather special "almost projection" $\boldsymbol{\rho}_{\delta}^{\star}$.
Proposition 6.2. For every $\bar{n}, Q \in \mathbb{N} \backslash\{0\}$ there are geometric constants $\delta_{0}, C>0$ with the following property. For every $\delta \in] 0, \delta_{0}\left[\right.$ there is $\boldsymbol{\rho}_{\delta}^{\star}: \mathbb{R}^{N(Q, \bar{n})} \rightarrow \mathcal{Q}=\boldsymbol{\xi}\left(\mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}}\right)\right)$ such that $\left|\boldsymbol{\rho}_{\delta}^{\star}(P)-P\right| \leq C \delta^{8^{-\bar{n} Q}}$ for all $P \in \mathcal{Q}$ and, for every $u \in W^{1,2}\left(\Omega, \mathbb{R}^{N}\right)$, the following holds:

$$
\begin{equation*}
\int\left|D\left(\boldsymbol{\rho}_{\delta}^{\star} \circ u\right)\right|^{2} \leq\left(1+C \delta^{8^{-\bar{n} Q-1}}\right) \int_{\left\{\operatorname{dist}(u, \mathcal{Q}) \leq \delta^{\bar{n} Q+1}\right\}}|D u|^{2}+C \int_{\left\{\operatorname{dist}(u, \mathcal{Q})>\delta^{\bar{n} Q+1}\right\}}|D u|^{2} . \tag{6.2}
\end{equation*}
$$

The proof of Proposition 6.2 is postponed to the next section. Here we show Theorem 6.1 and hence conclude the Theorems 1.4 and 1.6. Theorem 1.3 enters crucially in the argument when estimating the second summand of (6.2) for the regularization of $\boldsymbol{\xi} \circ f$.
6.1. Regularization by convolution. Here we construct the competitor.

Proposition 6.3. Let $\beta_{1} \in\left(0, \frac{1}{2 m}\right)$ and $T$ be an area minimizing current satisfying $A s$ sumption 1.1 in $\mathbf{C}_{4}$. Let $f$ be its $E^{\beta_{1}}$-Lipschitz approximation. Then, there exist constants
$\bar{\varepsilon}_{12}, \gamma_{12}, C>0$ and a subset of radii $B \subset[9 / 8,2]$ with $|B|>1 / 2$ with the following properties. If $\mathbf{E}\left(T, \mathbf{C}_{4}\right) \leq \bar{\varepsilon}_{12}$, for every $\sigma \in B$, there exists a $Q$-valued function $g \in \operatorname{Lip}\left(B_{\sigma}, \mathcal{A}_{Q}\right)$ such that

$$
\left.g\right|_{\partial B_{\sigma}}=\left.f\right|_{\partial B_{\sigma}}, \quad \operatorname{Lip}(g) \leq C E^{\beta_{1}}, \quad \operatorname{spt}(g(x)) \subset \Sigma \quad \forall x \in B_{\sigma}
$$

and

$$
\begin{equation*}
\int_{B_{\sigma}}|D g|^{2} \leq \int_{B_{\sigma} \cap K}|D f|^{2}+C E^{\gamma_{12}}\left(E+\mathbf{A}^{2}\right) \tag{6.3}
\end{equation*}
$$

Proof. By Remark 1.5 we assume that $\Psi(0)=0,\|D \Psi\|_{0} \leq C\left(E^{1 / 2}+\mathbf{A}\right)$ and $\left\|D^{2} \Psi\right\|_{0} \leq C$ A. Since $|D f|^{2} \leq C \mathbf{d}_{T} \leq C E^{2 \beta_{1}} \leq 1$ on $K$, by Theorem 1.3 there is $q_{1}=2 p_{1}>2$ such that

$$
\begin{equation*}
\||D f|\|_{L^{q_{1}}\left(K \cap B_{2}\right)}^{2} \leq C E^{1-1 / p_{1}}\left(E+\mathbf{A}^{2}\right)^{1 / p_{1}} \leq C\left(E+\mathbf{A}^{2}\right) \tag{6.4}
\end{equation*}
$$

Given two (vector-valued) functions $h_{1}$ and $h_{2}$ and two radii $0<\bar{r}<r$, we denote by $\operatorname{lin}\left(h_{1}, h_{2}\right)$ the linear interpolation in $B_{r} \backslash \bar{B}_{\bar{r}}$ between $\left.h_{1}\right|_{\partial B_{r}}$ and $\left.h_{2}\right|_{\partial B_{\bar{r}}}$. More precisely, if $(\theta, t) \in \mathbb{S}^{m-1} \times[0, \infty)$ are spherical coordinates, then

$$
\operatorname{lin}\left(h_{1}, h_{2}\right)(\theta, t)=\frac{r-t}{r-\bar{r}} h_{2}(\theta, t)+\frac{t-s}{r-\bar{r}} h_{1}(\theta, t)
$$

Next, let $\delta>0$ and $\varepsilon>0$ be two parameters and let $1<r_{1}<r_{2}<r_{3}<2$ be three radii, all to be chosen later. To keep the notation simple, we will write $\boldsymbol{\rho}^{\star}$ in place of $\boldsymbol{\rho}_{\delta}^{\star}$. Let $\varphi \in C_{c}^{\infty}\left(B_{1}\right)$ be a standard (nonnegative!) mollifier. We also use the notation $f(x)=$ $\left(f_{1}(x), f_{2}(x)\right) \in \mathcal{A}_{Q}\left(\mathbb{R}^{\bar{n}} \times \mathbb{R}^{l}\right)$ meaning that $f(x)=\sum_{i} \llbracket\left(f_{1}^{i}(x), f_{2}^{i}(x)\right) \rrbracket$ with $\left(f_{1}^{i}(x), f_{2}^{i}(x)\right) \in$ $\mathbb{R}^{\bar{n}} \times \mathbb{R}^{l}$ and the maps $f_{1}$ and $f_{2}$ are then given by $f_{j}(x)=\sum_{i} \llbracket f_{j}^{i}(x) \rrbracket$. This does not create confusion in "ordering the sheets": since the points $f^{i}(x)$ belong to $\Sigma$ we have indeed the relation $f_{2}^{j}(x)=\Psi\left(x, f_{1}^{j}(x)\right)$. We moreover set $f^{\prime}:=\boldsymbol{\xi} \circ f_{1}$. Recall the map $\boldsymbol{\rho}$ of [14, Theorem 2.1] and define:

$$
g^{\prime}:= \begin{cases}\sqrt{E} \boldsymbol{\rho} \circ \operatorname{lin}\left(\frac{f^{\prime}}{\sqrt{E}}, \boldsymbol{\rho}^{\star}\left(\frac{f^{\prime}}{\sqrt{E}}\right)\right) & \text { in } B_{r_{3}} \backslash B_{r_{2}}  \tag{6.5}\\ \sqrt{E} \boldsymbol{\rho} \circ \operatorname{lin}\left(\boldsymbol{\rho}^{\star}\left(\frac{f^{\prime}}{\sqrt{E}}\right), \boldsymbol{\rho}^{\star}\left(\frac{f^{\prime}}{\sqrt{E}} * \varphi_{\varepsilon}\right)\right) & \text { in } B_{r_{2}} \backslash B_{r_{1}} \\ \sqrt{E} \boldsymbol{\rho}^{\star}\left(\frac{f^{\prime}}{\sqrt{E}} * \varphi_{\varepsilon}\right) & \text { in } B_{r_{1}}\end{cases}
$$

Finally set $g_{1}:=\boldsymbol{\xi}^{-1} \circ g^{\prime}$ and $g:=\sum_{i} \llbracket\left(g_{1}^{i}, \Psi\left(x, g_{1}^{i}\right)\right) \rrbracket$. We claim that, for $\sigma:=r_{3}$ in a suitable set $B \subset[9 / 8,2]$ with $|B|>1 / 2$, we can choose $r_{2}=r_{3}-s$ and $r_{1}=r_{2}-s$ so that $g$ satisfies the conclusion of the proposition. Some computations will be simplified taking into account that our choice of the parameters will imply the following inequalities:

$$
\begin{equation*}
\delta^{2 \cdot 8^{-\bar{n} Q}} \leq s, \quad \varepsilon \leq s \quad \text { and } \quad E^{1-2 \beta_{1}} \leq \varepsilon^{m} \tag{6.6}
\end{equation*}
$$

We start noticing that clearly $\left.g\right|_{\partial B_{r_{3}}}=\left.f\right|_{\partial B_{r_{3}}}$. As for the Lipschitz constant, it suffices to estimate the Lipschitz constant of $g^{\prime}$. This can be easily done observing that:

$$
\begin{cases}\operatorname{Lip}\left(g^{\prime}\right) \leq C \operatorname{Lip}\left(f^{\prime} * \varphi_{\varepsilon}\right) \leq C \operatorname{Lip}\left(f^{\prime}\right) \leq C E^{\beta_{1}} & \text { in } B_{r_{1}}, \\ \operatorname{Lip}\left(g^{\prime}\right) \leq C \operatorname{Lip}\left(f^{\prime}\right)+C \frac{\left\|f^{\prime}-f^{\prime} * \varphi_{\varepsilon}\right\|_{L^{\infty}}}{s} \leq C\left(1+\frac{\varepsilon}{s}\right) \operatorname{Lip}\left(f^{\prime}\right) \leq C E^{\beta_{1}} & \text { in } B_{r_{2}} \backslash B_{r_{1}}, \\ \operatorname{Lip}\left(g^{\prime}\right) \leq C \operatorname{Lip}\left(f^{\prime}\right)+C E^{1 / 2} \frac{\delta^{8} \bar{n} Q}{s} \leq C E^{\beta_{1}}+C E^{1 / 2} \leq C E^{\beta_{1}} & \text { in } B_{r_{3}} \backslash B_{r_{2}} .\end{cases}
$$

In the first inequality of the last line we have used that, since $\mathcal{Q}$ is a cone, $E^{-1 / 2} f^{\prime}(x) \in \mathcal{Q}$ for every $x$ : therefore $\left|\rho^{\star}\left(f^{\prime} / E^{1 / 2}\right)-f^{\prime} / E^{1 / 2}\right| \leq C \delta^{8^{-\bar{n} Q}}$. We pass now to estimate the Dirichlet energy of $g$.

Step 1. Energy in $B_{r_{3}} \backslash B_{r_{2}}$. By Section 1.5, the energy of the first component $g_{1}$ coincides with the (classical!) Dirichlet energy of $g^{\prime}$. By Proposition 6.2, $\left|\rho^{\star}(P)-P\right| \leq C \delta^{8^{-\bar{n} Q}}$ for all $P \in \mathcal{Q}$. Thus, elementary estimates on the linear interpolation give

$$
\begin{align*}
\int_{B_{r_{3}} \backslash B_{r_{2}}}\left|D g^{\prime}\right|^{2} \leq & \frac{C E}{\left(r_{3}-r_{2}\right)^{2}} \int_{B_{r_{3}} \backslash B_{r_{2}}}\left|\frac{f^{\prime}}{\sqrt{E}}-\boldsymbol{\rho}^{\star}\left(\frac{f^{\prime}}{\sqrt{E}}\right)\right|^{2}+C \int_{B_{r_{3}} \backslash B_{r_{2}}}\left|D f^{\prime}\right|^{2} \\
& +C \int_{B_{r_{3} \backslash B_{r_{2}}}\left|D\left(\boldsymbol{\rho}^{\star} \circ f^{\prime}\right)\right|^{2} \leq C \int_{B_{r_{3}} \backslash B_{r_{2}}}\left|D f^{\prime}\right|^{2}+C E s^{-1} \delta^{2 \cdot 8^{-\bar{n} Q}}} . \tag{6.7}
\end{align*}
$$

As for $g_{2}$, we compute $D g_{2}^{i}(x)=D_{x} \Psi\left(x, g_{1}^{i}(x)\right)+D_{u} \Psi\left(x, g_{1}^{i}(x)\right) D g_{1}^{i}(x)$ and so

$$
\begin{equation*}
\int_{B_{r_{3}} \backslash B_{r_{2}}}\left|D g_{2}\right|^{2} \leq C s\left(E+\mathbf{A}^{2}\right) \tag{6.8}
\end{equation*}
$$

where we used the estimate $\left\|D g_{2}\right\|_{0} \leq C\|D \Psi\|_{0} \leq C\left(E^{1 / 2}+\mathbf{A}\right)$.
Step 2. Energy in $B_{r_{2}} \backslash B_{r_{1}}$. Here, using the same interpolation inequality and a standard estimate on convolutions of $W^{1,2}$ functions, we get

$$
\begin{align*}
& \int_{B_{r_{2}} \backslash B_{r_{1}}}\left|D g^{\prime}\right|^{2} \leq C \int_{B_{r_{2}+\varepsilon} \backslash B_{r_{1}-\varepsilon}}\left|D f^{\prime}\right|^{2}+\frac{C}{\left(r_{2}-r_{1}\right)^{2}} \int_{B_{r_{2}} \backslash B_{r_{1}}}\left|f^{\prime}-\varphi_{\varepsilon} * f^{\prime}\right|^{2} \\
\leq & C \int_{B_{r_{2}+\varepsilon} \backslash B_{r_{1}-\varepsilon}}\left|D f^{\prime}\right|^{2}+C \varepsilon^{2} s^{-2} \int_{B_{3}}\left|D f^{\prime}\right|^{2} \leq C \int_{B_{r_{2}+\varepsilon} \backslash B_{r_{1}-\varepsilon}}\left|D f^{\prime}\right|^{2}+C \varepsilon^{2} E s^{-2} . \tag{6.9}
\end{align*}
$$

Similarly, for the second component we have that

$$
\begin{equation*}
\int_{B_{r_{2}+\varepsilon} \backslash B_{r_{1}-\varepsilon}}\left|D g_{2}\right|^{2} \leq C\left(\mathbf{A}^{2}+E\right) s \tag{6.10}
\end{equation*}
$$

Step 3. Energy in $B_{r_{1}}$. Define $Z:=\left\{\operatorname{dist}\left(\frac{f^{\prime}}{\sqrt{E}} * \varphi_{\varepsilon}, \mathcal{Q}\right)>\delta^{\bar{n} Q+1}\right\}$ and use (6.2) to get

$$
\begin{equation*}
\int_{B_{r_{1}}}\left|D g^{\prime}\right|^{2} \leq\left(1+C \delta^{8^{-\bar{n} Q-1}}\right) \int_{B_{r_{1}} \backslash Z}\left|D\left(f^{\prime} * \varphi_{\varepsilon}\right)\right|^{2}+C \int_{Z}\left|D\left(f^{\prime} * \varphi_{\varepsilon}\right)\right|^{2}=: I_{1}+I_{2} \tag{6.11}
\end{equation*}
$$

We consider $I_{1}$ and $I_{2}$ separately. For $I_{1}$ we first observe the elementary inequality

$$
\begin{align*}
\left\|D\left(f^{\prime} * \varphi_{\varepsilon}\right)\right\|_{L^{2}}^{2} \leq\left\|\left|D f^{\prime}\right| * \varphi_{\varepsilon}\right\|_{L^{2}}^{2} \leq & \left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right) * \varphi_{\varepsilon}\right\|_{L^{2}}^{2}+\left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K^{c}}\right) * \varphi_{\varepsilon}\right\|_{L^{2}}^{2} \\
& +2\left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right) * \varphi_{\varepsilon}\right\|_{L^{2}}\left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K^{c}}\right) * \varphi_{\varepsilon}\right\|_{L^{2}} \tag{6.12}
\end{align*}
$$

where $K^{c}$ is the complement of $K$ in $B_{3}$. Recalling $r_{1}+\varepsilon \leq r_{1}+s=r_{2}$ we estimate the first summand in (6.12) as follows:

$$
\begin{equation*}
\left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right) * \varphi_{\varepsilon}\right\|_{L^{2}\left(B_{r_{1}}\right)}^{2} \leq \int_{B_{r_{1}+\varepsilon}}\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right)^{2} \leq \int_{B_{r_{2} \cap K} \cap}\left|D f^{\prime}\right|^{2} \tag{6.13}
\end{equation*}
$$

To treat the other terms recall that $\operatorname{Lip}\left(f^{\prime}\right) \leq C E^{\beta_{1}}$ and $\left|K^{c}\right| \leq C E^{1-2 \beta_{1}}$ :

$$
\begin{equation*}
\left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K^{c}}\right) * \varphi_{\varepsilon}\right\|_{L^{2}\left(B_{r_{1}}\right)}^{2} \leq C E^{2 \beta_{1}}\left\|\mathbf{1}_{K^{c}} * \varphi_{\varepsilon}\right\|_{L^{2}}^{2} \leq C E^{2 \beta_{1}}\left\|\mathbf{1}_{K^{c}}\right\|_{L^{1}}^{2}\left\|\varphi_{\varepsilon}\right\|_{L^{2}}^{2} \leq \frac{C E^{2-2 \beta_{1}}}{\varepsilon^{m}} \tag{6.14}
\end{equation*}
$$

Putting (6.13) and (6.14) in (6.12) and recalling $E^{1-2 \beta_{1}} \leq \varepsilon^{m}$ and $\int\left|D f^{\prime}\right|^{2} \leq C E$, we get

$$
\begin{equation*}
I_{1} \leq \int_{B_{r_{2}} \cap K}\left|D f^{\prime}\right|^{2}+C \delta^{8-\bar{n} Q-1} E+C \varepsilon^{-m / 2} E^{3 / 2-\beta_{1}} \tag{6.15}
\end{equation*}
$$

For what concerns $I_{2}$, first we argue as for $I_{1}$, splitting in $K$ and $K^{c}$, to deduce that

$$
\begin{equation*}
I_{2} \leq C \int_{Z}\left(\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right) * \varphi_{\varepsilon}\right)^{2}+C \varepsilon^{-m / 2} E^{3 / 2-\beta_{1}} \tag{6.16}
\end{equation*}
$$

Then, regarding the first summand in (6.16), we note that

$$
\begin{equation*}
|Z| \delta^{2 \bar{n} Q+2} \leq \int_{B_{r_{1}}}\left|\frac{f^{\prime}}{\sqrt{E}} * \varphi_{\varepsilon}-\frac{f^{\prime}}{\sqrt{E}}\right|^{2} \leq C \varepsilon^{2} \tag{6.17}
\end{equation*}
$$

Since $\left|D f^{\prime}\right| \leq|D f|$ (and recalling that $q_{1}=2 p_{1}>2$ ), we use (6.4) to obtain

$$
\begin{align*}
\int_{Z}\left(\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right) * \varphi_{\varepsilon}\right)^{2} & \leq|Z|^{\frac{p_{1}-1}{p_{1}}}\left\|\left(\left|D f^{\prime}\right| \mathbf{1}_{K}\right) * \varphi_{\varepsilon}\right\|_{L^{q_{1}}}^{2} \leq C\left(\frac{\varepsilon}{\delta^{\bar{n} Q+1}}\right)^{\frac{2\left(p_{1}-1\right)}{p_{1}}}\left\|\mid D f^{\prime}\right\|_{L^{q_{1}}(K)}^{2} \\
& \leq C\left(\frac{\varepsilon}{\delta^{\bar{n} Q+1}}\right)^{\frac{2\left(p_{1}-1\right)}{p_{1}}}\left(E+\mathbf{A}^{2}\right) \tag{6.18}
\end{align*}
$$

Gathering all the estimates together, (6.11), (6.15), (6.16) and (6.18) give

$$
\begin{equation*}
\int_{B_{r_{1}}}\left|D g^{\prime}\right|^{2} \leq \int_{B_{r_{2}} \cap K}\left|D f^{\prime}\right|^{2}+C\left(E \delta^{8-\bar{n} Q-1}+\frac{E^{3 / 2-\beta_{1}}}{\varepsilon^{m / 2}}+\left(E+\mathbf{A}^{2}\right)\left(\frac{\varepsilon}{\delta^{\bar{n} Q+1}}\right)^{\frac{2\left(p_{1}-1\right)}{p_{1}}}\right) \tag{6.19}
\end{equation*}
$$

On the other hand, for what concerns $g_{2}$ we can estimate as follows

$$
\begin{align*}
& \int_{B_{r_{1}}}\left|D g_{2}\right|^{2}=\int_{B_{r_{1}}}\left|D f_{2}\right|^{2}+\sum_{i} \int_{B_{r_{1}}}\left(D g_{2}^{i}-D f_{2}^{i}\right) \cdot\left(D g_{2}^{i}+D f_{2}^{i}\right) \\
\leq & \int_{B_{r_{1} \cap K}}\left|D f_{2}\right|^{2}+\int_{B_{r_{1} \backslash K} \backslash}\left|D f_{2}\right|^{2}+C\left(\mathbf{A}+E^{\frac{1}{2}}\right) \sum_{i} \int_{B_{r_{1}}}\left|D g_{2}^{i}-D f_{2}^{i}\right| \tag{6.20}
\end{align*}
$$

We already observed that $\left|D f_{2}\right| \leq C\left(\mathbf{A}+E^{1 / 2}\right)$, leading to the estimate $\int_{K^{c}}\left|D f_{2}\right|^{2} \leq$ $C\left(\mathbf{A}^{2}+E\right)\left|K^{c}\right| \leq C\left(\mathbf{A}^{2}+E\right) E^{1-2 \beta_{1}}$. As for the latter summand we compute

$$
\begin{aligned}
\left|D g_{2}^{i}-D f_{2}^{i}\right| \leq & \left|D_{x} \Psi\left(x, g_{1}^{i}\right)-D_{x} \Psi\left(x, f_{1}^{i}\right)\right| \\
& +\left|D_{u} \Psi\left(x, g_{1}^{i}(x)\right) D g_{1}^{i}\right|+\left|D_{u} \Psi\left(x, f_{1}^{i}(x)\right) D f_{1}^{i}\right| \\
\leq & C \mathbf{A} \mathcal{G}\left(g_{1}, f_{1}\right)+C\left(\mathbf{A}+E^{1 / 2}\right) E^{\beta_{1}}
\end{aligned}
$$

We next estimate $\left\|\mathcal{G}\left(g_{1}, f_{1}\right)\right\|_{\infty} \leq C g^{\prime \prime}-f^{\prime} \|_{\infty}$ and

$$
\begin{aligned}
\left\|g^{\prime}-f^{\prime}\right\|_{\infty} & \leq C \sqrt{E}\left(\left\|\boldsymbol{\rho}^{*}\left(\frac{f^{\prime}}{\sqrt{E}} * \varphi_{\varepsilon}\right)-\boldsymbol{\rho}^{*}\left(\frac{f^{\prime}}{\sqrt{E}}\right)\right\|_{\infty}+\left\|\boldsymbol{\rho}^{*}\left(\frac{f^{\prime}}{\sqrt{E}}\right)-\frac{f^{\prime}}{\sqrt{E}}\right\|_{\infty}\right) \\
& \leq C \operatorname{Lip}\left(\boldsymbol{\rho}^{*}\right)\left\|f^{\prime} * \varphi_{\varepsilon}-f^{\prime}\right\|_{L^{\infty}}+C E^{1 / 2} \delta^{8^{-\bar{n} Q}} \leq C E^{\beta_{1}} \varepsilon+C E^{1 / 2} \delta^{8^{-\bar{n} Q}} \leq C E^{\beta_{1}}
\end{aligned}
$$

We therefore conclude

$$
\begin{equation*}
\int_{B_{r_{1}}}\left|D g_{2}\right|^{2} \leq \int_{B_{r_{1} \cap K}}\left|D f_{2}\right|^{2}+C\left(\mathbf{A}^{2}+E\right) E^{\beta_{1}} \tag{6.21}
\end{equation*}
$$

Final estimate. Since $|D g|^{2}=\left|D g^{\prime}\right|^{2}+\left|D g_{2}\right|^{2}$, summing (6.7), (6.8), (6.9), (6.10), (6.19) and (6.21) (and recalling $\varepsilon<s$ ), we conclude

$$
\begin{aligned}
\int_{B_{r_{3}}}|D g|^{2} \leq & \int_{B_{r_{1} \cap K}}|D f|^{2}+C \int_{B_{r_{1}+3 s} \backslash B_{r_{1}-s}}\left|D f^{\prime}\right|^{2}+C\left(\mathbf{A}+E^{2}\right)\left(s+E^{\beta_{1}}\right) \\
& +C E\left(\delta^{8^{-\bar{n} Q-1}}+\frac{\varepsilon^{2}}{s^{2}}+\frac{\delta^{2 \cdot 8^{-\bar{n} Q}}}{s}+\frac{E^{1 / 2-\beta_{1}}}{\varepsilon^{m / 2}}+\left(1+\mathbf{A}^{2} E^{-1}\right)\left(\frac{\varepsilon}{\delta^{\bar{n} Q+1}}\right)^{\frac{2\left(p_{1}-1\right)}{p_{1}}}\right) .
\end{aligned}
$$

We set $\varepsilon=E^{a}, \delta=E^{b}$ and $s=E^{c}$, where

$$
a=\frac{1-2 \beta_{1}}{2 m}, \quad b=\frac{1-2 \beta_{1}}{4 m(\bar{n} Q+1)} \quad \text { and } \quad c=\frac{1-2 \beta_{1}}{8^{\bar{n} Q} 4 m(\bar{n} Q+1)} .
$$

This choice respects (6.6). Assume $E$ is small enough so that $s \leq \frac{1}{16}$. Now, if $C>0$ is a sufficiently large constant, there is a set $B^{\prime} \subset\left[9 / 8, \frac{29}{16}\right]$ with $\left|B^{\prime}\right|>1 / 2$ such that,

$$
\int_{B_{r_{1}+3 s} \backslash B_{r_{1}-s}}\left|D f^{\prime}\right|^{2} \leq C s \int_{B_{2}}\left|D f^{\prime}\right|^{2} \leq C E^{1+c} \quad \text { for every } r_{1} \in B^{\prime}
$$

Indeed by integrating in polar coordinates and by Fubini's Theorem we have that

$$
\begin{aligned}
\int_{\frac{9}{8}}^{\frac{29}{16}} d r \int_{B_{r+3 s} \backslash B_{r-s}}\left|D f^{\prime}\right|^{2} & =\int_{\frac{9}{8}}^{\frac{29}{16}} d r \int_{r-s}^{r+3 s} d t \int_{\partial B_{t}}\left|D f^{\prime}\right|^{2} d \mathcal{H}^{n-1} \\
& \leq 4 s \int_{\frac{9}{8}-s}^{2} d t \int_{\partial B_{t}}\left|D f^{\prime}\right|^{2} d \mathcal{H}^{n-1} \leq 4 s \int_{B_{2}}\left|D f^{\prime}\right|^{2}
\end{aligned}
$$

from which the conclusion follows for $C$ big enough:

$$
\begin{aligned}
& \left|\left\{r \in\left[\frac{9}{8}, \frac{29}{16}\right]: \int_{B_{r+3 s} \backslash B_{r-s}}\left|D f^{\prime}\right|^{2} \geq C s \int_{B_{2}}\left|D f^{\prime}\right|^{2}\right\}\right| \\
& \quad \leq \frac{1}{C s \int_{B_{2}}\left|D f^{\prime}\right|^{2}} \int_{\frac{9}{8}}^{\frac{29}{16}} d r \int_{B_{r+3 s} \backslash B_{r-s}}\left|D f^{\prime}\right|^{2} \leq \frac{4}{C}<\frac{1}{8} .
\end{aligned}
$$

For $\sigma=r_{3} \in B=2 s+B^{\prime}$ we then conclude, for some $\gamma\left(\beta_{1}, \bar{n}, n, m, Q\right)>0$,

$$
\int_{B_{\sigma}}|D g|^{2} \leq \int_{B_{\sigma} \cap K}|D f|^{2}+C E^{\gamma}\left(E+\mathbf{A}^{2}\right)
$$

6.2. Proof of Theorem 6.1. Choose $\beta_{1}=\frac{1}{4 m}$ and consider the set $B \subset[9 / 8,2]$ given in Proposition 6.3. Using the coarea formula and the isoperimetric inequality (the argument and the map $\varphi$ are the same in the proof of Theorem 4.2 and that of Proposition 5.4), we find $s \in B$ and an integer rectifiable current $R$ such that

$$
\partial R=\left\langle T-\mathbf{G}_{f}, \varphi, s\right\rangle \quad \text { and } \quad \mathbf{M}(R) \leq C E^{\frac{2 m-1}{2 m-2}} .
$$

Since $\left.g\right|_{\partial B_{s}}=\left.f\right|_{\partial B_{s}}$ and $g$ takes values in $\Sigma$, we can use $g$ in place of $f$ in the estimates and, arguing as before (see e.g. the proof of Theorem 5.4), we get, for a suitable $\gamma>0$ :

$$
\begin{equation*}
\|T\|\left(\mathbf{C}_{s}\right) \leq Q\left|B_{s}\right|+\int_{B_{s}} \frac{|D g|^{2}}{2}+C E^{1+\gamma} \stackrel{(6.3)}{\leq} Q\left|B_{s}\right|+\int_{B_{s} \cap K} \frac{|D f|^{2}}{2}+C E^{\gamma}\left(E+\mathbf{A}^{2}\right) . \tag{6.22}
\end{equation*}
$$

On the other hand, by Taylor's expansion in Remark 4.4,

$$
\begin{align*}
\|T\|\left(\mathbf{C}_{s}\right) & =\|T\|\left(\left(B_{s} \backslash K\right) \times \mathbb{R}^{n}\right)+\left\|\mathbf{G}_{f}\right\|\left(\left(B_{s} \cap K\right) \times \mathbb{R}^{n}\right) \\
& \geq\|T\|\left(\left(B_{s} \backslash K\right) \times \mathbb{R}^{n}\right)+Q\left|K \cap B_{s}\right|+\int_{K \cap B_{s}} \frac{|D f|^{2}}{2}-C E^{1+\gamma} . \tag{6.23}
\end{align*}
$$

Hence, from (6.22) and (6.23), we get $\mathbf{e}_{T}\left(B_{s} \backslash K\right) \leq C E^{\gamma}\left(E+\mathbf{A}^{2}\right)$.
This is enough to conclude the proof. Indeed, let $A \subset B_{9 / 8}$ be a Borel set. Using the higher integrability of $|D f|$ in $K$ (see (6.4)) and possibly selecting a smaller $\gamma>0$, we get

$$
\begin{aligned}
\mathbf{e}_{T}(A) & \leq \mathbf{e}_{T}(A \cap K)+\mathbf{e}_{T}(A \backslash K) \leq \int_{A \cap K} \frac{|D f|^{2}}{2}+C E^{\gamma}\left(E+\mathbf{A}^{2}\right) \\
& \leq C|A \cap K|^{\frac{p_{1}-1}{p_{1}}}\left(\int_{A \cap K}|D f|^{q_{1}}\right)^{2 / q_{1}}+C E^{\gamma}\left(E+\mathbf{A}^{2}\right) \\
& \leq C|A|^{\frac{p_{1}-1}{p_{1}}}\left(E+\mathbf{A}^{2}\right)+C E^{\gamma}\left(E+\mathbf{A}^{2}\right)
\end{aligned}
$$

6.3. Proofs of Theorems 1.4 and 1.6. As usual we assume, w.l.o.g., $r=1$ and $x=0$. Choose $\beta_{11}<\min \left\{\frac{1}{2 m}, \frac{\gamma_{11}}{2\left(1+\gamma_{11}\right)}\right\}$, where $\gamma_{11}$ is the constant in Theorem 6.1. Let $f$ be the $E^{\beta_{11}}$-Lipschitz approximation of $T$. Clearly (1.4) follows directly from Proposition 2.2 if $\gamma_{1}<\beta_{11}$. Set next $A:=\left\{\mathbf{m e}_{T}>2^{-m} E^{2 \beta_{11}}\right\} \cap B_{9 / 8}$. By Proposition 2.2, $|A| \leq C E^{1-2 \beta_{11}}$. If $\varepsilon_{1}$ is sufficiently small, apply (2.1) and estimate (6.1) to $A$ to conclude:

$$
\left|B_{1} \backslash K\right| \leq C E^{-2 \beta_{11}} \mathbf{e}_{T}(A) \leq C E^{\gamma_{11}-2 \beta_{11}\left(1+\gamma_{11}\right)}\left(E+\mathbf{A}^{2}\right)
$$

By our choice of $\gamma_{11}$ and $\beta_{11}$, this gives (1.5) for some positive $\gamma_{1}$. Finally, set $S=\mathbf{G}_{f}$. Recalling the strong Almgren estimate (6.1) and the Taylor expansion in Remark 4.4, we conclude: for every $0<\sigma \leq 1$

$$
\begin{gathered}
\left|\|T\|\left(\mathbf{C}_{\sigma}\right)-Q \sigma^{m} \omega_{m}-\int_{B_{\sigma}} \frac{|D f|^{2}}{2}\right| \leq \mathbf{e}_{T}\left(B_{\sigma} \backslash K\right)+\mathbf{e}_{S}\left(B_{\sigma} \backslash K\right)+\left|\mathbf{e}_{S}\left(B_{\sigma}\right)-\int_{B_{\sigma}} \frac{|D f|^{2}}{2}\right| \\
\leq C E^{\gamma_{11}}\left(E+\mathbf{A}^{2}\right)+C\left|B_{\sigma} \backslash K\right|+C \operatorname{Lip}(f)^{2} \int_{B_{\sigma}}|D f|^{2} \leq C E^{\gamma_{1}}\left(E+\mathbf{A}^{2}\right)
\end{gathered}
$$

The $L^{\infty}$ bound follows from Proposition 2.2 recalling that, by Remark 1.5, we can assume $\|D \Psi\|_{0} \leq C\left(E^{1 / 2}+\mathbf{A}\right)$. Finally, Theorem 1.6 is a special case of Theorem 4.2, since the map $f$ in Theorem 1.4 is the $E^{\gamma_{1}}$-Lipschitz approximation of $T$.

## 7. The "Almost" Projections $\boldsymbol{\rho}_{\delta}^{\star}$

In this section we show the existence of the maps $\boldsymbol{\rho}_{\delta}^{\star}$ in Proposition 6.2. Compared to the original ones introduced by Almgren, our $\boldsymbol{\rho}_{\delta}^{\star}$ 's have the advantage of depending on a single parameter. Our proof is different from Almgren's and gives more explicit estimates, relying heavily on the following simple corollary of Kirszbraun's Theorem.

Lemma 7.1. Let $f: \Omega \subset \mathbb{R}^{N_{1}} \rightarrow C \subset \mathbb{R}^{N_{2}}$ be a Lipschitz function and assume that $C$ is closed and convex. Then, there is an extension $\hat{f}$ of $f$ to the whole $\mathbb{R}^{N_{1}}$ which preserves the Lipschitz constant and takes values in $C$.

To prove Lemma 7.1 it suffices to take the map $\tilde{f}$ of the classical statement of Kirszbraun's theorem (see [20, Theorem 2.10.43]) which takes values in $\mathbb{R}^{N_{2}}$ and compose it with the orthogonal projection $\pi_{C}$ onto the convex closed set $C$, which is a 1-Lipschitz map in $\mathbb{R}^{N_{2}}$.
Proof of Proposition 6.2. The proof consists of four parts: the first one is a detailed description of the set $\mathcal{Q}$, whereas the remaining three give a rather explicit construction in this order:
(1) first we specify $\boldsymbol{\rho}_{\delta}^{\star}$ on $\mathcal{Q}$ : the resulting map will be called $\boldsymbol{\rho}^{\text {b }}$;
(2) then we extend it to a map $\boldsymbol{\rho}^{\sharp}$ on $\mathcal{Q}_{\delta^{n Q+1}}$, the $\delta^{n Q+1}$-neighborhood of $\mathcal{Q}$; $\boldsymbol{\rho}^{\sharp}$ will satisfy $\operatorname{Lip}\left(\boldsymbol{\rho}^{\sharp}\right) \leq 1+C \delta^{8^{-\bar{n} Q-1}}$ and $\left|\boldsymbol{\rho}^{\sharp}(P)-P\right| \leq C \delta^{8^{-\bar{n} Q}}$ for every $p \in \mathcal{Q}$;
(3) we then extend it to all $\mathbb{R}^{N}$ keeping its Lipschitz constant bounded.
(3) follows easily from (2): we consider $\boldsymbol{\xi}^{-1} \circ \boldsymbol{\rho}^{\sharp}: \mathcal{Q}_{\delta^{\bar{n} Q+1}} \rightarrow \mathcal{A}_{Q}$ and a Lipschitz extension $h: \mathbb{R}^{N} \rightarrow \mathcal{A}_{Q}$ of it with $\operatorname{Lip}(h) \leq C$, using [14, Theorem 1.7]. Our map is then $\boldsymbol{\rho}_{\delta}^{\star}:=\boldsymbol{\xi} \circ h$. Then (6.2) is an easy consequence of (2), (3) and the chain rule.

The description of $\mathcal{Q}$ and the proofs of (1) and (2) are given in the next subsections.
From now on we use $n$ instead of $\bar{n}$ to simplify the notation.
7.1. Conical simplicial structure of $\mathcal{Q}$. We first prove that $\mathcal{Q}$ is the union of families $\left\{\mathcal{F}_{i}\right\}_{i=0}^{n Q}$ of sets, the " $i$-dimensional faces" of $\mathcal{Q}$, with the following properties:
(p1) $\mathcal{Q}=\cup_{i} \cup_{F \in \mathcal{F}_{i}} F$;
(p2) $\mathcal{F}:=\cup_{i} \mathcal{F}_{i}$ is a collection of finitely many disjoint sets;
(p3) each face $F \in \mathcal{F}_{i}$ is a convex open $i$-dimensional cone, where open means that for every $x \in F$ there exists an $i$-dimensional disk $D$ with $x \in D \subset F$;
(p4) for each $F \in \mathcal{F}_{i}, \bar{F} \backslash F$ is the union of some elements of $\cup_{j<i} \mathcal{F}_{j}$.
(p5) for each $i<k \leq n Q$ and for each $F \in \mathcal{F}_{i}$, there exists $G \in \mathcal{F}_{k}$ such that $F \subset \bar{G}$.
Remark 7.2. With a small abuse of notation $\partial F$ will denote $\bar{F} \backslash F$ for any $F \in \mathcal{F}$.
So, $\mathcal{F}_{0}=\{0\} ; \mathcal{F}_{1}$ consists of finitely many half-lines meeting at 0 , i.e. of sets of type $l_{v}=\{\lambda v: \lambda \in] 0,+\infty[ \}$ for $v \in \mathbb{S}^{N-1} ; \mathcal{F}_{2}$ consists of finitely many 2-dimensional "infinite
triangles" delimited by pairs of half lines $l_{v_{1}}, l_{v_{2}} \in \mathcal{F}_{1}$ and by $\{0\}$; and so on. To prove this statement, first of all we recall the construction of $\boldsymbol{\xi}$ (see [14, 2.1.2]). After selecting a suitable finite collection of non zero vectors $\left\{e_{k}\right\}_{k=1}^{h}$ (in general $h>n$ ), we define the linear map $L: \mathbb{R}^{n Q} \rightarrow \mathbb{R}^{N}$ with $N:=h Q>n Q$ given by

$$
L\left(P_{1}, \ldots, P_{Q}\right):=(\underbrace{P_{1} \cdot e_{1}, \ldots, P_{Q} \cdot e_{1}}_{w^{1}}, \underbrace{P_{1} \cdot e_{2}, \ldots, P_{Q} \cdot e_{2}}_{w^{2}}, \ldots, \underbrace{P_{1} \cdot e_{h}, \ldots, P_{Q} \cdot e_{h}}_{w^{h}})
$$

Then, we consider the map $O: \mathbb{R}^{N} \rightarrow \mathbb{R}^{N}$ which maps ( $w^{1} \ldots, w^{h}$ ) into the vector $\left(v^{1}, \ldots, v^{h}\right)$ where each $v^{i}$ is obtained from $w^{i}$ ordering its components in increasing order. Note that the composition $O \circ L:\left(\mathbb{R}^{n}\right)^{Q} \rightarrow \mathbb{R}^{N}$ is now invariant under the action of the symmetric group $\mathscr{P}_{Q}$. Therefore, $\boldsymbol{\xi}$ is simply the induced map on $\mathcal{A}_{Q}=\left(\mathbb{R}^{n}\right)^{Q} / \mathscr{P}_{Q}$ and $\mathcal{Q}=\boldsymbol{\xi}\left(\mathcal{A}_{Q}\right)=O(V)$ where $V:=L\left(\mathbb{R}^{n Q}\right)$. Moreover, since the vectors $e_{i}$ 's span $\mathbb{R}^{n}$ (cf. [14, 2.1.2]), the map $L$ is injective and thus $V$ is an $n Q$-dimensional subspace.

Consider the following equivalence relation $\sim$ on $V$ :

$$
\left(w^{1}, \ldots, w^{h}\right) \sim\left(z^{1}, \ldots, z^{h}\right) \quad \text { if } \quad\left\{\begin{array}{l}
w_{j}^{i}=w_{k}^{i}  \tag{7.1}\\
w_{j}^{i}>w_{k}^{i}
\end{array} \Leftrightarrow z_{j}^{i}=z_{k}^{i}>z_{k}^{i} \quad \forall i, j, k\right.
$$

where $w^{i}=\left(w_{1}^{i}, \ldots, w_{Q}^{i}\right)$ and $z^{i}=\left(z_{1}^{i}, \ldots, z_{Q}^{i}\right)$ : if $w \sim z$, then $O$ rearranges their components with the same permutation. We let $\mathcal{E}$ denote the set of corresponding equivalence classes in $V$ and $\mathcal{C}:=\left\{L^{-1}(E): E \in \mathcal{E}\right\}$. The following fact is an obvious consequence of definition (7.1):

$$
L(P) \sim L(S) \quad \text { if and only if } \quad L\left(P_{\pi(1)}, \ldots, P_{\pi(Q)}\right) \sim L\left(S_{\pi(1)}, \ldots, S_{\pi(Q)}\right) \quad \forall \pi \in \mathscr{P}_{Q}
$$

Thus, $\pi(C) \in \mathcal{C}$ for every $C \in \mathcal{C}$ and every $\pi \in \mathscr{P}_{Q}$. Since $\boldsymbol{\xi}$ is injective and is induced by $O \circ L$, it follows that, for every pair $E_{1}, E_{2} \in \mathcal{E}$, either $O\left(E_{1}\right)=O\left(E_{2}\right)$ or $O\left(E_{1}\right) \cap O\left(E_{2}\right)=\emptyset$. Therefore, the family $\mathcal{F}:=\{O(E): E \in \mathcal{E}\}$ is a partition of $\mathcal{Q}$.

Clearly, each $E \in \mathcal{E}$ is a convex cone. Let $i$ be its dimension and $D$ any $i$-dimensional disk $D \subset E$. Denote by $x$ the center of $D$ and let $y$ be any other point of $E$. Then, by (7.1), the point $z=y-\varepsilon(x-y)=(1+\varepsilon) y-\varepsilon x$ belongs as well to $E$ for any $\varepsilon>0$ sufficiently small. The convex envelope of $D \cup\{z\}$, which is contained in $E$, contains in turn an $i$-dimensional disk centered in $y$ : therefore $E$ is an open convex cone. Since $\left.O\right|_{E}$ is a linear injective map, $F=O(E)$ is an open convex cone of dimension $i$. Therefore, $\mathcal{F}$ satisfies (p1)-(p3).

Next notice that, having fixed $w \in E$, a point $z$ belongs to $\bar{E} \backslash E$ if and only if
(1) $w_{j}^{i} \geq w_{k}^{i}$ implies $z_{j}^{i} \geq z_{k}^{i}$ for every $i, j$ and $k$;
(2) there exists $r, s$ and $t$ such that $w_{s}^{r}>w_{t}^{r}$ and $z_{s}^{r}=z_{t}^{r}$.

Thus, if $d$ is the dimension of $E, \partial E:=\bar{E} \backslash E$ (cf. Remark 7.2) is the union of some elements of $\cup_{j<d} \mathcal{E}_{j}$, where with $\mathcal{E}_{j}$ we denote the $j$-dimensional elements of $\mathcal{E}$. Observe that, since $O$ is continuous, we must have $\bar{F} \supset O(\bar{E})$. On the other hand, if $x \in \bar{F}$ and $x_{k} \rightarrow x$ is a sequence contained in $F$, then there is a sequence $\left\{y_{k}\right\} \subset E$ with $O\left(y_{k}\right)=x_{k}$. By the definition of $O$ the sequence $\left\{y_{k}\right\}$ is bounded and hence, up to subsequence, we can assume that it converges to $y \in \bar{F}$ : thus $O(y)=x$ and $O(\bar{E})=\bar{F}$. On the other hand, for
equivalence classes $E_{1}, E_{2}$ of different dimension we necessarily have $O\left(E_{1}\right) \cap O\left(E_{2}\right)=\emptyset$. Thus $O(\partial E) \cap O(E)=\emptyset$, i.e. $\partial F=O(\partial E)$, which shows ( $\mathrm{p}_{4}$ ).

For what concerns $\left(\mathrm{p}_{5}\right)$ we show first that if $L(P)=z \in E \in \mathcal{E}$ is such that $z_{j}^{i} \neq z_{k}^{i}$ for all $i$ and for all $j \neq k$, then $O(E) \in \mathcal{F}_{n Q}$. Indeed, if $t<{ }^{1} / 4 \min _{i, j \neq k}\left|z_{j}^{i}-z_{k}^{i}\right|$, then $L(P+v) \in E$ for every $v \in B_{t}(0) \subset \mathbb{R}^{n Q}$, i.e. $E$ is an $(n Q)$-dimensional convex cone. Therefore it follows that for every $F \in \mathcal{F}_{i}$ with $i<n Q$ there exists $G \in \mathcal{F}_{n Q}$ such that $F \subset \bar{G}$. To show this claim it is enough to prove that, if $F=O(E)$ and $L(P)=z \in E$, then $z$ is the limit of points $w \in V$ such that $w_{j}^{i} \neq w_{k}^{i}$ for all $i, j, k$, which can be easily proved by a simple perturbation argument. Next, we argue inductively on $k$ : knowing that $F \in \mathcal{F}_{i}$ is contained in $\bar{G}$ for some $G \in \mathcal{F}_{k}$ with $k>i+1$, we show that there is $H \in \mathcal{F}_{k-1}$ such that $F \subset \bar{H}$. Observe indeed that $F \subset \partial G=\bar{G} \backslash G$ and that, for dimensional reasons, $\bar{G} \backslash G$ must be contained in the closure of those $H \in \mathcal{F}_{k-1}$ such that $H \subset \bar{G}$. Let $H \in \mathcal{F}_{k-1}$ be such that $F \cap \bar{H} \neq \emptyset$. Consider $E, K \in \mathcal{E}$ such that $F=O(E)$ and $H=O(K)$. Let $x \in E$ such that $O(x) \in F \cap \bar{H}$ and $z \in K$. We then must have that $x_{k}^{i} \geq x_{j}^{i}$ whenever $z_{k}^{i}>z_{j}^{i}$ and that $x_{k}^{i}=x_{j}^{i}$ whenever $z_{k}^{i}=z_{j}^{i}$. By the very definition of $\sim$, the same property holds even if we replace $x$ with another element $\xi \in E$. Therefore the open segment $] \xi, z[$ must be contained in $K$, which in turn implies that $\xi \in \bar{K}$. Thus we conclude $F \subset \bar{H}$.
7.2. Construction of $\boldsymbol{\rho}^{b}$. The main building block in the construction of $\boldsymbol{\rho}^{b}$ is given by the following lemma.

Lemma 7.3. For $\tau \in] 0, \frac{1}{4}\left[\right.$ and any $D \in \mathbb{N} \backslash\{0\}$ consider the map $\Phi_{\tau}: \mathbb{R}^{D} \rightarrow \mathbb{R}^{D}$ defined by:

$$
\Phi_{\tau}(x)= \begin{cases}0 & \text { if }|x| \leq \tau \\ \sqrt{\tau} \frac{|x|-\tau}{\sqrt{\tau}-\tau} \frac{x}{|x|} & \text { if } \tau \leq|x| \leq \sqrt{\tau} \\ x & \text { if }|x| \geq \sqrt{\tau}\end{cases}
$$

Then $\left|\Phi_{\tau}(x)-x\right| \leq \tau$ and $\operatorname{Lip}\left(\Phi_{\tau}\right) \leq 1+2 \sqrt{\tau}$.
Proof. The proofs of the two claims are straightforward computations. First $\Phi_{\tau}(x)=x$ if $|x| \geq \sqrt{\tau}$ and $\left|\Phi_{\tau}(x)-x\right|=|x| \leq \tau$ if $|x| \leq \tau$. For $\tau \leq|x| \leq \sqrt{\tau}$ we compute

$$
\left|\Phi_{\tau}(x)-x\right|=\left|\frac{\sqrt{\tau}(|x|-\tau)}{\sqrt{\tau}-\tau}-|x|\right|=\tau \frac{\sqrt{\tau}-|x|}{\sqrt{\tau}-\tau} \leq \tau
$$

Next we show that $\left|D \Phi_{\tau}(x) \cdot v\right| \leq(1+2 \sqrt{\tau})|v|$ at any point of differentiability. This inequality obviously imply the claimed Lipschitz constant estimate because $\Phi_{\tau}$ is Lipschitz and its domain of definition is a convex set. The inequality is, moreover, obvious when $|x|<\tau$ and $|x|>\sqrt{\tau}$. For $\tau<|x|<\sqrt{\tau}$, we can compute

$$
D \Phi_{\tau}(x)=\frac{1-\frac{\tau}{|x|}}{1-\sqrt{\tau}} \operatorname{Id}+\frac{\frac{\tau}{|x|}}{1-\sqrt{\tau}} \frac{x}{|x|} \otimes \frac{x}{|x|} .
$$

The matrix is symmetric with positive eigenvalues (because $|x|>\tau$ ) and the maximal eigenvalue is $(1-\sqrt{\tau})^{-1} \leq 1+2 \sqrt{\tau}$, thereby proving our claim.
7.2.1. Special coordinates, conical sections and separation. Let $S_{k}$ be the $k$-dimensional skeleton of $\mathcal{Q}$, i.e. the union of $F \in \mathcal{F}_{k}$ and denote by $\left(S_{k}\right)_{\sigma}$ its $\sigma$-neighborhood $\{x$ : $\left.\operatorname{dist}\left(x, S_{k}\right)<\sigma\right\}$. Incidentally, $\left(S_{k}\right)_{\sigma}$ contains $\left(S_{i}\right)_{\sigma}$ for every $i<k$.

Definition 7.4 (Coordinates and conical sections). Fix any face $F \in \mathcal{F}_{k}$ and introduce Cartesian coordinates $(y, z) \in \mathbb{R}^{k} \times \mathbb{R}^{N-k}$ in such a way that $F \subset \mathbb{R}^{k} \times\{0\}$. For a positive constant $\tilde{c}$ consider the cone $\mathscr{C}(F):=\left\{(y, z) \in \mathcal{Q}:(y, 0) \in F,|z| \leq \tilde{c} \operatorname{dist}\left((y, 0), S_{k-1}\right)\right\}$. For any $p=(y, 0) \in F$ we set $V_{p}:=\left(\{y\} \times \mathbb{R}^{N-k}\right) \cap \mathscr{C}(F)$.

Note that, if $\tilde{c}$ is sufficiently small, we will have the following property

$$
\mathscr{C}(F) \cap \mathscr{C}(G) \neq \emptyset \quad \Longrightarrow \quad \text { either } F \subset \bar{G} \text { or } G \subset \bar{F}
$$

For every constants $a, b>0, k=1 \ldots, n Q-1$ and $F \in \mathcal{F}_{k}$, we fix coordinates as in Definition 7.4 and denote by $F_{a, b}$ the sets

$$
F_{a, b}:=\left\{(y, z):|z| \leq a,(y, 0) \in F \backslash\left(S_{k-1}\right)_{b}\right\}
$$

For the faces $F \in \mathcal{F}_{n Q}$ of maximal dimension and for every $a>0, F_{\star, a}$ denotes the set $F_{\star, a}:=F \backslash\left(S_{n Q-1}\right)_{a}$. The following lemma is an obvious corollary of the linear simplicial and conical structures of $\mathcal{Q}$.

Lemma 7.5. There is a constant $\bar{c}>0$ (independent of $a, b$ below) with the following property. Assume $F$ and $G$ are two distinct $k$ dimensional faces.

- If $k=n Q, a>0, x \in F_{\star, a}$ and $x^{\prime} \in G_{\star, a}$, then $\left|x-x^{\prime}\right| \geq \bar{c} a$;
- If $k<n Q, b / a>\bar{c}^{-1}, x \in F_{a, b}$ and $x^{\prime} \in G_{a, b}$, then $\left|x-x^{\prime}\right| \geq \bar{c} b$.

Moreover, if $F \in \mathcal{F}_{k}, H \in \mathcal{F}_{i}$ with $i>k$ and $F \not \subset \partial H$ (cf. Remark 7.2), then $\left|x-x^{\prime}\right| \geq \bar{c} a$ for every $x \in H$ and $x^{\prime} \in F \backslash\left(S_{k-1}\right)_{a}$.
7.2.2. The domains $\operatorname{Dom}\left(f_{k}\right)$. Next we choose constants $c_{k}:=\delta^{8^{-n Q+k}}$. If $\delta$ is small enough, each family $\left\{F_{2 \sqrt{c_{k}}, c_{k-1}^{2}}\right\}_{F \in \mathcal{F}_{k}}$ with $k<n Q$ is made by pairwise disjoint sets, which are at least $\bar{c} c_{k-1}^{2}$ far apart, where $\bar{c}$ is the constant of Lemma 7.5, and it holds $F_{2 \sqrt{c_{k}}, c_{k-1}^{2}} \subset$ $\mathscr{C}(F) \subset \mathcal{Q}$. We are ready to define the map $\rho^{b}:=\left.\boldsymbol{\rho}^{\star}\right|_{\mathcal{Q}}$ inductively "from the top to the bottom". More precisely we will define a family of maps $\left\{f_{k}\right\}_{k \in\{0, \ldots, n Q\}}$ on domains $\operatorname{Dom}\left(f_{k}\right) \subset \mathcal{Q}$ starting from $f_{n Q}$ and ending with $f_{0}=\boldsymbol{\rho}^{\text {b }}$. We first explicitly define $\operatorname{Dom}\left(f_{k}\right):=\mathcal{Q} \backslash\left(S_{k-1}\right)_{c_{k-1}}$ for $k>0$ and $\operatorname{Dom}\left(f_{0}\right)=\mathcal{Q}$, and in order to simplify our notation we then agree that $c_{-1}=\delta^{8^{-n Q-1}}$ and $S_{-1}=\left(S_{-1}\right)_{c_{-1}}=\emptyset$. Note that $\operatorname{Dom}\left(f_{k+1}\right) \not \subset$ $\operatorname{Dom}\left(f_{k}\right)$. It is obvious that

$$
\begin{equation*}
\operatorname{Dom}\left(f_{k}\right)=\left(\operatorname{Dom}\left(f_{k+1}\right) \cup \bigcup_{F \in \mathcal{F}_{k}} F_{2 \sqrt{c_{k}}, c_{k-1}^{2}}\right) \backslash\left(S_{k-1}\right)_{c_{k-1}} . \tag{7.2}
\end{equation*}
$$

Indeed, if $x \in \operatorname{Dom}\left(f_{k}\right) \backslash \operatorname{Dom}\left(f_{k+1}\right)$ we then must have $\operatorname{dist}\left(x, S_{k}\right)<c_{k}$ and $\operatorname{dist}\left(x, S_{k-1}\right) \geq$ $c_{k-1}$. Let $q \in S_{k}$ be such that $|x-q|<c_{k}$. Since $\operatorname{dist}\left(x, S_{k-1}\right) \geq c_{k-1}>c_{k}$, the point $q$ must necessarily belong to a $k$-dimensional face $F$. Fix coordinates as in Definition 7.4. If $x=(y, z)$, we then obviously have $|z|<c_{k} \leq 2 \sqrt{c_{k}}$. On the other hand $\operatorname{dist}\left((y, 0), S_{k-1}\right) \geq$ $\operatorname{dist}\left(x, S_{k-1}\right)-|z| \geq c_{k-1}-c_{k}>c_{k-1}^{2}$. This shows that $x \in F_{2 \sqrt{c_{k}}, c_{k-1}^{2}}$.
7.2.3. The maps $f_{k}$. On $\operatorname{Dom}\left(f_{n Q}\right)$ we define $f_{n Q}=\operatorname{Id}$ and specify next the procedure to define $f_{k}$ knowing $f_{k+1}$. Along the procedure we claim inductively the following.
Assumption 7.6 (Inductive step). The map $f_{k+1}$ has the following three properties.
$\left(\mathrm{a}_{k+1}\right) \operatorname{Lip}\left(f_{k+1}\right) \leq 1+C c_{k+1}^{1 / 2}$ and $\left|f_{k+1}(x)-x\right| \leq C c_{k+1}$.
$\left(\mathrm{b}_{k+1}\right)$ Consider $i \leq k+1$, an $i$-dimensional face $F$, the cone $\mathscr{C}(F)$ in Definition 7.4 and the corresponding coordinates. Then, $f_{k+1}$ factorizes on $\operatorname{Dom}\left(f_{k+1}\right) \cap \mathscr{C}(F)$ as

$$
\begin{equation*}
f_{k+1}(y, z)=\left(y, h_{k+1}^{F}(y, z)\right) \in \mathbb{R}^{i} \times \mathbb{R}^{N-i} \tag{7.3}
\end{equation*}
$$

(c $\left.c_{k+1}\right)$ For every $G \in \mathcal{F}_{i}$ with $i \geq k+1$, $f_{k+1}$ maps $\operatorname{Dom}\left(f_{k+1}\right) \cap\{x: \operatorname{dist}(x, G)<\delta\}$ into $\bar{G}$. Moreover the restriction of $f_{k+1}$ to $G_{c_{i}, c_{k}}$ is the orthogonal projection onto $G$.

The constants involved depend on $k$ but not on the parameter $\delta$ and since the process is iterated finitely many times, we will not keep track of such dependence. Note that $f_{n Q}$ satisfies $\left(\mathrm{a}_{n Q}\right)$, $\left(\mathrm{b}_{n Q}\right)$ and $\left(\mathrm{c}_{n Q}\right)$ trivially, because it is the identity map. Given $f_{k+1}$ we next show how to construct $f_{k}$. For every $p \in G \in \mathcal{F}_{k}$ with $p \notin\left(S_{k-1}\right)_{c_{k-1}^{2}}$, set coordinates as in Definition 7.4 and consider the cone $W_{p}:=\left\{(y, z) \in V_{p}:|z| \leq 2 \sqrt{c_{k}}\right\}$. Let now $\Phi_{\tau}$ be the map of Lemma 7.3 with $\tau=2 c_{k}$. The function $f_{k}$ is defined in $W_{p}$ by

$$
f_{k}(x)=f_{k}(y, z):=\left(y, h_{k}^{F}(y, z)\right):= \begin{cases}(y, 0) & \text { for }|z| \leq \tau / 2=c_{k}  \tag{7.4}\\ \left(y, \Phi_{\tau}\left(h_{k+1}^{F}(y, z)\right)\right) & \text { otherwise }\end{cases}
$$

If $q \in \operatorname{Dom}\left(f_{k}\right)$ does not belong to any $W_{p}$ as above, then we set $f_{k+1}(q)=f_{k}(q)$.
Observe that the definition above gives values to $f_{k}$ on a set which is larger than $\operatorname{Dom}\left(f_{k}\right)$ : this will be useful to carry on some of the estimates, but we insist that Assumption 7.6 will only be checked on $\operatorname{Dom}\left(f_{k}\right)$.
7.2.4. Well-definition and continuity. Consider a point $q \in \operatorname{Dom}\left(f_{k}\right)$. If $q$ is not contained in $F_{2 \sqrt{c_{k}}, c_{k-1}^{2}}$ for some $k$-dimensional face, then by (7.2) it is contained in the domain of $f_{k+1}$ and thus $f_{k}(q)$ is defined. If $q$ is contained in $F_{2 \sqrt{c_{k}}, c_{k-1}^{2}}$ for some $k$-dimensional face, then $q$ belongs to some $W_{p}$ as above. Let $q=(y, z)$. If $|z| \leq c_{k}$, then $f_{k}(q)$ is defined; otherwise, since $\operatorname{dist}\left(q, S_{k}\right) \geq c_{k}$, we infer that $q \in \operatorname{Dom}\left(f_{k+1}\right)$ and $f_{k}(q)$ is also defined.

As for the continuity, fix $(y, z) \in W_{p} \cap \operatorname{Dom}\left(f_{k}\right)$ with $p=(y, 0) \in F \in \mathcal{F}_{k}$. If $|z|=c_{k}$, then by $\left(\mathrm{a}_{k+1}\right)$ we have $\left|h_{k+1}^{F}(y, z)\right| \leq|z|+C c_{k+1} \leq \tau / 2+C \tau^{8}$. For $\delta$ sufficiently small this obviously implies $\left|h_{k+1}^{F}(y, z)\right| \leq \tau$ and thus, by the definition of $\Phi_{\tau}, \Phi_{\tau}\left(h_{k+1}^{F}(y, z)\right)=0$. On the other hand, if $|z|=2 \sqrt{c_{k}}$, then $\left|h_{k+1}^{F}(y, z)\right| \geq|z|-C c_{k+1}=2 \sqrt{c_{k}}-C c_{k}^{8} \geq \sqrt{2 c_{k}}$ and thus $\Phi_{\tau}\left(h_{k+1}^{F}(y, z)\right)=h_{k+1}^{F}(y, z)$. Therefore under this assumption we have $f_{k+1}(q)=f_{k}(q)$.

We next check that $f_{k}$ maps $\operatorname{Dom}\left(f_{k}\right)$ into $\mathcal{Q}$. This is true by induction where $f_{k}$ coincides with $f_{k+1}$. Fix therefore a point $q$ in some $W_{p} \cap \operatorname{Dom}\left(f_{k}\right)$ with $p \in F \in \mathcal{F}_{k}$ and let $G$ be the $i$-dimensional face containing $q$ with $i>k$. Then, $f_{k+1}(q)$ belongs to a face $\bar{G}$, by Assumption 7.6. By the estimate in $\left(\mathrm{a}_{k+1}\right)$ and the assumption $\left(\mathrm{b}_{k+1}\right)$, the face $G$ must intersect $\mathscr{C}(F)$ and thus $F \subset \bar{G}$. Observe that, by the properties of $\Phi_{\tau}$ and by the inductive assumption $\left(\mathrm{b}_{k+1}\right), f_{k}(q)$ is mapped in the segment joining $f_{k+1}(q)$ and $q$ and thus must belong to $\bar{G}$.
7.2.5. The inductive conclusions $\left(c_{k}\right)$ and $\left(b_{k}\right)$. The first claim of $\left(c_{k}\right)$ is simple to prove: as noticed, if a point $q \in \operatorname{Dom}\left(f_{k}\right)$ belongs also to $\operatorname{Dom}\left(f_{k+1}\right)$, then $f_{k}$ maps it into the closure of the face containing $q$. If the point is not contained in $\operatorname{Dom}\left(f_{k+1}\right)$, then it must be contained in the $c_{k}$-neighborhood of some $k$-dimensional face $F$ and hence it is mapped into $F$ : when this happens $F$ is a portion of the boundary of the face containing $q$. Next, fix a face $G \in \mathcal{F}_{i}$. If $i=k$, by the very definition of $f_{k}$, we have that the restriction of $f_{k}$ to $\operatorname{Dom}\left(f_{k}\right) \cap G_{c_{k}, c_{k-1}}$ is the orthogonal projection onto $G$. If $i>k$, we actually have that $f_{k}=f_{k+1}$ on $\operatorname{Dom}\left(f_{k}\right) \backslash\left(S_{k}\right)_{2 \sqrt{c_{k}}} \supset \mathcal{Q} \backslash\left(S_{i-1}\right)_{c_{k-1}}$.

Fix now an $i$-dimensional face $L$ with $i \leq k$, consider coordinates $\mathbb{R}^{i} \times \mathbb{R}^{n-i}$ as in Definition 7.4 and the corresponding $\mathscr{C}(L)$. If $q=(y, 0) \in L$, the condition $\left(\mathrm{b}_{k}\right)$ is equivalent to saying that $V_{q} \cap \operatorname{Dom}\left(f_{k}\right)$ gets mapped into $\{(y, 0)\} \times \mathbb{R}^{N-i}$. Fix a point $\tilde{q} \in V_{q}$. If $f_{k+1}(\tilde{q})=f_{k}(\tilde{q})$ there is nothing to prove. Otherwise it turns out that there is a $k$-dimensional face $F$ such that $\tilde{q} \in \mathscr{C}(F)$. But then we necessarily have $L \subset \bar{F}$. So, set coordinates $\mathbb{R}^{i} \times \mathbb{R}^{k-i} \times \mathbb{R}^{n-k}$ so that at the same time $L \subset \mathbb{R}^{i} \times\{0\} \times\{0\}$ and $F \subset \mathbb{R}^{i} \times \mathbb{R}^{k-i} \times\{0\}$. Thus, $(y, 0,0)$ is the coordinate of $q$ and $(y, z, w)$ that of $\tilde{q}$. According to our definition of $f_{k}, f_{k}(\tilde{q})=\left(y, z, w^{\prime}\right)$ for some $w^{\prime}$, which indeed implies the desired claim.
7.2.6. $C^{0}$ estimate. Observe that, for every $x$ where $f_{k}$ coincides with $f_{k+1}$, we have $\mid f_{k}(x)-$ $x \mid \leq C c_{k+1} \leq C c_{k}^{8}$. Instead, for any point $x$ where $f_{k}$ is newly defined, we distinguish the following two cases: either $x=(y, z)$ with $|z| \leq c_{k}$, in which case $\left|f_{k}(x)-x\right| \leq c_{k}$; or $x=(y, z)$ with $|z|>c_{k}$, and then by the estimates of Lemma 7.3 and the triangle inequality we have

$$
\left|f_{k}(x)-x\right| \leq\left|f_{k+1}(x)-f_{k}(x)\right|+\left|f_{k+1}(x)-x\right| \leq C c_{k+1}+\tau \leq C c_{k+1}+2 c_{k}
$$

7.2.7. Lipschitz estimate. We fix $x, x^{\prime} \in \operatorname{Dom}\left(f_{k}\right)$ and, apart from the trivial one $f_{k}(x)=$ $f_{k+1}(x)$ and $f_{k}\left(x^{\prime}\right)=f_{k+1}\left(x^{\prime}\right)$, we distinguish three cases.
Case 1: $x, x^{\prime} \in G_{2 \sqrt{c_{k}}, c_{k-1}^{2}}$ for some $k$-dimensional face $G$. Choosing coordinates as in Definition 7.4, we set $x=(y, z)$ and $x^{\prime}=\left(y^{\prime}, z^{\prime}\right)$. If both $|z|,\left|z^{\prime}\right| \leq \frac{\tau}{2}$, then $\left|f_{k}(x)-f_{k}\left(x^{\prime}\right)\right|=$ $\left|y-y^{\prime}\right| \leq\left|x-x^{\prime}\right|$. If $|z| \geq \frac{\tau}{2}$ and $\left|z^{\prime}\right| \geq \frac{\tau}{2}$, then

$$
\begin{aligned}
\left|f_{k}(x)-f_{k}\left(x^{\prime}\right)\right|^{2} & \leq\left|y-y^{\prime}\right|^{2}+\left(1+2 \tau^{1 / 2}\right)^{2}\left|h_{k+1}^{F}(y, z)-h_{k+1}^{F}\left(y^{\prime}, z^{\prime}\right)\right|^{2} \\
& \leq\left(1+2 \tau^{1 / 2}\right)^{2}\left(\left|y-y^{\prime}\right|^{2}+\left|h_{k+1}^{F}(y, z)-h_{k+1}^{F}\left(y^{\prime}, z^{\prime}\right)\right|^{2}\right) \\
& =\left(1+\sqrt{2 c_{k}}\right)^{2}\left|f_{k+1}(x)-f_{k+1}\left(x^{\prime}\right)\right|^{2} \leq\left(1+\sqrt{2 c_{k}}\right)^{2}\left(1+C \sqrt{c_{k+1}}\right)^{2}\left|x-x^{\prime}\right|^{2} .
\end{aligned}
$$

If $|z| \leq \frac{\tau}{2}$ and $\left|z^{\prime}\right|>\frac{\tau}{2}$, let $\tilde{z}$ be the point with $|\tilde{z}|=\frac{\tau}{2}$ on the segment joining $z$ and $z^{\prime}$, and $\tilde{x}=(y, \tilde{z})$. Observe that $f_{k}(\tilde{x})=f_{k}(x)=(y, 0)$ and that $\left|\tilde{x}-x^{\prime}\right|^{2}=\left|y-y^{\prime}\right|^{2}+$ $\left|z^{\prime}-\tilde{z}\right|^{2} \leq\left|y-y^{\prime}\right|^{2}+\left|z-z^{\prime}\right|^{2} \leq\left|x-x^{\prime}\right|^{2}$. On the other hand we have just shown $\left|f_{k}\left(x^{\prime}\right)-f_{k}(\tilde{x})\right| \leq\left(1+C c_{k}^{1 / 2}\right)\left|x^{\prime}-\tilde{x}\right|$.
Case 2: $x \in F_{2 \sqrt{c_{k}}, c_{k-1}^{2}}, x^{\prime} \in G_{2 \sqrt{c_{k}}, c_{k-1}^{2}}$ for distinct $F, G \in \mathcal{F}_{k}$. By Lemma 7.5, $\left|x-x^{\prime}\right| \geq$ $\bar{c} c_{k-1}^{2} \geq \bar{c} c_{k}^{1 / 4}$. On the other hand, we also have, by the $C^{0}$ estimate,

$$
\left|f_{k}(x)-f_{k}\left(x^{\prime}\right)\right| \leq\left|x-x^{\prime}\right|+C c_{k} \leq\left(1+C c_{k}^{3 / 4}\right)\left|x-x^{\prime}\right|
$$

Case 3: $x \in G_{2 \sqrt{c_{k}}, c_{k-1}^{2}}$ for some $k$-dimensional face $G$ and $f_{k}\left(x^{\prime}\right)=f_{k+1}\left(x^{\prime}\right)$. Without loss of generality we assume

- $G \in \mathcal{F}_{k}$;
- $x^{\prime} \notin G_{2 \sqrt{c_{k}},}, c_{k-1}^{2}$;
- $x^{\prime} \in H$ for some face $H$ (of dimension $i>k$ ).

We have two possibilities.
Case 3a: $G \not \subset \bar{H}$. Consider the closed set $\tilde{G}:=G \backslash\left(S_{k-1}\right)_{c_{k-1}^{2}}$. By Lemma $7.5 \operatorname{dist}\left(x^{\prime}, \tilde{G}\right) \geq$ $\bar{c} c_{k-1}^{2}$ and thus $\left|x-x^{\prime}\right| \geq \bar{c} c_{k-1}^{2}-2 \sqrt{c_{k}} \geq \frac{\bar{c}}{2} c_{k-1}^{2}$. We can therefore argue as in Case 2.
Case 3b: $G \subset \bar{H}$. We then have two possibilities. The first is that $x \in \operatorname{Dom}\left(f_{k+1}\right)$. Since $f_{k}\left(x^{\prime}\right)=f_{k+1}\left(x^{\prime}\right)$, we have $\left|f_{k}\left(x^{\prime}\right)-x^{\prime}\right| \leq C c_{k+1}=C c_{k}^{8}$. We use the coordinates of Definition 7.4 and $\left(\mathrm{a}_{k+1}\right)$ to conclude $f_{k}\left(x^{\prime}\right)=f_{k+1}\left(y^{\prime}, z^{\prime}\right)=\left(y^{\prime \prime}, z^{\prime \prime}\right)$ with $\left|z^{\prime \prime}\right| \geq\left|z^{\prime}\right|-C c_{k}^{8} \geq$ $2 \sqrt{c_{k}}-C c_{k}^{8} \geq \sqrt{2 c_{k}}$. We can therefore write $f_{k}\left(x^{\prime}\right)=\left(y^{\prime \prime}, \Phi_{\tau}\left(z^{\prime \prime}\right)\right)$ (because $\left.\Phi_{\tau}\left(z^{\prime \prime}\right)=z^{\prime \prime}\right)$ and, hence, recalling $f_{k}(x)=\left(y, \Phi_{\tau}\left(h_{k+1}^{F}(y, z)\right)\right)$ and $f_{k+1}(x)=\left(y, h_{k+1}^{F}(y, z)\right)$,

$$
\begin{aligned}
& \left|f_{k}\left(x^{\prime}\right)-f_{k}(x)\right|^{2} \leq\left|y-y^{\prime \prime}\right|^{2}+(1+2 \sqrt{\tau})^{2}\left|h_{k+1}^{F}(y, z)-z^{\prime \prime}\right|^{2} \\
\leq & (1+2 \sqrt{\tau})^{2}\left|f_{k+1}(x)-f_{k+1}\left(x^{\prime}\right)\right|^{2}
\end{aligned}
$$

We therefore conclude $\left|f_{k}\left(x^{\prime}\right)-f_{k}(x)\right| \leq\left(1+C \tau^{1 / 2}\right)\left|x^{\prime}-x\right| \leq\left(1+C c_{k}^{1 / 2}\right)\left|x^{\prime}-x\right|$.
The second possibility is that $x$ is not in the domain of definition of $f_{k+1}$. In that case $x$ is at distance $c_{k}$ from $G$ and thus $\left|x-x^{\prime}\right| \geq \sqrt{c_{k}}$. We then conclude that $\left|f_{k}(x)-f_{k}\left(x^{\prime}\right)\right| \leq$ $\left|x-x^{\prime}\right|+C c_{k} \leq\left(1+C \sqrt{c_{k}}\right)\left|x-x^{\prime}\right|$.
7.2.8. Summary. After $n Q$ steps, we get a function $f_{0}=\boldsymbol{\rho}^{b}: \mathcal{Q} \rightarrow \mathcal{Q}$ which satisfies

$$
\begin{align*}
& \operatorname{Lip}\left(\boldsymbol{\rho}^{b}\right) \leq 1+C \delta^{8^{-n Q-1}} \quad \text { and } \quad\left|\boldsymbol{\rho}^{b}(x)-x\right| \leq C \delta^{8^{-n Q}}  \tag{7.5}\\
& \boldsymbol{\rho}^{b}(\{x: \operatorname{dist}(x, F) \leq \delta\}) \subset \bar{F} \text { for every } F \in \mathcal{F}_{k}  \tag{7.6}\\
& \boldsymbol{\rho}^{\boldsymbol{b}}: F_{\delta, c_{0}^{1 / 8}} \rightarrow F \text { is the orthogonal projection on } F \text { for every } F \in \mathcal{F}_{k} . \tag{7.7}
\end{align*}
$$

7.3. The extension $\boldsymbol{\rho}^{\sharp}$ of $\boldsymbol{\rho}^{b}$ to $\mathcal{Q}_{\delta^{n Q+1}}$. Next we extend the map $\boldsymbol{\rho}^{\boldsymbol{b}}: \mathcal{Q} \rightarrow \mathcal{Q}$ to the $\delta^{n Q+1}$-neighborhood of $\mathcal{Q}$, keeping the estimate (7.5). We first observe that, since the number of all the faces is finite, when $\delta$ is small enough, there exists a constant $C=C(N)$ with the following property. Consider two distinct faces $F$ and $H$ in $\mathcal{F}_{i}$. If $x, y$ are two points contained, respectively, in $F_{\delta^{i+1}} \backslash \cup_{j<i} \cup_{G \in \mathcal{F}_{j}} G_{\delta^{j+1}}$ and $H_{\delta^{i+1}} \backslash \cup_{j<i} \cup_{G \in \mathcal{F}_{j}} G_{\delta^{j+1}}$, then

$$
\begin{equation*}
\operatorname{dist}(x, y) \geq C \delta^{i} \tag{7.8}
\end{equation*}
$$

Similarly if $F \in \mathcal{F}_{l}$ and $H \in \mathcal{F}_{i}$ with $l<i$ and $F \not \subset \bar{H}$, then for every $x \in F_{\delta^{l+1}}$ and $y \in H_{\delta^{i+1}} \backslash \cup_{j<i} \cup_{G \in \mathcal{F}_{j}} G_{\delta^{j+1}}$ it holds

$$
\begin{equation*}
\operatorname{dist}(x, y) \geq C \delta^{i} \tag{7.9}
\end{equation*}
$$

The extension $\boldsymbol{\rho}^{\sharp}$ is defined inductively, but this time "from the bottom to the top". The first extension $g_{0}$ is identically 0 on $B_{\delta}(0)$ (note that this is feasible because $\boldsymbol{\rho}^{b} \equiv 0$ in
$\left.B_{\delta}(0) \cap \mathcal{Q}\right)$. Now we come to the inductive step. Suppose we have an extension $g_{\ell}$ of $\boldsymbol{\rho}^{b}$, defined on the union of the $\delta^{\ell+1}$-neighborhoods of the $\ell$-skeletons $S_{\ell}$, for $\ell \in\{0, \ldots, k\}$, i.e.

$$
\mathbf{L}_{k}:=\mathcal{Q} \cup B_{\delta}(0) \cup \bigcup_{\ell=1}^{k} \bigcup_{F \in \mathcal{F}_{\ell}} F_{\delta^{\ell+1}}
$$

Assume inductively that $\operatorname{Lip}\left(g_{k}\right) \leq 1+C \delta^{8^{-n Q-1}}$ and assume that $g_{k}$ maps any $\delta^{j+1}$ _ neighborhood of any $j$-dimensional face into its closure, when $j \leq k$. Then, we define the extension of $g_{k}$ to $\mathbf{L}_{k+1}$ in the following way. For every face $F \in \mathcal{F}_{k+1}$, we set

$$
g_{k+1}:=\left\{\begin{array}{lll}
\rho^{b} & \text { on } & \mathcal{Q}  \tag{7.10}\\
g_{k} & \text { on } & \left(S_{k}\right)_{\delta^{k+1}} \cap F_{\delta^{k+2}}, \\
\mathbf{p}_{F} & \text { on } & \left\{x \in \mathbb{R}^{N}: \mathbf{p}_{F}(x) \in F_{\delta, 1}\right\} \cap F_{\delta^{k+2}}
\end{array}\right.
$$

where $\mathbf{p}_{F}$ stands for the orthogonal projection on $F$ (recall that by (7.7) $\rho^{b}=\mathbf{p}_{F}$ on $F \cap$ $F_{\delta^{k+2}, 1}$ ). Consider now a face $F$ as above and $U(F)$ the union of all the $\delta^{j+1}$-neighborhoods of the $j$-dimensional faces which belong to $\bar{F}$. As defined above, $g_{k+1}$ maps a portion of $U(F)$ into $\bar{F}$. We can use Lemma 7.1 to extend $g_{k+1}$ to $U(F)$ keeping the same Lipschitz constant, which we now compute. This constant is obviously smaller than $1+C \delta^{8^{-n Q-1}}$ on the domain $\left(\left(S_{k}\right)_{\delta^{k+1}} \cap F_{\delta^{k+2}}\right) \cup F$ by inductive hypothesis. The same constant is 1 on $\{x \in$ $\left.\mathbb{R}^{N}: \mathbf{p}_{F}(x) \in F_{\delta, 1}\right\} \cap F_{\delta^{k+2}}$. Consider now a point $x \in\left\{x \in \mathbb{R}^{N}: \mathbf{p}_{F}(x) \in F_{\delta, 1}\right\} \cap F_{\delta^{k+2}}$ and a point $y \in F \cup\left(\left(S_{k}\right)_{\delta^{k+1}} \cap F_{\delta^{k+2}}\right)$. If $y \notin\left(S_{k}\right)_{c_{0}^{1 / 8}}$, then necessarily $y \in F$ and we then have

$$
\left|g_{k+1}(x)-g_{k+1}(y)\right|=\left|\mathbf{p}_{F}(x)-y\right|=\left|\mathbf{p}_{F}(x)-\mathbf{p}_{F}(y)\right| \leq|x-y| .
$$

Otherwise we have $|x-y| \geq 1-c_{0}^{1 / 8}=1-\delta^{8^{-n Q-1}}$ and we can write

$$
\begin{aligned}
\left|g_{k+1}(x)-g_{k+1}(y)\right| & \leq\left|g_{k+1}(x)-y\right|+C c_{0} \leq|x-y|+\delta^{k+2}+C c_{0} \\
& \leq\left(1+\frac{\delta^{k+2}+C c_{0}}{1-C c_{0}^{1 / 8}}\right)|x-y| \leq\left(1+C \delta^{8^{-n Q-1}}\right)|x-y|
\end{aligned}
$$

Note that, if $x \in U\left(F_{1}\right) \cap U\left(F_{2}\right)$ for two distinct $F_{1}, F_{2} \in \mathcal{F}_{k+1}$, then $x \in \mathbf{L}_{k}$. Thus, the map $g_{k+1}$ is continuous. We next bound the global Lipschitz constant of $g_{k+1}$. Indeed consider points $x \in U\left(F_{1}\right) \backslash U\left(F_{2}\right)$ and $y \in U\left(F_{2}\right) \backslash U\left(F_{1}\right)$ for two distinct $F_{i} \in \mathcal{F}_{k+1}$. Since by (7.8) and (7.9) $|x-y| \geq C \delta^{k+1}$, we easily see that

$$
\begin{aligned}
\left|g_{k+1}(x)-g_{k+1}(y)\right| \leq & \left|g_{k+1}(x)-g_{k+1}\left(\mathbf{p}_{F_{1}}(x)\right)\right|+\left|g_{k+1}\left(\mathbf{p}_{F_{1}}(x)\right)-g_{k+1}\left(\mathbf{p}_{F_{2}}(y)\right)\right| \\
& +\left|g_{k+1}\left(\mathbf{p}_{F_{2}}(y)\right)-g_{k+1}(y)\right| \\
\leq & 2\left(1+C \delta^{8^{-n Q-1}}\right) \delta^{k+2}+\left|\boldsymbol{\rho}^{b}\left(\mathbf{p}_{F_{1}}(x)\right)-\boldsymbol{\rho}^{b}\left(\mathbf{p}_{F_{2}}(y)\right)\right| \\
\leq & 2\left(1+C \delta^{8^{-n Q-1}}\right) \delta^{k+2}+\left(1+C \delta^{8^{-n Q-1}}\right)\left|\mathbf{p}_{F_{1}}(x)-\mathbf{p}_{F_{2}}(y)\right| \\
\leq & 2\left(1+C \delta^{8^{-n Q-1}}\right) \delta^{k+2}+\left(1+C \delta^{8^{-n Q-1}}\right)\left(|x-y|+2 \delta^{k+2}\right) \\
\leq & \left(1+C \delta^{8^{-n Q-1}}\right)|x-y| .
\end{aligned}
$$

Next, consider the case $x \in \mathcal{Q} \backslash U(F), y \in U(F)$. If $|x-y| \geq \delta^{k+1}$, we can then argue as above and (considering that $g_{k+1}(x)=\boldsymbol{\rho}^{\mathrm{b}}(x)$ ) we bound

$$
\begin{aligned}
& \left|g_{k+1}(x)-g_{k+1}(y)\right| \leq\left(1+C \delta^{8^{-n Q-1}}\right) \delta^{k+2}+\left|\boldsymbol{\rho}^{b}(x)-\boldsymbol{\rho}^{b}\left(\mathbf{p}_{F}(y)\right)\right| \\
& \leq\left(1+C \delta^{8^{-n Q-1}}\right)\left(\delta^{k+2}+\left|x-\mathbf{p}_{F}(y)\right|\right) \leq\left(1+C \delta^{8^{-n Q-1}}\right)\left(\delta^{k+2}+|x-y|+\delta^{k+2}\right) \\
& \leq\left(1+C \delta^{8-n Q-1}\right)|x-y| .
\end{aligned}
$$

We therefore assume $|x-y| \leq \delta^{k+1}$. Observe also that, if $y \notin\left\{x \in \mathbb{R}^{N}: \mathbf{p}_{F}(x) \in\right.$ $\left.F_{\delta, 1}\right\} \cap F_{\delta^{k+2}}$, then $g_{k+1}(y)=g_{k}(y)$ and since $g_{k+1}(x)=\rho^{b}(x)=g_{k}(x)$, we know the Lipschitz bound by inductive assumption. We therefore conclude that $x \in F_{\delta^{k+2}+\delta^{k+1}, 1-\delta^{k+1}}$. Assuming $\delta_{0}$ small enough, $\delta^{k+2}+\delta^{k+1} \leq \delta$ and $1-\delta^{k+1} \geq \delta^{8^{-n Q-1}}=c_{0}^{1 / 8}$, therefore $x \in F_{\delta, c_{0}^{1 / 8}}$. By (7.7) we then have $\left|g_{k+1}(x)-g_{k+1}(y)\right|=\left|\mathbf{p}_{F}(x)-\mathbf{p}_{F}(y)\right| \leq|x-y|$.

Since $\mathcal{Q}$ and the union of the $U\left(F_{i}\right)$ is the domain of definition of $g_{k+1}$, this shows $\operatorname{Lip}\left(\underline{g_{k+1}}\right) \leq 1+C \delta^{8^{-n Q-1}}$. Note that by construction we also have that $U(F)$ is mapped into $\bar{F}$, which is the other inductive hypothesis.

After making the step above $n Q$ times we arrive to a map $g_{n Q}$ which extends $\boldsymbol{\rho}^{b}$ and is defined in a $\delta^{n Q+1}$-neighborhood of $\mathcal{Q}$. This is the map $\rho^{\sharp}$.

## 8. Persistence of $Q$-points: Proof of Theorem 1.7

Proof of Theorem 1.7. As usual, by scaling and translating we assume $x=0$ and $r=1$. According to [14, Theorem 3.9], there are constants $\bar{C}(m, n, Q), \kappa(m, n, Q)>0$ such that

$$
\begin{equation*}
\sup _{x \neq y \in B_{1 / 2}} \frac{\mathcal{G}(w(x), w(y))}{|y-x|^{\kappa}} \leq \bar{C}(\operatorname{Dir}(w))^{\frac{1}{2}} \quad \text { for any Dir-minimizer } w: B_{1} \rightarrow \mathcal{A}_{Q}\left(\mathbb{R}^{n}\right) \tag{8.1}
\end{equation*}
$$

The final choice of $\bar{s}$ will be specified at the very end, but for the moment we impose $\bar{s}<\frac{1}{4}$. Fix now $s<\bar{s}$ and $C^{\star}$ as in the statement and assume by contradiction that, no matter how small we choose $\hat{\varepsilon}>0$, there are a current $T$ and a submanifold $\Sigma$ as in Theorem 1.4 and a point $(p, q) \in \mathbf{C}_{1 / 2}$ satisfying:
(a) $E:=\mathbf{E}\left(T, \mathbf{C}_{4}\right)<\hat{\varepsilon}$ and $\mathbf{A}^{2} \leq C^{\star} E$;
(b) $\Theta(T,(p, q))=Q$;
(c) the $E^{\gamma_{1}}$-approximation $f$ (which is the map of Theorem 1.4) violates (1.9), that is

$$
\begin{equation*}
\int_{B_{s}(p)} \mathcal{G}(f, Q \llbracket \boldsymbol{\eta} \circ f \rrbracket)^{2}>\hat{\delta} s^{m} E . \tag{8.2}
\end{equation*}
$$

Set $\bar{\delta}=\frac{1}{4}$ and fix $\bar{\eta}>0$ (whose choice will be specified later). By (a), for a suitably small $\hat{\varepsilon}$ we can apply Theorem 1.6 in the coordinates of Remark 1.5: we let $u$ be the corresponding Dir-minimizer and $w=(u, \Psi(x, u))$. If $\bar{\eta}$ and $\hat{\varepsilon}$ are suitably small, we have

$$
\int_{B_{s}(p)} \mathcal{G}(w, Q \llbracket \boldsymbol{\eta} \circ w \rrbracket)^{2} \geq \frac{3 \hat{\delta}}{4} s^{m} E
$$

and $\sup \{\operatorname{Dir}(f), \operatorname{Dir}(w)\} \leq C E$ (here we use Remark 4.5). Thus there is $\bar{p} \in B_{s}(p)$ with $\mathcal{G}(w(\bar{p}), Q \llbracket \boldsymbol{\eta} \circ w(\bar{p}) \rrbracket)^{2} \geq \frac{3 \hat{\delta}}{4 \omega_{m}} E$ and, by (8.1), we conclude

$$
\begin{equation*}
g(x):=\mathcal{G}(w(x), Q \llbracket \boldsymbol{\eta} \circ w(x) \rrbracket) \geq\left(\frac{3 \hat{\delta}}{4 \omega_{m}} E\right)^{1 / 2}-2(C E)^{1 / 2} \bar{C} \bar{s}^{\kappa} \geq\left(\frac{\hat{\delta}}{2} E\right)^{1 / 2} \tag{8.3}
\end{equation*}
$$

where we assume that $\bar{s}$ is chosen small enough in order to satisfy the last inequality. Setting $h(x):=\mathcal{G}(f(x), Q \llbracket \boldsymbol{\eta} \circ f(x) \rrbracket)$, we recall that we have

$$
\int_{B_{s}(p)}|h-g|^{2} \leq C \bar{\eta} E
$$

Consider therefore the set $A:=\left\{h>\left(\frac{\hat{\delta}}{4} E\right)^{1 / 2}\right\}$. If $\bar{\eta}$ is sufficiently small, we can assume that $\left|B_{s}(p) \backslash A\right|<\frac{1}{8}\left|B_{s}\right|$. Further, define $\bar{A}:=A \cap K$, where $K$ is the set of Theorem 1.4. Assuming $\hat{\varepsilon}$ is sufficiently small we ensure $\left|B_{s}(p) \backslash \bar{A}\right|<\frac{1}{4}\left|B_{s}\right|$. Let $N$ be the smallest integer such that $N \frac{\hat{\delta} E}{64 Q s} \geq \frac{s}{2}$. Set $\sigma_{i}:=s-i \frac{\hat{\delta} E}{64 Q s}$ for $i \in\{0,1 \ldots, N\}$ and consider, for $i \leq N-1$, the annuli $\mathcal{C}_{i}:=B_{\sigma_{i}}(p) \backslash B_{\sigma_{i+1}}(p)$. If $\hat{\varepsilon}$ is sufficiently small, we can assume that $N \geq 2$ and $\sigma_{N} \geq \frac{s}{4}$. For at least one of these annuli we must have $\left|\bar{A} \cap \mathcal{C}_{i}\right| \geq \frac{1}{2}\left|\mathcal{C}_{i}\right|$. We then let $\sigma:=\sigma_{i}$ be the corresponding outer radius and we denote by $\mathcal{C}$ the corresponding annulus.

Consider now a point $x \in \mathcal{C} \cap \bar{A}$ and let $T_{x}$ be the slice $\langle T, \mathbf{p}, x\rangle$. Since $\bar{A} \subset K$, for a.e. $x \in \bar{A}$ we have $T_{x}=\sum_{i=1}^{Q} \llbracket\left(x, f_{i}(x)\right) \rrbracket$. Moreover, there exist $i$ and $j$ such that $\left|f_{i}(x)-f_{j}(x)\right|^{2} \geq \frac{1}{Q} \mathcal{G}(f(x), \llbracket \boldsymbol{\eta} \circ f(x) \rrbracket)^{2} \geq \frac{\hat{\delta}}{4 Q} E$ (recall that $x \in \bar{A} \subset A$ ). When $x \in \mathcal{C}$ and the points $(x, y)$ and $(x, z)$ belong both to $\mathbf{B}_{\sigma}((p, q))$, we must have

$$
|y-z|^{2} \leq 4\left(\sigma^{2}-\left(\sigma-\frac{\hat{\delta} E}{64 Q s}\right)^{2}\right) \leq \frac{\sigma \hat{\delta} E}{8 Q s} \leq \frac{\hat{\delta} E}{8 Q}
$$

Thus, for $x \in \bar{A} \cap \mathcal{C}$ at least one of the points $\left(x, f_{i}(x)\right)$ is not contained in $\mathbf{B}_{\sigma}((p, q))$. We conclude therefore

$$
\begin{align*}
\|T\|\left(\mathbf{C}_{\sigma}(p) \backslash \mathbf{B}_{\sigma}((p, q))\right) & \geq|\mathcal{C} \cap \bar{A}| \geq \frac{1}{2}|\mathcal{C}|=\frac{\omega_{m}}{2}\left(\sigma^{m}-\left(\sigma-\frac{\hat{\delta} E}{64 Q s}\right)^{m}\right) \\
& \geq \frac{\omega_{m}}{2} \sigma^{m}\left(1-\left(1-\frac{\hat{\delta} E}{64 Q s \sigma}\right)^{m}\right) \tag{8.4}
\end{align*}
$$

Recall that, for $\tau$ sufficiently small, $(1-\tau)^{m} \leq 1-\frac{m \tau}{2}$. Since $\sigma \geq \frac{s}{4}$, if $\hat{\varepsilon}$ is chosen sufficiently small we can therefore conclude

$$
\begin{equation*}
\|T\|\left(\mathbf{C}_{\sigma}(p) \backslash \mathbf{B}_{\sigma}(p)\right) \geq \frac{\omega_{m} \sigma^{m} \hat{\delta} E}{256 Q s \sigma} \geq \frac{\omega_{m}}{1024 Q} \hat{\delta} E \sigma^{m-2}=c_{0} \hat{\delta} E \sigma^{m-2} \tag{8.5}
\end{equation*}
$$

Next, by Theorem 1.4 and Theorem 1.6,

$$
\begin{equation*}
\|T\|\left(\mathbf{C}_{\sigma}(p)\right) \leq Q \omega_{m} \sigma^{m}+C E^{1+\gamma_{1}}+\bar{\eta} E+\int_{B_{\sigma}(p)} \frac{|D w|^{2}}{2} \tag{8.6}
\end{equation*}
$$

Moreover, as shown in [14, Section 3.3] (cf. [14, Proposition 3.10]), we have

$$
\begin{equation*}
\int_{B_{\sigma}(p)}|D w|^{2} \leq\|D \Psi\|^{2} \sigma^{m}+C \int_{B_{\sigma}(p)}|D u|^{2} \leq C\left(1+C^{\star}\right) E \sigma^{m}+C \operatorname{Dir}(u) \sigma^{m-2+2 \kappa} \tag{8.7}
\end{equation*}
$$

(for some constants $\kappa$ and $C$ depending only on $m, n$ and $Q$; in fact the exponent $\kappa$ is the one of (8.1)). Combining (8.5), (8.6) and (8.7), we conclude

$$
\begin{equation*}
\|T\|\left(\mathbf{B}_{\sigma}((p, q))\right) \leq Q \omega_{m} \sigma^{m}+\left(\bar{\eta}+C\left(1+C^{\star}\right) \sigma^{m}\right) E+C E^{1+\gamma_{1}}+C E \sigma^{m-2+2 \kappa}-c_{0} \sigma^{m-2} \hat{\delta} E \tag{8.8}
\end{equation*}
$$

Next, by the monotonicity formula, $\rho \mapsto \exp \left(C \mathbf{A}^{2} \rho^{2}\right) \rho^{-m}\|T\|\left(\mathbf{B}_{\rho}((p, q))\right)$ is a monotone function (indeed, the usual monotonicity formula of the theory of varifolds with bounded mean curvature gives the monotonicity of $\rho \mapsto \exp (C \mathbf{A} \rho) \rho^{-m}\|T\|\left(\mathbf{B}_{\rho}((p, q))\right)$, cf. [29, Theorem 17.6]); the slight improvement needed in this proof follows from minor modifications of the usual argument but, since we have not been able to find a reference, we provide a proof in Lemma A. 1 in the appendix). Using $\mathbf{A}^{2} \leq C^{\star} E, \Theta(T,(p, q))=Q$ and the Taylor expansion of the exponential, we conclude

$$
\begin{equation*}
\|T\|\left(\mathbf{B}_{\sigma}((p, q))\right) \geq Q \omega_{m} \sigma^{m}-C C^{\star} E \sigma^{m+2} \tag{8.9}
\end{equation*}
$$

Combining (8.8) and (8.9) we conclude

$$
\begin{equation*}
C\left(1+C^{\star}\right) \sigma^{2}+\left(\bar{\eta}+C E_{1}^{\gamma}\right) \sigma^{2-m}+C \sigma^{2 \kappa} \geq c_{0} \hat{\delta} \tag{8.10}
\end{equation*}
$$

Recalling that $\sigma \leq s<\bar{s}$, we can, finally, specify $\bar{s}$ : it is chosen so that $C\left(1+C^{\star}\right) \bar{s}^{2}+C \bar{s}^{2 \kappa}$ is smaller than $\frac{c_{0}}{2} \hat{\delta}$. Combined with (8.3) this choice of $\bar{s}$ depends, therefore, only upon $\hat{\delta}$. (8.10) becomes then

$$
\begin{equation*}
\left(\bar{\eta}+C E^{\gamma_{1}}\right) \sigma^{2-m} \geq \frac{c_{0}}{2} \hat{\delta} \tag{8.11}
\end{equation*}
$$

Next, recall that $\sigma \geq \frac{s}{4}$. We then choose $\hat{\varepsilon}$ so that $\left(\bar{\eta}+C \hat{\varepsilon}^{\gamma_{1}}\right)\left(\frac{s}{4}\right)^{2-m} \leq \frac{c_{0}}{4} \hat{\delta}$. This choice is incompatible with (8.11), thereby reaching a contradiction: for this choice of the parameter $\hat{\varepsilon}$ (which in fact depends only upon $\hat{\delta}$ and $s$ ) the conclusion of the Theorem, i.e. (1.9), must then be valid.

## Appendix A. Monotonicity formula

Lemma A.1. There is a constant $C$ depending only on $m, n$ and $\bar{n}$ with the following property. If $\Sigma \subset \mathbb{R}^{m+n}$ is a $C^{2}(m+\bar{n})$-dimensional submanifold with $\left\|A_{\Sigma}\right\|_{\infty} \leq \mathbf{A}$ and $T$ an m-dimensional integer-rectifiable current supported in $\Sigma$ which is stationary in $\Sigma$, then for every $\xi \in \Sigma$ the function $\rho \mapsto \exp \left(C \mathbf{A}^{2} \rho^{2}\right) \rho^{-m}\|T\|\left(\mathbf{B}_{\rho}(\xi)\right)$ is monotone on the interval $] 0, \bar{\rho}\left[\right.$, where $\bar{\rho}:=\min \left\{\operatorname{dist}(x, \operatorname{spt}(\partial T)),(C \mathbf{A})^{-1}\right\}$.
Proof. The argument is a minor variant of the classical proof of the monotonicity formula for varifolds with bounded mean curvature due to Allard (cf. [2]). Here the stronger hypothesis that $T$ is stationary in a $C^{2}$-submanifold allows a better estimate of the relevant error term. Without loss of generality assume $\xi=0$, let $s \in] 0, \bar{\rho}\left[\right.$ and $\varphi \in C_{c}^{1}(]-1,1[)$ with $\varphi \equiv 1$ in a neighborhood of 0 . For each $x \in \Sigma$ let $\mathbf{p}_{x}: \mathbb{R}^{m+n} \rightarrow T_{x} \Sigma$ be the orthogonal projection onto the tangent space to $\Sigma$ in $x$ and consider the vector field $X_{s}(x):=\varphi\left(\frac{|x|}{s}\right) \mathbf{p}_{x}(x)$. Note that $X_{s}$ is tangent to $\Sigma$ and thus $\delta T\left(X_{s}\right)=0$. In order to compute $\delta T\left(X_{s}\right)$, consider
at $\|T\|$-a.e. $x \in \operatorname{spt}(T)$ an orthonormal frame $e_{1}, \ldots, e_{m}$ with $e_{1} \wedge \ldots \wedge e_{m}=\vec{T}$. It turns out that

$$
\delta T\left(X_{s}\right)=\int \operatorname{div}_{\vec{T}} X_{s} d\|T\|=\int \sum_{i}\left\langle D_{e_{i}} X_{s}, e_{i}\right\rangle d\|T\|
$$

Next, at any $x \in \Sigma$ let $\nu_{1}, \ldots, \nu_{l}(l=n-\bar{n})$ be an orthonormal frame orthogonal to $\Sigma$. Since $\mathbf{p}_{x}(x)=x-\sum_{j}\left\langle x, \nu_{j}\right\rangle \nu_{j}$ and $\left\langle e_{i}, \nu_{j}\right\rangle=0$, we compute:

$$
\operatorname{div}_{\vec{T}} X_{s}(x)=\underbrace{\sum_{i}\left[D_{e_{i}}\left(\varphi\left(\frac{|x|}{s}\right)\right)\left\langle x, e_{i}\right\rangle+\varphi\left(\frac{|x|}{s}\right)\left\langle D_{e_{i}} x, e_{i}\right\rangle\right]}_{I}-\underbrace{\varphi\left(\frac{|x|}{s}\right) \sum_{i, j}\left\langle x, \nu_{j}\right\rangle\left\langle D_{e_{i}} \nu_{j}, e_{i}\right\rangle}_{I I} .
$$

$I$ is the usual expression appearing in the proof of the standard monotonicity formula for stationary varifolds. If we use the notation $r$ for the function $x \mapsto|x|$ and $\nabla^{\perp} r$ for the orthogonal projection on the orthogonal complement of $\operatorname{Span}\left\{e_{1}, \ldots, e_{m}\right\}$, we find $I=m \varphi\left(\frac{r}{s}\right)+\frac{r}{s} \varphi^{\prime}\left(\frac{r}{s}\right)\left(1-\left|\nabla^{\perp} r\right|^{2}\right)$ (see for instance [11, (2.2)]). In order to bound $I I$, we first observe that $\left\langle D_{e_{i}} \nu_{j}, e_{i}\right\rangle=-\left\langle A\left(e_{i}, e_{i}\right), \nu_{j}\right\rangle$. Next, since $r \leq(C \mathbf{A})^{-1}$, if $C$ is chosen sufficiently large we can assume that the geodesic segment of $\Sigma$ connecting 0 and $x$ has length $\ell<2 r$. Denote by $\gamma:[0, \ell] \rightarrow \Sigma$ a parametrization by arc-length of such a segment. Then,

$$
\begin{equation*}
\left\langle x, \nu_{j}(x)\right\rangle=\int_{0}^{\ell}\left\langle\dot{\gamma}(\sigma), \nu_{j}(\gamma(\ell))\right\rangle d \sigma=\int_{0}^{\ell} \underbrace{\left\langle\dot{\gamma}(\sigma),\left[\nu_{j}(\gamma(\ell))-\nu_{j}(\gamma(\sigma))\right]\right\rangle}_{g(\sigma)} d \sigma \tag{A.1}
\end{equation*}
$$

and observe that

$$
\left|g^{\prime}(\sigma)\right| \leq 2\left|\frac{d}{d \sigma} \dot{\gamma}(\sigma)\right|+\left|\left\langle\dot{\gamma}(\sigma), \frac{d}{d \sigma} \nu_{j}(\gamma(\sigma))\right\rangle\right| \leq 3|A(\dot{\gamma}(\sigma), \dot{\gamma}(\sigma))| .
$$

Since $g(\ell)=0$, integrating the latter inequality we conclude $|g(\sigma)| \leq 3 \ell \mathbf{A} \leq 6 r \mathbf{A}$, which in turn, together with (A.1), gives $\left|x \cdot \nu_{j}(x)\right| \leq 12 r^{2} \mathbf{A}$.

Putting all estimates together, we achieve the inequality $|I I| \leq C \varphi\left(\frac{r}{s}\right) r^{2} \mathbf{A}^{2}$. From here on we can follow the usual strategy leading to the monotonicity formula (cf. [29] or [11, Proof of Theorem 2.1]): letting the test function $\varphi$ converge from below to the indicator function of $]-1,1[$, after few manipulations we achieve the inequality

$$
\frac{d}{d s} \frac{\|T\|\left(\mathbf{B}_{s}\right)}{s^{m}} \geq-C \mathbf{A}^{2} s \frac{\|T\|\left(\mathbf{B}_{s}\right)}{s^{m}}
$$

which leads to the desired claim.
Remark A.2. The proof can be easily extended to varifolds which are stationary in $\Sigma$. In fact the argument above can be considerably shortened using directly the Monotonicity Formula of Section 5 in [2].

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