JACOBIANS OF $W^{1,p}$ HOMEOMORPHISMS, CASE $p = \lfloor n/2 \rfloor$

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ABSTRACT. We investigate a known problem whether a Sobolev homeomorphism between domains in \mathbb{R}^n can change sign of the Jacobian. The only case that remains open is when $f \in W^{1,[n/2]}$, $n \geq 4$. We prove that if $n \geq 4$, and a sense-preserving homeomorphism f satisfies $f \in W^{1,[n/2]}$, $f^{-1} \in W^{1,n-[n/2]-1}$ and either f is Hölder continuous on almost all spheres of dimension [n/2], or f^{-1} is Hölder continuous on almost all spheres of dimensions n-[n/2]-1, then the Jacobian of f is non-negative, $J_f \geq 0$, almost everywhere. This result is a consequence of a more general result proved in the paper. Here [x] stands for the greatest integer less than or equal to x.

In memoriam: Boqdan Bojarski (1931-2018)

1. Introduction and results

A diffeomorphism f between domains in \mathbb{R}^n has either positive or negative Jacobian $J_f = \det Df$. Recall that domains are open and connected. We say that a diffeomorphism is sense preserving (sense reversing) if its Jacobian is positive (negative). More generally, we say that a homeomorphism between domains is sense preserving (sense reversing) if it has local topological degree 1 (-1) at every point of its domain, see Section 2.5. Every homeomorphism is either sense preserving or sense reversing. It easily follows from the topological properties of the degree that if f is a sense preserving (reversing) homeomorphism of domains in \mathbb{R}^n , f is differentiable at x and $J_f(x) \neq 0$, then $J_f(x) > 0$ ($J_f(x) < 0$). In particular, if a homeomorphism is differentiable almost everywhere, then $J_f \geq 0$ a.e. or $J_f \leq 0$ a.e. However, it may happen that a homeomorphism is differentiable a.e., but its Jacobian equals zero a.e. For elementary constructions, see [40] and references therein.

Let us assume now that a homeomorphism f between domains in \mathbb{R}^n is in the Sobolev space $W_{\text{loc}}^{1,p}$, $p \geq 1$. If p > n-1, then f is differentiable a.e., [21, Corollary 2.25], and therefore $J_f \geq 0$ a.e. or $J_f \leq 0$ a.e. Another approach, based on the topological degree, allows one to extend this result to $p \geq n-1$. However, the method completely fails when $1 \leq p < n-1$. Note also that, if $1 \leq p < n$, then there are very pathological examples of Sobolev homeomorphisms with the Jacobian equal zero a.e. The first such example was constructed by Hencl [20], see also [5, 8].

In 2001, Hajłasz asked a question whether a Sobolev $W_{\text{loc}}^{1,p}$, $1 \leq p < n-1$, homeomorphism between domains in \mathbb{R}^n can change sign of the Jacobian. That is, whether there is a homeomorphism such that $J_f > 0$ on a set of positive measure and $J_f < 0$ on a set of positive measure.

Hencl and Malý [24] proved the following two results:

Theorem 1. Let $\Omega \subset \mathbb{R}^n$, $n \leq 3$, be a domain and let $f \in W^{1,1}_{loc}(\Omega, \mathbb{R}^n)$ be a sense preserving homeomorphism. Then $J_f \geq 0$ a.e.

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Theorem 2. Let $\Omega \subset \mathbb{R}^n$, $n \geq 4$, be a domain and let $f \in W^{1,p}_{loc}(\Omega,\mathbb{R}^n)$, p > [n/2] be a sense preserving homeomorphism. Then $J_f \geq 0$ a.e.

Here [x] stands for the greatest integer less than or equal to x. Also, by a homeomorphism $f:\Omega\to\mathbb{R}^n$ we mean a homeomorphism onto the image. Since for p>[n/2] we have the embedding into the Lorentz space $L^p_{\mathrm{loc}}\subset L^{[n/2],1}_{\mathrm{loc}}$, Theorem 2 is a special case of a more general result of Hencl and Malý:

Theorem 3. Let $\Omega \subset \mathbb{R}^n$, $n \geq 4$, be a domain and let $f: \Omega \to \mathbb{R}^n$ be a sense preserving homeomorphism such that $Df \in L^{[n/2],1}_{loc}$. Then $J_f \geq 0$ a.e.

On the other hand, Hencl and his collaborators, [4, 26], constructed the following surprising example:

Theorem 4. If $n \ge 4$ and $1 \le p < \lfloor n/2 \rfloor$, then there is a homeomorphism $f \in W^{1,p}((-1,1)^n, \mathbb{R}^n)$ such that $J_f > 0$ on a set of positive measure and $J_f < 0$ on a set of positive measure. Moreover, f has the Lusin property.

Recall that the Lusin property means that the sets of Lebesgue measure zero are mapped to sets of Lebesgue measure zero.

The result of [26] provides such a homeomorphism for $n \geq 4$ with p = 1, and the general case is obtained in [4]. See also [16, 17] for related examples of approximately differentiable homeomorphisms.

Theorems 1, 2, and 4 leave only the borderline case open.

Question 1. Let $\Omega \subset \mathbb{R}^n$, $n \geq 4$, be a domain. Does there exist a homeomorphism $f \in W_{\text{loc}}^{1,[n/2]}(\Omega,\mathbb{R}^n)$ such that $J_f > 0$ on a set of positive measure and $J_f < 0$ on a set of positive measure?

The main result of the paper answers this question in the negative under some additional assumptions.

Theorem 5. Let $\Omega \subset \mathbb{R}^n$, $n \geq 4$, be a domain and let $f \in W^{1,[n/2]}_{loc}(\Omega,\mathbb{R}^n)$ be a sense preserving homeomorphism such that $f^{-1} \in W^{1,n-[n/2]-1}_{loc}(f(\Omega),\mathbb{R}^n)$. Assume also that one of the two conditions is satisfied:

- (a) f maps almost all spheres of dimension [n/2] to sets of $\mathcal{H}^{[n/2]+1}$ -measure zero,
- (b) f^{-1} maps almost all spheres of dimension $n \lfloor n/2 \rfloor 1$ to sets of $\mathcal{H}^{n-\lfloor n/2 \rfloor}$ -measure zero.

Then $J_f \geq 0$ a.e.

Here and in what follows, by \mathcal{H}^k we shall denote the k-dimensional Hausdorff measure.

Remark 6. The space of k-dimensional spheres in \mathbb{R}^n can be parameterized by the product $G(k+1,n)\times\mathbb{R}^n\times(0,\infty)$, where G(k+1,n) is the Grassmannian of (k+1)-dimensional subspaces in \mathbb{R}^n . Indeed, $(V,x,r)\in G(k+1,n)\times\mathbb{R}^n\times(0,\infty)$ defines a sphere centered at x, of radius r and parallel to V. Since there is a natural measure on $G(k+1,n)\times\mathbb{R}^n\times(0,\infty)$, it makes sense to talk about almost every k-dimensional sphere in \mathbb{R}^n .

Remark 7. The classes of bi-Sobolev homeomorphisms, i.e., homeomorphisms such that f and f^{-1} belong to Sobolev spaces, have been investigated for example in [6, 8, 22, 23, 25, 33, 34, 36].

The corollaries listed below show particular situations when the condition (a) or (b) is satisfied.

Corollary 8. Let $\Omega \subset \mathbb{R}^n$, $n \geq 4$, be a domain and let $f \in W^{1,[n/2]}_{loc}(\Omega,\mathbb{R}^n)$ be a sense preserving homeomorphism such that $f^{-1} \in W^{1,n-[n/2]-1}_{loc}(f(\Omega),\mathbb{R}^n)$. Assume also that f or f^{-1} is Hölder continuous. Then $J_f \geq 0$ a.e.

Remark 9. In fact, it suffices to assume f is Hölder continuous on almost all $\lfloor n/2 \rfloor$ -dimensional spheres or f^{-1} is Hölder continuous on almost all $n - \lfloor n/2 \rfloor - 1$ dimensional spheres; the proof remains the same.

Corollary 10. Let $\Omega \subset \mathbb{R}^{2m}$, $m \geq 2$, be a domain in the even dimensional space and let $f \in W^{1,m}_{loc}(\Omega,\mathbb{R}^{2m})$ be a sense preserving homeomorphism such that $f^{-1} \in W^{1,m-1+\varepsilon}_{loc}(f(\Omega),\mathbb{R}^m)$. Then $J_f \geq 0$ a.e.

Corollary 11. Let $\Omega \subset \mathbb{R}^{2m}$, $m \geq 2$, be a domain in the even dimensional space and let $f \in W^{1,m}_{loc}(\Omega,\mathbb{R}^{2m})$ be a sense preserving homeomorphism such that $Df^{-1} \in L^{m-1,1}_{loc}$. Then $J_f \geq 0$ a.e.

Remark 12. In Corollaries 10 and 11, we restrict the setting to even dimensions, because a corresponding result in odd dimensions would be a consequence of Theorems 2 and 3 respectively.

The corollaries easily follow from Theorem 5.

Proof of Corollary 8. If $f \in W^{1,k}(\mathbb{S}^k, \mathbb{R}^m)$, then there is a set $E \subset \mathbb{S}^k$ of measure zero such that the complement of this set is the union of the sets such that on each of these sets f is Lipschitz continuous (see a discussion around (2) below) and hence the Hausdorff dimension of $f(\mathbb{S}^k \setminus E)$ is at most k. According to a theorem of Malý and Martio [30, Theorem C], [42, Theorem 1], if $f \in W^{1,k}(\mathbb{S}^k, \mathbb{R}^m)$ is Hölder continuous, then it maps sets of measure zero to sets of \mathcal{H}^k -measure zero so $\mathcal{H}^k(E) = 0$ and hence the Hausdorff dimension of $f(\mathbb{S}^k)$ is at most k. Let now f be as in Corollary 8. Assume that f is Hölder continuous. According to the Fubini theorem for Sobolev functions (Lemma 26), f restricted to almost all spheres [n/2]-dimensional spheres is a Hölder continuous map in $W^{1,[n/2]}$ so the image of almost every such sphere has Hausdorff dimension at most [n/2] and hence its $\mathcal{H}^{[n/2]+1}$ -measure is zero, so condition (a) from Theorem 5 is satisfied and the result follows. Similarly, if f^{-1} is Hölder continuous, the condition (b) is satisfied and the result follows.

Proof of Corollaries 10 and 11. Since $L^{m-1+\varepsilon}_{loc} \subset L^{m-1,1}_{loc}$, Corollary 10 follows from Corollary 11. According to [23, Theorem C], mappings $f: \mathbb{S}^k \to \mathbb{R}^m$ with the weak derivative in $L^{k,1}$ map sets of measure zero to sets of \mathcal{H}^k -measure zero, so exactly the same argument as in the proof of Corollary 8 yields that $\mathcal{H}^{k+1}(f(\mathbb{S}^k)) = 0$ and then, again as in the proof of Corollary 8, the result follows.

The main idea in the proofs of Theorems 1 and 2 is to use the linking number. If $n \geq 4$ and $p > \lfloor n/2 \rfloor$, one can find linked spheres in Ω of dimensions less than p. This allows one to use the Sobolev embedding theorem on the linked spheres to control the topological linking number in terms of the Sobolev norm of the mapping. Since a sense preserving homeomorphism maps linked spheres onto linked topological spheres with the same linking number, one can use this fact to prove that the Jacobian of a sense preserving map cannot be negative on a set of positive measure. A similar argument is used when $n \leq 3$.

The proof of our Theorem 5 is based on a similar idea. However, we cannot use the Sobolev embedding theorem on spheres, because now $p = k = \lfloor n/2 \rfloor$ equals to the dimension of one of the

linked spheres. This causes many technical problems and in order to handle them, we need to assume Sobolev regularity of the inverse map.

Although Theorem 5 gives an answer to Question 1 only in a very special case, the main motivation behind Theorem 5 was to modify the technique of the linking number so it could be used in the limiting case, in which we do not have the Sobolev embedding on spheres. We believe that if the answer to Question 1 is in the negative, the proof should be based on the linking number technique and we hope that, with further modifications, our new technique can lead to the negative answer to Question 1 in full generality, for all $n \ge 4$. However, we do not know yet how to do it and we are not even sure what the final answer to Question 1 is.

The paper is structured as follows. In Section 2 we collect basic tools that are used in the proof of Theorem 5. Some of the tools collected there are known, but some other are new and of independent interest. In Section 3 we recall the proof of Theorems 1 and 2. This helps to understand the main idea of our proof and to see what are the additional difficulties we have to face. In the last Section 4 we prove Theorem 5. We put a lot of effort to make the paper self-contained and accessible to those who are new to this area of research.

Notation in the paper is quite standard. The Lebesgue measure of a set $A \subset \mathbb{R}^n$ is denoted by |A|. By \mathcal{H}^k we denote the k-dimensional Hausdorff measure. A k-dimensional open ball centered at a point x, with radius r, is denoted by $\mathbb{B}^k(x,r)$, and \mathbb{B}^k denotes the open unit ball in k dimensions. Similarly, \mathbb{S}^k denotes the unit k-dimensional sphere. The surface measure on \mathbb{S}^k is denoted by $d\sigma(x)$. Open half-space will be denoted by $\mathbb{R}^{n+1}_+ = \mathbb{R}^n \times (0,\infty)$. $W^{1,p}$ is the Sobolev space of functions $f \in L^p$ with $\nabla f \in L^p$. The Lorentz space is denoted by $L^{p,q}$. We do not recall the definition of this space since it does not play any role in our proofs. The Jacobian of a mapping $f \colon \mathbb{R}^n \to \mathbb{R}^n$ is denoted by $L^p = \det Df$. A domain is an open and connected set. The integral average is denoted by

$$\int_{E} f \, dx = \frac{1}{|E|} \int_{E} f \, dx.$$

By C we denote a generic constant whose value may change in a single string of estimates. Writing C = C(n, m) we will indicate that the constant C depends on n and m only.

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A few days before completion of this work we learned the sad news that Professor Bogdan Bojarski had passed away. He was the PhD advisor of Piotr Hajłasz and an inspiration for both of us. We mourn his passing, and we dedicate this paper with deep respect to his memory.

2. Preliminaries

In this section we collect some basic facts that are used in the proof of the main result. We present the results in a slightly more general form than we actually need, because they might be useful for some other applications.

2.1. Chain rule. The main result of [8] (see also [33]) provides an example of a surjective homeomorphism $f:(0,1)^n\to(0,1)^n, n\geq 3$, such that $f\in W^{1,1}, f^{-1}\in W^{1,1}$ and $J_f=0$ a.e., $J_{f^{-1}}=0$ a.e. Note that $f^{-1}\circ f=\mathrm{Id}$, but the chain rule

$$Df^{-1}(f(x))Df(x) = \operatorname{Id}$$

cannot be satisfied on a set of positive measure because of the vanishing Jacobians. In fact, f maps the set of full measure to the set of measure zero where Df^{-1} is not defined. The situation is different if we assume that $J_f \neq 0$. Namely, we have

Lemma 13. Let $U, V \subset \mathbb{R}^n$ be open sets. Assume that $f \in W^{1,1}_{loc}(U, V)$, $g \in W^{1,1}_{loc}(V, \mathbb{R})$ and $g \circ f \in W^{1,1}_{loc}(U, \mathbb{R})$. Then

(1)
$$D(g \circ f)(x) = Dg(f(x))Df(x)$$

for almost all points x in the set $\{x \in U : J_f(x) \neq 0\}$.

Remark 14. In particular, the result says that Dg(f(x)) is well defined at almost all points x such that $J_f(x) \neq 0$.

Proof. The set where the Jacobian is different than zero splits into two sets where the Jacobian is positive and negative, respectively. Thus it suffices to show that (1) is satisfied almost everywhere in the set where the Jacobian is positive,

$$X = \{x \in U : J_f(x) > 0\},\$$

because a similar argument can be applied to the set where the Jacobian is negative.

It is well known [1, 3, 18] that $u \in W^{1,1}(\mathbb{R}^n)$ satisfies the pointwise inequality

(2)
$$|u(x) - u(y)| \le C(n)|x - y|(\mathcal{M}|Du|(x) + \mathcal{M}|Du|(y)) \quad \text{a.e.,}$$

where $\mathcal{M}|Du|(x) = \sup_{r>0} \int_{\mathbb{B}^n(x,r)} |Du|(y) \, dy$ is the Hardy-Littlewood maximal function. Hence for each t>0, u restricted to the set $\{\mathcal{M}|Du|\leq t\}$ is Lipschitz continuous. This implies that \mathbb{R}^n can be decomposed into a set of measure zero and countably many sets such that on each of these sets u is Lipschitz continuous. This fact and a partition of unity argument implies that U can be decomposed into Borel sets

$$(3) U = N_o \cup \bigcup_{i=1}^{\infty} K_i$$

such that $|N_o| = 0$ and $f|_{K_i}$ is Lipschitz continuous. We need to use here a partition of unity argument, because f is defined in U, while (2) applies to functions defined on \mathbb{R}^n .

It remains to show that (1) is satisfied at almost all points of the set $X \cap K_i$ for each i = 1, 2, ...Let f_i be a Lipschitz extension of $f|_{X \cap K_i}$ to all of \mathbb{R}^n (see [10, Theorem 3.1]). According to the Rademacher theorem, [10, Theorem 3.2], Df_i exists a.e. Also $Df_i = Df$ a.e. in $X \cap K_i$. Indeed, $f - f_i = 0$ in $X \cap K_i$ so $D(f - f_i) = 0$ a.e. in $X \cap K_i$ by [10, Theorem 4.4(iv)]. Let

$$W_i = \{x \in X \cap K_i : Df_i(x) \text{ exists and } Df_i(x) = Df(x)\}.$$

Since $|(X \cap K_i) \setminus W_i| = 0$, it remains to show that (1) is satisfied at almost all points of the set W_i . Note that $J_{f_i} > 0$ on W_i . According to [10, Lemma 3.3] we can decompose the set W_i into a family of pairwise disjoint Borel sets

$$W_i = \bigcup_{j=1}^{\infty} S_j$$

such that $f_i|_{S_j}$ is bi-Lipschitz for each j=1,2,..., and it remains to prove (1) at almost all points of S_j . If $|S_j|=0$, the result is obvious, so we can assume that $|S_j|>0$ and hence $f(S_j)=f_i(S_j)$ has positive measure, too.

Since $g \in W_{loc}^{1,1}$, we have a decomposition

$$f(S_j) = M_o \cup \bigcup_{k=1}^{\infty} E_k, \quad |M_o| = 0, \quad g|_{E_k} \text{ is Lipschitz continuous.}$$

Indeed, we have a decomposition of V similar to (3) and then we take intersections with the set $f(S_j)$. We can also assume that $|E_k| > 0$ for all k, as otherwise we could add sets E_k of measure zero to the set M_o . Since the mapping f is bi-Lipschitz on S_j , we have that $|(f|_{S_j})^{-1}(M_o)| = 0$ and it remains to prove that (1) is satisfied at almost all points of $Z_{jk} = (f|_{S_j})^{-1}(E_k)$ for k = 1, 2, ... Let g_k be a Lipschitz extension of $g|_{E_k}$ to all of \mathbb{R}^n . Then g_k is differentiable a.e. and $Dg_k = Dg$ almost everywhere in E_k . Since $f|_{S_j}$ is bi-Lipschitz, the preimage $(f|_{S_j})^{-1}$ of the set of points in E_k where g_k is not differentiable has measure zero. This and the classical chain rule for differentiable functions imply that $g \circ f = g_k \circ f_i$ in Z_{jk} and

$$D(g_k \circ f_i)(x) = Dg_k(f_i(x))Df_i(x) = Dg(f(x))Df(x) \quad \text{a.e. in } Z_{jk}.$$

Since $g_k \circ f_i$ is Lipschitz continuous and it coincides with $g \circ f$ in Z_{jk} , it follows that $D(g_k \circ f_i) = D(g \circ f)$ almost everywhere in Z_{jk} .

As an immediate corollary we obtain the following result that will be used in the proof of Theorem 5, see also [9, Theorem 1.1 and 1.3], [12, Lemma 2.1], [21].

Corollary 15. Assume that $\Omega \subset \mathbb{R}^n$ is a domain and $f \in W^{1,1}_{loc}(\Omega,\mathbb{R}^n)$ is a homeomorphism such that $f^{-1} \in W^{1,1}_{loc}(f(\Omega),\mathbb{R}^n)$. Then

(4)
$$(Df(x))^{-1} = Df^{-1}(f(x))$$

almost everywhere in the set where $J_f \neq 0$.

In particular, Corollary 15 applies to homeomorphisms described in Theorem 5.

Corollary 16. Let $U \subset \mathbb{R}^n$ be open and let $f \in W^{1,1}_{loc}(U,\mathbb{R}^n)$ be continuous. If a compact set $K \subset \{x \in U : J_f(x) \neq 0\}$ has positive measure, then the set f(K) has positive measure.

Remark 17. In general, continuous mappings (even homeomorphisms) may map measurable sets to non-measurable sets. This is why we assume that K is compact to guarantee measurability of the set f(K).

Proof. This is a corollary of the proof of Lemma 13 and we assume the same notation as in the proof of Lemma 13. In particular we assume that the sets W_i and S_j are defined in the same way.

Let $K \subset \{x \in U : J_f(x) \neq 0\}$ be compact and of positive measure. Since f is continuous, f(K) is compact and hence measurable. One of the sets $K \cap \{J_f > 0\}$ or $K \cap \{J_f < 0\}$ has positive measure. Without loss of generality we may assume that the set $K \cap \{J_f > 0\}$ has positive measure.

The sets W_i constructed in the proof of Lemma 13 cover almost all points of the set $\{J_f > 0\}$ so $|K \cap W_i| > 0$ for some i. Since W_i is the union of sets S_j , $|K \cap S_j| > 0$ for some j. The mapping $f|_{S_j}$ is bi-Lipschitz and it follows that $f(K \cap S_j)$ is measurable and of positive measure. Hence also f(K) has positive measure.

2.2. Blow-up technique. In this section we describe a *blow-up technique* (Lemma 21) that is often used in the study of partial differential equations. This technique has also been used in [24]. Later, we generalize the blow-up technique to a simultaneous blow-up for a homeomorphism and its inverse, Lemma 25. This result will be used in the proof of Theorem 5.

All results of this section are local in nature, so they are true for functions and mappings defined on domains in \mathbb{R}^n and not necessarily on all of \mathbb{R}^n . However, for simplicity of notation we decided to formulate the results on \mathbb{R}^n .

We will need the following two classical lemmata, the first of which is due to Lebesgue.

Lemma 18 (The Lebesgue Differentiation Theorem). If $f \in L^p_{loc}(\mathbb{R}^n)$, $1 \leq p < \infty$, then

(5)
$$\oint_{\mathbb{B}^n(x,r)} |f(y) - f(x)|^p dy \xrightarrow{r \to 0} 0 for a.e. x \in \mathbb{R}^n.$$

The points $x \in \mathbb{R}^n$ where (5) holds true are called *p-Lebesgue points* of f.

The second, due to Calderón and Zygmund, [7, Theorem 12], is also an immediate consequence of [10, Theorem 6.2].

Lemma 19. If $f \in W^{1,p}_{loc}(\mathbb{R}^n)$, $1 \leq p < \infty$, then

(6)
$$\int_{\mathbb{B}^n(x,r)} \left| \frac{f(y) - f(x) - Df(x) \cdot (y - x)}{r} \right|^p dy \xrightarrow{r \to 0} 0 \quad \text{for a.e. } x \in \mathbb{R}^n.$$

Note that the above lemmata immediately generalize to the case of vector valued functions $f \in W^{1,p}_{loc}(\mathbb{R}^n,\mathbb{R}^k)$, since it suffices to apply them to components of f. In particular we have that if $f \in W^{1,p}_{loc}(\mathbb{R}^n,\mathbb{R}^k)$, then

(7)
$$\int_{\mathbb{B}^n(x,r)} |Df(y) - Df(x)|^p dy \xrightarrow{r \to 0} 0 \text{ for a.e. } x \in \mathbb{R}^n.$$

Definition 20. Let $f \in W^{1,p}_{loc}(\mathbb{R}^n, \mathbb{R}^k)$, $1 \leq p < \infty$. We say that $x \in \mathbb{R}^n$ is a *p-good point* for f if both of the integrals (6) and (7) converge to zero.

Clearly, almost all points of \mathbb{R}^n are p-good points for $f \in W^{1,p}_{loc}$.

The basic blow-up technique is described by the following lemma. It allows us to regard f almost as a linear map near any p-good point.

Lemma 21. Let $f \in W^{1,p}_{loc}(\mathbb{R}^n, \mathbb{R}^k)$. For a p-good point $x_o \in \mathbb{R}^n$ and r > 0 we define

$$f_r(x) = \frac{f(x_o + rx) - f(x_o)}{r}$$
 and $f_0(x) = Df(x_o)x$.

Then f_r converges to the linear map f_0 in the norm of $W^{1,p}(\mathbb{B}^n,\mathbb{R}^k)$ as $r \to 0$, where $\mathbb{B}^n = \mathbb{B}^n(0,1)$ is the unit ball.

Proof. Let $x_o \in \mathbb{R}^n$ be a p-good point for f. Note that $Df_0(x) = Df(x_o)$. We have

$$\int_{\mathbb{B}^n} |f_r(x) - f_0(x)|^p \, dx = \int_{\mathbb{B}^n(x_o, r)} \left| \frac{f(y) - f(x_o) - Df(x_o) \cdot (y - x_o)}{r} \right|^p \, dy \to 0$$

as $r \to 0$, and

$$\oint_{\mathbb{B}^n} |Df_r(x) - Df_0(x)|^p \, dx = \oint_{\mathbb{B}^n(x_o, r)} |Df(y) - Df(x_o)|^p \, dy \to 0 \quad \text{as } r \to 0.$$

The rest of the section is devoted to a simultaneous blow-up for a homomorphism and its inverse.

Lemma 22. Let $f \in W^{1,p}_{\mathrm{loc}}(\mathbb{R}^n,\mathbb{R}^n), \ 1 \le p < \infty,$ be a homeomorphism such that $f^{-1} \in \mathbb{R}^n$ $W_{\mathrm{loc}}^{1,q}(\mathbb{R}^n,\mathbb{R}^n), \ 1 \leq q < \infty.$ Then almost all points of the set $\{x \in \mathbb{R}^n : J_f(x) \neq 0\}$ have the following three properties satisfied simultaneously

- (a) x is a p-good point for f,
- (b) f(x) is a q-good point for f^{-1} , (c) $(Df(x))^{-1} = (Df)^{-1}(f(x))$.

Proof. A homeomorphic image of a Lebesgue measurable set need not be measurable, but a homeomorphic image of a Borel set is Borel, so we need to work with Borel sets.

Let $A \subset \mathbb{R}^n$ be a Borel set of q-good points for f^{-1} such that $|\mathbb{R}^n \setminus A| = 0$. Then $f^{-1}(A)$ is Borel and hence measurable. Almost all points of the set $\{J_f \neq 0\} \cap f^{-1}(A)$ have properties (a) and (b) and in order to show that almost all points of the set $\{J_f \neq 0\}$ have properties (a) and (b) it suffices to show that the set

$$X = \{x \in \mathbb{R}^n : J_f(x) \neq 0\} \setminus f^{-1}(A)$$

has measure zero. Suppose to the contrary that |X| > 0. Let $K \subset X$ be a compact set of positive measure. Then $f(K) \subset \mathbb{R}^n \setminus A$, and according to Corollary 16, f(K) has positive measure. This is, however, impossible, since $\mathbb{R}^n \setminus A$ has measure zero.

We proved that almost all points of the set $\{J_f \neq 0\}$ have properties (a) and (b). Now it follows from Corollary 15 that almost all points of the set $\{J_f \neq 0\}$ have all three properties (a), (b) and (c).

The next lemma is easy to prove.

Lemma 23. Let f be as in Lemma 22 and let $A \in GL(n)$ be a non-degenerate linear transformation on \mathbb{R}^n . If a point $x_o \in \{J_f \neq 0\}$ satisfies conditions (a), (b) and (c), then x_o also satisfies conditions (a), (b), (c) for a homeomorphism $g = A \circ f$.

Remark 24. Whether a point x_0 satisfies conditions (a), (b) and (c) for the mapping f depends on the choice of representatives of Df and Df^{-1} . More precisely, it depends on how the values of $Df(x_o)$ and $Df^{-1}(f(x_o))$ are defined. However, we proved that no matter how we choose representatives of Df and Df^{-1} , almost all points will satisfy (a), (b), (c). If a point x_o satisfies conditions (a), (b) and (c) for the mapping f, then we will prove that x_o satisfies the same conditions for g, provided the representatives of Dg and Dg^{-1} are such that

$$Dg(x_o) = ADf(x_o), \quad Dg^{-1}(g(x_o)) = Df^{-1}(f(x_o))A^{-1},$$

but we can make a choice of such representatives without any harm being done.

Sketch of a proof. The proof that x_o is good for g is straightforward. The proof that it is also good for $g^{-1} = f^{-1} \circ A^{-1}$ follows from the change of variables $\Phi(y) = A^{-1}(y)$. Then the averages over the balls will became averages over scaled and translated ellipsoids $A^{-1}\mathbb{B}^n$ and it remains to observe (a well known fact) that if averages at (6) and (7) (for g^{-1} in place of f) over the balls \mathbb{B}^n converge to zero, then also the averages over the ellipsoids $A^{-1}\mathbb{B}^n(g(x_o), r)$ converge to zero. Indeed, the average over an ellipsoid can be estimated from above by the average over a larger ball that contains $A^{-1}\mathbb{B}^n(g(x_o), r)$, with a uniform constant that does not depend on the diameter of the ellipsoid. Finally the condition (c) for $g = A \circ f$ is a consequence of linear algebra and the choice of representatives of Dg and Dg^{-1} (see Remark 24). We leave details to the reader.

Let f and x_o are as in Lemma 23. If $J_f(x_o) > 0$ and $A = Df(x_o)^{-1}$, then

$$Dg(x_o) = Dg^{-1}(g(x_o)) = \mathrm{Id}.$$

If $J_f(x_o) < 0$ and $A = \mathcal{R}Df(x_o)^{-1}$, where

$$\mathcal{R} = \mathrm{Diag}(1, 1, \dots, 1, -1) = \begin{pmatrix} 1 & 0 & \dots & 0 & 0 \\ 0 & 1 & \dots & 0 & 0 \\ \vdots & \vdots & \ddots & \vdots & \vdots \\ 0 & 0 & \dots & 1 & 0 \\ 0 & 0 & \dots & 0 & -1 \end{pmatrix},$$

then

$$Dg(x_o) = Dg^{-1}(g(x_o)) = \mathcal{R},$$

because $\mathcal{R} = \mathcal{R}^{-1}$.

In both cases the linear transformation A has positive determinant.

The next lemma describes the simultaneous blow-up in the case of negative Jacobian. This is what we will need in the proof of Theorem 5. In the case of positive Jacobian one can easily formulate a similar result with \mathcal{R} replaced by Id.

Lemma 25. Let $f \in W^{1,p}_{loc}(\mathbb{R}^n,\mathbb{R}^n)$, $1 \leq p < \infty$, be a homeomorphism such that $f^{-1} \in W^{1,q}_{loc}(\mathbb{R}^n,\mathbb{R}^n)$, $1 \leq q < \infty$. Then for almost every point x_o of the set $\{x \in \mathbb{R}^n : J_f(x) < 0\}$, there is a linear transformation A with positive determinant such that the homeomorphism $g = A \circ f$ satisfies

$$\lim_{r \to 0^+} \|g_r - \mathcal{R}\|_{W^{1,p}(\mathbb{B}^n,\mathbb{R}^n)} = \lim_{r \to 0^+} \|(g_r)^{-1} - \mathcal{R}\|_{W^{1,q}(\mathbb{B}^n,\mathbb{R}^n)} = 0,$$

where

$$g_r(x) = \frac{g(x_o + rx) - g(x_o)}{r}, \quad r > 0.$$

Proof. Almost every point of the set $\{J_f < 0\}$ satisfies conditions (a), (b) and (c) of Lemma 22 for f. Fix such a point x_o . Then, by Lemma 23, x_o satisfies conditions (a), (b) and (c) for $g = A \circ f$. Choosing $A = \mathcal{R}(Df(x_o))^{-1}$, we have that

$$Dg(x_o) = Dg^{-1}(g(x_o)) = \mathcal{R}.$$

Since $J_f(x_o) < 0$, it follows that det A > 0. We identify \mathcal{R} with the linear transformation $x \mapsto \mathcal{R}x$. Since x_o is a p-good point for g and $g(x_o)$ is a q-good point for g^{-1} , Lemma 21 implies that

$$\lim_{r \to 0^+} \|g_r - \mathcal{R}\|_{W^{1,p}(\mathbb{B}^n,\mathbb{R}^n)} = \lim_{r \to 0^+} \|(g^{-1})_r - \mathcal{R}\|_{W^{1,q}(\mathbb{B}^n,\mathbb{R}^n)} = 0,$$

where $(g^{-1})_r$ is the blow up of g^{-1} at $g(x_o)$. It remains to observe that $(g^{-1})_r = (g_r)^{-1}$, which easily follows from the definitions.

2.3. Fubini's theorem for Sobolev spaces. The fact that almost all slices f_x , $f_{i,x}$ in the lemma below belong to the Sobolev space $W^{1,p}$ is well known. It is often called Fubini's theorem for Sobolev spaces. On the other hand, facts (8) and (9) are not so well known.

Lemma 26. Let $f, f_i \in W^{1,p}((0,1)^n)$, $1 \le p < \infty$, with $f_i \to f$ as $i \to \infty$. Denote the points of the cube $(0,1)^n$ by

$$(x,y) \in (0,1)^{n-\ell} \times (0,1)^{\ell} = (0,1)^n$$

and define

$$f_x(y) = f(x, y), \quad f_{i,x}(y) = f_i(x, y).$$

Then for almost all $x \in (0,1)^{n-\ell}$ we have $f_x, f_{i,x} \in W^{1,p}((0,1)^{\ell})$ and there is a subsequence $f_{i,j}$ such that for almost all $x \in (0,1)^{n-\ell}$

(8)
$$\lim_{j \to \infty} \|f_{i_j,x} - f_x\|_{W^{1,p}((0,1)^{\ell})} = 0.$$

Moreover, for any $\varepsilon > 0$, there is a compact set $K \subset (0,1)^{n-\ell}$ such that $|(0,1)^{n-\ell} \setminus K| < \varepsilon$ and

(9)
$$\lim_{j \to \infty} \sup_{x \in K} \|f_{i_j,x} - f_x\|_{W^{1,p}((0,1)^{\ell})} = 0.$$

Proof. The fact that $f_x \in W^{1,p}((0,1)^\ell)$ for almost all $x \in (0,1)^{n-\ell}$ is an easy consequence of the classical Fubini theorem applied to a sequence of smooth functions approximating f in $W^{1,p}((0,1)^n)$. We leave details to the reader. Similarly, we prove that for every $i \in \mathbb{N}$, $f_{i,x} \in W^{1,p}((0,1)^\ell)$ for almost all $x \in (0,1)^{n-\ell}$. Since we have countably many functions f, f_i , $i \in \mathbb{N}$, there is a set $A \subset (0,1)^{n-\ell}$ of full measure such that $f_x, f_{i,x} \in W^{1,p}$ for all $x \in A$ and all $i \in \mathbb{N}$. We have

$$||f_{i} - f||_{W^{1,p}((0,1)^{n})}^{p} = \int_{(0,1)^{n-\ell}} \int_{(0,1)^{\ell}} |f_{i}(x,y) - f(x,y)|^{p} + |Df_{i}(x,y) - Df(x,y)|^{p} dy dx$$

$$= \int_{(0,1)^{n-\ell}} ||f_{i,x} - f_{x}||_{W^{1,p}((0,1)^{\ell})}^{p} dx \to 0 \quad \text{as } i \to \infty.$$

In other words, $F_i \to 0$ in $L^1((0,1)^{n-\ell})$. Therefore, there is a subsequence F_{i_j} such that $F_{i_j}(x) \to 0$ for almost all $x \in (0,1)^{n-\ell}$, which is (8). Moreover, according to Egoroff's theorem [10, Theorem 1.16], for any $\varepsilon > 0$ there is a compact set $K \subset (0,1)^{n-\ell}$ such that $|(0,1)^{n-\ell} \setminus K| < \varepsilon$ and $F_{i_j} \to 0$ uniformly on K, which is (9).

The above result allows for a lot of flexibility and instead of applying Fubini's theorem to the products of cubes, we can apply it for example to $\mathbb{B}^{n-\ell} \times \mathbb{S}^{\ell}$, as described in the next result.

Lemma 27. Let $f_r, f_0 \in W^{1,p}(\mathbb{B}^{n-\ell} \times \mathbb{S}^\ell, \mathbb{R}^n)$, $1 \leq p < \infty$, $0 < r < r_0$, be a family of mappings such that $f_r \to f_0$ in $W^{1,p}$ as $r \to 0^+$. If $r_i \searrow 0$ is a sequence decreasing to zero, then there is a subsequence r_{i_j} such that $f_j = f_{r_{i_j}}$ satisfies:

For almost all $x \in \mathbb{B}^{n-\ell}$

$$f_j|_{\{x\}\times\mathbb{S}^\ell}\to f_0|_{\{x\}\times\mathbb{S}^\ell}\quad in\ W^{1,p}(\{x\}\times\mathbb{S}^\ell)\ as\ j\to\infty,$$

and for any $\varepsilon > 0$ there is a compact subset of the unit ball $K \subset \mathbb{B}^{n-\ell}$ such that $|\mathbb{B}^{n-\ell} \setminus K| < \varepsilon$ and

$$\sup_{x \in K} \|f_j - f_0\|_{W^{1,p}(\{x\} \times \mathbb{S}^\ell)} \to 0 \quad as \ j \to \infty.$$

2.4. Traces and extensions. The following lemma is known, but not very well known.

Proposition 28. For $n \ge 1$ and p > 1, there is a bounded linear extension operator

$$E: W^{1,p}(\mathbb{R}^n) \to W^{1,q} \cap C^{\infty}(\mathbb{R}^{n+1}_+), \quad where \ q = \frac{(n+1)p}{n}.$$

In other words, $W^{1,p}(\mathbb{R}^n)$ continuously embeds into the trace space $W^{1-\frac{1}{q},q}(\mathbb{R}^n)$ of $W^{1,q}(\mathbb{R}^{n+1}_+)$.

Corollary 29. For $n \geq 2$, there is a bounded linear extension operator

$$E: W^{1,n}(\mathbb{R}^n) \to W^{1,n+1} \cap C^{\infty}(\mathbb{R}^{n+1}_+).$$

In other words, $W^{1,n}(\mathbb{R}^n)$ continuously embeds into the trace space $W^{1-\frac{1}{n+1},n+1}(\mathbb{R}^n)$ of $W^{1,n+1}(\mathbb{R}^{n+1}_+)$.

Remark 30. Corollary 29 fails when n = 1, see [2], [29, Exercise 14.36] and [39, Proposition 4].

Remark 31. Proposition 28 was proved in [2]. It also follows from Theorem 14.32, Remark 14.35 and Proposition 14.40 in [29] (first edition). Corollary 29 was also proved in [15, Lemma 14] as a consequence of Theorem 2.5.6, Theorem 2.7.1, Proposition 2.3.2.2(8), Theorem 2.5.7 and 2.5.7(9) in [41]. All of the arguments listed here are difficult. Below we present an elementary and unpublished proof of Proposition 28 due to Jan Malý.

Proof (due to Jan Malý). In the proof we need the following result.

Lemma 32. If $F \in L^1_{loc}(\mathbb{R}^{n+1}_+)$, $f \in L^p(\mathbb{R}^n)$, p > 1, $n \ge 1$, and

(10)
$$|F(x,t)| \le C \int_{\mathbb{B}^n(x,t)} |f(y)| \, dy \quad for \ (x,t) \in \mathbb{R}^n \times (0,\infty) = \mathbb{R}^{n+1}_+,$$

then

$$||F||_{L^q(\mathbb{R}^{n+1}_+)} \le C||f||_{L^p(\mathbb{R}^n)}, \quad where \ q = \frac{(n+1)p}{n}.$$

Proof. Since the right-hand side of (10) is bounded, up to the constant C, by the maximal function $\mathcal{M}f(x)$, we have

(11)
$$\int_{0}^{r} |F(x,t)|^{q} dt \leq Cr(\mathcal{M}f)^{q}(x) \quad \text{for } r > 0.$$

On the other hand, the inequality

$$|F(x,t)| \le C \left(\int_{\mathbb{B}^n(x,t)} |f(y)|^p \, dy \right)^{1/p} \le Ct^{-\frac{n}{p}} ||f||_p$$

yields

(12)
$$\int_{r}^{\infty} |F(x,t)|^{q} dt \le C \|f\|_{p}^{q} r^{1-\frac{nq}{p}} = C \|f\|_{p}^{q} r^{-n} \quad \text{for } r > 0.$$

We used here the fact that

$$\frac{nq}{p} = n + 1 > 1.$$

Now if we choose r > 0 such that

$$r(\mathcal{M}f)^q(x) = ||f||_p^q r^{-n}, \quad \text{i.e.,} \quad r = \left(\frac{||f||_p}{\mathcal{M}f(x)}\right)^{\frac{p}{n}},$$

then the right-hand sides of (11) and (12) are equal to $C||f||_p^{p/n}(\mathcal{M}f(x))^p$. Adding integrals at (11) and (12) yields

$$\int_{0}^{\infty} |F(x,t)|^{q} dt \le C \|f\|_{p}^{p/n} (\mathcal{M}f(x))^{p}$$

and Fubini's theorem along with boundedness of the maximal operator in L^p , p > 1, give

$$\int_{\mathbb{R}^{n+1}_+} |F(x,t)|^q dt dx \le C \|f\|_p^{p/n} \|f\|_p^p = C \|f\|_p^q.$$

Now we can complete the proof of Proposition 28.

Let
$$\varphi \in C_0^{\infty}(\mathbb{B}^n)$$
, $\varphi \geq 0$, $\int_{\mathbb{B}^n} \varphi(x) dx = 1$, and $\varphi_t(x) = t^{-n} \varphi(x/t)$.

For $f \in L^1_{loc}(\mathbb{R}^n)$ we define

$$(Ef)(x,t) = (f * \varphi_t)(x) = \int_{\mathbb{R}^n} f(x - ty)\varphi(y) \, dy, \quad (x,t) \in \mathbb{R}^{n+1}_+.$$

Clearly, properties of the convolution guarantee that $Ef \in C^{\infty}(\mathbb{R}^{n+1}_+)$ and a simple change of variables yields

$$|(Ef)(x,t)| = \left| \int_{\mathbb{B}^n(x,t)} f(y)\varphi_t(x-y) \, dy \right| \le C \int_{\mathbb{B}^n(x,t)} |f(y)| \, dy,$$

where the constant C depends on φ only. Therefore, if $f \in L^p(\mathbb{R}^n)$, Lemma 32 gives the estimate (13) $||Ef||_{L^q(\mathbb{R}^{n+1})} \leq C||f||_{L^p(\mathbb{R}^n)}.$

Assume now that $u \in W^{1,p}(\mathbb{R}^n)$. We have

(14)
$$|\nabla(Eu)(x,t)| = |\nabla_{(x,t)}(Eu)(x,t)| \le C \int_{\mathbb{B}^n(x,t)} |\nabla u(y)| \, dy.$$

Indeed,

$$|\nabla_x(Eu)(x,t)| = \left| \int_{\mathbb{B}^n} (\nabla u)(x-ty)\varphi(y) \, dy \right| \le C \int_{\mathbb{B}^n(x,t)} |\nabla u(y)| \, dy,$$

and

$$\left|\frac{\partial}{\partial t}(Eu)(x,t)\right| = \left|\int\limits_{\mathbb{B}^n} (\nabla u)(x-ty)\cdot (-y)\varphi(y)\,dy\right| \leq C\int\limits_{\mathbb{R}^n(x,t)} |\nabla u(y)|\,dy,$$

because $|-y| \leq 1$ for $y \in \mathbb{B}^n$. Now (14) and Lemma 32 imply that

$$\|\nabla(Eu)\|_{L^q(\mathbb{R}^{n+1}_+)} \le C\|\nabla u\|_{L^p(\mathbb{R}^n)}.$$

This estimate and (13) for f = u complete the proof.

A localized version of Corollary 29 gives the following result:

Lemma 33. Let $n \geq 2$. Then there is an extension operator

$$E: W^{1,n}(\partial(\mathbb{S}^n \times [0,1])) \to W^{1,n+1} \cap C^{\infty}(\mathbb{S}^n \times (0,1))$$

such that if $h = f_0$ on $\mathbb{S}^n \times \{0\}$ and $h = f_1$ on $\mathbb{S}^n \times \{1\}$, then traces (in the Sobolev sense) satisfy $\operatorname{Tr} Eh(\cdot,0) = f_0$, $\operatorname{Tr} Eh(\cdot,1) = f_1$, and

$$||Eh||_{W^{1,n+1}(\mathbb{S}^n\times(0,1))} \le C(n) \left(||f_0||_{W^{1,n}(\mathbb{S}^n)} + ||f_1||_{W^{1,n}(\mathbb{S}^n)} \right).$$

The next result seems new. Its proof uses some ideas from the proof of [19, Proposition 3.3].

Proposition 34. Let $n \geq 2$, $m \geq n+1$ and let $f \in W^{1,n} \cap C^0(\mathbb{S}^n, \mathbb{R}^m)$, $g \in C^{\infty}(\mathbb{S}^n, \mathbb{R}^m)$. Then there is a function $H \in C^0(\mathbb{S}^n \times [0,1], \mathbb{R}^m) \cap C^{\infty}(\mathbb{S}^n \times [0,1])$ such that H(x,0) = g(x), H(x,1) = f(x) and

$$\mathcal{H}^{n+1}(H(\mathbb{S}^n \times [0,1))) \le C \|f - g\|_{W^{1,n}(\mathbb{S}^n)} (\|f - g\|_{W^{1,n}(\mathbb{S}^n)} + \|Dg\|_{L^{n+1}(\mathbb{S}^n)})^n,$$

where the constant C depends on n and m only and \mathcal{H}^{n+1} denotes the Hausdorff measure.

Remark 35. Fix $g \in C^{\infty}(\mathbb{S}^n, \mathbb{R}^m)$. Let $f, f_k \in W^{1,n} \cap C^0(\mathbb{S}^n, \mathbb{R}^m)$ and let H, H_k be the homotopies for f and f_k constructed as in Proposition 34. If $f_k \to f$ uniformly on \mathbb{S}^n as $k \to \infty$ (we do not require convergence in the Sobolev norm), then $H_k \to H$ uniformly on $\mathbb{S}^n \times [0,1]$. Indeed, the homotopies are defined by the formula (16) and the extension operator Eh is defined through the averaging and multiplication by a cut-off function, and such a construction is continuous in the uniform norm.

Remark 36. The proposition has a clear geometric interpretation. The image of a continuous mapping $f \in W^{1,n} \cap C^0(\mathbb{S}^n, \mathbb{R}^m)$ can be very large. It can even fill a ball in \mathbb{R}^m . However, if f is very close in the $W^{1,n}$ norm to a fixed smooth map g, then there is a homotopy between f and g such that the \mathcal{H}^{n+1} -volume of the image of the homotopy, except the endpoint where the homotopy equals f, is very small. We hope that this result might be useful for other applications.

In the proof we will need the following estimate.

Lemma 37. If A and B are two $n \times n$ -matrices, then

(15)
$$|\det(A+B)| \le C(n)(|A|^n + |B|^n),$$

where $|\cdot|$ stands for the Hilbert-Schmidt norm of a matrix.

Proof. Since the determinant is continuous and homogeneous of order n, we immediately get that

$$|\det A| = |A|^n \left| \det \left(\frac{A}{|A|} \right) \right| \le \Lambda |A|^n$$
, where $\Lambda = \sup\{|\det M| : |M| \le 1\}$.

Then (15) follows from the triangle inequality and the standard convexity estimate $(a+b)^n \le 2^{n-1}(a^n+b^n)$ for $a,b\ge 0$.

Proof of Proposition 34. Let

$$h(x) = \begin{cases} f - g & \text{on } \mathbb{S}^n \times \{1\}, \\ 0 & \text{on } \mathbb{S}^n \times \{0\}, \end{cases}$$

and define

(16)
$$H(x,t) = (Eh)(x,t) + g(x) \text{ for } (x,t) \in \mathbb{S}^n \times [0,1],$$

where E is the extension operator from Lemma 33. Since the extension Eh is continuous (smooth) up to the boundary if the function on the boundary is continuous (smooth), we conclude that $H \in C^0(\mathbb{S}^n \times [0,1], \mathbb{R}^m) \cap C^\infty(\mathbb{S}^n \times [0,1])$ and H(x,0) = g(x), H(x,1) = f(x).

According to the area formula [10] we have that

(17)
$$\mathcal{H}^{n+1}(H(\mathbb{S}^n \times [0,1))) = \int_0^1 \int_{\mathbb{S}^n} J_H(x,t) \, d\sigma(x) \, dt,$$

where

$$J_H(x,t) = \sqrt{\det\left((DH)^T(DH)\right)}.$$

Denote the derivative in $\mathbb{S}^n \times [0,1]$ by $D=(D_x,\partial_t)$, where D_x is the derivative on \mathbb{S}^n . Then

$$DH(x,t) = (D_x(Eh) + D_x g, \partial_t(Eh)).$$

The Cauchy-Binet formula [10, Sect. 3.2.1, Theorem 4], the Laplace expansion along the last column $(\partial_t(Eh))_I$, and Lemma 37 yield the following estimate, where the sum is taken over all $I = (i_1, \ldots, i_{n+1}), 1 \le i_1 < \ldots < i_{n+1} \le m$:

$$J_{H} = \sqrt{\sum_{I} \left(\det \left((D_{x}(Eh))_{I} + (D_{x}g)_{I}, (\partial_{t}(Eh))_{I} \right) \right)^{2}} \le C|\partial_{t}(Eh)|(|D_{x}(Eh)|^{n} + |D_{x}g|^{n}).$$

Therefore (17) gives

$$\mathcal{H}^{n+1}(H(\mathbb{S}^{n} \times [0,1))) \leq C \|Eh\|_{W^{1,n+1}(\mathbb{S}^{n} \times (0,1))}^{n+1}$$

$$+ C \left(\int_{0}^{1} \int_{\mathbb{S}^{n}} |D_{x}g|^{n+1} \right)^{\frac{n}{n+1}} \left(\int_{0}^{1} \int_{\mathbb{S}^{n}} |\partial_{t}(Eh)|^{n+1} \right)^{\frac{1}{n+1}}$$

$$\leq C \|f - g\|_{W^{1,n}(\mathbb{S}^{n})}^{n+1} + \|Dg\|_{L^{n+1}(\mathbb{S}^{n})}^{n} \|f - g\|_{W^{1,n}(\mathbb{S}^{n})}.$$

2.5. The local degree and the linking number. To keep the paper self contained, we present here a short introduction to the theory of local degree. The standard references are the books of Fonseca and Gangbo [11] and Outerelo and Ruiz [35]. Then, at the end of the section, we discuss the linking number.

Throughout this section, we assume that $\Omega \subset \mathbb{R}^n$ is a domain. By $C^1(\overline{\Omega}, \mathbb{R}^n)$ we denote the set of all these mappings from $\overline{\Omega}$ to \mathbb{R}^n which admit an extension to a C^1 mapping on some open $U \supset \overline{\Omega}$.

We begin by defining the local degree for C^1 mappings at their regular values.

Definition 38 ([11, Definition 1.2]). Assume $\phi \in C^1(\overline{\Omega}, \mathbb{R}^n)$. Let $p \in \mathbb{R}^n$ be a regular value of ϕ and $p \notin \phi(\partial\Omega)$. We define the local degree of ϕ at p with respect to Ω as

$$\deg(\phi, \Omega, p) = \sum_{x \in \phi^{-1}(p)} \operatorname{sgn} J_{\phi}(x).$$

Moreover, if $p \notin \phi(\overline{\Omega})$, we set $\deg(\phi, \Omega, p) = 0$.

It turns out that $deg(\phi, \Omega, \cdot)$ is constant on connected components of $\mathbb{R}^n \setminus \phi(\partial\Omega)$.

Proposition 39 ([11, Proposition 1.8]). Let V be a connected component of $\mathbb{R}^n \setminus \phi(\partial\Omega)$ and assume $p_1, p_2 \in V$ are regular values of ϕ . Then $\deg(\phi, \Omega, p_1) = \deg(\phi, \Omega, p_2)$.

Proposition 39 allows us to define the local degree also at critical points of ϕ , as long as they are not in the image of $\partial\Omega$.

Definition 40 ([11, Definition 1.9]). Assume Ω and ϕ are as in Definition 38 and let $p \in \phi(\Omega) \setminus \phi(\partial\Omega)$ be a critical value of ϕ . Let p_1 be any regular value of ϕ such that p and p_1 lie in the same connected component of $\mathbb{R}^n \setminus \partial\Omega$. We set $\deg(\phi, \Omega, p) = \deg(\phi, \Omega, p_1)$.

Note that by Sard's Lemma such p_1 always exists; Proposition 39 shows that the above definition does not depend on the choice of p_1 .

The local degree is a C^1 homotopy invariant:

Proposition 41 ([11, Theorem 1.12]). If $H : \overline{\Omega} \times [0,1] \to \mathbb{R}^n$ is a C^1 mapping such that $H(\cdot,0) = \phi(\cdot)$, $H(\cdot,1) = \psi(\cdot)$ and $p \notin H(\partial\Omega \times [0,1])$, then $\deg(\phi,\Omega,p) = \deg(\psi,\Omega,p)$.

Proposition 41 allows us to extend the notion of local degree to continuous mappings:

Definition 42 ([11, Definition 1.18]). If $\phi \in C(\overline{\Omega}, \mathbb{R}^n)$ and $p \notin \phi(\partial\Omega)$, we set $\deg(\phi, \Omega, p) = \deg(\psi, \Omega, p)$, where ψ is any mapping in $C^1(\overline{\Omega}, \mathbb{R}^n)$ such that $\sup_{x \in \overline{\Omega}} |\phi(x) - \psi(x)| < \operatorname{dist}(p, \phi(\partial\Omega))$.

One easily checks that the above definition is independent on the choice of ψ : if we choose $\psi_1, \psi_2 \in C^1(\overline{\Omega}, \mathbb{R}^n)$ satisfying $\sup_{x \in \overline{\Omega}} |\phi(x) - \psi_i(x)| < \operatorname{dist}(p, \phi(\partial \Omega)), i = 1, 2$, then $p \notin H(\partial \Omega \times [0, 1])$, where H is the standard homotopy between ψ_1 and ψ_2 , $H(x, t) = (1 - t)\psi_1(x) + t\psi_2(x)$, and thus, by Proposition 41, the local degrees of ψ_1 and ψ_2 at p are the same.

Remark 43. In fact, a standard smoothing argument shows that the local degree for continuous maps at a point p, given by the above definition, is a homotopy invariant (without the C^1 assumption), as long as the homotopy H does not map any points of $\partial\Omega$ into p, i.e. $p \notin H(\partial\Omega \times [0,1])$.

We shall need the following deep facts on the local degree.

Proposition 44 (Multiplication theorem, [11, Theorem 2.10]). Assume $\Omega \subset \mathbb{R}^n$ is a domain, $\phi \in C(\overline{\Omega}, \mathbb{R}^n)$, $V \subset \mathbb{R}^n$ is a domain containing $\phi(\overline{\Omega})$ and $\psi \in C(\overline{V}, \mathbb{R}^n)$. Let $D = V \setminus \phi(\partial \Omega)$ and denote by D_i the connected components of D. Then for any $p \notin (\psi \circ \phi)(\partial \Omega) \cup \psi(\partial V)$ we have $p \notin \psi(\partial D_i)$, and the following formula holds:

(18)
$$\deg(\psi \circ \phi, \Omega, p) = \sum_{i} \deg(\psi, D_i, p) \deg(\phi, \Omega, q_i)$$

for arbitrary $q_i \in D_i$.

Proposition 45. If ψ and ϕ are as in Proposition 44 and additionally ψ and ϕ are homeomorphisms, then

a) for any $p \notin (\psi \circ \phi)(\partial \Omega) \cup \psi(\partial V)$

(19)
$$\deg(\psi \circ \phi, \Omega, p) = \deg(\psi, \phi(\Omega), p) \deg(\phi, \Omega, \psi^{-1}(p)),$$

b) for every $q \in \phi(\Omega)$ we have either

$$\deg(\phi, \Omega, q) = \deg(\phi^{-1}, \phi(\Omega), \phi^{-1}(q)) = 1$$

or

$$deg(\phi, \Omega, q) = deg(\phi^{-1}, \phi(\Omega), \phi^{-1}(q)) = -1.$$

Proof. If $p \notin \psi(\phi(\Omega))$, then both sides of (19) are zero. Assume thus $p \in \psi(\phi(\Omega))$.

By Brouwer's Invariance of Domain Theorem [11, Theorem 3.30], $\phi(\partial\Omega) = \partial\phi(\Omega)$ and there is only one component D_i of $V \setminus \phi(\partial\Omega)$ such that $\psi(D_i)$ contains p, namely $D_i = \phi(\Omega)$, thus all the terms of the sum in (18) are zero, except the one with $D_i = \phi(\Omega)$. Also, we may choose $q_i = \psi^{-1}(p)$, which gives (19).

To prove b), take any domain Ω' such that $q \in \phi(\Omega')$ and $\overline{\Omega}' \subset \Omega$. Applying a) to $\phi|_{\overline{\Omega}'} : \overline{\Omega}' \to \mathbb{R}^n$ and $\psi = \phi^{-1} : \phi(\overline{\Omega}) \to \mathbb{R}^n$, we see that

$$1 = \deg(\operatorname{Id}, \Omega', \phi^{-1}(q)) = \deg(\phi^{-1}, \phi(\Omega'), \phi^{-1}(q)) \deg(\phi, \Omega', q),$$

thus $\deg(\phi, \Omega', q) = \deg(\phi^{-1}, \phi(\Omega'), \phi^{-1}(q)) = \pm 1$. What remains to prove is that $\deg(\phi, \Omega', q) = \deg(\phi, \Omega, q)$ and $\deg(\phi^{-1}, \phi(\Omega'), \phi^{-1}(q)) = \deg(\phi^{-1}, \phi(\Omega), \phi^{-1}(q))$. Let $\zeta \in C^1(\overline{\Omega}, \mathbb{R}^n)$ be such that

$$\sup_{\overline{\Omega}} |\phi - \zeta| < \operatorname{dist}(q, \phi(\partial \Omega')) < \operatorname{dist}(q, \phi(\partial \Omega)).$$

Then, by Definition 42, $\deg(\phi, \Omega, q) = \deg(\zeta, \Omega, q)$ and $\deg(\phi, \Omega', q) = \deg(\zeta, \Omega', q)$. However, for any $z \in \Omega \setminus \Omega'$,

$$|\zeta(z) - q| \ge |q - \phi(z)| - |\phi(z) - \zeta(z)| > \operatorname{dist}(q, \phi(\partial \Omega')) - \operatorname{dist}(q, \phi(\partial \Omega')) = 0,$$

thus $q \notin \zeta(\Omega \setminus \Omega')$ and applying Definition 38 we see that $\deg(\zeta, \Omega, q) = \deg(\zeta, \Omega', q)$. Calculations for ϕ^{-1} follow the same steps. This concludes the proof of b).

A homeomorphism $h: \Omega \to \mathbb{R}^n$ with degree +1 at all its values is called *sense*- or *orientation-preserving*; if the degree is -1 at all values, we call it *sense-reversing*.

As an immediate corollary of Proposition 45, b) and Proposition 39, we obtain that every homeomorphism of a domain is either sense-preserving or sense-reversing.

Corollary 46 (c.f. [11, Theorem 3.35]). Assume $\Omega \subset \mathbb{R}^n$ is open and connected and $h : \overline{\Omega} \to \mathbb{R}^n$ is a homeomorphism. Then either for every $p \in h(\Omega)$ we have $\deg(h, \Omega, p) = 1$ or for every $p \in h(\Omega)$ we have $\deg(h, \Omega, p) = -1$.

Corollary 47. If $h: \Omega \to \mathbb{R}^n$ is a sense-preserving (reversing) homeomorphism and $\Omega' \subset \Omega$, then $h|_{\Omega'}: \Omega' \to \mathbb{R}^n$ is also sense-preserving (reversing).

This is a corollary from the proof of Proposition 45, where we showed in the last step that $deg(\phi, \Omega, q) = deg(\phi, \Omega', q)$.

The terms sense-preserving and sense-reversing are justified by the following fact.

Proposition 48. Assume $\Omega \subset \mathbb{R}^n$ is a domain, $\phi \in C(\overline{\Omega}, \mathbb{R}^n)$ and $h : \overline{U} \to \mathbb{R}^n$ is a homeomorphism of a domain $U \supset \phi(\overline{\Omega})$.

• If h is sense-preserving, then for every $p \in U \setminus \phi(\partial\Omega)$ we have

$$deg(\phi, \Omega, p) = deg(h \circ \phi, \Omega, h(p)).$$

• If h is sense-reversing, then for every $p \in U \setminus \phi(\partial\Omega)$ we have

$$\deg(\phi, \Omega, p) = -\deg(h \circ \phi, \Omega, h(p)).$$

Proof. Assume h is sense-preserving and apply Proposition 44 to $h \circ \phi$: let $D = U \setminus \phi(\partial\Omega)$ and denote by D_i the connected components of D. Then $p \in D_j$ for exactly one $j \in \mathbb{N}$ and we have for any $q_i \in D_i$, $i \neq j$,

$$\begin{split} \deg(h \circ \phi, \Omega, h(p)) &= \deg(h, D_j, h(p)) \deg(\phi, \Omega, p) + \sum_{i \neq j} \deg(h, D_i, h(p)) \deg(\phi, \Omega, q_i) \\ &= \deg(\phi, \Omega, p), \end{split}$$

because $deg(h, D_j, h(p)) = +1$, while $deg(h, D_i, h(p)) = 0$ for $i \neq j$, since $h(p) \notin h(D_i)$.

The case of sense-reversing h is proved in exactly the same way.

Definition 49. If M and N are compact, connected, oriented, smooth n-dimensional manifolds without boundary and $\phi: M \to N$ is C^{∞} smooth, then we define

$$\deg(\phi, y) = \sum_{x \in \phi^{-1}(y)} \operatorname{sgn} J_{\phi}(x),$$

where $y \in N$ is a regular value of ϕ , and $J_{\phi}(x)$ is the determinant of the derivative $D\phi(x)$: $T_xM \to T_{\phi(x)}N$. It turns out that $\deg(\phi, y)$ does not depend on the choice of a regular value y. The common value of all $\deg(\phi, y)$ is denoted by $\deg \phi$ and is called the *degree of* ϕ .

One can prove that homotopic mappings have equal degrees. Since every continuous mapping is homotopic to a smooth one, one can extend the notion of degree to the class of all continuous mappings $\phi: M \to N$. For more details, see [32].

The following result relates the local degree of ϕ with the topological degree of ϕ restricted to the boundary.

Proposition 50 ([35, Proposition IV.4.6]). Let $\partial\Omega$ be a connected, compact and smooth manifold oriented by the outward normal vector ([35, Section II.7.7]) and assume $\phi \in C(\overline{\Omega}, \mathbb{R}^n)$ is such that $\phi|_{\partial\Omega}: \partial\Omega \to \mathbb{S}^{n-1}$. Then

$$\deg \phi|_{\partial\Omega} = \deg(\phi, \Omega, 0).$$

Our main purpose for introducing the local degree is to justify the properties of yet another invariant, the linking number.

The linking number is an important and well studied invariant in the theory of knots. It was introduced by Gauss in a short note of 1833 [13] (see also [38] for a nice historical account and modern interpretation): if γ_1 , γ_2 are two parameterized, non-intersecting, oriented curves in \mathbb{R}^3 , γ_1 , $\gamma_2 : \mathbb{S}^1 \to \mathbb{R}^3$, then the linking number $\ell(\gamma_1, \gamma_2)$ is defined (in modern notation) as the integral

(20)
$$\ell(\gamma_1, \gamma_2) = \frac{1}{4\pi} \int_{\text{SI} \times \text{SI}} \frac{\det(\gamma_1'(s), \gamma_2'(t), \gamma_1(s) - \gamma_2(t))}{|\gamma_1(s) - \gamma_2(t)|^3} \, ds \, dt.$$

The Gauss map for γ_1 , γ_2 is defined as

(21)
$$\mathbb{S}^1 \times \mathbb{S}^1 \ni (s,t) \xrightarrow{\psi} \frac{\gamma_1(s) - \gamma_2(t)}{|\gamma_1(s) - \gamma_2(t)|} \in \mathbb{S}^2.$$

It turns out (see e.g. [38]) that the linking number given by (20) is equal to the degree of the Gauss map: $\ell(\gamma_1, \gamma_2) = \deg(\psi)$. In colloquial terms, the linking number tells us how many times (counting directions) one curve winds around the other.

These definitions have been later generalized to pairs of non-intersecting manifolds in higher dimensions (see the paper by M. Kervaire, [28], who attributes the idea to A. Shapiro). Here, we

shall restrict ourselves to the case when the manifolds in question are two oriented spheres \mathbb{S}^k , \mathbb{S}^l , where k+l=n-1, continuously embedded in \mathbb{R}^n in such a way that the embedded spheres do not intersect. Then, if we denote the embeddings as before, in the case of curves, by $\gamma_1 : \mathbb{S}^k \to \mathbb{R}^n$ and $\gamma_2 : \mathbb{S}^l \to \mathbb{R}^n$, we can define the Gauss map of the two embeddings by the formula (21), (this time, however, $\psi : \mathbb{S}^k \times \mathbb{S}^l \to \mathbb{S}^{n-1}$), and the linking number of the two embedded (oriented) spheres, identified here with the embedding maps γ_1 and γ_2 , is defined again as $\ell(\gamma_1, \gamma_2) = \deg \psi$.

More generally, we can define, by the same formula, the linking number $\ell(\gamma_1, \gamma_2)$ between any two continuous maps $\gamma_1 : \mathbb{S}^k \to \mathbb{R}^n$ and $\gamma_2 : \mathbb{S}^l \to \mathbb{R}^n$, k+l=n-1, provided $\gamma_1(\mathbb{S}^k) \cap \gamma_2(\mathbb{S}^l) = \emptyset$.

The linking number is a homotopy invariant in the sense that if $\gamma_1, \tilde{\gamma}_1 : \mathbb{S}^k \to \mathbb{R}^n \setminus \gamma_2(\mathbb{S}^l)$ are two mappings of \mathbb{S}^k , which are homotopic in $\mathbb{R}^n \setminus \gamma_2(\mathbb{S}^l)$ (i.e. neither the two images $\gamma_1(\mathbb{S}^k)$, $\tilde{\gamma}_1(\mathbb{S}^k)$, nor the image of the homotopy between them intersects $\gamma_2(\mathbb{S}^l)$), then $\ell(\gamma_1, \gamma_2) = \ell(\tilde{\gamma}_1, \gamma_2)$. Indeed, if $\Gamma : \mathbb{S}^k \times [0, 1] \to \mathbb{R}^n$ is the homotopy between γ_1 and $\tilde{\gamma}_1$, the image of which is disjoint with the image of γ_2 , then $\Psi : \mathbb{S}^k \times \mathbb{S}^l \times [0, 1] \to \mathbb{S}^n$ given by the formula

$$(s,t,r) \xrightarrow{\Psi} \frac{\Gamma(s,r) - \gamma_2(t)}{|\Gamma(s,r) - \gamma_2(t)|}$$

is a homotopy between

$$\psi(s,t) = \frac{\gamma_1(s,t) - \gamma_2(s,t)}{|\gamma_1(s,t) - \gamma_2(s,t)|} \quad \text{and} \quad \tilde{\psi}(s,t) = \frac{\tilde{\gamma_1}(s,t) - \gamma_2(s,t)}{|\tilde{\gamma_1}(s,t) - \gamma_2(s,t)|},$$

and thus the linking numbers $\ell(\gamma_1, \gamma_2) = \deg \psi$ and $\ell(\tilde{\gamma_1}, \gamma_2) = \deg \tilde{\psi}$ are the same.

The following known invariance result is very hard to find in the literature, thus we present the proof here (essentially the same argument, in a less general situation, was given in [24]).

Proposition 51. Assume k + l = n - 1 and let $\gamma_1 : \mathbb{S}^k \to \mathbb{R}^n$ and $\gamma_2 : \mathbb{S}^l \to \mathbb{R}^n$ be continuous maps such that $\gamma_1(\mathbb{S}^k) \cap \gamma_2(\mathbb{S}^l) = \emptyset$. Let $h : \mathbb{R}^n \to \mathbb{R}^n$ be a homeomorphism. Then

- if h is sense-preserving, then $\ell(h \circ \gamma_1, h \circ \gamma_2) = \ell(\gamma_1, \gamma_2)$;
- if h is sense-reversing, then $\ell(h \circ \gamma_1, h \circ \gamma_2) = -\ell(\gamma_1, \gamma_2)$.

Proof. We begin by fixing some notation. Let $\bar{\gamma}_1 \in C(\mathbb{B}^{k+1}, \mathbb{R}^n)$ be any extension of γ_1 , i.e. $\bar{\gamma}_1|_{\partial\mathbb{B}^{k+1}} = \gamma_1$. Let $A = \mathbb{B}^{k+1} \times \mathbb{S}^l$ be a full torus, embedded smoothly in \mathbb{R}^n in a way that the orientation of the boundary of embedded A by the outward normal vector is consistent with the orientation of $\mathbb{S}^k \times \mathbb{S}^l$, and define $F: A \to \mathbb{R}^n$, $F(x, y) = \bar{\gamma}_1(x) - \gamma_2(y)$.

In the proof, we shall need a simple lemma, connecting the linking number with the degree of the non-normalized Gauss map F.

Lemma 52. With the above notation,

$$\ell(\gamma_1, \gamma_2) = \deg(F, A, 0).$$

Proof of Lemma 52. The claim would follow from Proposition 50, if F mapped ∂A to \mathbb{S}^{n-1} , not to \mathbb{R}^n (because $F = \gamma_1 - \gamma_2$ on ∂A). Note, however, that $F(\partial A) \subset \mathbb{R}^n \setminus \{0\}$. Set $d = \operatorname{dist}(\gamma_1(\mathbb{S}^k), \gamma_2(\mathbb{S}^l)) = \operatorname{dist}(F(\partial A), 0)$, and take $\phi : [0, \infty) \to [0, \infty)$ to be a smooth, positive function such that

$$\phi(s) = \begin{cases} 1 & \text{for } s \le d/2, \\ s & \text{for } s \ge d. \end{cases}$$

Then

$$h_t(z) = \frac{z}{\phi(t|z|)}, \quad t \in [0, 1],$$

is a homotopy connecting Id: $\mathbb{R}^n \to \mathbb{R}^n$ with a mapping which is identity on $\mathbb{B}^n(0, d/2)$ and a projection $z \mapsto z/|z|$ outside $\mathbb{B}^n(0, d)$. Obviously, $0 \notin (h_t \circ F)(\partial A)$ for any $t \in [0, 1]$, thus

$$\deg(F, A, 0) = \deg(h_0 \circ F, A, 0) = \deg(h_1 \circ F, A, 0).$$

However,

$$(h_1 \circ F)|_{\partial A}(x,y) = \frac{\gamma_1(x) - \gamma_2(y)}{|\gamma_1(x) - \gamma_2(y)|}$$

is the Gauss map whose degree, by definition, equals $\ell(\gamma_1, \gamma_2)$. Thus Proposition 50 yields

$$\ell(\gamma_1, \gamma_2) = \deg(h_1 \circ F)|_{\partial A} = \deg(h_1 \circ F, A, 0) = \deg(F, A, 0).$$

Assume now h is sense-preserving (the case of h sense-reversing is treated in the same way).

Let $G_t: A \times [0,1] \to \mathbb{R}^n$,

$$G_t(x,y) = h(t\bar{\gamma}_1(x)) - h(\gamma_2(y) - (1-t)\bar{\gamma}_1(x)).$$

Then $G_1(x,y) = (h \circ \bar{\gamma}_1)(x) - (h \circ \gamma_2)(y)$, thus by Lemma 52 applied to $h \circ \gamma_1$ and $h \circ \gamma_2$ in place of γ_1 and γ_2 ,

$$\ell(h \circ \gamma_1, h \circ \gamma_2) = \deg(G_1, A, 0).$$

We have

$$G_0(x,y) = h(0) - h(\gamma_2(y) - \bar{\gamma}_1(x)) = h(0) - h(-F(x,y)) = h(0) + (-\mathrm{Id}) \circ h \circ (-\mathrm{Id}) \circ F(x,y).$$

Note also that $G_t(x,y) = 0$ if and only if $\bar{\gamma}_1(x) = \gamma_2(y)$, which is not possible if $(x,y) \in \partial A$, thus $0 \notin G_t(\partial A)$ for any $t \in [0,1]$, and Remark 43 yields

$$\ell(h \circ \gamma_1, h \circ \gamma_2) = \deg(G_1, A, 0) = \deg(G_0, A, 0)$$

= \deg(h(0) + (-\text{Id}) \circ h \circ (-\text{Id}) \circ F, A, 0)
= \deg(F, A, 0) = \ell(\gamma_1, \gamma_2),

because $z \mapsto h(0) + (-\mathrm{Id}) \circ h \circ (-\mathrm{Id})(z)$ is a sense-preserving homeomorphism of \mathbb{R}^n which maps 0 to 0 and we can apply Proposition 48.

3. Proofs of Theorems 1 and 2

For n=1, the claim of Theorem 1 is obvious: a sense-preserving homeomorphism of an open subset of a real line is an increasing function, thus it is differentiable a.e. and its derivative is non-negative. In dimension n=2, every $W^{1,1}$ homeomorphism is again a.e. differentiable (see [31, 14]) and its weak Jacobian coincides with its classical one a.e. The sign of the latter reflects whether the homeomorphism preserves or reverses the local orientation, and thus for a sense preserving homeomorphism f we have $J_f \geq 0$ a.e.

Assume $n \geq 3$. To simplify the notation we shall write $\nu = n - 1 - \lfloor n/2 \rfloor$. We will argue by contradiction: assume that the set $\{J_f < 0\}$ has positive measure.

Pick a p-good point $x_o \in \{J_f < 0\}$ for f (in the sense of Definition 20) and consider the blow-up f_r of f at x_o :

$$f_r(x) = \frac{f(x_o + rx) - f(x_o)}{r}.$$

According to Lemma 21, $f_r \to f_0$ in $W^{1,p}(\mathbb{B}^n, \mathbb{R}^n)$ as $r \to 0$, where $f_0(x) = Df(x_0)x$. Note that f_0 is a linear, orientation reversing isomorphism.

We pick in \mathbb{B}^n two solid tori, i.e. smooth embeddings $\iota_1: \mathbb{B}^{[n/2]+1} \times \mathbb{S}^{\nu} \to \mathbb{B}^n$ and $\iota_2: \mathbb{B}^{\nu+1} \times \mathbb{S}^{[n/2]} \to \mathbb{B}^n$, such that for any $x \in \mathbb{B}^{[n/2]+1}$ and $y \in \mathbb{B}^{\nu+1}$ the embedded spheres $\mathbb{S}^{\nu}_x = \iota_1(\{x\} \times \mathbb{S}^{\nu})$ and $\mathbb{S}^{[n/2]}_y = \iota_2(\{y\} \times \mathbb{S}^{[n/2]})$ are linked, with linking number $\ell(\mathbb{S}^{\nu}_x, \mathbb{S}^{[n/2]}_y) = +1$ (see the next section for a particular construction).

By Lemma 27, we can choose a sequence $r_j \searrow 0$ and particular $x \in \mathbb{B}^{[n/2]+1}$ and $y \in \mathbb{B}^{\nu+1}$ such that $f_j = f_{r_j}$ converge to f_0 in $W^{1,p}$ on $\mathbb{S}^{\nu}_x \cup \mathbb{S}^{[n/2]}_y$. If n > 3, then $p > [n/2] \ge \nu$, and by the Sobolev-Morrey embedding theorem f_j converge to f_0 uniformly on $\mathbb{S}^{\nu}_x \cup \mathbb{S}^{[n/2]}_y$. If n = 3 and p = 1, the $W^{1,1}$ convergence of f_k on the sum of circles $\mathbb{S}^{\nu}_x \cup \mathbb{S}^{[n/2]}_y$ implies uniform convergence as well.

Since f_j converge to f_0 uniformly on $\mathbb{S}^{\nu}_x \cup \mathbb{S}^{[n/2]}_y$, f_j are homotopic to f_0 for j sufficiently large, and the image of each sphere in that homotopy does not intersect the image of the other (the image of \mathbb{S}^{ν}_x in that homotopy stays in a small tubular neighborhood of $f_0(\mathbb{S}^{\nu}_x)$, and, similarly, the image of $\mathbb{S}^{[n/2]}_y$ in that homotopy stays near $f_0(\mathbb{S}^{[n/2]}_y)$). Thus

$$\ell(f_j(\mathbb{S}_x^{\nu}), f_j(\mathbb{S}_y^{[n/2]})) = \ell(f_0(\mathbb{S}_x^{\nu}), f_0(\mathbb{S}_y^{[n/2]})) = -\ell(\mathbb{S}_x^{\nu}, \mathbb{S}_y^{[n/2]}) = -1,$$

since f_0 is a linear, orientation reversing homeomorphism.

However, each f_j , as a translation and rescaling of an orientation preserving homeomorphism f, is again an orientation preserving homeomorphism, thus Proposition 51 yields

$$\ell(f_j(\mathbb{S}_x^{\nu}), f_j(\mathbb{S}_y^{[n/2]})) = \ell(\mathbb{S}_x^{\nu}, \mathbb{S}_y^{[n/2]}) = +1,$$

which gives the desired contradiction.

Remark 53. Note that the proofs of Theorems 1 and 2 required only a few results from Section 2, namely Lemma 21, Lemma 27 and Proposition 50 (i.e. the whole Section 2.5). However, the proof of Theorem 5 will require the whole content of Section 2, that is, in addition to results needed for the proofs of Theorems 1 and 2, we will need Lemma 25 and Proposition 34.

4. Proof of Theorem 5.

Throughout the proof, we shall assume that the assumption a) holds, i.e. f maps almost all [n/2]-dimensional spheres into sets of [n/2] + 1 dimensional Hausdorff measure zero. The case when b) holds is treated in the same way. We simply need to exchange f and f^{-1} in the proof below.

We argue by contradiction: assume that the set $\{J_f < 0\}$ has positive measure.

To simplify the notation we shall write $\nu = n - 1 - [n/2]$. According to a local version of Lemma 25 for homeomorphisms on domains instead of \mathbb{R}^n , we can find x_o and a linear transformation $A \in GL(n)$ with $\det A > 0$ such that the sense preserving homeomorphism $g = A \circ f$ satisfies

$$\lim_{r \to 0^+} \|g_r - \mathcal{R}\|_{W^{1,[n/2]}(\mathbb{B}^n,\mathbb{R}^n)} = \lim_{r \to 0^+} \|(g_r)^{-1} - \mathcal{R}\|_{W^{1,\nu}(\mathbb{B}^n,\mathbb{R}^n)} = 0.$$

Let $\mathbb{B}_1 = \overline{\mathbb{B}}^{[n/2]+1}$, $\mathbb{B}_2 = \overline{\mathbb{B}}^{\nu+1}$ be closed unit balls. Note that the manifolds $\mathbb{B}_1 \times \mathbb{S}^{\nu}$ and $\mathbb{B}_2 \times \mathbb{S}^{[n/2]}$ have dimension n. Let $\iota_1 : \mathbb{B}_1 \times \mathbb{S}^{\nu} \hookrightarrow \mathbb{B}^n$ and $\iota_2 : \mathbb{B}_2 \times \mathbb{S}^{[n/2]} \hookrightarrow \mathbb{B}^n$ be smooth embeddings, smooth up to the boundary. According to Lemma 27 applied to the family $g_r \circ \iota_2$ and then for

the second time to the family $(g_r)^{-1} \circ \iota_1$, we can find a sequence $r_k \searrow 0$ such that the mappings $g_k := g_{r_k}$ satisfy:

There are compact sets $K_1 \subset \mathbb{B}_1$ and $K_2 \subset \mathbb{B}_2$ of positive measure such that

(22)
$$\lim_{k \to \infty} \sup_{x \in K_1} \|(g_k)^{-1} \circ \iota_1 - \mathcal{R} \circ \iota_1\|_{W^{1,\nu}(\{x\} \times \mathbb{S}^{\nu})} = 0$$

and

(23)
$$\lim_{k \to \infty} \sup_{y \in K_2} \|g_k \circ \iota_2 - \mathcal{R} \circ \iota_2\|_{W^{1,[n/2]}(\{y\} \times \mathbb{S}^{[n/2]})} = 0.$$

We shall define the embeddings ι_1 and ι_2 explicitly, to make sure that the embedded spheres $\iota_1(\{x\} \times \mathbb{S}^{\nu})$ and $\iota_2(\{y\} \times \mathbb{S}^{[n/2]})$ are linked with the linking number 1.

• for $x \in \mathbb{B}_1 \subset \mathbb{R}^{[n/2]+1}$, $\sigma \in \mathbb{S}^{\nu} \subset \mathbb{R}^{\nu+1}$, we set

$$\iota_1(x,\sigma) = (\frac{5+x_1}{10}\sigma_1, \dots, \frac{5+x_1}{10}\sigma_{\nu}, \frac{5+x_1}{10}\sigma_{\nu+1} - \frac{1}{4}, \frac{x_2}{10}\dots, \frac{x_{[n/2]+1}}{10}),$$

thus the image of ι_1 is the full torus with its core sphere lying in the hyperplane of the first $\nu + 1$ coordinates,

• for $y \in \mathbb{B}_2 \subset \mathbb{R}^{\nu+1}$ and $\rho \in \mathbb{S}^{[n/2]} \subset \mathbb{R}^{[n/2]+1}$, we set

$$\iota_2(y,\rho) = (\frac{y_1}{10}, \dots, \frac{y_{\nu}}{10}, \frac{5+y_{\nu+1}}{10}\rho_1 + \frac{1}{4}, \frac{5+y_{\nu+1}}{10}\rho_2, \dots, \frac{5+y_{\nu+1}}{10}\rho_{[n/2]+1}),$$

thus the image of ι_2 is the full torus with its core sphere lying in the hyperplane of the last $\lfloor n/2 \rfloor + 1$ coordinates.

To simplify the notation, as in the previous section we shall write $\mathbb{S}_x^{\nu} = \iota_1(\{x\} \times \mathbb{S}^{\nu}), \, \mathbb{S}_y^{[n/2]} = \iota_2(\{y\} \times \mathbb{S}^{[n/2]})).$

Since the assumption a) holds, i.e., f maps almost every sphere of dimension [n/2] to a set of ([n/2]+1)-Hausdorff measure zero, it easily follows that g_k , for every k, has the same property. We may thus assume, possibly shrinking K_2 , that for every $y \in K_2$ and any k the set $g_k(\mathbb{S}_y^{[n/2]})$ has ([n/2]+1)-Hausdorff measure zero. (Similarly, if b) holds, i.e., f^{-1} maps almost every sphere of dimension ν to a set of $(\nu+1)$ -Hausdorff measure zero, we can assume that for every $x \in K_1$ and any k the set $g_k^{-1}(\mathbb{S}_x^{\nu})$ has $(\nu+1)$ -Hausdorff measure zero.)

The images of ι_1 and ι_2 are two full, disjoint, linked tori, and for each $x \in \mathbb{B}_1$ and $y \in \mathbb{B}_2$, \mathbb{S}_x^{ν} and $\mathbb{S}_y^{[n/2]}$ are two linked spheres. We choose orientations of the spheres \mathbb{S}^{ν} and $\mathbb{S}^{[n/2]}$ so that the linking number equals

$$\ell(\mathbb{S}^{\nu}_{x},\mathbb{S}^{[n/2]}_{y}):=\ell\big(\iota_{1}|_{\{x\}\times\mathbb{S}^{\nu}},\iota_{2}|_{\{y\}\times\mathbb{S}^{[n/2]}}\big)=+1.$$

The first equality mans that for all $x \in \mathbb{B}_1$ and all $y \in \mathbb{B}_2$ we equip \mathbb{S}_x^{ν} and $\mathbb{S}_y^{[n/2]}$ with orientations so that the diffeomorphisms

$$\iota_1|_{\{x\}\times\mathbb{S}^{\nu}}:\{x\}\times\mathbb{S}^{\nu}\to\mathbb{S}^{\nu}_x\quad\text{and}\quad \iota_2|_{\{y\}\times\mathbb{S}^{[n/2]}}:\{y\}\times\mathbb{S}^{[n/2]}\to\mathbb{S}^{[n/2]}_y$$

are orientation preserving.

The linear transformation \mathcal{R} is the reflection in the last coordinate. Since the spheres $\mathbb{S}_y^{[n/2]}$ are centered at a point lying in $\mathbb{R}^{\nu+1} \times \{0\}$, the reflection \mathcal{R} preserves the center and hence $\mathcal{R}(\mathbb{S}_y^{[n/2]}) = \mathbb{S}_y^{[n/2]}$. However, the reflection \mathcal{R} changes the orientation of the sphere $\mathbb{S}_y^{[n/2]}$. More precisely, the mapping

(24)
$$\mathcal{R}: \mathbb{S}_y^{[n/2]} \to \mathbb{S}_y^{[n/2]} \quad \text{has degree } -1.$$

By $\mathcal{R}(\mathbb{S}_y^{[n/2]})$ we denote the sphere $\mathbb{S}_y^{[n/2]}$ with the opposite orientation, so the mapping $\mathcal{R}: \mathbb{S}_y^{[n/2]} \to \mathcal{R}(\mathbb{S}_y^{[n/2]})$ is orientation preserving. In particular, (24) implies that

(25)
$$\ell(\mathbb{S}_x^{\nu}, \mathcal{R}(\mathbb{S}_y^{[n/2]})) = -1.$$

On the other hand, the spheres \mathbb{S}_x^{ν} are centered at points lying in $\{0\} \times \mathbb{R}^{[n/2]+1}$ and the last coordinate of the center is $x_{[n/2]+1}/10$, where $x = (x_1, \dots, x_{[n/2]+1})$, so $\mathcal{R}(\mathbb{S}_x^{\nu}) = \mathbb{S}_{\tilde{x}}^{\nu}$ where $\tilde{x} = (x_1, \dots, x_{[n/2]}, -x_{[n/2]+1})$. Also, the orientation of $\mathcal{R}(\mathbb{S}_x^{\nu})$ is the same as that of $\mathbb{S}_{\tilde{x}}^{\nu}$. More precisely, the mapping $\mathcal{R}: \mathbb{S}_x^{\nu} \to \mathbb{S}_{\tilde{x}}^{\nu}$ is orientation preserving, so $\mathcal{R}(\mathbb{S}_x^{\nu})$ denotes the sphere $\mathbb{S}_{\tilde{x}}^{\nu}$ with the original orientation. In particular,

(26)
$$\ell(\mathcal{R}(\mathbb{S}_{x}^{\nu}), S_{y}^{[n/2]}) = \ell(\mathbb{S}_{\tilde{x}}^{\nu}, \mathbb{S}_{y}^{[n/2]}) = +1.$$

By Proposition 34 and Remark 35, for each $k \in \mathbb{N}$ we may define a continuous map

$$H_{1,k}: K_1 \times \mathbb{S}^{\nu} \times [0,1] \to \mathbb{R}^n$$

such that for each $x \in K_1$, $H_{1,k}(x,\cdot,\cdot)$ is a homotopy between

$$H_{1,k}(x,\cdot,1) = g_k^{-1} \circ \iota_1|_{\{x\} \times \mathbb{S}^{\nu}} \quad \text{and} \quad H_{1,k}(x,\cdot,0) = \mathcal{R} \circ \iota_1|_{\{x\} \times \mathbb{S}^{\nu}}.$$

Moreover, Proposition 34 along with (22) yields

$$\sup_{x \in K_1} \mathcal{H}^{\nu+1}(H_{1,k}(\{x\} \times \mathbb{S}^{\nu} \times [0,1)) \to 0 \quad \text{as } k \to \infty.$$

Therefore by taking a suitable subsequence of g_k^{-1} (still denoted by g_k^{-1}) we may require that

(27)
$$\sup_{x \in K_1} \mathcal{H}^{\nu+1}(H_{1,k}(\{x\} \times \mathbb{S}^{\nu} \times [0,1))) < \frac{1}{2^k}, \text{ for all } k.$$

Likewise, we define a continuous map

$$H_{2,k}: K_2 \times \mathbb{S}^{[n/2]} \times [0,1] \to \mathbb{R}^n$$

such that for all $y \in K_2$, $H_{2,k}(y,\cdot,\cdot)$ is a homotopy between

$$H_{2,k}(y,\cdot,1) = g_k \circ \iota_2|_{\{y\} \times \mathbb{S}^{[n/2]}}$$
 and $H_{2,k}(y,\cdot,0) = \mathcal{R} \circ \iota_2|_{\{y\} \times \mathbb{S}^{[n/2]}}.$

Again, Proposition 34 along with (23) yields

$$\sup_{y \in K_2} \mathcal{H}^{[n/2]+1}(H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1)) \to 0 \quad \text{as } k \to \infty.$$

Since for every $y \in K_2$ we have

$$\mathcal{H}^{[n/2]+1}(H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times \{1\}) = \mathcal{H}^{[n/2]+1}(g_k(\mathbb{S}_u^{[n/2]})) = 0,$$

by taking a suitable subsequence we may require that

(28)
$$\sup_{y \in K_2} \mathcal{H}^{[n/2]+1}(H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1])) < \frac{1}{2^k} \quad \text{for all } k.$$

Note that this is a stronger condition than (27) in the sense that now we have the estimate for the whole interval [0, 1], while in (27) we only have the estimate for the interval [0, 1).

We prove the following:

Lemma 54. For every $y \in K_2$, for almost every $x \in K_1$

(29)
$$\exists_{m \in \mathbb{N}} \ \forall_{k \geq m} \quad \mathbb{S}_x^{\nu} \cap H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1]) = \varnothing.$$

Also, for every $x \in K_1$, for almost every $y \in K_2$

$$\exists_{m \in \mathbb{N}} \ \forall_{k \ge m} \quad \mathbb{S}_y^{[n/2]} \cap H_{1,k}(\{x\} \times \mathbb{S}^{\nu} \times [0,1)) = \varnothing.$$

Proof of Lemma 54. We will only prove (29) since the proof of (30) follows from the same reasoning. Fix $y \in K_2$. We want to show that for almost all $x \in K_1$, (29) is satisfied.

We define a projection of the embedded full torus $\iota_1(\mathbb{B}_1 \times \mathbb{S}^{\nu})$ onto \mathbb{B}_1 by

$$\pi: \iota_1(\mathbb{B}_1 \times \mathbb{S}^{\nu}) \to \mathbb{B}_1, \qquad \pi(\iota(x, \sigma)) = x.$$

It is easy to see that π is 10-Lipschitz and hence it increases the Hausdorff measure $\mathcal{H}^{[n/2]+1}$ of a set at most by a constant factor $C(n) = 10^{[n/2]+1}$. Let

$$T_k(y) = H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1]) \cap \iota_1(\mathbb{B}_1 \times \mathbb{S}^{\nu})$$

be a part of the image of the homotopy $H_{2,k}(y,\cdot,\cdot)$ that is contained in the domain of the projection π . Estimate (28) and the fact that π increases the Hausdorff measure by at most C(n) imply

$$\mathcal{H}^{[n/2]+1}\left(\bigcup_{k=m}^{\infty} \pi(T_k(y))\right) < \sum_{k=m}^{\infty} C(n)2^{-k} = C(n)2^{-m+1} \to 0 \text{ as } m \to \infty,$$

so

$$\mathcal{H}^{[n/2]+1}\left(\bigcap_{m=1}^{\infty}\bigcup_{k=m}^{\infty}\pi(T_k(y))\right)=0.$$

To complete the proof of (29) it suffices to show that (29) is satisfied by all

(31)
$$x \in K_1 \setminus \bigcap_{m=1}^{\infty} \bigcup_{k=m}^{\infty} \pi(T_k(y)) = \bigcup_{m=1}^{\infty} \bigcap_{k=m}^{\infty} (K_1 \setminus \pi(T_k(y))).$$

Note that if $x \in K_1 \setminus \pi(T_k(y))$, then

$$\pi(\mathbb{S}_x^{\nu} \cap T_k(y)) \subset \pi(\mathbb{S}_x^{\nu}) \cap \pi(T_k(y)) = \{x\} \cap \pi(T_k(y)) = \emptyset,$$

so

(32)
$$\mathbb{S}_{x}^{\nu} \cap H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1]) = \mathbb{S}_{x}^{\nu} \cap T_{k}(y) = \varnothing.$$

Therefore if x belongs to the set (31), then

$$\exists_{m \in \mathbb{N}} \ \forall_{k \ge m} \quad x \in K_1 \setminus \pi(T_k(y)).$$

Since the condition $x \in K_1 \setminus \pi(T_k(y))$ implies (32), claim (29) follows.

To finish the proof of Theorem 5, we want to choose $x \in K_1$ and $y \in K_2$ in such a way that

i) there exists m such that

$$\mathbb{S}_{x}^{\nu} \cap \bigcup_{k=m}^{\infty} H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1]) = \varnothing,$$

i.e., the sphere \mathbb{S}_x^{ν} avoids, for all sufficiently large k, the image of the homotopy $H_{2,k}(y,\cdot,\cdot) \circ \iota_2^{-1}$ joining $g_k|_{\mathbb{S}_y^{[n/2]}}$ with $\mathcal{R}|_{\mathbb{S}_y^{[n/2]}}$, and simultaneously

ii) there exists m such that

$$\mathbb{S}_{y}^{[n/2]} \cap \bigcup_{k=m}^{\infty} H_{1,k}(\{x\} \times \mathbb{S}^{\nu} \times [0,1)) = \varnothing,$$

i.e., the sphere $\mathbb{S}_y^{[n/2]}$ avoid, for all sufficiently large k, the image of the homotopy $H_{1,k}(x,\cdot,\cdot)\circ\iota_1^{-1}$ joining $g_k^{-1}|_{\mathbb{S}_x^{\nu}}$ with $\mathcal{R}|_{\mathbb{S}_x^{\nu}}$, except possibly at the endpoint: we do not rule out yet that $\mathbb{S}_y^{[n/2]}\cap g_k^{-1}(\mathbb{S}_x^{\nu})\neq\varnothing$.

Assume we have chosen x and y satisfying conditions i) and ii) above. Then, for all sufficiently large k, $\mathbb{S}_x^{\nu} \cap g_k(\mathbb{S}_y^{[n/2]}) = \emptyset$, and since g_k^{-1} is a homeomorphism, this immediately implies that $\mathbb{S}_y^{[n/2]} \cap g_k^{-1}(\mathbb{S}_x^{\nu}) = \emptyset$. Therefore, for all sufficiently large k, the sphere $\mathbb{S}_y^{[n/2]}$ avoids the image of the whole homotopy joining $g_k^{-1}|_{\mathbb{S}_x^{\nu}}$ with $\mathcal{R}|_{\mathbb{S}_x^{\nu}}$, including the endpoint:

ii') there exists m such that

$$\mathbb{S}_y^{[n/2]} \cap \bigcup_{k=m}^{\infty} H_{1,k}(\{x\} \times \mathbb{S}^{\nu} \times [0,1]) = \varnothing.$$

Denote by $A_1 \subset K_1 \times K_2$ the set of all (x, y) satisfying the condition i) above. By Lemma 54, (29), for every $y_o \in K_2$ the section $A_1 \cap \{(x, y_o) : a \in K_1\}$ is of full measure, and thus, by Fubini's theorem, A_1 is of full measure in $K_1 \times K_2$, provided that A_1 is a measurable set. Similarly, the set $A_2 \subset K_1 \times K_2$ of these (x, y), which satisfy the condition ii), if measurable, is of full measure in $K_1 \times K_2$. We shall leave the issue of measurability of A_1 and A_2 and address it at the end of the proof. Since A_1 and A_2 are of full measure, their intersection is not empty and we can find x and y simultaneously satisfying the conditions i) and ii) and hence conditions i) and ii').

In particular, there is $x \in K_1$, $y \in K_2$ and $k \in \mathbb{N}$ such that

$$\mathbb{S}_{x}^{\nu} \cap H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1]) = \mathbb{S}_{y}^{[n/2]} \cap H_{1,k}(\{x\} \times \mathbb{S}^{\nu} \times [0,1]) = \emptyset.$$

We fix such a point $(x, y) \in K_1 \times K_2$ and we look at linking numbers of spheres and their images in g_k , g_k^{-1} and \mathcal{R} .

The mappings $g_k^{-1}|_{\mathbb{S}_x^{\nu}}$ and $\mathcal{R}|_{\mathbb{S}_x^{\nu}}$ are homotopic (with $H_{1,k}(x,\cdot,\cdot)\circ\iota_1^{-1}$ providing the homotopy), and the image of the homotopy does not intersect $\mathbb{S}_y^{[n/2]}$. This and (26) yield

$$+1 = \ell(\mathbb{S}_{\tilde{x}}^{\nu}, \mathbb{S}_{y}^{[n/2]}) = \ell(\mathcal{R}(\mathbb{S}_{x}^{\nu}), \mathbb{S}_{y}^{[n/2]}) = \ell(g_{k}^{-1}(\mathbb{S}_{x}^{\nu}), \mathbb{S}_{y}^{[n/2]})$$

and, since g_k is a sense preserving homeomorphism,

$$+1 = \ell(g_k^{-1}(\mathbb{S}_x^{\nu}), \mathbb{S}_y^{[n/2]}) = \ell(\mathbb{S}_x^{\nu}, g_k(\mathbb{S}_y^{[n/2]})).$$

Next, using the homotopy between $\mathcal{R}|_{\mathbb{S}_{u}^{[n/2]}}$ and $g_{k}|_{\mathbb{S}_{u}^{[n/2]}}$, given by $H_{2,k}(y,\cdot,\cdot)\circ\iota_{2}^{-1}$, we have

$$+1 = \ell(\mathbb{S}_x^{\nu}, g_k(\mathbb{S}_y^{[n/2]})) = \ell(\mathbb{S}_x^{\nu}, \mathcal{R}(\mathbb{S}_y^{[n/2]})) = -1,$$

where the last equality follows from (25). This gives the desired contradiction and finishes the proof, except for the set aside problem of measurability of the sets A_1 and A_2 .

Since the proof of measurability of both sets follows exactly the same scheme, we shall prove only that A_1 is measurable.

If we write

$$W_k = \{(x, y) \in K_1 \times K_2 : \mathbb{S}_x^{\nu} \cap H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0, 1]) = \emptyset\},\$$

then $A_1 = \bigcup_{m=1}^{\infty} \bigcap_{k=m}^{\infty} W_k$, thus to prove measurability of A_1 , it suffices to prove it for W_k .

Let

$$F_k = (\iota_1, H_{2,k}) \colon K_1 \times \mathbb{S}^{\nu} \times K_2 \times \mathbb{S}^{[n/2]} \times [0,1] \to \mathbb{R}^n \times \mathbb{R}^n$$

and denote by $\Delta = \{(x,x) : x \in \mathbb{R}^n\}$ the diagonal in $\mathbb{R}^n \times \mathbb{R}^n$. Then, since F_k is continuous, $K_1 \times \mathbb{S}^{\nu} \times K_2 \times \mathbb{S}^{[n/2]} \times [0,1]$ is compact and Δ is closed, the set $F_k^{-1}(\Delta) \subset \mathbb{R}^{2n+1}$ is compact. Let now

$$\Pi: K_1 \times \mathbb{S}^{\nu} \times K_2 \times \mathbb{S}^{[n/2]} \times [0,1] \to K_1 \times K_2$$

be the projection on the first and third factors. The set $\Pi(F_k^{-1}(\Delta))$ is a compact subset of $K_1 \times K_2$. We have

$$(x,y) \in \Pi(F_k^{-1}(\Delta))$$

$$\Leftrightarrow \text{ there exist } \sigma \in \mathbb{S}^{\nu}, \ \rho \in \mathbb{S}^{[n/2]}, \ t \in [0,1] \text{ with } (x,\sigma,y,\rho,t) \in F_k^{-1}(\Delta)$$

$$\Leftrightarrow \exists_{\sigma,\rho,t} \quad F_k(x,\sigma,y,\rho,t) \in \Delta$$

$$\Leftrightarrow \exists_{\sigma,\rho,t} \quad \iota_1(x,\sigma) = H_{2,k}(y,\rho,t)$$

$$\Leftrightarrow \mathbb{S}^{\nu}_x \cap H_{2,k}(\{y\} \times \mathbb{S}^{[n/2]} \times [0,1]) \neq \varnothing,$$

which shows that $W_k = (K_1 \times K_2) \setminus \Pi(F_k^{-1}(\Delta))$, and that W_k is an open (and thus measurable) subset of $K_1 \times K_2$. This concludes the proof of measurability of A_1 and the proof of Theorem 5.

Remark 55. Under the assumptions of Corollaries 10 and 11, the proof simplifies greatly. Recall that in these corollaries we assume n=2m, thus $\nu=n-[n/2]-1=m-1$. If we assume $f^{-1} \in W^{1,m-1+\varepsilon}$, then $g_k^{-1} \to \mathcal{R}$ in $W^{1,m-1+\varepsilon}$, and by the Morrey-Sobolev imbedding, on almost every sphere \mathbb{S}^{ν}_x this convergence is uniform (the same conclusion holds if we assume $Df^{-1} \in L^{m-1,1}_{loc}$). We can set the homotopy between $\mathcal{R}|_{\mathbb{S}^{\nu}_x}$ and $g_k^{-1}|_{\mathbb{S}^{\nu}_x}$ to be $H_1, k(x, \sigma, t) = tg_k^{-1}(\iota_1(x, \sigma) + (1-t)\mathcal{R}(\iota_1(x, \sigma))$; then for a.e. $x \in \mathbb{B}_1$ and sufficiently large k the whole image of the homotopy $H_1(x, \cdot, \cdot)$ lies in the torus $\iota_1(\mathbb{B}_1 \times \mathbb{S}^{[n/2]})$, thus for any $y \in \mathbb{B}_2$ this image does not intersect $\mathbb{S}^{[n/2]}_y$:

(34)
$$\exists_l \ \forall_{k>l} \quad \mathbb{S}_{y}^{[n/2]} \cap H_1(\{x\} \times \mathbb{S}^{\nu} \times [0,1]) = \varnothing.$$

Fix any $y \in \mathbb{B}_2$. Denoting by $\tilde{Y}_k(y)$ the projection of $\tilde{T}_k(y)$ onto the cross section $\iota_1(\mathbb{B}_1 \times \{\sigma_o\})$ (in analogy to $Y_k(y)$), we see that the (m+1)-dimensional Hausdorff measure of $\tilde{Y}_k(y)$ tends to zero. We can thus find $x \in \mathbb{B}_1$ such that (34) holds and $\iota_1(x,\sigma_o) \cap \tilde{Y}_k(y) = \emptyset$ for some k > l. Then \mathbb{S}_x^{ν} does not intersect the image of the homotopy joining $\mathcal{R}|_{\mathbb{S}_y^{[n/2]}}$ with $g_k|_{\mathbb{S}_y^{[n/2]}}$, except possibly at the endpoint – we have not ruled out that $\mathbb{S}_x^{\nu} \cap g_k(\mathbb{S}_y^{[n/2]}) \neq \emptyset$. We have thus found x and y satisfying conditions i) and ii) in the proof of Theorem 5 (although with x and y exchanged) and we may conclude the proof as it is done there.

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