

## Reshetnyak's Theorem and The Inner Distortion

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**Abstract:** We prove that a quasilight mapping of finite distortion with locally  $n$ -integrable weak partials and locally integrable inner distortion is discrete and open.

**Keywords:** Mapping of finite distortion, quasiregular mapping, Reshetnyak's theorem.

### 1. INTRODUCTION

We call  $f : \Omega \rightarrow \mathbb{R}^n$ ,  $n \geq 2$ , a mapping of finite distortion if  $f \in W_{\text{loc}}^{1,1}(\Omega, \mathbb{R}^n)$ ,  $J_f \in L_{\text{loc}}^1(\Omega)$ , and if there exists a measurable function  $K : \Omega \rightarrow [1, \infty)$  such that

$$|Df(x)|^n \leq K(x)J_f(x) \quad \text{a.e. } x \in \Omega.$$

Here  $|Df(x)|$  and  $J_f(x)$  are the operator norm and the Jacobian determinant of  $Df(x)$ , respectively. If  $K \in L^\infty(\Omega)$ ,  $f$  is called quasiregular, or a mapping of bounded distortion.

For a mapping of finite distortion  $f$ , the outer and inner distortion functions  $K_O$  and  $K_I$  are defined as

$$K_O(x) = \frac{|Df(x)|^n}{J_f(x)} \quad \text{and} \quad K_I(x) = \frac{|D^\#f(x)|^n}{J_f(x)^{n-1}},$$

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respectively, when  $0 < |Df(x)|$ ,  $J_f(x) < \infty$ , and  $K_O(x) = K_I(x) = 1$  otherwise. Here  $D^\sharp f(x)$  is the adjoint matrix of  $Df(x)$ . Then we have

$$K_I^{1/(n-1)}(x) \leq K_O(x) \leq K_I^{n-1}(x) \quad \text{a.e. } x \in \Omega.$$

In the late 1960s, Reshetnyak proved that a non-constant mapping of bounded distortion is always continuous, open and discrete. This theorem initiated the by now well-established theory of mappings of bounded distortion, see [13], [14], [6].

Recently, a lot of research has been done in order to find the sharp assumptions of Reshetnyak's theorem in the class of mappings of finite distortion, cf. [3], [4], [5], [7], [8], [11]. In this note we continue this line of research by giving a new partial result towards a conjecture of Iwaniec and Šverák [7].

**Theorem 1.1.** *Suppose that  $f : \Omega \rightarrow \mathbb{R}^n$ ,  $n \geq 2$ , is a quasilight mapping of finite distortion satisfying  $f \in W_{\text{loc}}^{1,n}(\Omega, \mathbb{R}^n)$  and  $K_I \in L_{\text{loc}}^1(\Omega)$ . Then  $f$  is discrete and open.*

By definition, a mapping  $f$  is called quasilight if the components of every point-inverse  $f^{-1}(y)$  are compact. The Iwaniec-Šverák conjecture is Theorem 1.1 without the quasilightness assumption. In [7] the conjecture is proved for  $n = 2$ . An example of Ball [2] shows that the integrability assumption on  $K_I$  cannot be relaxed in Theorem 1.1.

There are other partial results concerning the Iwaniec-Šverák conjecture, see [3], [4], [5] and [11]. The novelty in Theorem 1.1 lies in the fact that it only deals with the inner distortion; the previous results are proved under assumptions on the outer distortion function. In particular, Hencl and Malý [5] proved Theorem 1.1 assuming  $K_O \in L_{\text{loc}}^{n-1}(\Omega)$ , and Manfredi and Villamor [11] without the quasilightness assumption when  $K_O \in L_{\text{loc}}^p(\Omega)$  for some  $p > n - 1$ . It is clear that, when working with the inner distortion, one has to find methods different from those used in the above-mentioned works. We prove Theorem 1.1 by using the conformal modulus of  $(n - 1)$ -dimensional sets, the coarea formula, and elementary topological considerations. Also, we use several results concerning the theory of mappings of finite distortion. Another natural intermediate step towards the Iwaniec-Šverák conjecture would be the theorem of Manfredi and Villamor under the assumption  $K_I \in L_{\text{loc}}^p(\Omega)$  for some  $p > 1$  (instead of the assumption on  $K_O$ ), which we cannot prove. For closely related results on the

global invertibility properties of Sobolev mappings, see [2, Theorem 2] and [15, Corollary 2].

## 2. PRELIMINARIES

In this section we recall some known properties of mappings satisfying the assumptions of Theorem 1.1. First, let  $f : \Omega \rightarrow \mathbb{R}^n$  be a continuous map, and  $U \subset\subset \Omega$  open. Then the (local) topological degree  $\mu(y, f, U)$  is well-defined for every  $y \in \mathbb{R}^n \setminus f(\partial U)$ , see [14, I.4]. We will use the following facts:

$$(2.1) \quad \mu(y, f, U) = 0 \text{ if } y \notin f(U),$$

$$(2.2) \quad \mu(y, f, U) = \mu(v, f, U)$$

whenever  $y$  and  $v$  lie in the same component of  $\mathbb{R}^n \setminus f(\partial U)$ , and

$$(2.3) \quad \mu(y, f, U) = \sum_{i=1}^k \mu(y, f, U_i)$$

if both sides are well-defined, and if  $U_1, \dots, U_k$  are disjoint open sets satisfying

$$U \cap f^{-1}(y) \subset \bigcup_{i=1}^k U_i \subset U.$$

We call  $f$  sense-preserving if  $\mu(y, f, U) > 0$  whenever  $y \in f(U) \setminus f(\partial U)$ . Notice that if  $f$  is sense-preserving, then

$$\mu(y, f, U) \leq \mu(y, f, V)$$

whenever both sides are well-defined and  $U \subset V$ .

We say that  $f$  satisfies condition  $N$  if the  $n$ -measure  $|f(E)| = 0$  whenever  $|E| = 0$ . For mappings of finite distortion with locally  $n$ -integrable partials, we have

**Theorem 2.1** ([3, Theorem 1.3]). *Suppose that  $f \in W_{\text{loc}}^{1,n}(\Omega, \mathbb{R}^n)$  is a mapping of finite distortion. Then*

- (1)  $f$  has a continuous representative,
- (2)  $f$  is sense-preserving,
- (3)  $f$  satisfies condition  $N$ ,
- (4)  $f$  is differentiable almost everywhere in  $\Omega$ .

Part 3. implies that the change of variables formula holds for  $f$ . In fact, if  $U \subset\subset \Omega$  is open, we have

$$(2.4) \quad N(y, f, U) = \mu(y, f, U)$$

for almost every  $y \in \mathbb{R}^n \setminus f(\partial U)$ , see [5, Proposition 2]. Here

$$N(y, f, U) = \text{card}\{f^{-1}(y) \cap U\}.$$

Since  $K_O^{1/(n-1)} \leq K_I$  almost everywhere, [9, Theorem 1.2] implies

**Theorem 2.2.** *Suppose that  $f$  is as in Theorem 1.1. Then  $J_f(x) > 0$  for almost every  $x \in \Omega$ . In particular, if  $(A_i)$ ,  $A_i \subset \Omega$ , is a decreasing sequence of measurable sets so that  $|A_1| < \infty$  and  $\cap_i A_i \subset f^{-1}(y)$  for some  $y \in \mathbb{R}^n$ , then*

$$\int_{A_i} K_I \rightarrow 0 \quad \text{as } i \rightarrow \infty.$$

The following characterization of quasilightness will be useful in the sequel.

**Theorem 2.3** ([16, Theorem 3.1]). *A mapping  $f : \Omega \rightarrow \mathbb{R}^n$  is quasilight if and only if every point  $x \in \Omega$  has a neighborhood  $U \subset\subset \Omega$  such that  $f(x) \notin f(\partial U)$ .*

We call a mapping  $f$  light if every point-inverse  $f^{-1}(y)$  is totally disconnected. Hence a light mapping is quasilight in particular.

**Lemma 2.4** ([14, VI Lemma 5.6]). *If  $f : \Omega \rightarrow \mathbb{R}^n$  is continuous, light and sense-preserving, then  $f$  is discrete and open.*

By combining Theorem 2.1 and Lemma 2.4, we see that in order to prove Theorem 1.1 it suffices to show that  $f$  is light. We conclude this section with a topological lemma.

**Lemma 2.5.** *Let  $f$  be as in Theorem 1.1. Suppose that  $V \subset \mathbb{R}^n$  is homeomorphic to  $B(0, 1)$  and  $\bar{V}$  to  $\bar{B}(0, 1)$ , and that  $\emptyset \neq U \subset\subset \Omega$  is a component of  $f^{-1}(V)$ . Then  $f(\partial U) = \partial V$ , and  $f(U) = V$ .*

*Proof.* First,  $f(\partial U) \subset \partial V$  by the continuity of  $f$ . Hence, for every  $a \in f(U)$ ,  $\mu(a, f, U)$  is well-defined, and strictly positive by Theorem 2.1. By (2.1), there exists  $b \in \mathbb{R}^n$  such that  $\mu(b, f, U) = 0$ . Hence, by (2.2),  $f(\partial U)$  separates  $f(U)$  and  $b$ , and so  $f(\partial U) = \partial V$ . Also, if there exists a point  $p \in V \setminus f(U)$ , then  $\mu(p, f, U) = 0$ . But  $p$  and  $f(U)$  lie in the same component of  $\mathbb{R}^n \setminus f(\partial U) = \mathbb{R}^n \setminus \partial V$ .

Hence, by (2.2),  $\mu(p, f, U) = \mu(a, f, U) > 0$  whenever  $a \in f(U)$ . We conclude that  $f(U) = V$ .  $\square$

### 3. PREIMAGES OF RADIAL SEGMENTS

From now on we assume that  $f$  is as in Theorem 1.1. Recall from Section 2 that in order to prove Theorem 1.1 it suffices to show that  $f$  is light. We assume, in contrary, that there exists a point  $a \in \mathbb{R}^n$  such that some component of  $f^{-1}(a)$  has positive  $\mathcal{H}^1$ -measure. Without loss of generality,  $a = 0 \in f(\Omega)$ , and  $E$  is a component of  $f^{-1}(0)$  so that  $\mathcal{H}^1(E) > 0$ . Then Theorem 1.1 is proved if we can show that  $\mathcal{H}^1(E)$  has to be zero.

We denote the projection  $(x_1, \dots, x_n) \mapsto x_1$  by  $\text{pr}$ . By scaling and rotating, if necessary, we may assume that  $\mathcal{H}^1(\text{pr}(E)) = 1$ . By Theorem 2.3, there exists a domain  $G \subset\subset \Omega$  so that  $E \subset G$ , and a number  $M > 0$  so that  $|f(x)| \geq M$  for every  $x \in \partial G$ . Moreover, by Theorem 2.1, there exists  $m \in \mathbb{N}$  so that

$$(3.1) \quad \mu(y, f, G) = m \quad \text{for every } y \in B(0, M).$$

For  $0 < R < M$ , we denote the  $E$ -component of  $f^{-1}(B(0, R))$  by  $E_R$ . Then  $E_R \subset G$ . We define radial segments

$$I(R, \phi) = \{(t, \phi) : t \in (R/2, R)\}$$

in polar coordinates, and denote  $A_R = B(0, R) \setminus \overline{B(0, R/2)}$ , and  $U_R = E_R \cap f^{-1}(A_R)$ . The first main ingredient in the proof of Theorem 1.1 is the following.

**Proposition 3.1.** *There exists  $0 < M_0 < M$ , so that for each  $R < M_0$  there exist  $\phi_R \in S(0, 1)$  and  $a_R \in \mathbb{R}$ , so that if we denote*

$$L_R = (a_R - (4m)^{-1}, a_R + (4m)^{-1}),$$

then

$$L_R \subset \text{pr}(E) \quad \text{and} \quad \text{pr}^{-1}(L_R) \cap E_R \cap f^{-1}(I(R, \phi_R)) = \emptyset.$$

*Proof.* For  $R < M$ , define

$$h_R : U_R \rightarrow S(0, 1), \quad h_R(x) = \frac{f(x)}{|f(x)|}.$$

Then  $h_R^{-1}(\phi) = f^{-1}(I(R, \phi)) \cap E_R$  for every  $\phi \in S(0, 1)$ . Also, we have

$$|J_{n-1}h_R(x)| \leq \frac{|D^\#f(x)|}{|f(x)|^{n-1}} \quad \text{a.e. } x \in U_R,$$

for the  $(n - 1)$ -dimensional Jacobian of  $h_R$ . Then, the coarea formula (cf. [10]), and Hölder’s inequality yield

$$\begin{aligned} \int_{S(0,1)} \mathcal{H}^1(h_R^{-1}(\phi)) \, d\mathcal{H}^{n-1}(\phi) &= \int_{U_R} |J_{n-1}h_R| \leq \int_{U_R} \frac{|D^\#f|}{|f|^{n-1}} \\ (3.2) \qquad \qquad \qquad &= \int_{U_R} \frac{K_I^{1/n} J_f^{(n-1)/n}}{|f|^{n-1}} \\ &\leq \left( \int_{U_R} K_I \right)^{1/n} \left( \int_{U_R} \frac{J_f}{|f|^n} \right)^{(n-1)/n}. \end{aligned}$$

Since  $\mu(y, f, E_R) \leq m$  for every  $y \in B(0, R)$ , and  $U_R \subset E_R$ , the change of variables formula gives

$$(3.3) \qquad \int_{U_R} \frac{J_f}{|f|^n} \leq \int_{A_R} \frac{\mu(y, f, E_R)}{|y|^n} \, dy \leq m\omega_{n-1} \log 2.$$

Moreover, by Theorem 2.2,

$$(3.4) \qquad \int_{U_R} K_I \rightarrow 0 \quad \text{as } R \rightarrow 0.$$

Now, by combining (3.2), (3.3) and (3.4), we have: for every  $\epsilon > 0$  there exists  $k < M$  so that

$$(3.5) \qquad \int_{S(0,1)} \mathcal{H}^1(E_R \cap f^{-1}(I(R, \phi))) \, d\mathcal{H}^{n-1}(\phi) < \epsilon$$

for every  $R \leq k$ . Moreover, by slightly changing the set  $A_R$ , we see that (3.5) also holds for  $W_{R,\phi} = \overline{E_R} \cap f^{-1}(\overline{I(R, \phi)})$ .

Let  $R$  be as above. Next, we claim that, for each  $\phi \in S(0, 1)$ ,  $W_{R,\phi}$  consists of at most  $m$  components. Fix  $\phi$  and let  $\{J_i\}$ ,  $i = 1, \dots, N$  be a finite set of preimage components of  $\overline{I(R, \phi)}$  in  $\overline{E_R}$ . Denote by  $I_\delta$  the closed  $\delta$ -neighborhood of  $I(R, \phi)$ . Then  $I_\delta$  has  $N_\delta$  different preimage components  $\tilde{I}_\delta^j$  containing some  $J_i$ . When  $\delta$  is small enough,  $\tilde{I}_\delta^j \subset G$  for every  $j = 1, \dots, N_\delta$ . Then, by Lemma 2.5,  $f(\tilde{I}_\delta^j) = I_\delta$  for  $\delta$  small enough. Moreover, for  $\delta < \delta_0$  we have  $N_\delta = N$ . Then, if  $y \in I(R, \phi)$ , Theorem 2.1 and (2.3) yield

$$N \leq \sum \mu(y, f, \text{int } \tilde{I}_\delta^j) \leq \mu(y, f, G) = m.$$

This proves the claim.

Suppose that  $M_0 < M$  is small enough, so that (3.5) holds with  $\epsilon = \omega_{n-1}(100m)^{-1}$ . Then, in particular, for every  $R \leq M_0$  there exists  $\phi_R \in S(0, 1)$  such that

$$\mathcal{H}^1(\text{pr}(W_{R,\phi_R})) < (100m)^{-1}.$$

Moreover, we showed that  $W_{R,\phi_R}$  consists of at most  $m$  components. Now the proposition follows from our assumption  $\mathcal{H}^1(\text{pr}(E)) = 1$ .  $\square$

#### 4. MODULUS ESTIMATES AND THE PROOF OF THEOREM 1.1

In this section we prove Theorem 1.1, except for an upper bound for the conformal modulus of certain  $(n-1)$ -dimensional sets (Proposition 4.2). For a measurable function  $\omega \in L^1_{\text{loc}}(\Omega)$ ,  $\Omega \subset \mathbb{R}^n$ , and a family  $\Lambda = \{V_i : i \in I\}$  of Borel sets, set

$$M_\omega \Lambda = \inf_{\rho \in X(\Lambda)} \int_\Omega \omega \rho^{n/(n-1)},$$

where  $X(\Lambda)$  is the set of all Borel functions  $\rho : \Omega \rightarrow [0, \infty]$  satisfying

$$\int_{V_i} \rho d\mathcal{H}^{n-1} \geq 1$$

for every  $V_i \in \Lambda$  with  $\mathcal{H}^{n-1}(V_i) > 0$ . If  $\omega = 1$  almost everywhere in  $\Omega$ , we denote  $M_\omega$  by  $M$ .

Now fix  $R$  and  $a_R$  as in Proposition 3.1. Denote  $l = ((8m)^{-1}, (4m)^{-1})$ ,

$$\begin{aligned} V_t^+ &= E_R \cap \text{pr}^{-1}(\{a_R + t\}), & V_t^- &= E_R \cap \text{pr}^{-1}(\{a_R - t\}), \\ V_t &= V_t^+ \cup V_t^-, & Q_R^+ &= \{x \in V_t^+ : t \in l\}, & Q_R^- &= \{x \in V_t^- : t \in l\}, \end{aligned}$$

and

$$\Lambda_R = \{V_t : t \in l\}.$$

**Lemma 4.1.** *We have*

$$(16m)^{-n/(n-1)} \left( \int_{E_R} K_I \right)^{-1/(n-1)} \leq M_{K_I^{-1/(n-1)}} \Lambda_R \leq mMf(\Lambda_R).$$

*Proof.* Since  $f \in W^{1,n}(E_R, \mathbb{R}^n)$ , the restrictions of  $f$  to the components  $G_t^j$  of  $V_t$  belong to  $W^{1,n}(G_t^j, \mathbb{R}^n)$  for almost every  $t \in l$ . In particular, for those  $t$  the change of variables formula holds in  $V_t$ , see [12]. Also, Theorems 2.1 and 2.2 show that  $\mathcal{H}^{n-1}(f(V_t)) > 0$  for almost every  $t \in l$ .

Now fix  $\rho \in X(f(\Lambda_R))$ . Then, for almost every  $t \in I$ , the change of variables formula yields

$$(4.1) \quad \int_{V_t} (\rho \circ f) |D^\# f| d\mathcal{H}^{n-1} \geq \int_{fV_t} \rho d\mathcal{H}^{n-1} \geq 1,$$

i.e. the function  $\rho' : E_R \rightarrow [0, \infty]$ , defined as  $\rho'(x) = (\rho \circ f)(x) |D^\# f(x)|$  for  $x \in V_t$ ,  $t \in I$ , when (4.1) holds,  $\rho'(x) = \infty$  when  $x \in V_t$ ,  $t \in I$ , and (4.1) does not hold, and  $\rho'(x) = 0$  otherwise, belongs to  $X(\Lambda_R)$ . Now, by using the change of variables formula in  $E_R$ , with the fact that  $\mu(y, f, E_R) \leq m$  for every  $y \in B(0, R)$ , we have

$$\begin{aligned} \int_{E_R} (\rho')^{n/(n-1)} K_I^{-1/(n-1)} &= \int_{E_R} (\rho \circ f)^{n/(n-1)} |D^\# f|^{n/(n-1)} K_I^{-1/(n-1)} \\ &= \int_{E_R} (\rho \circ f)^{n/(n-1)} J_f \\ &\leq \int_{\mathbb{R}^n} \rho(y)^{n/(n-1)} \mu(y, f, E_R) dy \leq m \int_{\mathbb{R}^n} \rho^{n/(n-1)}. \end{aligned}$$

Since  $\rho \in X(f(\Lambda_R))$  is arbitrary, the second inequality in the lemma follows.

To prove the first inequality, fix  $g \in X(\Lambda_R)$ . Then, for every  $t \in I$ ,

$$1 \leq \int_{V_t^+} g d\mathcal{H}^{n-1} + \int_{V_t^-} g d\mathcal{H}^{n-1}.$$

By Fubini's theorem,

$$(8m)^{-1} \leq \int_{Q_R^+} g + \int_{Q_R^-} g,$$

so that one of the integrals, say the one over  $Q_R^+$ , is greater than  $(16m)^{-1}$ . Then, Hölder's inequality yields

$$(4.2) \quad (16m)^{-1} \leq \int_{Q_R^+} g K_I^{-1/n} K_I^{1/n} \leq \left( \int_{Q_R^+} g^{n/(n-1)} K_I^{-1/(n-1)} \right)^{(n-1)/n} \left( \int_{Q_R^+} K_I \right)^{1/n}.$$

Since  $g$  is arbitrary, (4.2) proves the first inequality in the lemma.  $\square$

In order to complete the proof of Theorem 1.1, we need an upper bound for  $Mf(\Lambda_R)$ .

**Proposition 4.2.**

$$Mf(\Lambda_R) \leq C,$$

where  $C > 0$  only depends on  $n$ .



We will prove Proposition 4.2 in Section 5. Assuming the proposition, Theorem 1.1 now follows: combining Lemma 4.1 with the proposition yields

$$(16m)^{-n/(n-1)} \left( \int_{E_R} K_I \right)^{-1/(n-1)} \leq mC,$$

where  $C$  does not depend on  $R$ . Thus,

$$\int_{E_R} K_I \geq T > 0,$$

with  $T$  independent of  $R$ . This contradicts Theorem 2.2, since

$$\bigcap_{R>0} E_R = E.$$

We conclude that  $\mathcal{H}^1(E) = 0$ , as desired.

## 5. PROOF OF PROPOSITION 4.2

We assume that  $n \geq 3$ . For  $n = 2$  the proposition is trivial. The idea for the proof is to show, using Proposition 3.1, that the sets  $f(V_t)$  separate  $I(R, \phi_R)$  and another “large” set in  $A_R$ . There are some technicalities, though, that slightly complicate matters.

Fix a point  $\xi \in \text{pr}^{-1}(a_R) \cap E$ , and denote by  $W$  the  $\xi$ -component of  $\mathbb{R}^n \setminus (V_{(8m)^{-1}} \cup \partial E_R)$ . Notice that, by the definition of  $V_t$ ,

$$(5.1) \quad \text{pr}(W) \subset (a_R - (8m)^{-1}, a_R + (8m)^{-1}).$$

**Lemma 5.1.** *For almost every  $r \in (R/2, R)$  there exist  $q_r \in W$  and a neighborhood  $U_r \subset W$  of  $q_r$  so that  $|p_r| = |f(q_r)| = r$  and*

$$(5.2) \quad f^{-1}(p_r) \cap U_r = \{q_r\}.$$

*Proof.* First, by Proposition 3.1, there exists a segment  $\alpha$  joining  $\partial E_R$  and  $\xi$  in  $W \cap \text{pr}^{-1}(a_R)$ . Fix a small  $\epsilon > 0$ . Then, for any  $x \in B^{n-1}(0, \epsilon)$ , we can choose a segment  $\alpha_x$  as follows: if  $\tilde{\alpha}$  is the line spanned by  $\alpha$ , then  $\tilde{\alpha}_x = \tilde{\alpha} + x$ ,  $x \in B^{n-1}(0, \epsilon) \subset H$ , where  $H \ni 0$  is the hyperplane orthogonal to  $\tilde{\alpha}$ . Moreover,  $\alpha_x$  is a segment in  $\tilde{\alpha}_x$  joining  $\partial E_R$  and  $B(\xi, \epsilon)$  in  $W$ . Choose  $\epsilon$  to be small enough, so that  $f(\alpha_x)$  connects  $S(0, R)$  and  $S(0, R/2)$  for every  $x \in B^{n-1}(0, \epsilon)$ .

By the definition of a mapping of finite distortion, and Theorems 2.1 and 2.2, there exists  $x_0 \in B^{n-1}(0, \epsilon)$  so that

- (1)  $f$  is absolutely continuous on  $\alpha_{x_0}$ ,
- (2)  $f$  is differentiable  $\mathcal{H}^1$ -almost everywhere on  $\alpha_{x_0}$ ,
- (3)  $J_f > 0$   $\mathcal{H}^1$ -almost everywhere on  $\alpha_{x_0}$ .

If  $f$  is differentiable at  $z \in \alpha_{x_0}$ , and  $J_f(z) > 0$ , then, for every  $\nu > 0$  small enough,

$$f(z) \notin f(S(z, \nu)).$$

Because this is true for almost every  $z \in \alpha_{x_0}$ , the absolute continuity of  $f$  on  $\alpha_{x_0}$  completes the proof.  $\square$

Denote by  $D$  the exceptional set in Lemma 5.1. For a radius  $r \in (R/2, R) \setminus D$ , denote  $\{\beta_r\} = S(0, r) \cap I(R, \phi_R)$ . By (5.1), Lemma 5.1 and (5.3) below,  $\beta_r \neq p_r$  for every  $r$ .

**Lemma 5.2.** *Let  $\kappa : [0, 1] \rightarrow S(0, r)$  be a one-to-one  $C^\infty$ -path such that  $\kappa(0) = p_r$  and  $\kappa(1) = \beta_r$ . Then, for every  $t \in ((8m)^{-1}, (4m)^{-1})$ ,*

$$\kappa((0, 1)) \cap f(V_t) \neq \emptyset.$$

*Proof.* Recall that

$$(5.3) \quad \text{pr}^{-1}((a_R - (4m)^{-1}, a_R + (4m)^{-1})) \cap E_R \cap f^{-1}(\beta_r) = \emptyset$$

by Proposition 3.1. For  $q_r$  and  $U_r$  as in Lemma 5.1, denote by  $\tilde{\kappa}$  the  $q_r$ -component of  $f^{-1}(\kappa([0, 1]))$ . By using Lemma 2.5 as below, we see that  $\tilde{\kappa} \neq \{q_r\}$ . Then, by (5.2), we find  $s \in (0, 1)$ , and a component  $\kappa'$  of  $f^{-1}(\kappa([s, 1]))$  so that  $\kappa' \cap U_r \neq \emptyset$  and  $\kappa' \subset \tilde{\kappa}$ .

We assume that  $\kappa' \cap V_t = \emptyset$ . Since  $f(\partial E_R) = S(0, R)$ , we conclude that  $\kappa'$  is compact. On the other hand,  $\beta_r = \kappa(1) \notin f(\kappa')$  by (5.3). Thus there exists  $t \in (s, 1)$  so that

$$(5.4) \quad t = \max\{\tau : \kappa(\tau) \in f(\kappa')\}.$$

Choose a point  $x_t \in f^{-1}(\kappa(t)) \cap \kappa'$ . By our assumption on  $\kappa'$ , the  $x_t$ -component of  $f^{-1}(\kappa(t))$  does not intersect  $V_t$ . Then there exists a ball  $B = B(\kappa(t), \epsilon)$  so that the  $x_t$ -component  $U_t$  of  $f^{-1}(B)$  does not intersect  $V_t$ . By Lemma 2.5  $f(U_t) = B$ , and since  $\kappa$  is  $C^\infty$ , applying Lemma 2.5 to the  $\epsilon$ -neighborhoods of  $\kappa((t - \delta, t + \delta))$  for small enough  $\delta$ , and the  $x_t$ -components of their preimages, shows that actually  $\kappa([t, t + \delta]) \subset f(\kappa')$ , contradicting (5.4). The proof is complete.  $\square$

**Lemma 5.3.** For every  $r \in (R/2, R) \setminus D$ , there exists a Borel function  $\rho_r : S(0, r) \rightarrow [0, \infty]$  so that, whenever  $t \in ((8m)^{-1}, (4m)^{-1})$ ,

$$(5.5) \quad \int_{S(0,r) \cap f(V_t)} \rho_r d\mathcal{H}^{n-2} \geq C_1/r,$$

and

$$(5.6) \quad \int_{S(0,r)} \rho_r^{n/(n-1)} d\mathcal{H}^{n-1} \leq C_2/r,$$

where the constants  $C_1, C_2 > 0$  only depend on  $n$ .

*Proof.* We first map  $S(0, r)$  onto  $S(e_n/2, 1/2)$  by a map  $T$  which is a composition of scaling, translation and rotation, so that  $T(\beta_r) = e_n$ . Then, if  $\rho : S(e_n/2, 1/2) \rightarrow [0, \infty]$  satisfies

$$(5.7) \quad \int_{(T \circ f)(V_t)} \rho d\mathcal{H}^{n-2} \geq C_1(n)$$

for all  $t \in ((8m)^{-1}, (4m)^{-1})$ , and

$$(5.8) \quad \int_{S(e_n/2, 1/2)} \rho^{n/(n-1)} d\mathcal{H}^{n-1} \leq C_2(n),$$

then the function  $\rho_r = r^{1-n}(\rho \circ T)$  satisfies (5.5) and (5.6). Hence it suffices to show (5.7) and (5.8).

If we map  $S(e_n/2, 1/2)$  onto  $\overline{\mathbb{R}^{n-1}}$  by the stereographic projection  $h$ ,

$$h(x) = e_n + (x - e_n)/|x - e_n|^2,$$

then  $e_n = T(\beta_r)$  gets mapped to  $\infty$ . We denote

$$a = (h \circ T)(p_r) \in \mathbb{R}^{n-1}.$$

We define  $\rho : \mathbb{R}^{n-1} \rightarrow [0, \infty]$ ,

$$\rho(x) = |x - a|^{2-n}(1 + |x|^2)^{n-2},$$

and denote  $Y_t = (h \circ T \circ f)(V_t)$ . Then we have to show that

$$(5.9) \quad \int_{Y_t} \frac{\rho(x)}{(1 + |x|^2)^{n-2}} d\mathcal{H}^{n-2}(x) = \int_{Y_t} |x - a|^{2-n} d\mathcal{H}^{n-2}(x) \geq C_1(n)$$

for all  $t \in ((8m)^{-1}, (4m)^{-1})$ , and

$$(5.10) \quad \int_{\mathbb{R}^{n-1}} \frac{\rho^{n/(n-1)}(x)}{(1 + |x|^2)^{n-1}} d\mathcal{H}^{n-1}(x) = \int_{\mathbb{R}^{n-1}} \frac{|x - a|^{1-n+1/(n-1)}}{(1 + |x|^2)^{1/(n-1)}} d\mathcal{H}^{n-1}(x) \leq C_2(n).$$

By Lemma 5.2, for every  $\alpha \in S^{n-2}(0, 1)$ , the half-line

$$I_\alpha = \{a + \alpha t : t > 0\}$$

intersects  $Y_t$ . For  $i \in \mathbb{Z}$ , denote  $A_i = B(a, 2^i) \setminus B(a, 2^{i-1})$ , and

$$\Phi_i = \{\alpha \in S^{n-2}(0, 1) : I_\alpha \cap A_i \cap Y_t \neq \emptyset\}.$$

Then a projection argument shows that

$$(5.11) \quad \int_{Y_t \cap A_i} |x - a|^{2-n} d\mathcal{H}^{n-2}(x) \geq C(n)\mathcal{H}^{n-2}(\Phi_i).$$

Since  $\sum_i \mathcal{H}^{n-2}(\Phi_i) = \omega_{n-2}$ , (5.9) follows by summing over  $i$ .

In order to prove (5.10), we first consider the case  $|a| > 1$ . We divide  $\mathbb{R}^{n-1}$  to  $N_1 = B^{n-1}(a, |a|/2)$ ,  $N_2 = B^{n-1}(0, |a|/2)$  and  $N_3 = \mathbb{R}^{n-1} \setminus (N_1 \cup N_2)$ . Then

$$\begin{aligned} \int_{N_1} \frac{|x - a|^{1-n+1/(n-1)}}{(1 + |x|^2)^{1/(n-1)}} d\mathcal{H}^{n-1}(x) &\leq C|a|^{\frac{-2}{n-1}} \int_{N_1} |x - a|^{1-n+1/(n-1)} d\mathcal{H}^{n-1}(x) \\ &\leq C|a|^{\frac{-1}{n-1}}, \end{aligned}$$

$$\begin{aligned} \int_{N_2} \frac{|x - a|^{1-n+1/(n-1)}}{(1 + |x|^2)^{1/(n-1)}} d\mathcal{H}^{n-1}(x) &\leq C|a|^{-2+\frac{1}{n-1}} \int_0^{|a|/2} \frac{r}{(1 + r^2)^{1/(n-1)}} dr \\ &\leq C|a|^{\frac{-1}{n-1}}, \end{aligned}$$

and, since  $10|x| \geq |x - a|$  for  $x \in N_3$ ,

$$\begin{aligned} \int_{N_3} \frac{|x - a|^{1-n+1/(n-1)}}{(1 + |x|^2)^{1/(n-1)}} d\mathcal{H}^{n-1}(x) &\leq C \int_{N_3} |x - a|^{1-n-1/(n-1)} d\mathcal{H}^{n-1}(x) \\ &\leq C|a|^{\frac{-1}{n-1}}. \end{aligned}$$

Combining the integrals proves (5.10) in the case  $|a| > 1$ . The case  $|a| \leq 1$  is similar, but now it suffices to consider the division  $\tilde{N}_1 = B^{n-1}(a, 3)$ ,  $\tilde{N}_2 = \mathbb{R}^{n-1} \setminus \tilde{N}_1$ .

□

Define  $\rho : A_R \rightarrow [0, \infty]$ ,  $\rho(x) = \rho_{|x|}(x)$ , where  $\rho_r$  is as in Lemma 5.3 for  $r \notin D$ , and  $\rho_r = 0$  otherwise. Since the restrictions of  $f$  to the components  $G_t^j$  of  $V_t$  belong to  $W^{1,n}(G_t^j, \mathbb{R}^n)$  for almost every  $t$ ,  $f(V_t)$  is countably  $(n - 1)$ -rectifiable

for  $t \in ((8m)^{-1}, (4m)^{-1}) \setminus Q$ , where  $\mathcal{H}^1(Q) = 0$ . Then, Lemma 5.3, and the coarea formula for rectifiable sets, cf. [1, Theorem 2.93 and Remark 2.94], yield

$$\int_{f(V_t)} \rho d\mathcal{H}^{n-1} \geq C(n) \int_{R/2}^R \int_{f(V_t) \cap S(0,r)} \rho d\mathcal{H}^{n-2} dr \geq C(n)$$

for every  $t \in ((8m)^{-1}, (4m)^{-1}) \setminus Q$ . Also, by Lemma 5.3,

$$\int_{A_R} \rho^{n/(n-1)} = \int_{R/2}^R \int_{S(0,r)} \rho^{n/(n-1)} d\mathcal{H}^{n-1} dr \leq C(n) \int_{R/2}^R \frac{dr}{r} \leq C(n).$$

By Theorem 2.1,  $M\{f(V_t) : t \in Q\} = 0$ . The proof of Proposition 4.2 is complete.

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