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# The fourth-order total variation flow in $\mathbb{R}^n$

Dedicated to Professor Neil Trudinger on the occasion of his 80th birthday

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## Abstract

We define rigorously a solution to the fourth-order total variation flow equation in  $\mathbb{R}^n$ . If  $n \geq 3$ , it can be understood as a gradient flow of the total variation energy in  $D^{-1}$ , the dual space of  $D_0^1$ , which is the completion of the space of compactly supported smooth functions in the Dirichlet norm. However, in the low dimensional case  $n \leq 2$ , the space  $D^{-1}$  does not contain characteristic functions of sets of positive measure, so we extend the notion of solution to a larger space. We characterize the solution in terms of what is called the Cahn-Hoffman vector field, based on a duality argument. This argument relies on an approximation lemma which itself is interesting.

We introduce a notion of calibrability of a set in our fourth-order setting. This notion is related to whether a characteristic function preserves its form throughout the evolution. It turns out that all balls are calibrable. However, unlike in the second-order total variation flow, the outside of a ball is calibrable if and only if  $n \neq 2$ . If  $n \neq 2$ , all annuli are calibrable, while in the case  $n = 2$ , if an annulus is too thick, it is not calibrable.

We compute explicitly the solution emanating from the characteristic function of a ball. We also provide a description of the solution emanating from any piecewise constant, radially symmetric datum in terms of a system of ODEs.

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Keywords: fourth-order; total variation flow; calibrability; subdifferential; radial solution.

## 1 Introduction

We consider the fourth-order total variation flow equation in  $\mathbb{R}^n$  of the form

$$u_t = -\Delta \operatorname{div} \frac{\nabla u}{|\nabla u|}. \quad (1.1)$$

We aim to give explicit description of its solutions emanating from piecewise constant radial data. However, it turns out that the definition of a solution is itself non-trivial since  $-\Delta$  does not have a bounded inverse on  $L^2(\mathbb{R}^n)$ . Our first goal is thus to provide a rigorous definition of a

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solution. Our second goal is to find explicit formula for the solution to (1.1) when the initial datum  $u(0, x) = u_0(x)$  is the characteristic function of a ball or an annulus. In other words,

$$u_0 = a_0 \mathbf{1}_{B_{R_0}} \quad \text{or} \quad u_0 = a_0 \mathbf{1}_{A_{R_0^1}^{R_0^0}} \quad a_0 \in \mathbb{R},$$

where  $\mathbf{1}_K$  is the characteristic function of a set  $K \subset \mathbb{R}^n$ , i.e.,

$$\mathbf{1}_K(x) = \begin{cases} 1, & x \in K \\ 0, & x \in \mathbb{R}^n \setminus K. \end{cases}$$

Here  $B_R$  denotes the open ball of radius  $R$  centered at  $0 \in \mathbb{R}^n$  and  $A_{R_0^1}^{R_0^0}$  denotes the annulus defined by  $A_{R_0^1}^{R_0^0} = B_{R_1} \setminus \overline{B_{R_0}}$ . Our major concern is whether or not the solution remains a characteristic function throughout the evolution. For example, in the case  $u_0 = a_0 \mathbf{1}_{B_{R_0}}$ , whether or not the solution  $u$  of (1.1) is of the form

$$u(t, x) = a(t) \mathbf{1}_{B_{R(t)}}$$

with a function  $a = a(t)$ . In other words, we are asking whether the speed  $u_t$  on the ball  $B_{R(t)}$  and on its complement are constant in the spatial variable. As in the second-order problem [3] (see also [16]), this leads to the notion of calibrability of a set. In the case of the second-order problem  $u_t = \operatorname{div}(\nabla u / |\nabla u|)$ , a ball and its complement are always calibrable and  $R(t) \equiv R_0$ , i.e. the ball does not expand nor shrink [3]. In our problem,  $R(t)$  may not be constant.

We first note that the definition of a solution itself is non-trivial. The fourth-order total variation flow has been mainly studied in the periodic setting [13], [11] or in a bounded domain with some boundary conditions [14]. Formally, it is a gradient flow of the total variation functional

$$TV(u) := \int_{\Omega} |\nabla u|$$

with respect to the inner product

$$(u, v)_{-1} = \int_{\Omega} u(-\Delta)^{-1} v$$

when  $\Omega$  is a domain in  $\mathbb{R}^n$  or a flat torus  $\mathbb{T}^n$ . In the periodic setting, i.e.  $\Omega = \mathbb{T}^n$  as in [13], [11], it is interpreted as a gradient flow in  $H_{av}^{-1}$  which is the dual space of  $H_{av}^1$ , the space of average-free  $H^1$  functions equipped with the inner product

$$(u, v)_1 = \int_{\Omega} \nabla u \cdot \nabla v.$$

For the homogeneous Dirichlet boundary condition with bounded  $\Omega$ ,  $H_{av}^{-1}$  is replaced by  $D^{-1}$ , the dual space of  $D_0^1 = D_0^1(\Omega)$ , which is the completion of  $C_c^\infty(\Omega)$  in the norm associated with the inner product  $(u, v)_1$ ; here  $C_c^\infty(\Omega)$  denotes the space of all smooth functions compactly supported in  $\Omega$ . By the Poincaré inequality, both  $H_{av}^1$  and  $D_0^1(\Omega)$  can be regarded as subspaces of  $L^2(\Omega)$ . However, if  $\Omega$  equals  $\mathbb{R}^n$ , the situation is more involved. If  $n \geq 3$ ,  $D_0^1(\mathbb{R}^n)$  is continuously and densely embedded in  $L^{2^*}(\mathbb{R}^n)$ , where  $2^* = np/(n-p)$  so that  $2^* = 2n/(n-2)$ , by the Sobolev inequality. In fact,

$$D_0^1(\mathbb{R}^n) = D^1(\mathbb{R}^n) \cap L^{2^*}(\mathbb{R}^n), \quad D^1(\mathbb{R}^n) = \{u \in L_{loc}^1(\mathbb{R}^n) \mid \nabla u \in L^2(\mathbb{R}^n)\}$$

see e.g. [10]. On the other hand, if  $n \leq 2$ ,  $D_0^1$  is isometrically identified with the quotient space  $\dot{D}^1(\mathbb{R}^n) := D^1(\mathbb{R}^n)/\mathbb{R}$ , when  $D^1(\mathbb{R}^n)$  is equipped with inner product  $(u, v)_1$  [10]. Thus, we need to be careful when  $n \leq 2$  because an element of  $D_0^1(\mathbb{R}^n)$  is determined only up to a constant. In any case,  $D_0^1(\mathbb{R}^n)$  is a Hilbert space with the scalar product

$$(u, v)_{D_0^1(\mathbb{R}^n)} = \int_{\mathbb{R}^n} \nabla u \cdot \nabla v.$$

Therefore, we can identify  $D_0^1(\mathbb{R}^n)$  with its dual space by means of the isometry

$$-\Delta: u \mapsto (u, \cdot)_{D_0^1(\mathbb{R}^n)}.$$

On the other hand, let us define a subspace  $\tilde{D}^{-1}(\mathbb{R}^n) \subset D_0^1(\mathbb{R}^n)'$  by

$$\begin{aligned} \tilde{D}^{-1}(\mathbb{R}^n) &= \left\{ w \mapsto \int_{\mathbb{R}^n} uw : u \in C_c^\infty(\mathbb{R}^n) \right\} \quad \text{if } n \geq 3, \\ \tilde{D}^{-1}(\mathbb{R}^n) &= \left\{ w \mapsto \int_{\mathbb{R}^n} uw : u \in C_{c,av}^\infty(\mathbb{R}^n) \right\} \quad \text{if } n = 1 \text{ or } n = 2, \end{aligned}$$

where

$$C_{c,av}^\infty(\mathbb{R}^n) = \left\{ u \in C_c^\infty(\mathbb{R}^n) : \int_{\mathbb{R}^n} u = 0 \right\}.$$

Then the closure  $D^{-1}(\mathbb{R}^n)$  of  $\tilde{D}^{-1}(\mathbb{R}^n)$  coincides with  $D_0^1(\mathbb{R}^n)'$  [10]. Note that the restriction to  $C_{c,av}^\infty(\mathbb{R}^n)$  in the definition of  $\tilde{D}^{-1}(\mathbb{R}^n)$  in  $n = 1, 2$  is necessary for the functionals to be well-posed on  $D^1(\mathbb{R}^n)/\mathbb{R}$ . In any case, since (by definition) the space of test functions  $\mathcal{D}(\mathbb{R}^n)$  is continuously and densely embedded in  $D_0^1(\mathbb{R}^n)$ , we also have a continuous embedding  $D^{-1}(\mathbb{R}^n) = D_0^1(\mathbb{R}^n)' \subset \mathcal{D}'(\mathbb{R}^n)$ . Throughout the paper, we will often drop  $(\mathbb{R}^n)$  in the notation for spaces of functions on  $\mathbb{R}^n$ , e. g.  $D^{-1} = D^{-1}(\mathbb{R}^n)$ .

In the first step, we give a rigorous definition of the total variation functional  $TV$  on  $D^{-1}$ . Then we calculate the subdifferential of  $TV$  in  $D^{-1}$  space. Since it is a homogeneous functional, we are able to apply a duality method [3] to characterize the subdifferential, provided that  $TV$  is well approximated by nice functions in  $D^{-1}$ . We know that  $C_{c,av}^\infty(\mathbb{R}^n)$  is dense in  $D^{-1}$  for  $n \leq 2$ ; see e. g. [10]. However, it is not immediately clear whether  $TV$  is simultaneously approximable. Fortunately, it turns out that for any  $w \in D^{-1}$ , there is a sequence  $w_k \in C_{c,av}^\infty(\mathbb{R}^n)$  which converges to  $w$  in  $D^{-1}$  and  $TV(w_k) \rightarrow TV(w)$  as  $k \rightarrow \infty$ . This approximation part is relatively involved since we have to use an efficient cut-off function. Using the approximation, we are able to characterize the subdifferential  $\partial_{D^{-1}}TV$  of  $TV$  in  $D^{-1}$  by adapting the argument in [3]. Namely, we have

$$\partial_{D^{-1}}TV(u) = \{v = \Delta \operatorname{div} Z \mid Z \in L^\infty(\mathbb{R}^n), |Z| \leq 1, -\langle u, \operatorname{div} Z \rangle = TV(u)\},$$

where  $\langle \cdot, \cdot \rangle$  denotes the canonical pairing of  $D^{-1}$  and  $D_0^1$ . A vector field  $Z$  corresponding to an element of the subdifferential is often called a *Cahn-Hoffman vector field*. The equation (1.1) should be interpreted as the gradient flow of  $TV$  in  $D^{-1}$ , i. e.

$$u_t \in -\partial_{D^{-1}}TV(u), \tag{1.2}$$

and its unique solvability for any initial datum  $u_0 \in D^{-1}$  is guaranteed by the classical theory of maximal monotone operators ([23], [6]). By our characterization of the subdifferential, we are able to give a more explicit definition of a solution which is consistent with that proposed in [14]. Namely, for  $u_0 \in D^{-1}$  with  $TV(u_0) < \infty$ , a function  $u \in C([0, T[, D^{-1})$  is a solution to (1.2) with  $u(0) = u_0$  if and only if there exists  $Z \in L^\infty(]0, T[ \times \mathbb{R}^n)$  satisfying  $\operatorname{div} Z \in L^2(0, T; D_0^1(\mathbb{R}^n))$  such that

$$u_t = -\Delta \operatorname{div} Z \quad \text{in } D^{-1}(\mathbb{R}^n) \tag{1.3}$$

$$|Z(t, x)| \leq 1 \quad \text{for a. e. } x \in \mathbb{R}^n \tag{1.4}$$

$$\langle u, \operatorname{div} Z \rangle = -TV(u) \tag{1.5}$$

for a. e.  $t \in ]0, T[$ . This is convenient for calculating explicit solutions.

Unfortunately, in  $n \leq 2$ , for a compactly supported square integrable function  $u_0$ , we know that  $u_0 \in D^{-1}$  if and only if  $u_0$  is average-free, i. e.  $\int_{\mathbb{R}^n} u_0 = 0$  (see Lemma 16). Thus, the characteristic function of any bounded, measurable set of positive measure does not belong to  $D^{-1}$ . We have

to extend a class of initial data  $u_0$  such that  $u_0 = \psi + w_0$  with  $w_0 \in D^{-1}$  while  $\psi$  is a fixed compactly supported  $L^2$  function. We consider a gradient flow  $u_t \in -\partial_{D^{-1}}TV(u)$  in the affine space  $\psi + D^{-1}$ . Since  $\partial_{D^{-1}}$  is a directional partial derivative in the direction of  $D^{-1}$ , it is more convenient to consider solutions to evolutionary variational inequality

$$\frac{1}{2} \frac{d}{dt} \|u(t) - g\|_{D^{-1}}^2 \leq TV(g) - TV(u(t)) \quad \text{for a. e. } t > 0 \quad (1.6)$$

for any  $g \in \psi + D^{-1}$  [2]. In the case  $\psi = 0$ , it is easy to show that the evolutionary variational inequality is equivalent to (1.2). Indeed, by definition of the subdifferential, (1.2) is equivalent to

$$(-u_t, g - u(t))_{D^{-1}} \leq TV(g) - TV(u(t))$$

for any  $g \in D^{-1}$ . The left-hand side equals  $(d/dt) (\|u - g\|^2/2)$ . Thus, the equivalence follows if  $\psi = 0$ .

It is easy to check that there is at most one solution to the evolutionary variational inequality (1.6). The solution  $u$  is constructed by solving

$$w_t \in -\partial_{D^{-1}}TV(w + \psi) \quad \text{with} \quad w(0) = w_0 = u_0 - \psi$$

and setting  $u = w + \psi$ . Characterization of the (directional) subdifferential is more involved since  $w \mapsto TV(w + \psi)$  is no more positively one-homogeneous. We identify the one dimensional space  $\{c\psi | c \in \mathbb{R}\}$  with  $\mathbb{R}$  and consider the Hilbert space  $E^{-1}$  defined as the orthogonal sum  $D^{-1} \oplus \mathbb{R}$ . We calculate the subdifferential by the duality method since  $TV$  is now positively one-homogeneous on  $E^{-1}$ . We then project this subdifferential onto  $D^{-1}$  to get a characterization of a (directional) subdifferential  $\partial_{D^{-1}}TV$ . We end up with a characterization of solution to (1.2) similar to (1.3)–(1.5), with (1.5) adjusted in a suitable way. If we also denote  $E_0^1 = D^1$  in  $n \leq 2$ ,  $E_0^1 = D_0^1$ ,  $E^{-1} = D^{-1}$  in  $n \geq 3$  and

$$\langle u, v \rangle_E = \begin{cases} \langle u, v \rangle & \text{if } n \geq 3, \\ \langle w, [v] \rangle + c \int \psi v, & \text{where } u = w + c\psi, \quad w \in D^{-1} \text{ if } n \leq 2, \end{cases} \quad (1.7)$$

for  $u \in E^{-1}$ ,  $v \in E_0^1$ , we end up with the following definition of solution

**Definition 1.** Assume that  $u_0 \in E^{-1}$ . We say that  $u \in C([0, \infty[, E^{-1})$  with  $u_t \in L_{loc}^2([0, \infty[, D^{-1})$  is a solution to (1.1) with initial datum  $u_0$  if there exists  $Z \in L^\infty([0, \infty[ \times \mathbb{R}^n)$  with  $\text{div } Z(t, \cdot) \in E_0^1$  for a. e.  $t > 0$  such that

$$u_t = -\Delta \text{div } Z \quad \text{in } D^{-1}(\mathbb{R}^n) \quad (1.8)$$

$$|Z(t, x)| \leq 1 \quad \text{for a. e. } x \in \mathbb{R}^n \quad (1.9)$$

$$-\langle u, \text{div } Z \rangle_E = TV(u) \quad (1.10)$$

for a. e.  $t > 0$ .

and associated well-posedness result

**Theorem 2.** Let  $u_0 \in E^{-1}$ . There exists a unique solution to (1.1) with initial datum  $u_0$ .

Our next problem is whether or not the speed of a characteristic function of a set is spatially constant inside and outside of the set. By the general theory ([6], [23]), the speed is determined by the minimal section (canonical restriction)  $\partial_{D^{-1}}^0TV$  of  $\partial_{D^{-1}}TV$ . In other words,  $\partial_{D^{-1}}^0TV(u) = v_0$  minimizes  $\|v\|_{D^{-1}}$  for  $v \in \partial_{D^{-1}}TV(u)$ , i.e.,

$$\partial_{D^{-1}}^0TV(u) := \arg \min \{ \|v\|_{D^{-1}} \mid v \in \partial_{D^{-1}}TV(u) \}.$$

To motivate the notion of calibrability, we consider a smooth function  $u$  such that

$$\bar{U} = \{x \in \mathbb{R}^n \mid u(x) = 0\},$$

where  $U$  is a smooth open set. Outside  $\bar{U}$ , we assume that  $\nabla u \neq 0$ . To fix the idea, we assume that  $\partial U$  has negative signature (orientation) in the sense that  $u < 0$  outside  $\bar{U}$ . By our specification of  $u$ , we see that

$$\partial_{D^{-1}}^0 TV(u) = \arg \min \left\{ \|\operatorname{div} Z\|_{D_0^1} \mid |Z| \leq 1 \text{ in } U, Z = \nabla u / |\nabla u| \text{ in } \bar{U}^c, \operatorname{div} Z \in D_0^1 \right\}.$$

Since  $\operatorname{div} Z$  is locally integrable,  $Z \cdot \nu$  does not jump across  $\partial U$ , where  $\nu$  is the exterior unit normal of  $\partial U$ . In this case,

$$Z \cdot \nu = Z \cdot \nabla u / |\nabla u| = -1 \quad \text{on } \partial U. \quad (1.11)$$

Since  $\nabla \operatorname{div} Z$  does not have a singular part,  $\operatorname{div} Z$  does not jump across  $\partial U$ . In this case,

$$\operatorname{div} Z = -\operatorname{div} \nu(x) \quad \text{on } \partial U. \quad (1.12)$$

However,  $v = \Delta \operatorname{div} Z$  may have a non-zero singular part concentrated on  $\partial U$  even if  $v = v_0$ , i.e.,  $v$  is the minimizer. This phenomenon is observed in [19], [20], [11] in a one-dimensional periodic setting. Different from the second-order problem, this causes expansion or shrinking of the ball when the solution  $u$  is of the form  $u(t, x) = a(t)\mathbf{1}_{B_{R(t)}}$ . If  $u > 0$  outside  $U$ , the minus in (1.11) (1.12) should be replaced by the plus.

If  $\partial_{D^{-1}}^0 TV(u)$  is constant on  $B_{R(t)}$  and  $(\overline{B_{R(t)}})^c$ , this property is preserved under the evolution, which leads us to definition of calibrability. We say that  $U$  (with negative signature) is *calibrable* if  $\Delta \operatorname{div} Z_0$  is constant on  $U$ , where  $Z_0$  belongs to

$$\arg \min \left\{ \|\nabla \operatorname{div} Z\|_{L^2(U)} \mid Z \text{ satisfies (1.11), (1.12) and } |Z| \leq 1, \text{ a.e. } x \in U \right\}. \quad (1.13)$$

Note that the value of  $\partial_{D^{-1}}^0 TV(u)$  on  $U$  is determined by  $U$  and its signature does not depend on particular value of  $u$ . In the second-order problem, we say that  $U$  (with negative signature) is calibrable if  $-\operatorname{div} \tilde{Z}_0$  is constant on  $U$  for

$$\tilde{Z}_0 = \arg \min \left\{ \|\operatorname{div} Z\|_{L^2(U)} \mid z \text{ satisfies (1.11) and } |Z| \leq 1, \text{ a.e. } x \in U \right\}$$

as in [16]. This is formally equivalent to *-calibrability* in [4], [3].

In our fourth-order problem, if one minimizes  $\|\nabla \operatorname{div} Z\|_{L^2(U)}$  under (1.12) but without the constraint  $|Z| \leq 1$ , a minimizer  $Z_1$  must satisfy

$$-\Delta w = \lambda \quad \text{in } U \quad (1.14)$$

$$w = -\operatorname{div} \nu \quad \text{on } \partial U \quad (1.15)$$

for  $w = \operatorname{div} Z_1$  with some constant  $\lambda$ . If  $U$  is bounded, this constant  $\lambda$  is determined by (1.11) since

$$\int_U w \, d\mathcal{L}^n = \int_U \operatorname{div} Z_1 \, d\mathcal{L}^n = \int_{\partial U} Z_1 \cdot \nu \, d\mathcal{H}^{n-1} = -\mathcal{H}^{n-1}(\partial U), \quad (1.16)$$

where  $\mathcal{H}^{n-1}$  denotes the  $n-1$  dimensional Hausdorff and  $\mathcal{L}^n$  denotes the Lebesgue measure in  $\mathbb{R}^n$ . For bounded  $U$ , we prove that if  $v = \Delta \operatorname{div} Z_2$  satisfies (1.14), (1.15), (1.16) for  $w = \operatorname{div} Z_2$ , and moreover it satisfies (1.11) and  $|Z_2| \leq 1$  in  $U$  a.e., then  $v$  is the minimizer and  $U$  is calibrable (Theorem 26). We call such vector field  $Z_2$  a calibration for  $U$ .

In the radially symmetric setting, it is not difficult to show that  $Z_0$  in (1.13) can be chosen in the form  $z(|x|) \frac{x}{|x|}$ . Indeed, if  $Z_0$  is belongs to the set of minimizers (1.13), then its rotational average  $\bar{Z}_0$  belongs to (1.13) as well, because averaging preserves (1.11), (1.12) and the inequality  $|Z| \leq 1$ . Since the angular part of  $\bar{Z}_0$  does not contribute to the divergence, it is possible to delete this part (Lemma 30). We thus conclude that there is an element of (1.13) of form  $Z(x) = z(|x|) \frac{x}{|x|}$ . Thus, the equation (1.14) can be written as the third-order ODE of the form

$$-r^{1-n} \left( r^{n-1} \left( r^{1-n} (r^{n-1} z)' \right)' \right)' = \lambda \quad (1.17)$$

since  $\operatorname{div} Z = r^{1-n}(r^{n-1}z)'$ . If  $U$  is  $B_R$  with negative signature, condition (1.11) implies

$$z(R) = -1. \quad (1.18)$$

Since  $\operatorname{div} Z = z' + (n-1)z/r$ , condition (1.12) implies that

$$z'(R) = 0. \quad (1.19)$$

Solving (1.17) under the assumption that  $z$  is smooth near zero under conditions (1.18), (1.19), we eventually get a unique solution (1.17)–(1.19) of the form

$$z(r) = \frac{1}{2} \left( \frac{r}{R} \right)^3 - \frac{3}{2} \frac{r}{R}, \quad \lambda = -\frac{n(n+2)}{R^3}$$

for all  $n \geq 1$ . It is easy to see that  $Z(x) = z(|x|) \frac{x}{|x|}$  satisfies the constraint  $|Z| \leq 1$  in  $B_R$ . We conclude that all balls are calibrable. More careful argument is necessary, but we are able to discuss calibrability of an annulus as well as a complement of a ball.

**Theorem 3.** (i) All balls are calibrable for all  $n \geq 1$ .

(ii) All complement of balls are calibrable except  $n = 2$ .

(iii) If  $n = 2$ , all complement of balls are not calibrable.

(iv) All annuli (with definite signature) are calibrable except in  $n = 2$ .

(v) For  $n = 2$ , there is  $Q_* > 1$  such that an annulus (with definite signature) is calibrable if and only if the ratio of the exterior radius over the interior radius is smaller than or equal to  $Q_*$ . In other words,  $A_{R_0}^{R_1}$  is calibrable if and only if  $R_1/R_0 \leq Q_*$ .

Theorem 3(v) is consistent with (iii) since  $R_1 \rightarrow \infty$  implies  $A_{R_0}^{R_1}$  converges to  $\overline{B_{R_0}}^c$ , a complement of the closure of the ball  $B_{R_0}$ . Note that in the case of an annulus, there is a possibility we take a signature which is different on the exterior boundary  $\partial B_{R_1}$  and the interior boundary  $\partial B_{R_0}$ . We also study such indefinite cases.

We now calculate an explicit solution of (1.1) starting from  $u_0 = a_0 \mathbf{1}_{B_{R_0}}$ . We first discuss the case  $n \neq 2$ . Since a ball and its complement is calibrable, the solution is of the form

$$u(t, x) = a(t) \mathbf{1}_{B_{R(t)}}. \quad (1.20)$$

We take the (radial) calibration  $Z_{in}$  in  $B_{R(t)}$  and  $Z_{out}$  in  $\mathbb{R}^n \setminus \overline{B_{R(t)}}$  and set

$$Z(x, t) = \begin{cases} Z_{in}(x), & x \in B_{R(t)} \\ Z_{out}(x), & x \in \mathbb{R}^n \setminus \overline{B_{R(t)}}. \end{cases}$$

Here  $Z_{out}(x) = z_{out}(|x|) \frac{x}{|x|}$  can be calculated as

$$z_{out}(r) = -\frac{n-1}{2} \left( \frac{r}{R} \right)^3 + \frac{n-3}{2} \left( \frac{r}{R} \right)^{1-n}$$

while, as we already discussed,  $z_{in}$  for  $Z_{in}(x) = z_{in}(|x|) \frac{x}{|x|}$  is of the form

$$z_{in}(r) = \frac{1}{2} \left( \frac{r}{R} \right)^3 - \frac{3}{2} \frac{r}{R}.$$

This  $Z$  satisfies (1.9) and (1.10), and moreover  $\operatorname{div} Z \in D_0^1$  for any  $t > 0$ . Moreover,  $\operatorname{div} Z$  is continuous across  $\partial B_{R(t)}$ . However,  $\nabla \operatorname{div} Z$  may jump across  $\partial B_{R(t)}$ . Actually,

$$-\Delta \operatorname{div} Z = \lambda \mathbf{1}_{B_{R(t)}} + \nu \cdot (\nabla \operatorname{div} Z_{in} - \nabla \operatorname{div} Z_{out}) \delta_{\partial B_{R(t)}},$$

where  $\delta_\Gamma(\varphi) = \int_\Gamma \varphi d\mathcal{H}^{n-1}$  or  $\delta_\Gamma = \mathcal{H}^{n-1} \llcorner \Gamma$  for a hypersurface  $\Gamma$  and  $\nu$  is the exterior unit normal of  $\partial B_{R(t)}$ , i.e.,  $\nu = x/R(t)$ . Here  $\lambda = -\frac{n(n+2)}{R^3}$ . By a direct calculation, the quantity  $\nu \cdot (\nabla \operatorname{div} Z_{in} - \nabla \operatorname{div} Z_{out}) = -\frac{n(n-4)}{R^2}$ . Since  $u_t = -\Delta \operatorname{div} Z$ , by

$$\partial_t (a \mathbf{1}_{B_R}) = \frac{da}{dt} \mathbf{1}_{B_R} + a \frac{dR}{dt} \delta_{\partial B_R},$$

we conclude that

$$\frac{da}{dt} = -\frac{n(n+2)}{R^3}, \quad \frac{dR}{dt} = -\frac{n(n-4)}{aR^2}.$$

Since

$$\frac{d}{dt}(aR^3) = -n(n+2) - 3n(n-4) = -n(4n-10),$$

an explicit form of a solution is given as

$$a(t) = a_0 \left(1 - \frac{n(4n-10)}{a_0 R_0^3} t\right)^{\frac{n+2}{4n-10}}, \quad R(t) = R_0 \left(1 - \frac{n(4n-10)}{a_0 R_0^3} t\right)^{\frac{n-4}{4n-10}}.$$

As noticed earlier, in the case  $n = 2$ , the complement of the disk is not calibrable. If  $u$  is a radially strictly decreasing function outside  $B_R$ , we expect  $Z_{out}(x) = -x/|x|$  for  $|x| > R(t)$ . In [14], it is proposed that a solution  $u$  to (1.1) must satisfy

$$u_t = -\Delta \operatorname{div} Z_{out}.$$

Since  $\operatorname{div} Z_{out} = -(n-1)/|x|^2$  and  $\nabla \operatorname{div} Z_{out} = \frac{(n-1)x}{|x|^3}$ , this implies

$$u_t(t, x) = -\frac{(n-1)(n-3)}{|x|^3}, \quad x \in (\overline{B_{R(t)}})^c = \mathbb{R}^n \setminus \overline{B_{R(t)}}. \quad (1.21)$$

In the case  $n = 2$ ,  $\nabla \operatorname{div} Z_{out} \in L^2\left(\left(\overline{B_{R(t)}}\right)^c\right)$  so  $Z_{out}$  is a Cahn-Hoffman vector field.

If we start with  $u_0 = a_0 \mathbf{1}_{B_{R_0}}$  with  $a_0 > 0$  for  $n = 2$ , the expected form of a solution is

$$u_t(t, x) = a(t) \mathbf{1}_{B_{R(t)}} + \frac{t}{|x|^3} \mathbf{1}_{\overline{B_{R(t)}}^c}, \quad (1.22)$$

where

$$\frac{da}{dt} = -\frac{2 \cdot 4}{R^3}, \quad \left(a(t) - \frac{t}{R(t)^3}\right) \frac{dR}{dt} = \frac{2 \cdot 2}{R^2}. \quad (1.23)$$

Analyzing this ODE system, we can deduce qualitative properties of the solution. Summing up our results yields

**Theorem 4.** *Let  $u_0 = a_0 \mathbf{1}_{B_{R_0}}$  with  $a_0 > 0$ .*

*If  $n \neq 3$ , then the solution  $u$  to (1.1) with initial datum  $u_0$  is of the form*

$$u(t, x) = a(t) \mathbf{1}_{B_{R(t)}} \quad \text{for } t < t_* = a_0 R_0^3 / (n(4n-10))$$

*and  $u(t, x) \equiv 0$  for  $t \geq t_*$ . (The time  $t_*$  is called the extinction time.) Moreover,  $a(t)$  is decreasing and  $a(t) \rightarrow 0$  as  $t \uparrow t_*$ .*

- (i)  $R(t)$  is increasing and  $R(t) \rightarrow \infty$  as  $t \uparrow t_*$  for  $n = 3$ .
- (ii)  $R(t) = R_0$  for  $n = 4$ .
- (iii)  $R(t)$  is decreasing and  $R(t) \rightarrow 0$  as  $t \uparrow t_*$  for  $n \geq 5$ .

*If  $n = 2$ , then the solution is not a characteristic function for  $t > 0$ . It is of the form (1.21) and moves by (1.23). In particular, there is no extinction time,  $R(t)$  is increasing and  $a(t)$  is decreasing. Moreover,  $R(t) \rightarrow \infty$  and  $a(t) \rightarrow 0$  as  $t \rightarrow \infty$ . The gap  $a(t) - \frac{t}{R(t)^3}$  is always positive.*

*If  $n = 1$ , then the solution is of the form  $u(t, x) = a(t) \mathbf{1}_{B_{R(t)}}$  for  $t > 0$ . Moreover,  $R(t)$  is increasing and  $a(t)$  is decreasing with  $R(t) \rightarrow \infty$  and  $a(t) \rightarrow 0$  as  $t \rightarrow \infty$ .*

We note that the infinite extinction time observed in  $n \leq 2$  is related to the fact that 0 is not an element of the affine space  $u_0 + D^{-1}$  where the flow lives if  $\int u_0 \neq 0$ . In [13], finite time extinction for solution to (1.1) is proved in a periodic setting for average zero initial data when the space dimension  $n \leq 4$ . Our result is unrelated to their result because we consider (1.1) in  $\mathbb{R}^n$ .

The formula (1.21) does not give a solution to (1.1) when  $n \geq 4$  since  $\nabla \operatorname{div} Z_{out}$  does not belong to  $L^2\left(\overline{(B_{R(t)})^c}\right)$ . In the case  $n = 3$ , this formula is consistent with our definition. If we consider  $u_0$  strictly radially decreasing for  $|x| > R_0$  and  $u_0(x) = u_*$  for  $|x| \leq R_0$ , then  $u_0$  does not belong to the domain of  $\partial_{D^{-1}}TV$  for  $n \geq 4$ . In other words, there is no Cahn-Hoffman vector field.

These results contrast with the second-order total variation flow

$$u_t = \operatorname{div}(\nabla u / |\nabla u|).$$

In the second-order problem, a ball and an annulus are always calibrable with their complements, see e.g. [3] or [16, Section 5]. Furthermore,  $u_t(t, \cdot)$  is a locally integrable function without singular part for  $t > 0$ . Thus, for example, the solution starting from  $u_0 = a_0 \mathbf{1}_{B_{R_0}}$  ( $a_0 > 0$ ) must be  $u(t, x) = a(t) \mathbf{1}_{B_{R_0}}$  with  $a(t) = -\lambda t + a_0$ , where  $\lambda$  is the Cheeger ratio, i.e.  $\lambda = \mathcal{H}^{n-1}(\partial B_{R_0}) / \mathcal{L}^n(B_{R_0})$ . In particular, the extinction time  $t_*$  equals  $t_* = a_0 / \lambda$ .

We conclude this paper by deriving a system of ODEs prescribing the solution in the case when the initial datum is a piecewise constant, radially symmetric function, which we call a stack. To be precise, we say that  $w \in E^{-1}$  is a stack if it is of the form

$$w = a^0 \mathbf{1}_{B_{R^0}} + a^1 \mathbf{1}_{A_{R^0}^{R^1}} + \dots + a^{N-1} \mathbf{1}_{A_{R^{N-2}}^{R^{N-1}}} + a^N \mathbf{1}_{\mathbb{R}^n \setminus B_{R^{N-1}}},$$

$0 < R^0 < R^1 < \dots < R^{N-1}$ ,  $a^k \in \mathbb{R}$ . In particular, we obtain

**Theorem 5.** *Let  $n \neq 2$  and let  $u_0$  be a stack. If  $u$  is the solution to (1.1), then  $u(t, \cdot)$  is a stack for  $t > 0$ .*

In the case  $n = 2$ , this result is no longer true, as evidenced by Theorem 4. However the solution can still be prescribed by a finite system of ODEs.

A total variation flow type equation

$$w_t = -\Delta \left( \operatorname{div} \frac{\nabla w}{|\nabla w|} + \beta \operatorname{div}(\nabla w |\nabla w|) \right) \quad (1.24)$$

was introduced by [30] to describe the height of crystal surface moved by relaxation dynamics below the roughening temperature, where  $\beta > 0$ . For this equation, characterization of the subdifferential of the corresponding energy was given by Y. Kashima in periodic setting [19], [20] and under Dirichlet condition on a bounded domain [20]. The speed of a facet (a flat part of the graph) is calculated for  $n = 1$  in [19] and for a ball with the Dirichlet condition under radial symmetry [20]. Different from the second-order problem, the speed of a facet is determined not only by the shape of facet. Also it has been already observed in [19], that the minimal section may have a delta part although the behavior of the corresponding solution was not studied there. A numerical computation was given in [22]. The equation (1.24) was derived as a continuum limit of models describing motion of steps on crystal surface as discussed in [26], where numerical simulation was given; see also [21].

In [7], a crystalline diffusion flow was proposed and calculated numerically. In a special case, it is of the form  $w_t = -\partial_x^2(W'(w_x))$ , where  $W$  is a piecewise linear convex function, when the curve is given as the graph of a function. This equation was analyzed in [12] in a class of piecewise linear (in space) solutions.

Fourth-order equations of type (1.1) were proposed for image denoising as an improvement over the second-order total variation flow. For example, the equation

$$w_t = -\Delta \operatorname{div}(\nabla w / |\nabla w|) + \lambda(f - w),$$

where  $f$  is an original image which is given and  $\lambda > 0$ , corresponds to the Osher-Solé-Vese model [27]. The well-posedness of this equation was proved by using the Galerkin method by [8].

For (1.1), an extinction time estimate was given in [13] for  $n = 1, 2, 3, 4$  in the periodic setting. It was extended to the Dirichlet problem in a bounded domain by [14]. In the review paper [11], it was proved that the solution  $u$  of (1.1) in  $n = 1$  may become discontinuous instantaneously even if the initial datum is Lipschitz continuous, because the speed may have a delta part.

There are a few numerical studies for (1.1) in the periodic setting. A duality-based numerical scheme which applies the forward-backward splitting has been proposed in [15]. A split Bregman method was adjusted to (1.1) and also (1.24) in [17]. In these methods, the singularity of the equation at  $\nabla u = 0$  is not regularized. However, all above studies deal with either periodic, Dirichlet or Neumann boundary condition for a bounded domain. It has never been rigorously studied in  $\mathbb{R}^n$ , although in [14] there are some preliminary calculations for radial solution in  $\mathbb{R}^n$ .

This paper is organized as follows. In Section 2, we discuss basic properties of the total variation on  $D^{-1}$ , notably we show strict density of  $C_{c,av}^\infty$ . In Section 3, we characterize the subdifferential of  $TV$  in  $D^{-1}$ , adjusting a duality argument in [3]. In Section 4, we give a rigorous definition of a solution to (1.1) in  $\mathbb{R}^n$ . In Section 5, we introduce the notion of calibrability. In Section 6, we discuss calibrability of rotationally symmetric sets in  $\mathbb{R}^n$ . In Section 7, we study solutions emanating from piecewise constant, radially symmetric data.

## 2 The total variation functional on $D^{-1}$

In this section, we give a rigorous definition of the total variation  $TV$  on  $D^{-1}$  and relate it to the usual total variation defined on  $L^1_{loc}$ . The main tool that we use here as well as in the following section is an approximation lemma, which for a given  $w \in D^{-1}$  produces a sequence of nice functions  $w_k \in D^{-1}$  that converges to  $w$  in  $D^{-1}$  and  $TV(w_k) \rightarrow TV(w)$ .

Let us denote

$$X_1 = \{ \psi \in C_c^\infty(\mathbb{R}^n, \mathbb{R}^n), \|\psi\|_{L^\infty(\mathbb{R}^n, \mathbb{R}^n)} \leq 1 \}.$$

We define  $TV: D^{-1}(\mathbb{R}^n) \rightarrow [0, \infty]$  by

$$TV(u) = \sup_{\psi \in X_1} \langle u, \operatorname{div} \psi \rangle.$$

Let us compare this definition with the usual total variation, which we denote here by  $\overline{TV}: L^1_{loc}(\mathbb{R}^n) \rightarrow [0, \infty]$ , defined by

$$\overline{TV}(u) = \sup_{\psi \in X_1} \int_{\mathbb{R}^n} u \operatorname{div} \psi.$$

First of all, as in the case  $\overline{TV}$ , we easily check that  $TV$  is lower semicontinuous with respect to the weak-\* (and, a fortiori, strong) convergence in  $D^{-1}(\mathbb{R}^n)$ . Indeed, if  $v_k \xrightarrow{*} v$  in  $D^{-1}(\mathbb{R}^n)$ ,

$$TV(v) = \sup_{\psi \in X_1} \{ \langle v, \operatorname{div} \psi \rangle \} = \sup_{\psi \in X_1} \liminf_{k \rightarrow \infty} \{ \langle v_k, \operatorname{div} \psi \rangle \} \leq \liminf_{k \rightarrow \infty} \sup_{\psi \in X_1} \{ \langle v_k, \operatorname{div} \psi \rangle \} = \liminf_{k \rightarrow \infty} TV(v_k).$$

In fact, we have

**Lemma 6.** *We have  $D(TV) \subset L^1_{loc}$ , and so  $D(TV) \subset D(\overline{TV})$  with  $TV$  and  $\overline{TV}$  coinciding on  $D(TV)$ . In particular, if  $n \geq 2$ ,  $D(TV) \subset L^{1^*}(\mathbb{R}^n)$ . If  $n = 1$ ,*

$$D(TV) \subset L^\infty_0(\mathbb{R}) = \{ w \in L^\infty(\mathbb{R}): \lim_{x \rightarrow \pm\infty} w(x) = 0 \}.$$

The proof of this fact is a consequence of the lemma below and we postpone it.

**Lemma 7.** *For any  $w \in D^{-1}(\mathbb{R}^n)$  there exists a sequence  $w_k \in C_{c,av}^\infty(\mathbb{R}^n)$  such that*

$$w_k \rightarrow w \text{ in } D^{-1}(\mathbb{R}^n)$$

and

$$TV(w_k) \rightarrow TV(w).$$

To prove it, we will use a special choice of cut-off function and associated variant of the Sobolev-Poincaré inequality. For  $R > 0$ , let us denote by  $\vartheta_R$  the element of minimal norm in  $D_0^1(\mathbb{R}^n)$  among those  $w \in D_0^1(\mathbb{R}^n)$  that satisfy  $w(x) = 1$  if  $|x| \leq \frac{R}{2}$ ,  $w(x) = 0$  if  $|x| \geq R$ . It is an easy exercise to show that for  $\frac{R}{2} \leq |x| \leq R$

$$\vartheta_R(x) = (2^{n-2} - 1)^{-1} \left( \left( \frac{|x|}{R} \right)^{2-n} - 1 \right) \text{ if } n \neq 2, \quad \vartheta_R(x) = \frac{\log \frac{R}{|x|}}{\log 2} \text{ if } n = 2.$$

In either case,

$$\nabla \vartheta_R(x) = C_n \frac{|x|^{-n} x}{R^{2-n}} \text{ if } \frac{R}{2} \leq |x| \leq R. \quad (2.1)$$

**Lemma 8.** *If  $p \in [1, n[$  and  $q \in [1, p^*]$ , then for all  $w \in C^1(\mathbb{R}^n)$ ,  $R > 0$  there holds*

$$\left\| w - \frac{\int \vartheta_R w}{\int \vartheta_R} \right\|_{L^q(B_R)} \leq C R^{1 + \frac{n}{q} - \frac{n}{p}} \|\nabla w\|_{L^p(B_R)} \quad (2.2)$$

with  $C = C(n, p)$  and

$$\|\nabla \vartheta_R\|_{L^p(\mathbb{R}^n)} = C R^{-\frac{1}{p}(n-1)(2-p)} \quad (2.3)$$

with a different  $C = C(n, p)$ .

*Proof.* Let  $v \in C^1(\mathbb{R}^n)$ . Following the proof of the standard Poincaré inequality by contradiction using Rellich-Kondrachov theorem, we obtain

$$\left\| v - \frac{\int \vartheta_1 v}{\int \vartheta_1} \right\|_{L^p(B_1)} \leq C \|\nabla v\|_{L^p(B_1)}.$$

Applying the Sobolev inequality in  $B_1$  to the function  $v - \frac{\int \vartheta_1 v}{\int \vartheta_1}$ , we upgrade this to

$$\left\| v - \frac{\int \vartheta_1 v}{\int \vartheta_1} \right\|_{L^q(B_1)} \leq C \|\nabla v\|_{L^p(B_1)} \quad (2.4)$$

Next, let  $v(x) = w(Rx)$  for a given  $w \in C^1(\mathbb{R}^n)$ . We observe that

$$\vartheta_1(x) = \vartheta_R(Rx) \quad \text{for } x \in \mathbb{R}^n$$

and so, by a change of variables  $x = y/R$ ,

$$\int \vartheta_1 = \frac{1}{R^n} \int \vartheta_R, \quad \int \vartheta_1 v = \frac{1}{R^n} \int \vartheta_R w.$$

Applying the same change of variables to both sides of (2.4) we conclude the proof of (2.2).

The proof of (2.3) is a matter of direct calculation.  $\square$

Let us now return to the proof of the approximation lemma.

*Proof of Lemma 7.* Given  $w \in D^{-1}$ , let

$$w_{\varepsilon, R} = \left( \varrho_\varepsilon * w - \frac{\int \vartheta_R \varrho_\varepsilon * w}{\int \vartheta_R} \right) \vartheta_R.$$

Equivalently, for  $\varphi \in D_0^1(\mathbb{R}^n)$ ,

$$\langle w_{\varepsilon, R}, \varphi \rangle = \left\langle w, \varrho_\varepsilon * \left( \left( \varphi - \frac{\int \vartheta_R \varphi}{\int \vartheta_R} \right) \vartheta_R \right) \right\rangle.$$

Denoting  $\tilde{w} = (-\Delta)^{-1}w$ ,

$$\begin{aligned} \langle w_{\varepsilon,R} - w, \varphi \rangle &= \int \nabla \varrho_\varepsilon * \tilde{w} \cdot \nabla \left( \left( \varphi - \frac{\int \vartheta_R \varphi}{\int \vartheta_R} \right) \vartheta_R - \varphi \right) + \int (\nabla \varrho_\varepsilon * \tilde{w} - \nabla \tilde{w}) \cdot \nabla \varphi \\ &= \int \nabla \varrho_\varepsilon * \tilde{w} \cdot (\vartheta_R - 1) \nabla \varphi + \int \nabla \varrho_\varepsilon * \tilde{w} \cdot \left( \varphi - \frac{\int \vartheta_R \varphi}{\int \vartheta_R} \right) \nabla \vartheta_R + \int (\nabla \varrho_\varepsilon * \tilde{w} - \nabla \tilde{w}) \cdot \nabla \varphi. \end{aligned}$$

We estimate the second term on the r. h. s. using the Poincaré inequality from Lemma 8, taking into account that the support of the integrand is contained in  $\bar{A}_R$ , where  $A_R = B_R \setminus \bar{B}_{R/2}$ ,

$$\begin{aligned} \left| \int \nabla \varrho_\varepsilon * \tilde{w} \cdot \left( \varphi - \frac{\int \vartheta_R \varphi}{\int \vartheta_R} \right) \nabla \vartheta_R \right| &\leq C \|\nabla \varrho_\varepsilon * \tilde{w} \mathbf{1}_{A_R}\|_{L^2(\mathbb{R}^n)} \left\| \varphi - \frac{\int \vartheta_R \varphi}{\int \vartheta_R} \right\|_{L^2(\mathbb{R}^n)} \|\nabla \vartheta_R\|_{L^\infty(\mathbb{R}^n)} \\ &\leq C \|\nabla \varrho_\varepsilon * \tilde{w} \mathbf{1}_{A_R}\|_{L^2(\mathbb{R}^n)} \|\nabla \varphi\|_{L^2(\mathbb{R}^n)}. \end{aligned}$$

Thus,

$$\begin{aligned} \|w_{\varepsilon,R} - w\|_{D^{-1}(\mathbb{R}^n)} &= \sup_{\|\varphi\|_{D^1_0(\mathbb{R}^n)} \leq 1} \langle w_{\varepsilon,R} - w, \varphi \rangle \\ &\leq \|(1 - \vartheta_R) \nabla \varrho_\varepsilon * \tilde{w}\|_{L^2(\mathbb{R}^n)} + C \|\nabla \varrho_\varepsilon * \tilde{w} \mathbf{1}_{A_R}\|_{L^2(\mathbb{R}^n)} + \|\nabla \varrho_\varepsilon * \tilde{w} - \nabla \tilde{w}\|_{L^2(\mathbb{R}^n)} \end{aligned}$$

and so

$$\lim_{\varepsilon \rightarrow 0^+} \lim_{R \rightarrow \infty} \|w_{\varepsilon,R} - w\|_{D^{-1}(\mathbb{R}^n)} = 0. \quad (2.5)$$

Next we estimate  $TV(w_{\varepsilon,R})$ . Due to lower semicontinuity of  $TV$ , we can assume without loss of generality that  $TV(w) < \infty$ . First, we note that  $\varrho_\varepsilon * w \in D^{-1}(\mathbb{R}^n) \cap C^\infty(\mathbb{R}^n)$  ([29], [28]) and

$$\|\varrho_\varepsilon * \psi\|_{L^\infty(\mathbb{R}^n, \mathbb{R}^n)} \leq \|\psi\|_{L^\infty(\mathbb{R}^n, \mathbb{R}^n)}.$$

Thus, for any  $\psi \in L^\infty(\mathbb{R}^n, \mathbb{R}^n)$ ,

$$TV(\varrho_\varepsilon * w) = \sup_{\psi \in X_1} \{\langle w, \operatorname{div} \varrho_\varepsilon * \psi \rangle\} \leq TV(w).$$

In particular, this implies that  $\nabla \varrho_\varepsilon * w \in L^1(\mathbb{R}^n)$  for  $\varepsilon > 0$  and  $\int |\nabla \varrho_\varepsilon * w| = TV(\varrho_\varepsilon * w) \leq TV(w)$ . By Lemma 8,

$$\begin{aligned} TV(w_{\varepsilon,R}) &= \int |\nabla w_{\varepsilon,R}| \leq \int \vartheta_R |\nabla \varrho_\varepsilon * w| + \int \left| \left( \varrho_\varepsilon * w - \frac{\int \vartheta_R \varrho_\varepsilon * w}{\int \vartheta_R} \right) \nabla \vartheta_R \right| \\ &\leq \int |\nabla \varrho_\varepsilon * w| + C \left\| \varrho_\varepsilon * w - \frac{\int \vartheta_R \varrho_\varepsilon * w}{\int \vartheta_R} \right\|_{L^{1^*}(\mathbb{R}^n)} \|\nabla \vartheta_R\|_{L^n(\mathbb{R}^n)} \\ &\leq \|\nabla \varrho_\varepsilon * w\|_{L^1(\mathbb{R}^n)} + C \|\nabla \varrho_\varepsilon * w\|_{L^1(\mathbb{R}^n)} R^{-\frac{(n-1)(n-2)}{n}} \leq \left( 1 + CR^{-\frac{(n-1)(n-2)}{n}} \right) TV(w). \quad (2.6) \end{aligned}$$

If  $n \geq 3$ , together with lower semicontinuity of  $TV$ , this yields

$$\lim_{(\varepsilon, R) \rightarrow (0, \infty)} TV(w_{\varepsilon,R}).$$

Taking into account (2.5), by a diagonal procedure we can select sequences  $(\varepsilon_k)$ ,  $(R_k)$  such that  $w_k := w_{\varepsilon_k, R_k}$  satisfies both requirements in the assertion. On the other hand, if  $n = 1$  or  $n = 2$ , (2.6) only implies uniform boundedness of  $TV(\nabla w_{\varepsilon,R})$ .

Let now  $n = 2$ . Since  $w_{\varepsilon,R}$  have compact support, uniform boundedness of  $TV(\nabla w_{\varepsilon,R})$  implies uniform bound on  $w_{\varepsilon,R}$  in  $L^2(\mathbb{R}^n)$  by the Sobolev inequality. As  $w_{\varepsilon,R}$  converges to  $\varrho_\varepsilon * w$  in

$D^{-1}(\mathbb{R}^n)$ , we have  $\varrho_\varepsilon * w \in L^2(\mathbb{R}^n)$ . This allows us to improve (2.6):

$$\begin{aligned} TV(w_{\varepsilon,R}) &= \int |\nabla w_{\varepsilon,R}| \leq \int \vartheta_R |\nabla \varrho_\varepsilon * w| + \int |\varrho_\varepsilon * w \nabla \vartheta_R| + \frac{\int \vartheta_R \varrho_\varepsilon * w}{\int \vartheta_R} \int |\nabla \vartheta_R| \\ &\leq \|\nabla \varrho_\varepsilon * w\|_{L^1(\mathbb{R}^2)} + C \|\varrho_\varepsilon * w \mathbf{1}_{A_R}\|_{L^2(\mathbb{R}^2)} \|\nabla \vartheta_R\|_{L^2(\mathbb{R}^2)} + C \frac{\|\vartheta_R\|_{D_0^1(\mathbb{R}^2)} \|\varrho_\varepsilon * w\|_{D^{-1}(\mathbb{R}^2)}}{\|\vartheta_R\|_{L^1(\mathbb{R}^2)}} \|\nabla \vartheta_R\|_{L^1(\mathbb{R}^2)} \\ &\leq TV(w) + C \|\varrho_\varepsilon * w \mathbf{1}_{A_R}\|_{L^2(\mathbb{R}^2)} + \frac{C}{R} \|\varrho_\varepsilon * w\|_{D^{-1}(\mathbb{R}^2)}. \end{aligned} \quad (2.7)$$

The r. h. s. of (2.7) converges to  $TV(w)$  as  $R \rightarrow \infty$  and we conclude as before.

Next, consider  $n = 1$ . In this case, finiteness of  $TV(w)$  implies that  $\varrho_\varepsilon * w \in L^\infty(\mathbb{R}^n)$  and there exist  $g_\varepsilon^\pm \in \mathbb{R}$  such that

$$\lim_{x \rightarrow \pm\infty} \varrho_\varepsilon * w(x) = g_\varepsilon^\pm.$$

Now, let  $\eta_R^\pm$  be the element of minimal norm in  $D_0^1(\mathbb{R})$  under constraints

$$\eta_R^\pm(\pm x) = 1 \text{ if } x \in [2R, 3R], \quad \eta_R^\pm(\pm x) = 0 \text{ if } x \notin [R, 4R].$$

(Clearly,  $\eta_R^\pm$  is a continuous, piecewise affine function.) We have

$$|\langle \varrho_\varepsilon * w, \eta_R^\pm \rangle| \leq \|\varrho_\varepsilon * w\|_{D^{-1}(\mathbb{R})} \|\eta_R^\pm\|_{D_0^1(\mathbb{R})} \rightarrow 0 \text{ as } R \rightarrow \infty.$$

On the other hand, since  $\eta_R^\pm$  are compactly supported and  $\varrho_\varepsilon * w$  coincides as distribution with a locally integrable function, we can calculate

$$\langle \varrho_\varepsilon * w, \eta_R^\pm \rangle = \int \varrho_\varepsilon * w \eta_R^\pm \rightarrow \infty \cdot g_\varepsilon^\pm \text{ as } R \rightarrow \infty,$$

so  $g_\varepsilon^\pm = 0$ . Therefore, we can estimate

$$\begin{aligned} TV(w_{\varepsilon,R}) &= \int |\nabla w_{\varepsilon,R}| \leq \int \vartheta_R |\nabla \varrho_\varepsilon * w| + \int |\varrho_\varepsilon * w \nabla \vartheta_R| + \frac{\int \vartheta_R \varrho_\varepsilon * w}{\int \vartheta_R} \int |\nabla \vartheta_R| \\ &\leq TV(w) + \frac{2}{R} \int_{A_R} |\varrho_\varepsilon * w| + 2 \frac{\int \vartheta_R \varrho_\varepsilon * w}{\int \vartheta_R}. \end{aligned} \quad (2.8)$$

Since we have shown that  $\varrho_\varepsilon * w(x) \rightarrow 0$  as  $x \rightarrow \pm\infty$ , the averages on the r. h. s. converge to 0 and we conclude as before.  $\square$

As a first application of the approximation lemma, we demonstrate Lemma 6 announced before.

*Proof of Lemma 6.* Let  $w \in D(TV)$  and let  $(w_k) \subset C_{c,av}^\infty(\mathbb{R}^n)$  be the sequence provided by Lemma 7. Let first  $n > 1$ . Since  $\nabla w_k$  is uniformly bounded in  $L^1(\mathbb{R}^n, \mathbb{R}^n)$ , by the Sobolev embedding  $w_k$  is uniformly bounded in  $L^{1^*}(\mathbb{R}^n)$ . Therefore,  $w \in L^{1^*}(\mathbb{R}^n)$ .

In case  $n = 1$ , since  $\nabla w_k$  are compactly supported and uniformly bounded in  $L^1(\mathbb{R})$ ,  $w_k$  is uniformly bounded in  $L^\infty(\mathbb{R})$ . From these two bounds it follows that  $w \in L^\infty(\mathbb{R})$ ,  $\nabla w \in M(\mathbb{R})$ ,  $w_k \rightarrow w$  in  $L_{loc}^1(\mathbb{R})$ . Identifying  $w$  with its semicontinuous representative, there exist  $g^\pm$  such that  $w(x) \rightarrow g^\pm$  as  $x \rightarrow \pm\infty$ . Suppose that  $g^+ \neq 0$  and let  $M$  be such that

$$\int_{]M, +\infty[} |\nabla w| < \frac{|g^+|}{2}.$$

In particular,  $|w(x)| > \frac{|g^+|}{2}$  for  $x > M$ . Take any  $x_0 > M$  such that  $w_k(x_0) \rightarrow w(x_0)$  and  $|\nabla w|(\{x_0\}) = 0$ . Then

$$\begin{aligned} \liminf_{k \rightarrow \infty} \int_{\mathbb{R}} |\nabla w_k| &\geq \liminf_{k \rightarrow \infty} \int_{]-\infty, x_0[} |\nabla w_k| + \liminf_{k \rightarrow \infty} \int_{]x_0, \infty[} |\nabla w_k| \\ &\geq \int_{]-\infty, x_0[} |\nabla w| + \liminf_{k \rightarrow \infty} |w_k(x_0)| = \int_{\mathbb{R}} |\nabla w| - \int_{]x_0, \infty[} |\nabla w| + |w(x_0)| > \int_{\mathbb{R}} |\nabla w| \end{aligned}$$

which contradicts strict convergence of  $w_k$ , hence  $g^+ = 0$ . Similarly we prove that  $g^- = 0$ .  $\square$

### 3 Subdifferential of the total variation

We give a characterization of the total variation in the space  $D^{-1} = D^{-1}(\mathbb{R}^n)$ . The basic idea of the proof is a duality argument, which has been carried out in the case of  $L^2$  subdifferentials. In the case of  $L^2$  setting, the idea goes back to the unpublished note of F. Alter and a detailed proof is given in [3]. Let  $\mathcal{E}$  be a functional on a real Hilbert space  $H$  equipped with an inner product  $(\cdot, \cdot)_H$ . The main idea is to characterize the subdifferential  $\partial\mathcal{E}$  by the polar  $\mathcal{E}^0$  of  $\mathcal{E} : H \rightarrow [-\infty, \infty]$  is defined by

$$\mathcal{E}^0(v) := \sup \{(u, v)_H \mid u \in H, \mathcal{E}(u) \leq 1\} = \sup \{(u, v)/\mathcal{E}(u) \mid u \in D(\mathcal{E}), \mathcal{E}(u) \neq 0\},$$

where  $D(\mathcal{E}) = \{u \in H \mid |\mathcal{E}(u)| < \infty\}$ . We first recall a lemma [3, Lemma 1.7].

**Lemma 9.** *Let  $\mathcal{E}$  be convex. Assume that  $\mathcal{E}$  is positively one-homogeneous, i.e.,*

$$\mathcal{E}(\lambda u) = \lambda \mathcal{E}(u)$$

for all  $\lambda > 0, u \in H$ . Then,  $v \in \partial\mathcal{E}(u)$  if and only if  $\mathcal{E}^0(v) \leq 1$  and  $(u, v)_H = \mathcal{E}(u)$ .

*Proof.* We give a proof for completeness. Since  $\mathcal{E}$  is convex and positively one-homogeneous,

$$\mathcal{E}(u + w) \leq \mathcal{E}(u) + \mathcal{E}(w). \quad (3.1)$$

Indeed,

$$\frac{1}{2}\mathcal{E}(u + w) = \mathcal{E}((u + w)/2) \leq (\mathcal{E}(u) + \mathcal{E}(w))/2.$$

“Only if” part. If  $v \in \partial\mathcal{E}(u)$ , then

$$\mathcal{E}(w) \geq (v, w)_H \quad \text{for all } w \in H \quad (3.2)$$

which is equivalent to  $\mathcal{E}^0(v) \leq 1$ . Indeed, by (3.1) and the definition of  $v$ , we see

$$\mathcal{E}(w) \geq \mathcal{E}(u + w) - \mathcal{E}(u) \geq (v, w)_H \quad (3.3)$$

for all  $w \in H$ .

The second identity  $(u, v)_H = \mathcal{E}(u)$  is nothing but the Euler equation for a homogeneous functional. We take  $w = u$  in (3.3) to get

$$\mathcal{E}(u) \geq (v, u)_H.$$

If we take  $w = -u/2$ , then

$$\frac{1}{2}\mathcal{E}(u) - \mathcal{E}(u) \geq -(v, u)_H/2,$$

which implies  $\mathcal{E}(u) \leq (v, u)_H$ . Thus  $\mathcal{E}(u) = (v, u)_H$ .

“If” part. Since  $\mathcal{E}^0(v) \leq 1$ , we see

$$\mathcal{E}(u+w) \geq (v, u+w)_H.$$

Thus,

$$\mathcal{E}(u+w) - \mathcal{E}(u) \geq (v, u+w)_H - \mathcal{E}(u) = (v, u+w)_H - (v, u)_H = (v, w)_H$$

by the Euler equation  $\mathcal{E}(u) = (v, u)_H$ , i. e.  $v \in \partial\mathcal{E}(u)$ .  $\square$

**Remark 10.** *By general theory of convex functionals, we know that*

$$(\mathcal{E}^0)^0 = \mathcal{E}$$

if  $\mathcal{E}$  is a lower semicontinuous, convex, positively one-homogeneous functional [3, Proposition 1.6].

This property is essential for the proof of

**Theorem 11.** *Let  $\Psi: D^{-1} \rightarrow [0, \infty]$  be defined by*

$$\Psi(v) = \inf \{ \|z\|_\infty \mid v = \Delta \operatorname{div} Z, Z \in L^\infty(\mathbb{R}^n), \operatorname{div} Z \in D_0^1 \}.$$

Then  $(TV)^0 = \Psi$ .

**Remark 12.** (i) *By definition,  $\Psi$  is a convex, lower semi-continuous, positively one-homogeneous function, so  $(\Psi^0)^0 = \Psi$ .*

(ii) *if  $\Psi(v) < \infty$ , the infimum is attained. Theorem 11 together with Lemma 9 implies the following characterization of the subdifferential of  $TV$ .*

**Theorem 13.** *An element  $v \in D^{-1}$  belongs to  $\partial TV(u)$  if and only if there is  $Z \in L^\infty(\mathbb{R}^n)$  with  $\operatorname{div} Z \in D_0^1$  such that*

$$(i) \quad |Z| \leq 1$$

$$(ii) \quad v = \Delta \operatorname{div} Z$$

$$(iii) \quad -\langle u, \operatorname{div} Z \rangle = TV(u).$$

*Proof.* By Lemma 9 and Theorem 11,

$$v \in \partial TV(u) \iff \Psi(v) \leq 1 \text{ and } (v, u)_{D^{-1}} = TV(u).$$

The property  $\Psi(v) \leq 1$  together with Remark 12(ii) implies (i), (ii) and  $\operatorname{div} Z \in D_0^1$ .

$$(v, u)_{D^{-1}} = \langle u, (-\Delta)^{-1}v \rangle = -\langle u, \operatorname{div} Z \rangle.$$

It is not difficult to check the converse.  $\square$

*Proof of Theorem 11.* The inequality  $TV^0 \leq \Psi$ :

We take  $v \in D^{-1}$  with  $\Psi(v) < \infty$ . By Remark 12(ii), there is  $Z \in L^\infty(\mathbb{R}^n)$  with  $v = \Delta \operatorname{div} Z$  with  $\operatorname{div} Z \in D_0^1$  such that  $\Psi(v) = \|Z\|_\infty$ . By Lemma 7, there is  $u_k \in C_{c,av}^\infty$  such that  $TV(u_k) \rightarrow TV(u)$ ,  $u_k \rightarrow u$  in  $D^{-1}$ . We observe that

$$\begin{aligned} (u_k, v)_{D^{-1}} &= \langle u_k, (-\Delta)^{-1}v \rangle = -\langle u_k, \operatorname{div} Z \rangle \\ &= \int_{\mathbb{R}^n} Z \cdot \nabla u_k dx \leq \|Z\|_\infty TV(u_k). \end{aligned}$$

Sending  $k \rightarrow \infty$ , we conclude that

$$(u, v)_{D^{-1}} \leq \|Z\|_\infty \quad \text{for all } u \in D^{-1} \text{ with } TV(u) \leq 1.$$

By definition of  $\Psi$ , this implies  $TV^0 \leq \Psi$ .

The inequality  $\Psi \leq TV^0$ :  
By definition,

$$\begin{aligned} TV(u) &= \sup \{ \langle u, -\operatorname{div} z \rangle \mid z \in C_c^\infty(\mathbb{R}^n), |z| \leq 1 \} \\ &= \sup \left\{ \frac{\langle u, -\operatorname{div} z \rangle}{\|z\|_\infty} \mid z \in C_c^\infty(\mathbb{R}^n), z \neq 0 \right\}. \end{aligned}$$

Since

$$\langle u, -\operatorname{div} z \rangle = \langle u, (-\Delta)^{-1} \Delta \operatorname{div} z \rangle = (u, \Delta \operatorname{div} z)_{D^{-1}},$$

we observe that

$$\begin{aligned} TV(u) &= \sup \left\{ \frac{(u, \Delta \operatorname{div} z)_{D^{-1}}}{\|z\|_\infty} \mid z \in C_c^\infty(\mathbb{R}^n), z \neq 0 \right\} \\ &\leq \sup \left\{ \frac{(u, \Delta \operatorname{div} z)_{D^{-1}}}{\Psi(\Delta \operatorname{div} z)} \mid z \in C_c^\infty(\mathbb{R}^n), \Psi(\Delta \operatorname{div} z) \neq 0 \right\} \\ &\leq \Psi^0(u). \end{aligned}$$

This implies that  $TV^0 \geq (\Psi^0)^0 = \Psi$ . □

## 4 Definition of solution

For a gradient flow of a convex functional, there is a general theory initiated by Y. Kōmura [23] and developed by H. Brézis [6] and others. It is summarized as follows.

**Proposition 14** ([6]). *Let  $H$  be a real Hilbert space. Let  $\mathcal{E}$  be a lower semicontinuous, convex functional on  $H$  with values in  $]-\infty, \infty]$ . Assume that  $D(\mathcal{E})$  is dense in  $H$ . Then, for any  $u_0 \in H$ , there exists a unique solution  $u \in C([0, \infty[, H)$  which is absolutely continuous in  $(\delta, T)$  (for any  $\delta < T < \infty$ ) satisfying*

$$\begin{cases} u_t \in -\partial\mathcal{E}(u) & \text{a.e. } t > 0 \\ u(0) = u_0. \end{cases}$$

Moreover,

$$\int_s^t \|u_\tau\|_H^2 d\tau \leq \mathcal{E}(u(s)) - \mathcal{E}(u(t)) \quad \text{for all } t \geq s > 0.$$

If  $\mathcal{E}(u_0) < \infty$ , then  $s = 0$  is allowed. In particular,  $u_t \in L^2(0, \infty; H)$ .

As in [2], this solution satisfies the evolutionary variational inequality

$$\frac{1}{2} \frac{d}{dt} \|u - f\|_H^2 \leq \mathcal{E}(f) - \mathcal{E}(u) \quad \text{a.e. } t > 0$$

for any  $f \in H$ . Indeed, by definition,  $u_t \in -\partial\mathcal{E}(u)$  is equivalent to saying

$$(-u_t, f - u)_H \leq \mathcal{E}(f) - \mathcal{E}(u)$$

for any  $f \in H$ . Since the left-hand side equals to

$$\frac{1}{2} \frac{d}{dt} \|u - f\|_H^2,$$

we obtain the evolutionary variational inequality. The evolutionary variational inequality is not only an equivalent formulation of the gradient flow  $u_t \in -\partial\mathcal{E}(u)$ , but also apply to a gradient flow of a metric space by replacing  $\|u - f\|_H$  by distance between  $u$  and  $f$ ; see [2] for the theory.

Since the subdifferential of  $TV$  in  $D^{-1}$  is now calculated, we are able to justify an explicit definition of a solution proposed in [14].

**Theorem 15.** Assume that  $u \in C([0, \infty[, D^{-1})$ . Then  $u$  is a solution of  $u_t \in -\partial_{D^{-1}}TV(u)$  with  $u_0 = u(0)$  in the sense of Proposition 14 if and only if there exists  $Z \in L^\infty([0, \infty[ \times \mathbb{R}^n)$  satisfying

$$\operatorname{div} Z \in L^2(\delta, \infty; D_0^1(\mathbb{R}^n)) \quad \text{for any } \delta > 0$$

such that for a.e.  $t \in ]0, \infty[$  there holds

$$u_t = -\Delta \operatorname{div} Z \quad \text{in } D^{-1}(\mathbb{R}^n),$$

$$|Z| \leq 1 \quad \mathcal{L}^n\text{-a.e.}$$

and

$$\langle u, \operatorname{div} Z \rangle = -TV(u).$$

(If  $TV(u_0) < \infty$ ,  $\delta = 0$  is allowed.)

*Proof.* The Theorem essentially follows from Theorem 13. We only need to justify that a Cahn-Hoffman vector field  $Z$  defined separately for every time instance by Theorem 13 can be chosen to be jointly measurable, i.e.  $Z \in L^\infty([0, \infty[ \times \mathbb{R}^n)$ . As in the second-order case [3], this can be done by recalling that  $u$  is a limit of time-discretizations.  $\square$

Unfortunately, in the case  $n \leq 2$ , the characteristic function  $1_A$  of a set  $A$  of positive measure is not in  $D^{-1}$  since  $\int 1_A dx \neq 0$ . We shall define a new space containing  $1_A$  as follows. We take a function  $\psi \in L^2(\mathbb{R}^n)$  with compact support such that  $\int_{\mathbb{R}^n} \psi = 1$ . We introduce a vector space

$$E_\psi^{-1} = \{w + c\psi \mid w \in D^{-1}(\mathbb{R}^n), c \in \mathbb{R}\}.$$

This space is independent of the choice of  $\psi$ . Indeed, let  $\psi_i \in L^2(\mathbb{R}^n)$  be compactly supported and  $\int \psi_i dx = 1$  ( $i = 1, 2$ ). An element  $w + c\psi_1 \in E_{\psi_1}^{-1}$  can be rewritten as

$$w + c\psi_1 = w + c(\psi_1 - \psi_2) + c\psi_2.$$

The next lemma implies  $q = c(\psi_1 - \psi_2) \in D^{-1}(\mathbb{R}^n)$  since  $\int q = 0$ . We then conclude that  $w + c\psi_1 \in E_{\psi_2}^{-1}$ .

**Lemma 16.** Assume that  $n \leq 2$ . A compactly supported function  $q \in L^2(\mathbb{R}^n)$  belongs to  $D^{-1}(\mathbb{R}^n)$  if and only if  $\int_{\mathbb{R}^n} q = 0$ .

*Proof.* If  $q \in D^{-1} \cap L^1$ , then

$$\int q = \langle q, [1] \rangle = 0,$$

where  $[1]$  stands for the element of  $D_0^1$  whose representatives are 1 as well as 0.

Now suppose that a compactly supported function  $q \in L^2(\mathbb{R}^n)$  satisfies  $\int_{\mathbb{R}^n} q = 0$ . Given a  $[\varphi] \in D_0^1$ , by the Poincaré inequality, we have for any  $R > 0$ , independently of the representative  $\varphi \in D^1$ ,

$$\left\| \varphi - \frac{1}{|B_R|} \int_{B_R} \varphi \right\|_{L^2(B_R)} \leq C_R \|\nabla \varphi\|_{L^2(B_R)} \leq C_R \|\varphi\|_{D_0^1}.$$

In particular,  $\varphi \in L_{loc}^2$ . Taking into account this and the assumption  $\int q = 0$ , we see that the linear functional

$$\langle q, [\varphi] \rangle = \int q \varphi \tag{4.1}$$

on  $D_0^1$  is well defined. Moreover, if  $R > 0$  is large enough that  $\operatorname{supp} q \subset B_R$ ,

$$\left| \int q \varphi \right| = \left| \int_{B_R} q \left( \varphi - \frac{1}{|B_R|} \int_{B_R} \varphi \right) \right| \leq \|q\|_{L^2(B_R)} \left\| \varphi - \frac{1}{|B_R|} \int_{B_R} \varphi \right\|_{L^2(B_R)} \leq C_R \|q\|_{L^2} \|\varphi\|_{D_0^1}.$$

Thus, the functional defined by (4.1) is bounded, i.e.  $q \in D^{-1}$ .  $\square$

Since  $E_\psi^{-1}$  is independent of the choice of  $\psi$ , we suppress  $\psi$  and denote this space by  $E^{-1}$ . In case  $n \geq 3$ , we will use notation  $E^{-1} = D^{-1}$ . We also denote  $E_0^1 = D^1$  if  $n \leq 2$ ,  $E_0^1 = D_0^1$  if  $n \geq 3$ . For  $u \in E^{-1}$ ,  $v \in E_0^1$ , we denote

$$\langle u, v \rangle_E = \begin{cases} \langle w, v \rangle & \text{if } n \geq 3, \\ \langle w, [v] \rangle + c \int \psi v & \text{if } n \leq 2, \end{cases} \quad (4.2)$$

where  $u = w + c\psi$ ,  $w \in D^{-1}$ ,  $\psi \in L_c^2$ ,  $\int \psi = 1$ . As before, we check that the value of  $\langle u, v \rangle_E$  does not depend on the choice of this decomposition.

We recall that if  $n \geq 3$ ,  $E_0^1 = D_0^1$  and  $E^{-1} = D^{-1}$  come with a Hilbert space structure. We also define inner products on  $E_0^1$ ,  $E^{-1}$  in case  $n \leq 2$  by

$$(v_1, v_2)_{E_0^1} := ([v_1], [v_2])_{D_0^1} + \int \psi v_1 \int \psi v_2,$$

$$(u_1, u_2)_{E^{-1}} := (w_1, w_2)_{D^{-1}} + c_1 c_2$$

for  $u_i = w_i + c_i \psi$ ,  $w_i \in D^{-1}$ ,  $c_i \in \mathbb{R}$  ( $i = 1, 2$ ). This gives an orthogonal decomposition

$$E^{-1} = D^{-1} \oplus \mathbb{R}.$$

We note that although the values of those products may depend on the choice of  $\psi$ , the topologies they induce on  $E_0^1$ ,  $E^{-1}$  do not. Formula (4.2) associates to any  $u \in E^{-1}$  a continuous linear functional on  $E_0^1$ . The resulting mapping is an isometric isomorphism between  $E^{-1}$  and the continuous dual to  $E_0^1$ .

We extend  $TV$  onto  $E^{-1}$  by defining

$$TV(u) := \sup_{\psi \in X_1} \langle u, -\operatorname{div} \psi \rangle_E.$$

As usual, we check that  $TV$  is a convex, weakly-\* (and strongly) lower semicontinuous functional. In particular, for a fixed  $g \in E^{-1}$ , the functional  $w \mapsto TV(w + g)$  is convex and lower semicontinuous on  $D^{-1}$ . We next give a definition of a solution of  $u_t \in -\partial TV(u)$  in the space  $E^{-1}$ . It turns out the idea of evolutionary variational inequality is very convenient since it is a flow in an affine space  $g + D^{-1}$  for some  $g \in E^{-1}$ .

**Definition 17.** Assume that  $u_0 \in E^{-1}$ . We say that  $u : [0, \infty[ \rightarrow E^{-1}$  is a solution to

$$u_t \in -\partial_{D^{-1}} TV(u) \quad (4.3)$$

in the sense of EVI (evolutionary variational inequality) with initial datum  $u_0$  if

- (i)  $u - g$  is absolutely continuous on  $[\delta, T]$  (for any  $\delta < T < \infty$ ) with values in  $D^{-1}$  and continuous up to zero and
- (ii)  $u - g$  satisfies the evolutionary variational inequality

$$\frac{1}{2} \frac{d}{dt} \|u(t) - g\|_{D^{-1}}^2 \leq TV(g) - TV(u(t)), \quad u(0) = u_0$$

holds for a. e.  $t > 0$ ,

provided that  $u_0 - g \in D^{-1}$  and  $g \in E^{-1}$ .

**Theorem 18.** For any  $u_0 \in E^{-1}$ , there exists a unique solution  $u$  of (4.3) in the sense of EVI.

*Proof.* The uniqueness part is standard [2]. We give a full proof for the reader's convenience and for completeness. Let  $u^i$  ( $i = 1, 2$ ) be a solution to (4.3) in the sense of EVI with initial datum  $u_0^i$  such that  $u_0^1 - u_0^2 \in D^{-1}$ . Since  $u^2$  is a solution,  $u^2(s) - u_0^2 \in D^{-1}$  for all  $s \geq 0$  by setting  $g = u_0^2$ .

This implies that  $u_0^1 - u^2(s) \in D^{-1}$  since  $u_0^1 - u_0^2 \in D^{-1}$ . Since  $u^1$  is a solution, we take  $g = u^2(s)$  and conclude that

$$\frac{1}{2} \frac{d}{dt} \|u^1(t) - u^2(s)\|_{D^{-1}}^2 \leq TV(u^2(s)) - TV(u^1(t)) \text{ for a.e. } t > 0, \quad \text{all } s > 0.$$

Interchanging  $u^1$  and  $u^2$ , we obtain

$$\frac{1}{2} \frac{d}{dt} \|u^2(t) - u^1(s)\|_{D^{-1}}^2 \leq TV(u^1(s)) - TV(u^2(t)) \text{ for a.e. } t > 0, \quad \text{all } s > 0.$$

Adding these two inequalities, we end up with

$$\frac{1}{2} \frac{\partial}{\partial t} \|u^1(t) - u^2(s)\|_{D^{-1}}^2 + \frac{1}{2} \frac{\partial}{\partial t} \|u^2(t) - u^1(s)\|_{D^{-1}}^2 \leq 0.$$

Evaluating the left-hand side at  $t = s$ , we observe that it equals

$$\frac{1}{2} \frac{d}{dt} \|u^1(t) - u^2(t)\|_{D^{-1}}^2.$$

We now conclude that  $\|u^1(t) - u^2(t)\|_{D^{-1}}^2$  is non-increasing so that

$$\|u^1(t) - u^2(t)\|_{D^{-1}} \leq \|u_0^1 - u_0^2\|_{D^{-1}} \quad \text{for all } t \geq 0. \quad (4.4)$$

This yields uniqueness of solution.

The existence is more involved. For  $u_0 = w_0 + g_0 \in E^{-1}$  with  $w_0 \in D^{-1}$ , we consider the gradient flow of the form

$$w_t \in -\partial_{D^{-1}} TV(w + g_0), \quad w(0) = w_0. \quad (4.5)$$

Applying Proposition 14, there is a unique solution  $w$  to (4.5) for  $w_0 \in D^{-1}$ . This solution satisfies the evolutionary variational inequality

$$\frac{1}{2} \frac{d}{dt} \|w - f\|_{D^{-1}}^2 \leq TV(f + g_0) - TV(w + g_0) \quad \text{for a.e. } t > 0$$

for any  $f \in D^{-1}$ . Setting  $u = w + g_0$ ,  $g = f + g_0$ , we end up with

$$\frac{1}{2} \frac{d}{dt} \|w - g\|_{D^{-1}}^2 \leq TV(g) - TV(u(t)).$$

Since  $g$  can be taken arbitrary such that  $u_0 - g \in D^{-1}$ , this shows that  $u$  is the solution of (4.3) in the sense of EVI; see condition (i) is easy by Proposition 14.  $\square$

From the proof of uniqueness we have a contraction property (4.4).

**Corollary 19.** *Let  $u^i$  be the solution to (4.3) in the sense of EVI with  $u^i(0) = u_0^i$  for  $i = 1, 2$  where  $u_0^i \in E^{-1}$ . Then*

$$\|u^1(t) - u^2(t)\|_{D^{-1}} \leq \|u_0^1 - u_0^2\|_{D^{-1}} \quad \text{for all } t \geq 0$$

provided that  $u_0^1 - u_0^2 \in D^{-1}$ .

It is non-trivial to characterize the subdifferential  $\partial_{D^{-1}} TV$ . For this purpose, we introduce a mapping  $I$  which plays a role analogous to  $-\Delta$  in  $n \geq 3$ .

**Lemma 20.** *Let  $n \leq 2$ . The mapping  $I: E_0^1 \rightarrow E^{-1}$  defined by*

$$I(f) = (-\Delta)[f] + \left( \int_{\mathbb{R}^n} f \psi \right) \psi$$

is an isometric isomorphism.

*Proof.* It is clear that  $I(f) \in E^{-1}$  and  $I$  is linear. For given  $u = w + c\psi \in E^{-1}$  with  $w \in D^{-1}$ ,  $c \in \mathbb{R}$ , there is  $\bar{f} \in D_0^1$  such that  $(-\Delta)\bar{f} = w$ . Since a representative  $f$  of  $\bar{f}$  is determined up to an additive constant, there is a unique representative  $f$  such that

$$\int_{\mathbb{R}^n} f\psi = c.$$

Thus, the mapping  $I$  is surjective. If  $I(f) = 0$ , then  $(-\Delta)[f] = 0$  so  $[f] = 0$ . Thus  $f$  is a constant. Since  $\int f\psi = 0$ , this constant must be zero, so  $f = 0$ . Thus,  $I$  is injective. Recalling our definitions of inner products on  $E_0^1$ ,  $E^{-1}$ , it is easy to check that  $I$  is an isometry.  $\square$

We have a characterization of the polar of  $TV$  in  $E^{-1}$  as in Theorem 11.

**Theorem 21.** *Let  $n \leq 2$ . Let  $\Psi$  be given by*

$$\Psi(v) = \inf \{ \|Z\|_\infty \mid v = I(-\operatorname{div} Z), Z \in L^\infty(\mathbb{R}^n), \operatorname{div} Z \in E_0^1 \}$$

for  $v \in E^{-1}$ . Then  $(TV)^0 = \Psi$ .

Admitting this fact, we are able to give a characterization of the subdifferential.

**Theorem 22.** *Let  $n \leq 2$ . An element  $v \in E^{-1}$  belongs to  $\partial_{E^{-1}}TV(u)$  if and only if there is  $Z \in L^\infty(\mathbb{R}^n)$  with  $\operatorname{div} Z \in E_0^1$  such that*

- (i)  $|Z| \leq 1$
- (ii)  $v = I(-\operatorname{div} Z)$
- (iii)  $-\langle u, \operatorname{div} Z \rangle_E = TV(u)$ .

*Proof of Theorem 22.* The proof parallels that of Theorem 13. By Lemma 9 and Theorem 21

$$v \in \partial TV(u) \iff \Psi(v) \leq 1 \text{ and } (u, v)_{E^{-1}} = TV(u).$$

The properties (i), (ii) together with  $\operatorname{div} Z \in E_0^1$  are equivalent to  $\Psi(v) \leq 1$ . Since

$$(u, v)_{E^{-1}} = (w, (-\Delta)[v])_{D^{-1}} + c \int_{\mathbb{R}^n} v\psi = \langle w, [v] \rangle + c \int_{\mathbb{R}^n} v\psi, \quad (4.6)$$

the Euler equation  $(u, v)_{E^{-1}} = TV(u)$  is equivalent to (iii).  $\square$

*Proof of Theorem 21.* The proof parallels that of Theorem 11. We first prove that

$$(u, v)_{E^{-1}} \leq \|Z\|_\infty \quad \text{for all } u \in E^{-1} \text{ with } TV(u) \leq 1$$

for  $v = I(-\operatorname{div} Z)$ . This implies  $TV^0 \leq \Psi$ . The estimate  $(u, v)_{E^{-1}} \leq \|Z\|_\infty$  formally follows from the identity (4.6). Indeed, by (4.6), we see

$$(u, v)_{E^{-1}} = -\langle w, [\operatorname{div} Z] \rangle - c \int_{\mathbb{R}^n} \psi \operatorname{div} Z.$$

If  $u$  is in  $C_c^\infty(\mathbb{R}^n)$ , then, by this formula, we obtain

$$(u, v)_{E^{-1}} = - \int_{\mathbb{R}^n} u \operatorname{div} Z \, dx = \int_{\mathbb{R}^n} \nabla u \cdot Z \, dx \leq \|Z\|_\infty TV(u).$$

By approximation, as in the proof of Theorem 11, we conclude the desired estimate.

The other inequality  $\Psi \leq TV^0$  follows from  $TV \leq \Psi^0$ . The proof of  $TV \leq \Psi^0$  is parallel to that of Theorem 11 by replacing  $\Delta \operatorname{div} Z$  by  $I(-\operatorname{div} Z)$  and the  $D^{-1}$  inner product by the  $E^{-1}$  inner product, respectively, if one notes the identity (4.6). Since  $I$  is an isometry,  $\Psi$  is lower semicontinuous, and we conclude that  $\Psi = TV^0$  by Remark 10.  $\square$

We have to be careful, since the  $E^{-1}$  gradient flow

$$u_t \in -\partial_{E^{-1}}TV(u)$$

does not correspond to the total variation flow  $u_t = (-\Delta) \operatorname{div}(\nabla u/|\nabla u|)$ . By Theorem 22(iii),  $Z = \nabla u/|\nabla u|$  if  $\nabla u \neq 0$ . Thus the  $E^{-1}$  gradient flow is formally of the form

$$u_t = (-\Delta) \operatorname{div}(\nabla u/|\nabla u|) + \psi \int_{\mathbb{R}^n} \psi \operatorname{div}(\nabla u/|\nabla u|) dx.$$

To recover the original total variation flow, we consider “partial” subdifferential in the direction of  $D^{-1}$ . Let  $P$  be the orthogonal projection from  $E^{-1}$  to  $D^{-1}$ . Then, by definition,

$$\partial_{D^{-1}}TV(w + c\psi) = P\partial_{E^{-1}}TV(u).$$

The equation

$$w_t \in -\partial_{D^{-1}}TV(w + c\psi)$$

is now formally of the form

$$u_t = (-\Delta) \operatorname{div}(\nabla u/|\nabla u|)$$

since  $c\psi$  is time-independent. Here is a precise statement.

**Theorem 23.** *Let  $n \leq 2$ . Consider the functional  $\mathcal{F} : w \mapsto TV(w + c\psi)$  on  $D^{-1}$  for a fixed  $c \in \mathbb{R}$  and  $\psi$ . Then,  $\partial_{D^{-1}}TV(w + c\psi) = P\partial_{E^{-1}}TV(u)$  for  $u = w + c\psi$ . In particular, an element  $v \in D^{-1}$  belongs to  $\partial_{D^{-1}}\mathcal{F}(w + c\psi)$  if and only if there is  $Z \in L^\infty(\mathbb{R}^n)$  with  $\operatorname{div} Z \in E_0^1$  such that*

- (i)  $|Z| \leq 1$ ,
- (ii)  $v = \Delta \operatorname{div} Z$ ,
- (iii)  $-\langle u, \operatorname{div} Z \rangle_E = TV(u)$ .

(In case  $n \leq 2$ , by  $\Delta \operatorname{div} Z$  we understand  $\Delta[\operatorname{div} Z]$ .) This characterization is important to calculate the solution of  $u_t = (-\Delta) \operatorname{div}(\nabla u/|\nabla u|)$  for  $n \leq 2$  explicitly. In fact, we recover the characterization of a solution in the sense of EVI as in Theorem 15. This amounts to Theorem 2.

## 5 The notion of calibrability

We are interested in sets where the speed of solution  $u_t$  is spatially constant. The speed is given as minus the minimal section of the subdifferential, i. e.

$$\partial_{D^{-1}}^0TV(u) := \arg \min \{ \|v\|_{D^{-1}} \mid v \in \partial_{D^{-1}}TV(u) \}.$$

Since  $\partial_{D^{-1}}TV(u)$  is closed and convex,  $\partial_{D^{-1}}^0TV(u)$  is uniquely determined if  $\partial_{D^{-1}}TV(u) \neq \emptyset$ . Since we have characterized the subdifferential, we end up with

$$\begin{aligned} \partial_{D^{-1}}^0TV(u) = \arg \min \{ \|v\|_{D^{-1}} \mid v = \Delta \operatorname{div} Z, Z \in L^\infty(\mathbb{R}^n), |Z| \leq 1, \\ \operatorname{div} Z \in E_0^1(\mathbb{R}^n), -\langle u, \operatorname{div} Z \rangle_E = TV(u) \}. \end{aligned}$$

Although the minimizer  $v$  is unique, the corresponding  $Z$  may not be unique. Let  $U$  be a smooth open set in  $\mathbb{R}^n$ . We consider a smooth function  $u$  such that

$$\bar{U} = \{x \in \mathbb{R}^n \mid u(x) = 0\}$$

and  $\partial_{D^{-1}}TV(u) \neq \emptyset$ . Such a closed set is often called a facet. Assume further that  $\nabla u \neq 0$  outside  $\bar{U}$ . Let  $Z$  be a vector field satisfying  $v = \Delta \operatorname{div} Z$  for  $v \in \partial TV(u)$ . It is easy to see that outside the facet  $\bar{U}$ ,

$$Z(x) = \nabla u(x) / |\nabla u(x)|$$

by  $-\langle u, \operatorname{div} Z \rangle_E = TV(u)$ . Since  $\|v\|_{D^{-1}} = \|\operatorname{div} Z\|_{D_0^1}$ , we see that

$$\partial_{D^{-1}}^0 TV(u) = \arg \min \left\{ \|\operatorname{div} Z\|_{D_0^1} \mid |Z| \leq 1 \text{ in } U, Z = \nabla u / |\nabla u| \text{ in } \bar{U}^c, \operatorname{div} Z \in E_0^1 \right\}.$$

Since  $\operatorname{div} Z$  is locally integrable, the normal trace is well-defined from inside as an element of  $L^\infty(\partial U)$  [3] and it must agree with that from outside, i.e.

$$\nu \cdot Z(x) = \nu(x) \cdot \nabla u / |\nabla u| = \nu(x) \cdot \chi \nu(x) = \chi(x),$$

where  $\nu(x)$  is the exterior unit normal of  $\partial U$  and

$$\chi(x) = \begin{cases} 1 & \text{if } u > 0 \text{ outside } \bar{U} \text{ near } x \in \partial U, \\ -1 & \text{otherwise.} \end{cases}$$

Since  $\operatorname{div} Z$  is in  $E_0^1$ , its trace from inside and outside must agree, i.e.,

$$\operatorname{div} Z(x) = \chi \operatorname{div} \nu(x).$$

Let  $Z_0$  be a minimizer corresponding to  $v = \partial_{D^{-1}}^0 TV(u)$ . Since the value  $Z_0$  outside  $\bar{U}$  is always the same, we consider its restriction on  $U$  and still denote by  $Z_0$ . Then,

$$Z_0 = \arg \min \left\{ \int_U |\nabla \operatorname{div} Z|^2 \mid |Z| \leq 1 \text{ in } U \quad \nu \cdot Z = \chi \text{ on } \partial U, \operatorname{div} Z = \chi \operatorname{div} \nu \text{ on } \partial U \right\}.$$

In the case  $u > 0$  outside  $\bar{U}$ ,  $\operatorname{div} Z$  on  $\partial U$  must equal to the minus mean curvature (the sum of all principal curvatures) in the direction of  $\nu$ . Although  $\operatorname{div} Z_0 \in D_0^1(\mathbb{R}^n)$  so that  $\nabla \operatorname{div} Z_0 \in L^2(\mathbb{R}^n)$ , the quantity  $\nabla \operatorname{div} Z_0$  may jump across  $\partial U$ . Thus  $\Delta \operatorname{div} Z_0$  may contain singular part which is a driving force to move the facet boundary ‘‘horizontally’’ during its evolution under the fourth-order total variation equation as observed in the previous section and earlier in [11]. In the second-order problem, the speed does not contain any singular part so the jump discontinuity does not move.

We are interested in a situation where  $\Delta \operatorname{div} Z_0$  is constant over  $U$ . In the spirit of [24], we call any continuous function  $\chi: \partial U \rightarrow \{-1, 1\}$  a *signature* for  $U$ .

**Definition 24.** Let  $U$  be a smooth open set in  $\mathbb{R}^n$  with signature  $\chi$ . We say that  $U$  is  $(D^{-1})$ -calibrable (with signature  $\chi$ ) if there exists a minimizer  $Z_0$  of

$$\int_U |\nabla \operatorname{div} Z|^2 \tag{5.1}$$

under the constraint

$$|Z| \leq 1 \text{ a. e. on } U \tag{5.2}$$

with boundary conditions

$$\nu \cdot Z = \chi, \quad \operatorname{div} Z = \chi \operatorname{div} \nu \quad \text{on } \partial U, \tag{5.3}$$

with the property that

$$\Delta \operatorname{div} Z_0 \text{ is constant over } U. \tag{5.4}$$

We call any such  $Z_0$  a  $(D^{-1})$ -calibration for  $U$  (with signature  $\chi$ ).

We shall study this variational problem. We first set  $w = \operatorname{div} Z$ ,  $w_0 = \operatorname{div} Z_0$ . Then,  $\int_U |\nabla \operatorname{div} Z|^2$  is nothing but the Dirichlet energy for  $w$ , i.e.,

$$e(w) = \int_U |\nabla w|^2$$

and  $w_0$  minimizes this  $e(w)$  under the constraint. Since  $\nu \cdot Z = \chi$ ,

$$\int_U w d\mathcal{L}^n = \int_{\partial U} \chi d\mathcal{H}^{n-1}.$$

The other boundary condition can be rewritten as

$$w = \chi \operatorname{div} \nu \quad \text{on } \partial U.$$

**Proposition 25.** *Let  $U$  be a smooth bounded domain in  $\mathbb{R}^n$ . Assume that  $Z_0$  is a calibration for  $U$  with signature  $\chi$ . Then,  $w_0 = \operatorname{div} Z_0$  must satisfy the Saint-Venant problem*

$$\begin{cases} -\Delta w = \lambda & \text{in } U \\ w = \chi \operatorname{div} \nu & \text{on } \partial U \end{cases} \quad (5.5)$$

with the constraint

$$\int_U w \, d\mathcal{L}^n = \int_{\partial U} \chi \, d\mathcal{H}^{n-1}, \quad (5.7)$$

where  $\lambda$  is some constant.

We are able to prove the converse.

**Theorem 26.** *Let  $U$  be a smooth bounded domain in  $\mathbb{R}^n$ . Assume that  $Z_* \in L^\infty(U)$  satisfies  $|Z_*| \leq 1$  in  $U$  and that  $w_* = \operatorname{div} Z_* \in D_0^1(U)$  is a solution to the Saint-Venant problem (5.5), (5.6), where  $\lambda$  is a constant. Assume that  $Z_*$  satisfies  $\nu \cdot Z_* = \chi$ . Then  $Z_*$  is a minimal Cahn-Hoffman vector field. In particular,  $U$  is calibrable with signature  $\chi$  and  $Z_*$  is a calibration.*

*Proof.* We first note that  $w_* = \operatorname{div} Z_*$  must satisfy (5.7). We consider the variational problem of  $e(w)$  under the Dirichlet condition (5.6) and the constraint (5.7). Since the problem is strictly convex, there is a unique minimizer  $\bar{w}$  in  $D_0^1(U)$ . By Lagrange's multiplier method,  $\bar{w}$  must satisfy (5.5) because of the constraint (5.7). (Actually, a weak solution  $\bar{w}$  of (5.5) is a smooth solution of (5.5), (5.6) by the standard regularity theory of linear elliptic partial differential equations [18, Chapter 6].) As we see below, the constant  $\lambda$  is uniquely determined by (5.5) and (5.7). For  $\bar{w}$ , there always exists  $Z \in C^\infty(\bar{U})$  such that

$$\operatorname{div} Z = \bar{w} \quad \text{in } U, \quad \nu \cdot Z = \chi \quad \text{on } \partial U. \quad (5.8)$$

Indeed, let  $p$  be a solution of the Neumann problem

$$\Delta p = \bar{w} \quad \text{in } U, \quad \nu \cdot \nabla p = \chi \quad \text{on } \partial U.$$

Such a solution  $p$  always exists since  $\bar{w}$  satisfies the compatibility condition (5.7) and it is smooth up to  $\bar{U}$ ; see e. g. [10, 18]. If we set  $Z = \nabla p$ , then  $Z$  satisfies the desired property (5.8). Thus, the minimum  $e(\bar{w})$  of the Dirichlet energy under the constraint (5.7) agrees with

$$\min \left\{ \int_U |\nabla \operatorname{div} Z|^2 \mid \nu \cdot Z = \chi, \operatorname{div} Z = \chi \operatorname{div} \nu \text{ on } \partial U \right\}.$$

Since  $w_* = \bar{w}$  and  $|Z_*| \leq 1$ , this shows that  $Z_*$  is a minimal Cahn-Hoffman vector field. Thus,  $U$  is calibrable and  $Z_*$  is a calibration.  $\square$

**Lemma 27.** *Let  $U$  be a smooth bounded domain in  $\mathbb{R}^n$ . Let  $w$  solve*

$$\begin{cases} -\Delta w = \lambda & \text{in } U \\ w = f & \text{on } \partial U \end{cases}$$

for  $\lambda \in \mathbb{R}$ ,  $f \in C(\partial U)$ . This solution can be written as

$$w = \lambda w_{\text{sv}} + h_f$$

where  $w_{\text{sv}}$  solves the Saint-Venant problem

$$\begin{cases} -\Delta w_{\text{sv}} = 1 & \text{in } U \\ w_{\text{sv}} = 0 & \text{on } \partial U \end{cases}$$

and  $h_f$  is the harmonic extension of  $f$  to  $U$ . In particular, if  $\int_U w = c$  is given then  $\lambda$  is uniquely determined by

$$\lambda \int_U w_{\text{sv}} + \int_U h_f = c,$$

since  $w_{\text{sv}} > 0$  in  $U$ .

The decomposition  $w = \lambda w_{sv} + h_f$  is rather clear. The property  $w_{sv} > 0$  in  $U$  follows from the maximum principle [18].

We now compare the definition of calibrability for the second-order problem.

**Definition 28.** *Let  $U$  be a smooth open set in  $\mathbb{R}^n$  with signature  $\chi$ . We say that  $\bar{U}$  is ( $L^2$ -)calibrable if there is a minimizer  $Z_0$  of*

$$\int_U |\operatorname{div} Z|^2$$

*under the constraint  $|Z| \leq 1$  in  $U$  and the boundary condition  $\nu \cdot Z = \chi$  with the property that  $\operatorname{div} Z_0$  is a constant over  $U$ .*

This definition is slightly weaker than the calibrability used in [1, 25], where  $a(t)1_U$  is a solution of the total variation flow in  $\mathbb{R}^n$  with some function  $a(t)$  of  $t$ ; see also [3]. This requires that  $\partial_{L^2}^0 TV(u)$  is constant not only on  $U$  but also  $U^c$ . Our definition follows from that of [5].

## 6 Calibrability of rotationally symmetric sets

**Definition 29.** *We say that a Lebesgue measurable subset  $U$  (defined up to a set of measure zero) of  $\mathbb{R}^n$  is a generalized annulus if  $U$  is non-empty, open, connected and rotationally symmetric, i. e. invariant under the linear action of  $SO(n)$  on  $\mathbb{R}^n$ .*

It is easy to see that any generalized annulus is a ball, an annulus, the complement of a ball or the whole space  $\mathbb{R}^n$ . In other words, any generalized annulus is of form

$$A_{R_0}^{R_1} = \{x \in \mathbb{R}^n : R_0 < |x| < R_1\}$$

with  $0 \leq R_0 < R_1 \leq \infty$ . We note that

$$A_0^R = B_R$$

as measurable sets for  $R > 0$ . In this section we will settle the question which generalized annuli are calibrable.

**Lemma 30.** *Let  $U$  be a generalized annulus. Suppose that  $U$  is calibrable with signature  $\chi$ . Then there exists a calibration  $\bar{Z}$  for  $(U, \chi)$  of form  $\bar{Z}(x) = z(|x|) \frac{x}{|x|}$ .*

*Proof.* Let  $Z$  be any calibration for  $(U, \chi)$ . Let  $\mu_n$  be the Haar measure on  $SO(n)$ . We define  $\bar{Z}$  as the average

$$\bar{Z}(x) = \int LZ(L^{-1}x) d\mu_n(L).$$

It is an exercise in vector calculus to check that  $\bar{Z}$  satisfies boundary conditions (5.3) and that  $\Delta \operatorname{div} \bar{Z}$  is a constant (equal to  $\Delta \operatorname{div} Z$ ) on  $U$ . By convexity,  $|\bar{Z}| \leq 1$  and  $\bar{Z}$  is also a minimizer of (5.1). Thus  $\bar{Z}$  is a calibration for  $(U, \chi)$ . By definition, it is invariant under rotations, i. e.

$$L\bar{Z}(L^{-1}x) = \bar{Z}(x)$$

for  $L \in SO(n)$ ,  $x \in \mathbb{R}^n$ . In the case  $n = 1$  this already shows that  $\bar{Z}$  is in the desired form. In higher dimensions, we consider the orthogonal decomposition

$$\bar{Z}(x) = \bar{Z}^\perp(x) + \bar{Z}^T(x) := \frac{x}{|x|} \otimes \frac{x}{|x|} \bar{Z}(x) + \left( I - \frac{x}{|x|} \otimes \frac{x}{|x|} \right) \bar{Z}(x).$$

Both  $\bar{Z}^\perp$  and  $\bar{Z}^T$  are invariant under rotations. In particular, for any given  $R > 0$ , the restriction of  $\bar{Z}^T$  to  $\mathbb{S}_R^{n-1}$  is an invariant tangent vector field on  $\mathbb{S}_R^{n-1}$ . Note that any such vector field is smooth. If  $n = 3$ , it follows by the hedgehog uncombability theorem [9, Proposition 7.15] that

$\bar{Z}^T \equiv 0$ . If  $n > 3$ , any vector field invariant on  $\mathbb{S}^{n-1}$  is in particular invariant on a sphere  $\mathbb{S}^2$  containing any given point in  $\mathbb{S}^{n-1}$ , so the same conclusion follows. Thus, we have

$$\bar{Z}(x) = \bar{Z}^\perp(x) = \frac{x}{|x|} \otimes \frac{x}{|x|} \bar{Z}(x) = \frac{x}{|x|} \cdot \bar{Z}(x) \frac{x}{|x|} =: \bar{z}(x) \frac{x}{|x|}.$$

By rotational invariance, we have  $\bar{z}(x) = z(|x|)$ , which concludes the proof.

We are left with the case  $n = 2$  in which there exists a one-dimensional space of invariant tangent fields on  $\mathbb{S}^{n-1} = \mathbb{S}^1$  spanned by  $e^T(x) := (x_2, -x_1)$ . Thus, we have

$$\bar{Z}^T(x) = z^T(|x|)e^T(x).$$

We calculate

$$\operatorname{div} \bar{Z}^T(x) = z^T(|x|) \operatorname{div} e^T(x) + (z^T)'(|x|) \frac{x}{|x|} \cdot e^T(x) = 0.$$

Thus, we can disregard  $\bar{Z}^T$  and choose  $\bar{Z}^\perp$  as our calibration, since it satisfies conditions (5.2)-(5.4) (recall that  $\bar{Z}^T$  and  $\bar{Z}^\perp$  are orthogonal) and  $\bar{Z}^T$  does not contribute to the value of (5.1). As before, we see that

$$\bar{Z}^\perp(x) = z(|x|) \frac{x}{|x|}.$$

□

Let  $U$  be a generalized annulus. By Lemma 30, if  $U$  is calibrable, then there exists a calibration  $Z$  for  $U$  of form  $Z = z(|x|) \frac{x}{|x|}$ . It follows from (5.5) that  $z$  needs to satisfy the ODE

$$-r^{1-n} \left( r^{n-1} \left( r^{1-n} (r^{n-1} z)' \right)' \right)' = \lambda. \quad (6.1)$$

The general solution to this ODE is

$$z(r) = c_0 r^3 + c_1 r^{3-n} + c_2 r + c_3 r^{1-n} \quad (6.2)$$

where  $c_0 = -\frac{\lambda}{2n(n+2)}$  if  $n \neq 2$  and

$$z(r) = c_0 r^3 + c_1 r \log r + c_2 r + c_3 r^{-1}. \quad (6.3)$$

where  $c_0 = -\frac{\lambda}{16}$  if  $n = 2$ . We will now try to find a calibration for  $U$  by solving a suitable boundary value problem for (6.1).

## 6.1 Balls

Let  $U = B_R(0)$ . To focus attention, we choose  $\chi = -1$  on  $\partial U$ . In this case, boundary conditions (5.3) lead to

$$z(R) = -1, \quad z'(R) = 0. \quad (6.4)$$

If  $n \geq 2$ , in order to satisfy the requirements  $|Z| \leq 1$  and  $\nabla \operatorname{div} Z$ , we need to restrict to  $c_1 = c_3 = 0$  in (6.2). We make the same choice also in case  $n = 1$ , as it leads to the right result. Then, applying (6.4) in (6.2) or (6.3), we obtain a system of two affine equations for two unknowns  $\lambda, c_2$ . We solve it obtaining

$$z(r) = \frac{1}{2} \left( \frac{r}{R} \right)^3 - \frac{3}{2} \frac{r}{R}, \quad (6.5)$$

$$\lambda = -\frac{n(n+2)}{R^3}. \quad (6.6)$$

We check that  $z$  satisfies  $|z| \leq 1$  on  $[0, R]$ , so  $Z$  is a calibration for  $B_R$ . Thus, all balls are calibrable in any dimension.

## 6.2 Complements of balls

Let  $U = \mathbb{R}^n \setminus B_R$ . For consistency with the previous case, we choose  $\chi = 1$  on  $\partial U$ . In this case, boundary conditions (5.3) also lead to

$$z(R) = -1, \quad z'(R) = 0. \quad (6.7)$$

Let us first assume that  $n \geq 3$ . In order to satisfy the requirement  $|Z| \leq 1$ , we need to restrict to  $\lambda = c_2 = 0$  in (6.2). Again, applying (6.4) in (6.2) leads to a system of two affine equations for two unknowns  $c_1, c_3$ . We solve it obtaining

$$z(r) = -\frac{n-1}{2} \left(\frac{r}{R}\right)^{3-n} + \frac{n-3}{2} \left(\frac{r}{R}\right)^{1-n}. \quad (6.8)$$

Again, we easily check that  $z$  satisfies  $|z| \leq 1$  on  $[0, R]$ , so  $Z$  is a calibration for  $B_R$ .

In the omitted cases  $n = 1, 2$ , requirement  $|Z| \leq 1$  implies  $\lambda = c_1 = c_2 = 0$  in (6.2). If  $n = 1$ , there exists  $z$  of such form satisfying (6.7):  $z(r) \equiv -1$ , consistently with (6.8). On the other hand, if  $n = 2$ , applying (6.7) to (6.2) with  $\lambda = c_1 = c_2 = 0$  leads to a contradiction.

Summing up, all complements of balls are calibrable if  $n \neq 2$ . On the other hand, if  $n = 2$  all complements of balls turn out not to be calibrable.

## 6.3 Annuli

Let now  $U = A_{R_0}^{R_1} = B_{R_1} \setminus B_{R_0}$ ,  $0 < R_0 < R_1$ . In this case  $\partial U$  has two connected components, so there exist two distinct choices of signature: constant and non-constant. Let us first consider the former. To focus attention, we choose  $\chi \equiv -1$ . Then, boundary conditions (5.3) take form

$$z(R_0) = 1, \quad z(R_1) = -1, \quad z'(R_0) = z'(R_1) = 0. \quad (6.9)$$

Applying (6.9) to (6.2) or (6.3) leads to a system of four affine equations with four unknowns. In the case  $n \neq 2$ , the solution is

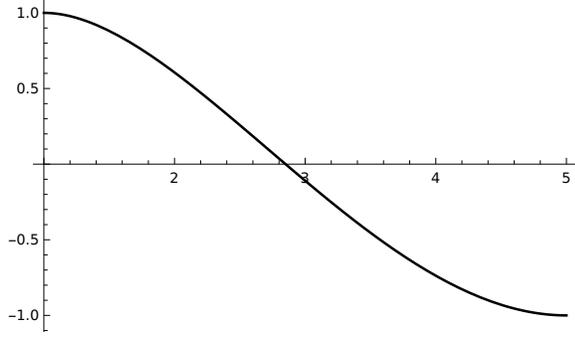
$$\begin{aligned} c_0 &= \frac{-(R_1 - R_0)R_1^n R_0^n ((n-1)(n-2)(R_1 + R_0)^2 - 2R_0 R_1) + 2R_1^3 R_0^{2n} - 2R_0^3 R_1^{2n}}{R_1 R_0 (R_1^n R_0^n (n^2 (R_1^2 - R_0^2)^2 + 8R_1^2 R_0^2) - 4R_1^2 R_0^{2n+2} - 4R_0^2 R_1^{2n+2})}, \\ c_1 &= \frac{(R_1 + R_0) (R_1^n (2(n-1)R_1^2 + (n+2)R_1 R_0 - (n+2)R_0^2) + R_0^n ((n+2)R_1^2 - (n+2)R_1 R_0 - 2(n-1)R_0^2))}{2(n^2 - 4)R_1^3 R_0^3 + 4R_1^{3-n} R_0^{n+3} + 4R_1^{n+3} R_0^{3-n} - n^2 R_1^5 R_0 - n^2 R_1 R_0^5}, \\ c_2 &= \frac{(R_1 - R_0)R_1^n R_0^n (6R_1^2 R_0^2 + (n-1)nR_1^3 R_0 + (n-1)nR_1^4 + (n-1)nR_1 R_0^3 + (n-1)nR_0^4) - 6R_1^3 R_0^{2n+2} + 6R_0^3 R_1^{2n+2}}{R_1 R_0 (R_1^n R_0^n (n^2 (R_1^2 - R_0^2)^2 + 8R_1^2 R_0^2) - 4R_1^2 R_0^{2n+2} - 4R_0^2 R_1^{2n+2})}, \\ c_3 &= \frac{(R_1 + R_0) (R_0^2 R_1^n (-2(n-3)R_1^2 - nR_1 R_0 + nR_0^2) - R_1^2 R_0^n (nR_1^2 - R_0(nR_1 + 2(n-3)R_0)))}{2(n^2 - 4)R_1^3 R_0^3 + 4R_1^{3-n} R_0^{n+3} + 4R_1^{n+3} R_0^{3-n} - n^2 R_1^5 R_0 - n^2 R_1 R_0^5}. \end{aligned} \quad (6.10)$$

This can be rewritten in a form emphasizing homogeneity:

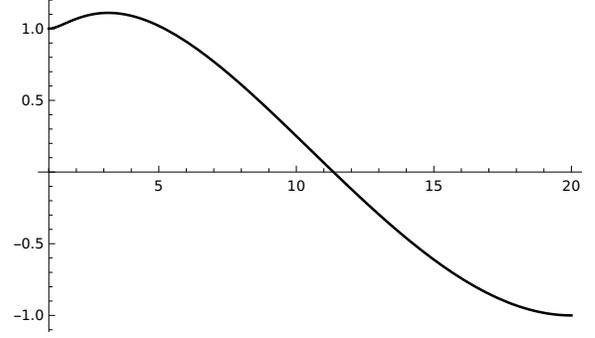
$$\begin{aligned} c_0 &= \frac{-(Q-1)Q^n ((n-1)(n-2)(Q+1)^2 - 2Q) + 2Q^3 - 2Q^{2n}}{Q^{n-2} (n^2 (Q^2 - 1)^2 + 8Q^2) - 4 - 4Q^{2n}} R_1^{-3}, \\ c_1 &= \frac{(Q+1) (Q^n (2(n-1)Q^2 + (n+2)Q - (n+2)) + ((n+2)Q^2 - (n+2)Q - 2(n-1)))}{2(n^2 - 4)Q^n + 4 + 4Q^{2n} - n^2 Q^{n+2} - n^2 Q^{n-2}} R_1^{n-3}, \\ c_2 &= \frac{(Q-1)Q^n (6Q^2 + (n-1)nQ^3 + (n-1)nQ^4 + (n-1)nQ + (n-1)n) - 6Q^3 + 6Q^{2n+2}}{Q^n (n^2 (Q^2 - 1)^2 + 8Q^2) - 4Q^2 - 4Q^{2n+2}} R_1^{-1}, \\ c_3 &= \frac{(Q+1) (Q^n (-2(n-3)Q^2 - nQ + n) - Q^2 (nQ^2 - (nQ + 2(n-3))))}{2(n^2 - 4)Q^{n+2} + 4Q^2 + 4Q^{2n+2} - n^2 Q^{n+4} - n^2 Q^n} R_1^{n-1}, \end{aligned} \quad (6.11)$$

where we denoted  $Q = R_1/R_0$ . We can further simplify it to

$$\begin{aligned} c_0 &= \frac{2Q^3(Q^{2n-3} - 1) + (Q-1)Q^n ((n-1)(n-2)(Q+1)^2 - 2Q)}{4(Q^n - 1)^2 - n^2(Q^2 - 1)^2 Q^{n-2}} R_1^{-3}, \\ c_1 &= \frac{(Q+1)(2(n-1)(Q^{n+2} - 1) + (n+2)Q(Q-1)(Q^{n-1} + 1))}{4(Q^n - 1)^2 - n^2(Q^2 - 1)^2 Q^{n-2}} R_1^{n-3}, \\ c_2 &= -\frac{6Q(Q^{2n-1} - 1) + (Q-1)Q^{n-2}(6Q^2 + n(n-1)(1+Q)(1+Q^3))}{4(Q^n - 1)^2 - n^2(Q^2 - 1)^2 Q^{n-2}} R_1^{-1}, \\ c_3 &= -\frac{(Q+1)(2(n-3)(Q^n - 1) + nQ(Q-1)(Q^{n-3} + 1))}{4(Q^n - 1)^2 - n^2(Q^2 - 1)^2 Q^{n-2}} R_1^{n-1}. \end{aligned} \quad (6.12)$$



(a)  $R_0 = 1, R_1 = 5.$



(b)  $R_0 = 1, R_1 = 20.$

Figure 1: Plots of  $z$  for an annulus with constant signature for two different values of  $Q$  in case  $n = 2$ .

We need to check whether condition  $|Z| \leq 1$  is satisfied. We calculate

$$\begin{aligned} z''(r) &= 6c_0r + (n-3)(n-2)c_1r^{1-n} + n(n-1)c_3r^{-n-1} \\ &= r^{-n-1}(6c_0r^{n+2} + (n-3)(n-2)c_1r^2 + n(n-1)c_3) =: r^{-n-1}w(r). \end{aligned} \quad (6.13)$$

Using the form (6.12), we can check that  $c_0 > 0, c_1 > 0$  for all  $Q > 1$ . Therefore,  $w$  has at most one zero on the half-line  $r > 0$ . Consequently,  $z''$  has at most one zero, so  $z$  has at most one inflection point. Taking into account (6.9),  $z$  cannot have a local extremum on  $]R_0, R_1[$ . Thus,  $|z| \leq 1$  on  $]R_0, R_1[$  and  $Z$  is a valid calibration.

In the case  $n = 2$ , the solution is

$$\begin{aligned} c_0 &= \frac{R_1^2 - R_0^2 + 2R_1R_0 \log(R_1/R_0)}{4R_1(R_1 - R_0)R_0(-R_1^2 + R_0^2 + (R_1^2 + R_0^2) \log(R_1/R_0))} \\ c_1 &= \frac{-R_1^3 - 3R_1^2R_0 - 3R_1R_0^2 - R_0^3}{2R_1R_0(-R_1^2 + R_0^2 + (R_1^2 + R_0^2) \log(R_1/R_0))} \\ c_2 &= \frac{-3R_1^4 + 3R_0^4 + 2(R_1^4 - R_1^3R_0 + R_1^2R_0^2 - 3R_1R_0^3) \log(R_1) + 2(3R_1^3R_0 - R_1^2R_0^2 + R_1R_0^3 - R_0^4) \log(R_0)}{4R_1(R_1 - R_0)R_0(-R_1^2 + R_0^2 + (R_1^2 + R_0^2) \log(R_1/R_0))} \\ c_3 &= \frac{R_1(-3R_1^2R_0 + 3R_0^3 + 2(R_1^2R_0 - R_1R_0^2 + R_0^3) \log(R_1/R_0))}{4(R_1 - R_0)(-R_1^2 + R_0^2 + (R_1^2 + R_0^2) \log(R_1/R_0))} \end{aligned}$$

which can be rewritten (again, denoting  $Q = R_1/R_0$ ) as

$$\begin{aligned} c_0 &= \frac{Q^2(Q^2 - 1 + 2Q \log Q)}{4(Q-1)(-Q^2 + 1 + (Q^2 + 1) \log Q)} R_1^{-3}, \\ c_1 &= \frac{-(Q+1)^3}{2(-Q^2 + 1 + (Q^2 + 1) \log Q)} R_1^{-1}, \\ c_2 &= \frac{-3(Q^4 - 1) - 2(3Q^3 - Q^2 + Q - 1) \log Q}{4(Q-1)(-Q^2 + 1 + (Q^2 + 1) \log Q)} R_1^{-1} + \frac{(Q+1)^3}{2(-Q^2 + 1 + (Q^2 + 1) \log Q)} R_1^{-1} \log R_1, \\ c_3 &= \frac{-3Q^2 + 3 + 2(Q^2 - Q + 1) \log Q}{4(Q-1)(-Q^2 + 1 + (Q^2 + 1) \log Q)} R_1. \end{aligned} \quad (6.14)$$

As before, we calculate the second derivative of  $z$ :

$$z''(r) = 6c_0r + c_1r^{-1} + 2c_3r^{-3} = r^{-3}(6c_0r^4 + c_1r^2 + 2c_3) =: r^{-3}w(r).$$

The polynomial  $w$  has at most 2 positive roots, and so does  $z''$ . By (6.9), at least one of them belongs to  $]R_0, R_1[$ . Furthermore, since  $c_0 > 0$  for  $Q > 1$ ,  $z''(r)$  is positive for large values of  $r$ . Taking into account these observations, we deduce that  $|z| \leq 1$  on  $[R_0, R_1]$  if and only if  $z''(R_0) \leq 0$ . This inequality is equivalent to

$$m(Q) := \log Q - \frac{(Q^2 - 1)(2Q - 1)}{Q(Q^2 - 2Q + 3)} \leq 0.$$

We compute

$$m(1) = 0, \quad \lim_{Q \rightarrow +\infty} m(Q) = +\infty, \quad m'(Q) = \frac{(Q-3)(Q-1)(Q+1)^3}{Q^2(Q^2-2Q+3)^2}. \quad (6.15)$$

We observe that  $m$  has exactly one zero  $Q_*$  on  $]1, +\infty[$ , and  $m(Q) \leq 0$  if and only if  $Q \leq Q_*$ . Therefore,  $Z$  is a valid calibration for  $A_{R_0}^{R_1}$  with signature  $-1$  if and only if  $R_1/R_0 \leq Q_*$ . By (6.15) it is evident that  $Q_* > 3$ . Numerical computation using Wolfram Mathematica shows that  $Q_* \approx 9.7$ . Thus,  $A_{R_0}^{R_1}$  with constant signature is calibrable if and only if  $R_1/R_0 \leq Q_*$ . This concludes the proof of Theorem 3.

Now, let us consider non-constant signature. We assume that  $\chi = 1$  on  $\partial B_{R_0}$  and  $\chi = -1$  on  $\partial B_{R_1}$ . This choice leads to

$$z(R_0) = -1, \quad z(R_1) = -1, \quad z'(R_0) = z'(R_1) = 0. \quad (6.16)$$

If  $n \neq 2$ , the solution to the resulting affine system is

$$\begin{aligned} c_0 &= \frac{-2R_1^3 R_0^{2n} - 2R_0^3 R_1^{2n} + R_0^n R_1^n (R_0 + R_1) \left( (n-2)(n-1)R_0^2 - 2((n-3)n+1)R_0 R_1 + (n-2)(n-1)R_1^2 \right)}{R_0 R_1 \left( R_0^n R_1^n \left( n^2 (R_0^2 - R_1^2)^2 + 8R_0^2 R_1^2 \right) - 4R_1^2 R_0^{2n+2} - 4R_0^2 R_1^{2n+2} \right)} \\ c_1 &= -\frac{-3nR_1 R_0^{n+2} + 2(n-1)R_0^{n+3} + (n+2)R_1^3 R_0^n + (n+2)R_0^3 R_1^n - 3nR_0 R_1^{n+2} + 2(n-1)R_1^{n+3}}{4R_0^{3-n} R_1^{3-n} (R_0^n - R_1^n)^2 - n^2 R_0 R_1 (R_0^2 - R_1^2)^2} \\ c_2 &= \frac{6R_1^3 R_0^{2n+2} + 6R_0^3 R_1^{2n+2} - R_0^n R_1^n (R_0 + R_1) \left( (n-1)nR_0^4 - (n-1)nR_0^3 R_1 - (n-1)nR_0 R_1^3 + (n-1)nR_1^4 + 6R_0^2 R_1^2 \right)}{R_0 R_1 \left( R_0^n R_1^n \left( n^2 (R_0^2 - R_1^2)^2 + 8R_0^2 R_1^2 \right) - 4R_1^2 R_0^{2n+2} - 4R_0^2 R_1^{2n+2} \right)} \\ c_3 &= \frac{(R_0 - R_1) \left( R_0^2 R_1^n (n(R_0 - R_1)(R_0 + 2R_1) + 6R_1^2) - R_1^2 R_0^n (-2(n-3)R_0^2 + nR_0 R_1 + nR_1^2) \right)}{4R_0^{3-n} R_1^{3-n} (R_0^n - R_1^n)^2 - n^2 R_0 R_1 (R_0^2 - R_1^2)^2} \end{aligned} \quad (6.17)$$

which we rewrite as

$$\begin{aligned} c_0 &= \frac{-2Q^3 - 2Q^{2n} + Q^n(1+Q) \left( (n-2)(n-1) - 2((n-3)n+1)Q + (n-2)(n-1)Q^2 \right)}{Q^{n-2} \left( n^2(1-Q^2)^2 + 8Q^2 \right) - 4 - 4Q^{2n}} R_1^{-3} \\ c_1 &= -\frac{-3nQ + 2(n-1) + (n+2)Q^3 + (n+2)Q^n - 3nQ^{n+2} + 2(n-1)Q^{n+3}}{4(1-Q^n)^2 - n^2 Q^{n-2} (1-Q^2)^2} R_1^{-3} \\ c_2 &= \frac{6Q^3 + 6Q^{2n+2} - Q^n(1+Q) \left( (n-1)n - (n-1)nQ - (n-1)nQ^3 + (n-1)nQ^4 + 6Q^2 \right)}{Q^n \left( n^2(1-Q^2)^2 + 8Q^2 \right) - 4Q^2 - 4Q^{2n+2}} R_1^{-1} \\ c_3 &= \frac{(1-Q) \left( Q^n (n(1-Q)(1+2Q) + 6Q^2) - Q^2 (-2(n-3) + nQ + nQ^2) \right)}{4Q^2(1-Q^n)^2 - n^2 Q^n (1-Q^2)^2} R_1^{-1} \end{aligned} \quad (6.18)$$

We note that in case  $n = 1$  the solution reduces to

$$c_0 = 0, \quad c_1 = 0, \quad c_2 = 0, \quad c_3 = -1,$$

while in case  $n = 3$  it reduces to

$$c_0 = 0, \quad c_1 = -1, \quad c_2 = 0, \quad c_3 = 0.$$

In both of these cases  $z$  is constant and we have  $\lambda = c_0 = 0$ . On the other hand, if  $n \geq 4$ , we can check that  $c_0 > 0$  for  $Q > 1$ . Recalling (6.13), we observe that  $z''$  has at most two zeros on the positive half-line and  $z''(r) > 0$  for large values of  $r$ . On the other hand, by (6.16), if  $z$  has  $N$  local extrema on  $]R_0, R_1[$ , it needs to have at least  $N + 1$  inflection points. We deduce from these conditions that  $z$  has exactly one local maximum and no local minima, and therefore  $z \geq -1$  on  $]R_0, R_1[$ . It remains to check whether  $z \leq 1$  on  $]R_0, R_1[$ . Let now

$$f(r) = r^{1-n} (r^{n-1} z(r))' = f'(r) + (n-1) \frac{f(r)}{r}.$$

Then, by (6.1), (6.16),  $f$  is a solution to the second-order elliptic problem

$$\mathcal{A}f = \lambda, \quad f(R_0) = -\frac{n-1}{R_0}, \quad f(R_1) = -\frac{n-1}{R_1},$$

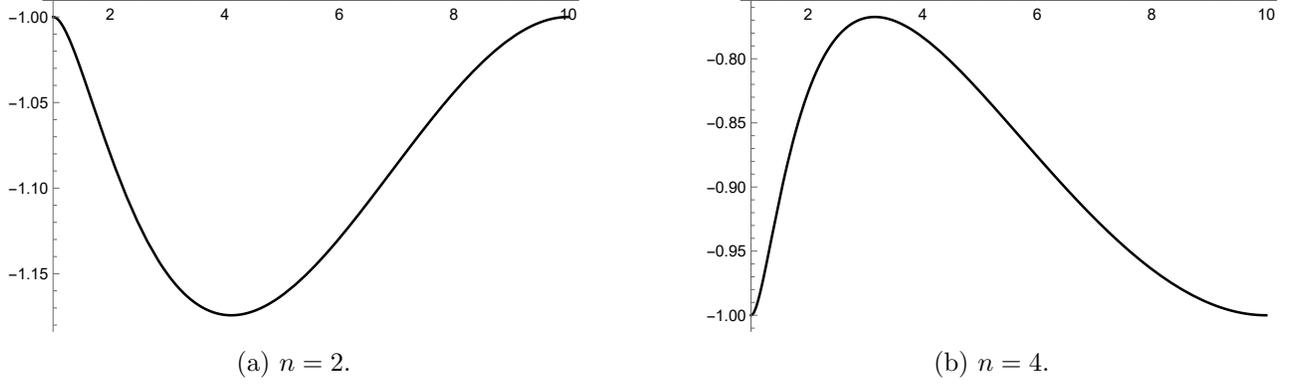


Figure 2: Plots of  $z$  for an annulus with non-constant signature for two values of  $n$ , with  $R_0 = 1$ ,  $R_1 = 10$ .

where

$$\mathcal{A}f = -r^{1-n}(r^{n-1}f'(r))' = -f''(r) - (n-1)\frac{f'(r)}{r}.$$

Since  $c_0 \geq 0$ , we have  $\lambda < 0$  for  $Q > 1$ . By the classical weak maximum principle [18, Theorem 3.1.],

$$\max_{[R_0, R_1]} f = \max_{\{R_0, R_1\}} f = -\frac{n-1}{R_0}.$$

Now, if  $z$  has a local maximum at  $r_0$ , then  $z'(r_0) = 0$ , so  $f(r_0) = \frac{z(r_0)}{r_0}$ . Consequently,

$$\frac{z(r_0)}{r_0} \leq -\frac{n-1}{R_0} < 0,$$

so  $z < 0$  on  $[R_0, R_1]$ . Thus, if  $n \neq 2$ , all annuli with non-constant signature are calibrable.

We move to the case  $n = 2$ . Now, the solution to the affine system for coefficients of  $z$  is

$$\begin{aligned} c_0 &= \frac{-R_0^2 - 2R_0R_1 \log(R_1/R_0) + R_1^2}{4R_0R_1(R_0+R_1)(-R_0^2+R_1^2 - (R_0^2+R_1^2) \log(R_1/R_0))}, \\ c_1 &= \frac{(R_0-R_1)^3}{2R_0R_1(-R_0^2+R_1^2 - (R_0^2+R_1^2) \log(R_1/R_0))}, \\ c_2 &= \frac{(R_0-R_1)(R_0+R_1)(R_0^3+4R_0R_1+R_1^2) - 2R_0(R_0^3+R_0^2R_1+R_0R_1^2+3R_1^3) \log(R_0) + 2R_1(3R_0^3+R_0^2R_1+R_0R_1^2+R_1^3) \log(R_1)}{4R_0R_1(R_0+R_1)(-R_0^2+R_1^2 - (R_0^2+R_1^2) \log(R_1/R_0))}, \\ c_3 &= \frac{R_0R_1(2(R_0^2+R_0R_1+R_1^2) \log(R_1/R_0)) + 3(R_0-R_1)(R_0+R_1)}{4(R_0+R_1)(-R_0^2+R_1^2 - (R_0^2+R_1^2) \log(R_1/R_0))} \end{aligned}$$

or equivalently

$$\begin{aligned} c_0 &= \frac{Q^2(-1-2Q \log Q + Q^2)}{4(1+Q)(-1+Q^2-(1+Q^2) \log Q)} R_1^{-3}, \\ c_1 &= \frac{(1-Q)^3}{2(-1+Q^2-(1+Q^2) \log Q)} R_1^{-1}, \\ c_2 &= \frac{(1-Q)(1+Q)(1+4Q+Q^2) + 2(1+Q+Q^2+3Q^3) \log Q}{4(1+Q)(-1+Q^2-(1+Q^2) \log Q)} R_1^{-1} + \frac{(Q-1)^3}{2(-1+Q^2-(1+Q^2) \log Q)} R_1^{-1} \log R_1, \\ c_3 &= \frac{2(1+Q+Q^2) \log Q + 3(1-Q)(1+Q)}{4(1+Q)(-1+Q^2-(1+Q^2) \log Q)} R_1. \end{aligned}$$

We can check that in this case  $c_0 < 0$  for  $Q > 1$ . By the same argument as in the previous case, we show that  $z < -1$  in  $]R_0, R_1[$ , so it does not define a valid calibration. Thus, in the case  $n = 2$  all annuli with non-constant signature are not calibrable.

## 7 Explicit solutions

### 7.1 Balls

In this section, our goal is to provide explicit description of solutions to (1.1) emanating from the characteristic function of a ball

$$u_0 = a_0 \mathbf{1}_{B_{R_0}}. \quad (7.1)$$

In the case of second-order total variation flow, the solutions with initial datum (7.1) are known to be of form

$$u(t) = a(t) \mathbf{1}_{B_{R_0}}$$

with finite extinction time, i. e. there exists  $t_* > 0$  such that  $a(t) = 0$  for  $t \geq t_*$ . In the fourth order case, based on the treatment of case  $n = 1$  in [11], we would expect the solutions to have the form

$$u(t) = a(t) \mathbf{1}_{B_{R(t)}}, \quad (7.2)$$

at least until an extinction time beyond which  $u(t, \cdot) \equiv 0$ . This intuition turns out to be correct in every dimension except  $n = 2$ .

Let first  $n \geq 3$ . As we have checked in section 6, in this case both balls and complements of balls are calibrable. Thus, as long as the solution is of form (7.2) in time instance  $t \geq 0$ , we expect a valid Cahn-Hoffman vector field  $Z$  to be given by

$$Z(x) = \begin{cases} Z_{in}(x) & \text{if } |x| \in [0, R[ \\ Z_{out}(x) & \text{if } |x| > R, \end{cases} \quad (7.3)$$

where  $Z_{in}$  is the calibration we constructed for a ball  $B_R$  and  $Z_{out}$  is the calibration we constructed for the complement of that ball, recall:

$$\lambda = -\frac{n(n+2)}{R^3}, \quad (7.4)$$

$$Z_{in}(x) = \frac{1}{2} \left( \frac{|x|}{R} \right)^3 \frac{x}{|x|} - \frac{3}{2} \frac{x}{R}, \quad Z_{out}(x) = -\frac{n-1}{2} \left( \frac{|x|}{R} \right)^{3-n} \frac{x}{|x|} + \frac{n-3}{2} \left( \frac{|x|}{R} \right)^{1-n} \frac{x}{|x|}. \quad (7.5)$$

We further calculate:

$$\begin{aligned} \operatorname{div} Z_{in}(x) &= \frac{n+2}{2} \frac{|x|^2}{R^3} - \frac{3n}{2} \frac{1}{R}, & \operatorname{div} Z_{out}(x) &= -(n-1) \frac{|x|^{2-n}}{R^{3-n}}, \\ \nabla \operatorname{div} Z_{in}(x) &= (n+2) \frac{x}{R^3}, & \nabla \operatorname{div} Z_{out}(x) &= (n-1)(n-2) \frac{|x|^{-n} x}{R^{3-n}}. \end{aligned}$$

It is straightforward to check that  $\operatorname{div} Z \in D^1(\mathbb{R}^n) \cap L^{2^*}(\mathbb{R}^n) = D_0^1(\mathbb{R}^n)$ . Next, we deduce

$$\begin{aligned} u_t &= -\Delta \operatorname{div} Z \\ &= -\Delta \operatorname{div} Z_{in} \mathcal{L}^n \llcorner_{B_R} - \Delta \operatorname{div} Z_{out} \mathcal{L}^n \llcorner_{\mathbb{R}^n \setminus B_R} + \frac{x}{|x|} \cdot (\nabla \operatorname{div} Z_{in} - \nabla \operatorname{div} Z_{out}) \mathcal{H}^{n-1} \llcorner_{\partial B_R} \\ &= -\frac{n(n+2)}{R^3} \mathcal{L}^n \llcorner_{B_R} - \frac{n(n-4)}{R^2} \mathcal{H}^{n-1} \llcorner_{\partial B_R}. \end{aligned} \quad (7.6)$$

Then, using the identity  $\frac{d}{dt} \int_{\mathbb{R}^n} u = \int_{\mathbb{R}^n} u_t$ , we obtain (recall notation (7.2))

$$a(t) \mathcal{H}^{n-1}(\partial B_{R(t)}) \frac{dR}{dt} = a(t) \frac{d}{dt} \mathcal{L}^n(B_{R(t)}) = -\frac{n(n-4)}{R^2} \mathcal{H}^{n-1}(\partial B_{R(t)}).$$

Summing up, evolution of initial datum (7.1) is given by (7.2) with  $a, R$  satisfying

$$\frac{da}{dt} = -\frac{n(n+2)}{R^3}, \quad \frac{dR}{dt} = -\frac{n(n-4)}{R^2 a}. \quad (7.7)$$

This system can be explicitly solved by noticing that

$$\frac{d}{dt}(aR^3) = -n(n+2) - 3n(n-4) = -n(4n-10)$$

and therefore

$$aR^3 = a_0R_0^3 - n(4n-10)t$$

along trajectories. The solution is

$$a(t) = a_0 \left(1 - \frac{n(4n-10)}{a_0R_0^3}t\right)^{\frac{n+2}{4n-10}}, \quad R(t) = R_0 \left(1 - \frac{n(4n-10)}{a_0R_0^3}t\right)^{\frac{n-4}{4n-10}}. \quad (7.8)$$

We note that the solution satisfies

$$\left(\frac{a}{a_0}\right)^{n-4} = \left(\frac{R}{R_0}\right)^{n+2}$$

along trajectories. (This "first integral" could also have been used to solve the system (7.7).) Let us point out a few observations concerning the solutions:

- the extinction time is equal to  $t_* = \frac{a_0R_0^3}{n(4n-10)}$ ,
- if  $n = 3$ ,  $R(t)$  is increasing and  $R(t) \rightarrow +\infty$  as  $t \rightarrow t_*^-$ ,
- if  $n = 4$ ,  $R(t) = R_0$  is constant,
- in higher dimensions,  $R(t)$  is decreasing and  $R(t) \rightarrow 0$  as  $t \rightarrow t_*^-$ .

In the case  $n = 2$  we were able to exhibit a calibration for the ball  $B_R$ , but not for its complement. Another possible ansatz on the Cahn-Hoffman vector field of form (7.3) is one where  $Z_{in}$  is the calibration we constructed for  $B_R$  and  $Z_{out}$  is the choice considered in [14]:

$$Z_{in}(x) = \frac{1}{2} \left(\frac{|x|}{R}\right)^3 \frac{x}{|x|} - \frac{3}{2} \frac{x}{R}, \quad Z_{out}(x) = -\frac{x}{|x|}. \quad (7.9)$$

We calculate

$$\operatorname{div} Z_{out} = -\frac{(n-1)}{|x|}, \quad \nabla \operatorname{div} Z_{out} = \frac{(n-1)x}{|x|^3}, \quad (7.10)$$

hence

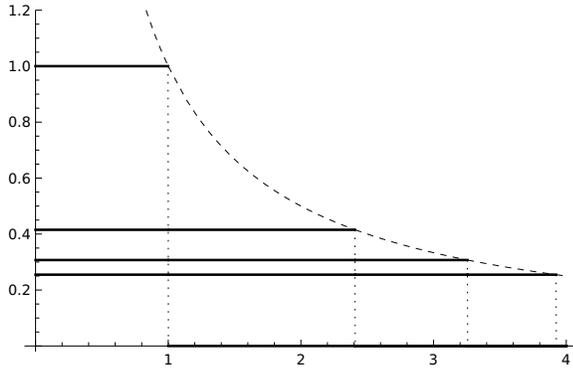
$$u_t(t, x) = -\frac{(n-1)(n-3)}{|x|^3} \quad \text{in } \mathcal{D}'(\mathbb{R}^n \setminus \overline{B_{R(t)}}). \quad (7.11)$$

If  $n \geq 4$ , this would lead to  $u(t)$  being radially strictly increasing for positive  $t$  and large values of  $|x|$ , which would be at odds with our choice of  $Z_{out}$ . In fact, if  $n \geq 4$ ,  $\operatorname{div} Z \notin D_0^1(\mathbb{R}^n)$  for any  $Z$  of this form. However, in smaller dimensions this ansatz remains a viable option. If  $n = 3$ , it leads to the same solution as before. On the other hand, if  $n = 2$ , we obtain a solution which is not of form (7.2). Instead, we are led to assume

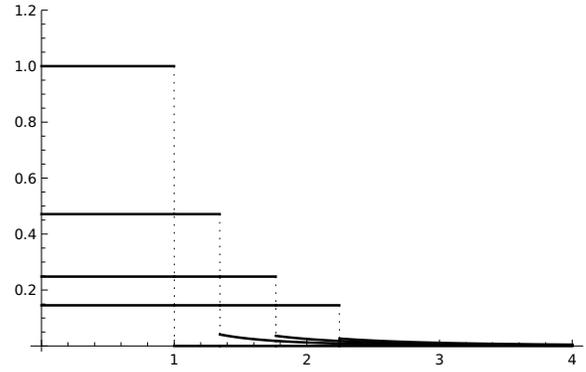
$$u(t, x) = a(t)\mathbf{1}_{B_{R(t)}} + \frac{t}{|x|^3}\mathbf{1}_{\mathbb{R}^2 \setminus B_{R(t)}}. \quad (7.12)$$

We have:

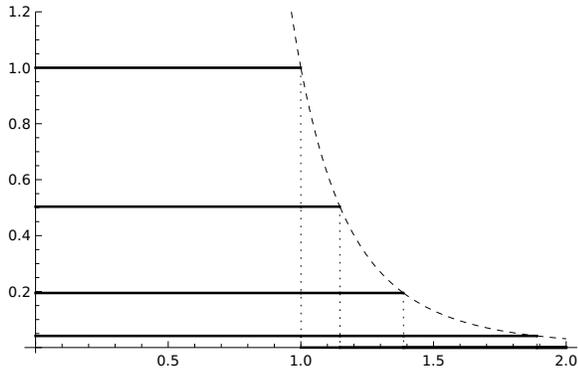
$$\begin{aligned} \operatorname{div} Z_{in}(x) &= 2\frac{|x|^2}{R^3} - 3\frac{1}{R}, & \operatorname{div} Z_{out}(x) &= -\frac{1}{|x|}, \\ \nabla \operatorname{div} Z_{in}(x) &= 4\frac{x}{R^3}, & \nabla \operatorname{div} Z_{out}(x) &= \frac{x}{|x|^3}, \\ u_t &= -\Delta \operatorname{div} Z = -\frac{8}{R^3}\mathcal{L}^2 \llcorner_{B_R} + \frac{1}{|x|^3}\mathcal{L}^2 \llcorner_{\mathbb{R}^2 \setminus B_R} + \frac{3}{R^2}\mathcal{H}^1 \llcorner_{\partial B_R} \end{aligned}$$



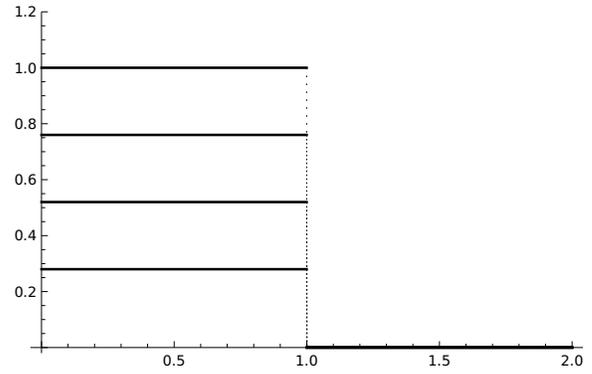
(a) Case  $n = 1$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.8, t = 1.6, t = 2.4$ .



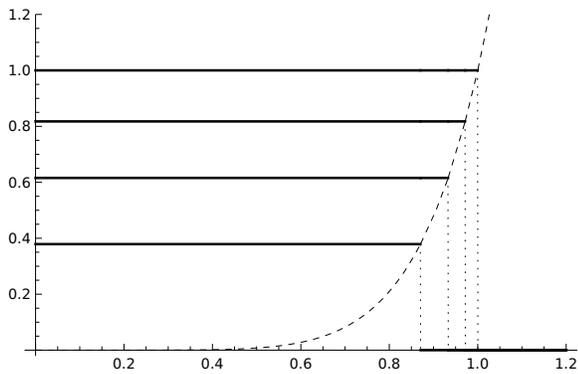
(b) Case  $n = 2$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.1, t = 0.2, t = 0.3$ .



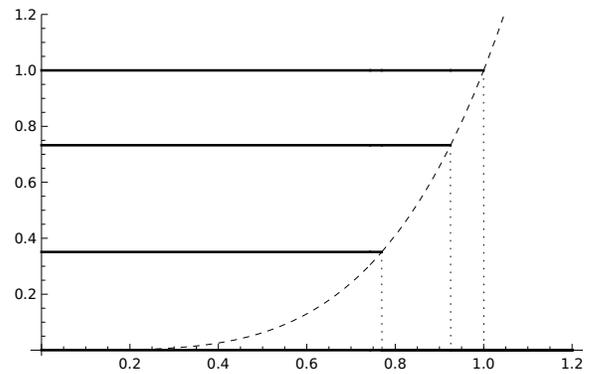
(c) Case  $n = 3$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.04, t = 0.08, t = 0.12$ .



(d) Case  $n = 4$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.01, t = 0.02, t = 0.03$ .



(e) Case  $n = 5$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.005, t = 0.01, t = 0.015$ .



(f) Case  $n = 6$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.005, t = 0.01, t = 0.015$ .

Figure 3: Plots of the solution  $u(t, x)$  emanating from the characteristic function of the unit ball as a function of  $|x|$  for chosen values of  $t$ .

and, recalling (7.12),

$$\left(a(t) - \frac{t}{R(t)^3}\right) \mathcal{H}^1(\partial B_{R(t)}) \frac{dR}{dt} = \left(a(t) - \frac{t}{R(t)^3}\right) \frac{d}{dt} \mathcal{L}^2(B_{R(t)}) = \frac{3}{R(t)^2} \mathcal{H}^1(\partial B_{R(t)}).$$

Thus, we arrive at ODE system

$$\frac{da}{dt} = -\frac{8}{R^3}, \quad \frac{dR}{dt} = \frac{3R}{aR^3 - t}. \quad (7.13)$$

This system is not autonomous, but it can be integrated by noticing that along trajectories

$$\frac{d}{dt} (aR^3 - t) = \frac{9aR^3}{aR^3 - t} - 8 - 1 = \frac{9t}{aR^3 - t}$$

and so

$$aR^3 = \sqrt{a_0^2 R_0^6 + 9t^2} + t.$$

This implies, first of all, that

$$a(t) > \frac{t}{R(t)^3}$$

for all  $t > 0$  and the form of solution (7.12) is preserved as long as the solution does not vanish. Furthermore, we can rewrite the system (7.13) in decoupled form

$$\frac{d}{dt} \log a = \frac{-8}{\sqrt{a_0^2 R_0^6 + 9t^2} + t}, \quad \frac{d}{dt} \log R^3 = \frac{9}{\sqrt{a_0^2 R_0^6 + 9t^2}}. \quad (7.14)$$

These equations can be explicitly integrated:

$$a(t) = a_0 \frac{a_0^4 R_0^{12} + 8a_0^2 R_0^{12} t^2}{\left(\sqrt{a_0^2 R_0^6 + 9t^2} + 3t\right)^2 \left(a_0^2 R_0^6 + 6t^2 + 2t\sqrt{a_0^2 R_0^6 + 9t^2}\right)},$$

$$R(t) = R_0 \sqrt{1 + 6t \frac{3t + \sqrt{a_0^2 R_0^6 + 9t^2}}{a_0^2 R_0^6}}.$$

We observe that the solutions exist globally and

$$\lim_{t \rightarrow \infty} a(t) = 0, \quad \lim_{t \rightarrow \infty} R(t) = \infty.$$

In particular  $u$  stays in the form (7.12) for all  $t > 0$ .

Finally we consider  $n = 1$ . In this case, both ansätze considered before lead to the same solution:

$$Z_{in}(x) = \frac{1}{2} \left(\frac{x}{R}\right)^3 - \frac{3}{2} \frac{x}{R}, \quad Z_{out}(x) = -\operatorname{sgn} x \quad (7.15)$$

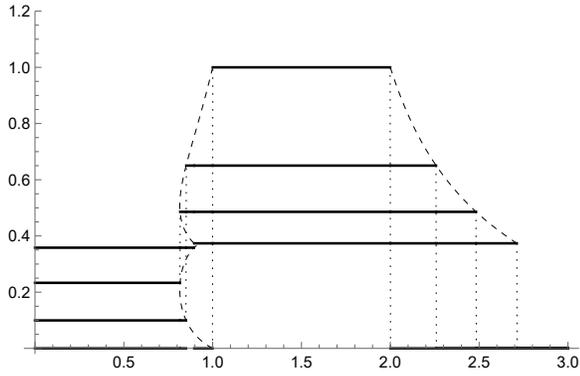
which coincides with (7.5). Repeating the calculations following (7.5), we obtain a solution of form (7.2) satisfying (7.8), i. e.

$$u(t) = a(t) \mathbf{1}_{B_{R(t)}}, \quad a(t) = a_0 \left(1 + \frac{6}{a_0 R_0^3} t\right)^{-\frac{1}{2}}, \quad R(t) = R_0 \left(1 + \frac{6}{a_0 R_0^3} t\right)^{\frac{1}{2}}. \quad (7.16)$$

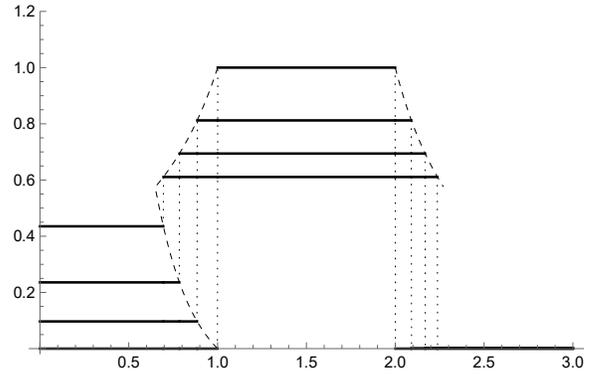
Note that now, as opposed to the case  $n \geq 3$ , the coefficient multiplying  $t$  is positive. Like in  $n = 2$ , the extinction time is infinite and we have

$$\lim_{t \rightarrow \infty} a(t) = 0, \quad \lim_{t \rightarrow \infty} R(t) = \infty.$$

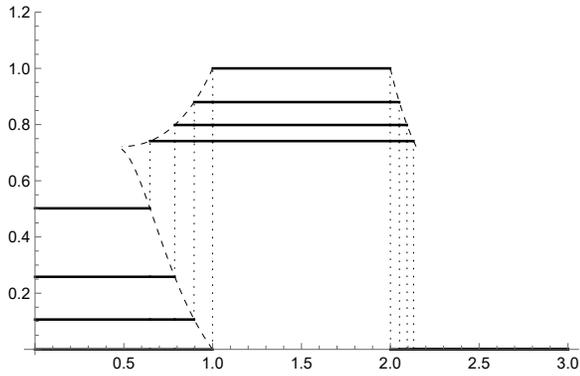
This concludes the proof of Theorem 4.



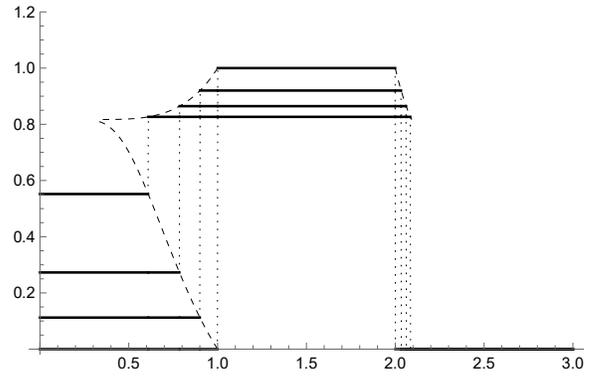
(a) Case  $n = 1$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.025, t = 0.5, t = 0.075$ .



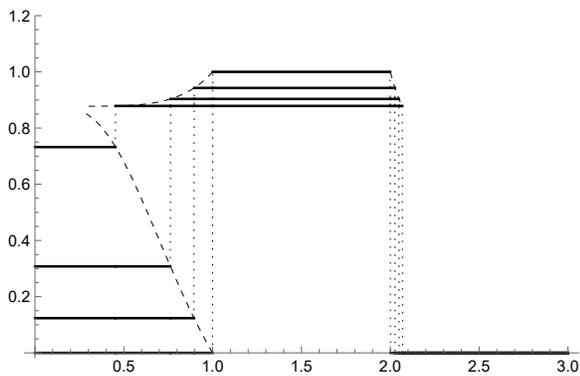
(b) Case  $n = 2$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.01, t = 0.02, t = 0.03$ .



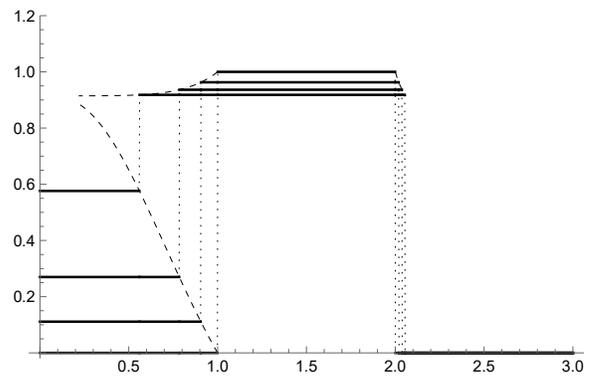
(c) Case  $n = 3$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.006, t = 0.012, t = 0.018$ .



(d) Case  $n = 4$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.004, t = 0.008, t = 0.012$ .



(e) Case  $n = 5$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.003, t = 0.006, t = 0.009$ .



(f) Case  $n = 6$ . Solid lines: plots of  $u(t, \cdot)$  for  $t = 0, t = 0.002, t = 0.004, t = 0.006$ .

Figure 4: Plots of the solution  $u(t, x)$  emanating from the characteristic function of annulus  $A_{R^1}^{R^2}$  as a function of  $|x|$  for chosen values of  $t$  with  $R^1 = 1, R^2 = 2$ .

## 7.2 Stacks

Using the calibrations we constructed for generalized annuli, we will now derive a system of ODEs locally prescribing the solution emanating from any piecewise constant, radially symmetric datum (a *stack*).

**Definition 31.** Let  $w \in D(TV)$ . We say that  $w$  is a stack if there exists a number  $N \in \mathbb{N}$  and sequences  $0 < R^0 < R^1 < \dots < R^{N-1}$ ,  $a^0, a^1, \dots, a^N$  with  $a^k \in \mathbb{R}$  such that

$$w = a^0 \mathbf{1}_{B_{R^0}} + a^1 \mathbf{1}_{A_{R^0}^{R^1}} + \dots + a^{N-1} \mathbf{1}_{A_{R^{N-2}}^{R^{N-1}}} + a^N \mathbf{1}_{\mathbb{R}^n \setminus B_{R^{N-1}}}.$$

Suppose first that  $n \neq 2$ , in which case all connected components of level sets of any stack  $w$  are calibrable. Let  $u_0$  be a stack

$$u_0 = a_0^0 \mathbf{1}_{B_{R_0^0}} + a_0^1 \mathbf{1}_{A_{R_0^0}^{R_0^1}} + \dots + a_0^{N-1} \mathbf{1}_{A_{R_0^{N-2}}^{R_0^{N-1}}} + a_0^N \mathbf{1}_{\mathbb{R}^n \setminus B_{R_0^{N-1}}}, \quad (7.17)$$

where  $a^{k-1} \neq a^k$  for  $k = 1, \dots, N$ ,  $a_0^N = 0$ . We expect that if  $u$  is the solution emanating from  $u_0$ , then  $u(t, \cdot)$  is a stack of form

$$u(t, \cdot) = a^0(t) \mathbf{1}_{B_{R^0(t)}} + a^1(t) \mathbf{1}_{A_{R^0(t)}^{R^1(t)}} + \dots + a^{N-1}(t) \mathbf{1}_{A_{R^{N-2}(t)}^{R^{N-1}(t)}} + a^N(t) \mathbf{1}_{\mathbb{R}^n \setminus B_{R^{N-1}(t)}}, \quad (7.18)$$

with  $a^N(t) = 0$  for all  $t > 0$ , and that  $a^{k-1} \neq a^k$ ,  $k = 1, \dots, N$  for small  $t$ . We construct a Cahn-Hoffman vector field  $Z(t, \cdot)$  for  $u(t, \cdot)$  by pasting together calibrations  $Z^k$  for  $B_{R^0(t)}$ ,  $A_{R^k(t)}^{R^{k+1}(t)}$ ,  $\mathbb{R}^n \setminus B_{R^{N-1}(t)}$  with suitable choice of signatures. We have

$$\begin{aligned} u_t = -\Delta \operatorname{div} Z &= -\Delta \operatorname{div} Z^0 \mathcal{L}^n \llcorner_{B_{R^0}} - \sum_{k=1}^n \Delta \operatorname{div} Z^k \mathcal{L}^n \llcorner_{A_{R^{k-1}}^{R^k}} - \Delta \operatorname{div} Z^N \mathcal{L}^n \llcorner_{\mathbb{R}^n \setminus B_{R^N}} \\ &+ \sum_{k=0}^n \frac{x}{|x|} \cdot (\nabla \operatorname{div} Z^k - \nabla \operatorname{div} Z^{k+1}) \mathcal{H}^{n-1} \llcorner_{S_{R^k}}. \end{aligned} \quad (7.19)$$

We denote

$$\begin{aligned} &\left. \frac{x}{|x|} \cdot (\nabla \operatorname{div} Z^k - \nabla \operatorname{div} Z^{k+1}) \right|_{S_{R^k}} \\ &= z_{rr}^k(R^k) - z_{rr}^{k+1}(R^k) + \frac{n-1}{R^k} (z_r^k(R^k) - z_r^{k+1}(R^k)) - \frac{n-1}{(R^k)^2} (z^k(R^k) - z^{k+1}(R^k)) \\ &= z_{rr}^k(R^k) - z_{rr}^{k+1}(R^k) =: d^k. \end{aligned}$$

The values of  $d^k$  are functions of  $R^0, \dots, R^{N-1}$ . Assuming that  $R^k$  are regular enough and  $\varepsilon, |t-s|$  are small enough, we have

$$\frac{d}{dt} \int_{A_{R^k(s)-\varepsilon}^{R^k(s)+\varepsilon}} u = \int_{A_{R^k(s)-\varepsilon}^{R^k(s)+\varepsilon}} u_t,$$

whence

$$(a^k(t) - a^{k-1}(t)) \mathcal{H}^{n-1}(S_{R^k(t)}) \frac{dR^k}{dt} = (a^k(t) - a^{k-1}(t)) \frac{d}{dt} \mathcal{L}^n(A_{R^k(s)-\varepsilon}^{R^k(t)}) = d^k(t) \mathcal{H}^{n-1}(S_{R^k(t)}).$$

Further, for  $k = 0, \dots, N$ , we denote by  $\lambda^k$  the value of  $-\Delta \operatorname{div} Z^k(t, \cdot)$  which is constant since  $Z^k$  is a calibration. Then, we can write down the system of ODEs for  $a^k$  and  $R^k$ :

$$\frac{da^k}{dt} = \lambda^k \text{ for } k = 0, \dots, N, \quad \frac{dR^k}{dt} = \frac{d^k}{a^k - a^{k+1}} \text{ for } k = 0, \dots, N-1. \quad (7.20)$$

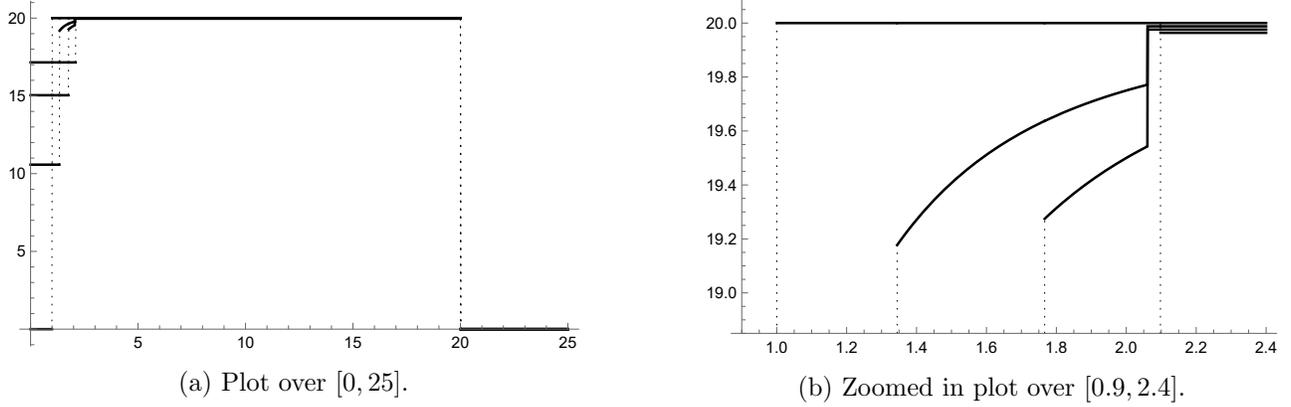


Figure 5: Plots of the solution  $u(t, x)$  emanating from the characteristic function of annulus  $A_{R_1}^{R_2}$  with  $R_0 = 1$ ,  $R_1 = 20$  in  $n = 2$  as a function of  $|x|$  for  $t = 0$ ,  $t = 2$ ,  $t = 4$ ,  $t = 6$ .

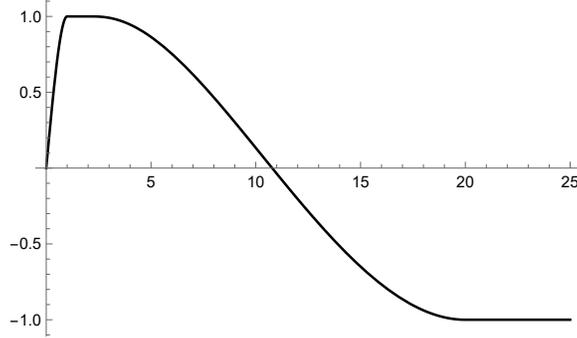


Figure 6: Plot of the Cahn-Hoffman vector field  $Z(t, x)$  for the characteristic function of annulus  $A_{R_1}^{R_2}$  with  $R_0 = 1$ ,  $R_1 = 20$  in  $n = 2$  as a function of  $|x|$ .

Let  $c_0^k$  denote  $c_0$  given by (6.17) if  $\text{sgn}(a^{k+1} - a^k) = \text{sgn}(a^k - a^{k-1})$  or by (6.10) if  $\text{sgn}(a^{k+1} - a^k) \neq \text{sgn}(a^k - a^{k-1})$ , with  $R^{k+1}$  and  $R^k$  in place of  $R_1$  and  $R_0$ . Then, we have

$$\lambda^0 = \text{sgn}(a^1 - a^0) \frac{n(n+2)}{(R^0)^3}, \quad \lambda^k = 2n(n+2) \text{sgn}(a^{k+1} - a^k) c_0^k \text{ for } k = 1, \dots, N-1, \quad \lambda^N = 0. \quad (7.21)$$

We observe that in a neighborhood of any initial datum  $R_0^0, \dots, R_0^{N-1}, a_0^0, \dots, a_0^N$ ,  $R_0^k < R_0^{k+1}$ ,  $a_0^k \neq a_0^{k+1}$ , the r.h.s. of (7.20) is regular in  $R^0, \dots, R^{N-1}, a^0, \dots, a^N$ , so locally the system has a unique solution. Unique solvability fails when a time instance  $t > 0$  is reached such that  $a^k(t) = a^{k+1}(t)$ ,  $R^k(t) = R^{k+1}(t)$  or  $R_0 = 0$ . In such case  $u(t, \cdot)$  is again a stack with a smaller  $N$ , and we can restart our procedure. This concludes the proof of Theorem 5.

Next we deal with the remaining case of dimension  $n = 2$ . In this case, our attempt to obtain a radial calibration failed for complements of balls and for some annuli. Again, let  $u_0$  be a stack of form (7.17). For  $k = 1, \dots, N$ , let  $\sigma^k = \text{sgn}(a_0^k - a_0^{k-1})$ . We assume the following ansatz on the solution  $u$  and the associated field  $Z$  for small  $t > 0$ :

$$u(t, \cdot) = a^0(t) \text{ on } B_{R^0(t)}, u(t, \cdot) = a^k(t) \text{ on } A_{\max(R^{k-1}(t), R^k(t)/Q_*)}^{R^k(t)} \text{ if } \sigma^{k+1} \neq \sigma^k, \quad k = 1, \dots, N-1, \quad (7.22)$$

$$Z(t, x) = \sigma^k \frac{x}{|x|} \text{ on } A_{R^{k-1}(t)}^{R^k(t)} \text{ if } \sigma^{k+1} \neq \sigma^k \text{ or on } A_{R^{k-1}(t)}^{R^k(t)/Q_*} \text{ if } \sigma^{k+1} = \sigma^k, \quad k = 1, \dots, N-1,$$

$$Z(t, x) = \sigma^{k+1} \frac{x}{|x|} \text{ on } \mathbb{R}^n \setminus B_{R^{N-1}(t)}. \quad (7.23)$$

We complete the definition of a Cahn-Hoffman field  $Z$  consistent with (7.22), (7.23) by pasting the calibrations  $Z^k$  with suitable choice of signatures into the gaps left in (7.23). This leads to

$$u_t(t, \cdot) = \lambda^0(t) \text{ in } \mathcal{D}'(B_{R^0(t)}),$$

$$u_t(t, x) = \lambda^k(t) \text{ in } \mathcal{D}'\left(A_{\max(R^{k-1}(t), R^k(t)/Q_*)}^{R^k(t)}\right),$$

$$u_t(t, x) = \frac{\sigma^k}{|x|^3} \text{ in } \mathcal{D}'\left(A_{R^{k-1}(t)}^{R^k(t)/Q_*}\right) \text{ if } \sigma^k \neq \sigma^{k+1} \text{ or in } \mathcal{D}'\left(A_{R^{k-1}(t)}^{R^k(t)}\right) \text{ if } \sigma^k = \sigma^{k+1}, \quad k = 1, \dots, N-1,$$

$$u_t(t, x) = \frac{\sigma^N}{|x|^3} \text{ in } \mathcal{D}'(\mathbb{R}^2 \setminus B_{R^N(t)}).$$

Moreover,  $u_t(t, \cdot) \in M(\mathbb{R}^2)$  and

$$u_t \llcorner S_{R^k} = \frac{x}{|x|} \cdot ((\nabla \operatorname{div} Z)^- - (\nabla \operatorname{div} Z)^+) \mathcal{H}^1 \llcorner S_{R^k} =: d^k$$

for  $k = 0, \dots, N-1$ , where  $(\nabla \operatorname{div} Z)^\pm$  are the one sided limits as  $|x| \rightarrow (R^k)^\pm$ . The values of  $d^k$  are functions of  $R^0, \dots, R^{N-1}$ . Reasoning as in the case  $n \neq 2$ , the evolution of  $R^k$  is governed by equations

$$\frac{dR^k}{dt} = \frac{d^k}{u(t, x)|_{|x|=(R^k)^-} - u(t, x)|_{|x|=(R^k)^+}}. \quad (7.24)$$

The values  $u(t, x)|_{|x|=(R^k)^+}$  are either prescribed by ODEs

$$\frac{da^k}{dt} = \lambda^k \quad (7.25)$$

with  $\lambda^k$  functions of  $R^0, \dots, R^{N-1}$  in calibrable regions where  $u(t, x) = a^k(t)$ , or explicitly determined by  $u_t(t, x) = \sigma^k/|x|^3$  in bending regions. It is important to note that in the case  $\sigma^k \neq \sigma^{k+1}$ ,  $R^{k-1} \leq R^k/Q_*$  the functions  $d^k, \lambda^k$  do not depend on  $R^{k-1}$ . Thus, one can first solve a part of the system (7.24), (7.25) for the outer annuli, then calculate  $u$  in the bending region (without knowing a priori its inner boundary) and move on to solving innermore parts of (7.24), (7.25). This way, finding the solution is indeed again reduced to solving a system of ODEs.

## References

- [1] F. Alter, V. Caselles and A. Chambolle, A characterization of convex calibrable sets in  $\mathbb{R}^N$ , *Math. Ann.*, **332** (2005), 329–366.
- [2] L. Ambrosio, N. Luigi and G. Savaré, *Gradient flows in metric spaces and in the space of probability measures*, Second edition, Lectures in Mathematics ETH Zürich, Birkhäuser Verlag, Basel, 2008.
- [3] F. Andreu-Vaillo, V. Caselles and J. M. Mazón, *Parabolic quasilinear equations minimizing linear growth functionals*, Progress in Mathematics, 223, Birkhäuser Verlag, Basel, 2004.
- [4] G. Bellettini, V. Caselles and M. Novaga, The total variation flow in  $\mathbb{R}^N$ , *J. Differential Equations*, **184** (2002), 475–525.
- [5] G. Bellettini, M. Novaga and M. Paolini, Characterization of facet breaking for nonsmooth mean curvature flow in the convex case, *Interfaces Free Bound.*, **3** (2001), 415–446.

- [6] H. Brézis, *Opérateurs maximaux monotones et semi-groupes de contractions dans les espaces de Hilbert*, North-Holland Mathematics Studies, No. 5, Notas de Matemática, No. 50, North-Holland Publishing Co., Amsterdam-London; American Elsevier Publishing Co., Inc., New York, 1973.
- [7] W. C. Carter, A. R. Roosen, J. W. Cahn and J. E. Taylor, Shape evolution by surface diffusion and surface attachment limited kinetics on completely faceted surfaces, *Acta Metall. Mater.*, **43** (1995), 4309–4323.
- [8] C. M. Elliott and S. A. Smitheman, Analysis of the TV regularization and  $H^{-1}$  fidelity model for decomposing an image into cartoon plus texture, *Commun. Pure Appl. Anal.*, **6** (2007), 917–936.
- [9] W. Fulton, *Algebraic topology. A first course*, Graduate Texts in Mathematics, 153, Springer-Verlag, New York, 1995.
- [10] G. P. Galdi, *An introduction to the mathematical theory of the Navier-Stokes equations. Steady-state problems*, Second edition, Springer Monographs in Mathematics, Springer, New York, 2011.
- [11] M.-H. Giga and Y. Giga, Very singular diffusion equations: second and fourth order problems, *Japanese J. Ind. Appl. Math.*, **27** (2010), 323–345.
- [12] M.-H. Giga and Y. Giga, Crystalline surface diffusion flow for graph-like curves, *Hokkaido University Preprint Series in Mathematics*, **1142** (2022).
- [13] Y. Giga and R. V. Kohn, Scale-invariant extinction time estimates for some singular diffusion equations, *Discrete Contin. Dyn. Syst.*, **30** (2011), 509–535.
- [14] Y. Giga, H. Kuroda and H. Matsuoka, Fourth-order total variation flow with Dirichlet condition: Characterization of evolution and extinction time estimates, *Adv. Math. Sci. Appl.*, **24** (2014), 499–534.
- [15] Y. Giga, M. Muszkieta and P. Rybka, A duality based approach to the minimizing total variation flow in the space  $H^{-s}$ , *Jpn. J. Ind. Appl. Math.*, **36** (2019), 261–286.
- [16] Y. Giga and N. Požár, Motion by crystalline-like mean curvature: a survey, *Hokkaido University Preprint Series in Mathematics*, **1140** (2021).
- [17] Y. Giga and Y. Ueda, Numerical computations of split Bregman method for fourth order total variation flow, *J. Comput. Phys.*, **405** (2020), 109114.
- [18] D. Gilbarg and N. S. Trudinger, *Elliptic partial differential equations of second order*, Reprint of the 1998 edition, Classics in Mathematics, Springer-Verlag, Berlin, 2001.
- [19] Y. Kashima, A subdifferential formulation of fourth order singular diffusion equations, *Adv. Math. Sci. Appl.*, **14** (2004), 49–74.
- [20] Y. Kashima, Characterization of subdifferentials of a singular convex functional in Sobolev spaces of order minus one, *J. Funct. Anal.*, **262** (2012), 2833–2860.
- [21] R. V. Kohn, Surface relaxation below the roughening temperature: some recent progress and open questions, In: *Nonlinear partial differential equations*, Abel Symp., 7, Springer, Heidelberg, 2012, 207–221.
- [22] R. V. Kohn and H. M. Versieux, Numerical analysis of a steepest-descent PDE model for surface relaxation below the roughening temperature, *SIAM J. Numer. Anal.*, **48** (2010), 1781–1800.
- [23] Y. Kōmura, Nonlinear semi-groups in Hilbert space, *J. Math. Soc. Japan*, **19** (1967), 493–507.
- [24] M. Łasica, S. Moll and P. B. Mucha, Total variation denoising in  $l^1$  anisotropy. *SIAM J. Imaging Sci.*, **10** (2017), 1691–1723.
- [25] G. P. Leonardi, An overview on the Cheeger problem, In: *New trends in shape optimization*, Internat. Ser. Numer. Math., 166, Birkhäuser/Springer, Cham, 2015, 117–139.

- [26] I. V. Odisharia, Simulation and analysis of the relaxation of a crystalline surface, Ph.D. thesis, New York University, New York 2006.
- [27] S. Osher, A. Solé and L. Vese, Image decomposition and restoration using total variation minimization and the  $H^{-1}$  norm, *Multiscale Model. Simul.*, **1** (2003), 349–370.
- [28] W. Rudin, *Functional analysis*, Second edition, International Series in Pure and Applied Mathematics, McGraw-Hill, Inc., New York, 1991.
- [29] L. Schwartz, *Théorie des distributions*, Nouvelle édition, entièrement corrigée, refondue et augmentée, Publications de l’Institut de Mathématique de l’Université de Strasbourg, IX-X Hermann, Paris 1966.
- [30] H. Spohn, Surface dynamics below the roughening transition, *J. Phys. I France*, **3** (1993), 69–81.

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