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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 Ω 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 Ω , 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 Ω . By the Poincaré inequality, both H_{av}^1 and $D_0^1(\Omega)$ can be regarded as subspaces of $L^2(\Omega)$. However, if Ω 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 ψ 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 ∂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 ∂U , where ν is the exterior unit normal of ∂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 ∂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 ∂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 λ . If U is bounded, this constant λ 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 ∂B_{R_1} and the interior boundary ∂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 Γ and ν 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 λ 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 ϑ_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 (ε_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 η_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, η_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 η_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 ∇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 ∇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 $\Psi: D^{-1} \rightarrow [0, \infty]$ by 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, Ψ 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) $|Z| \leq 1$
- (ii) $v = \Delta \operatorname{div} Z$
- (iii) $-\langle u, \operatorname{div} Z \rangle = TV(u)$.

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

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

By definition of Ψ , 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 (δ, 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 ψ . 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_ψ^{-1} is independent of the choice of ψ , we suppress ψ 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 ψ , 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 Ψ 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, Ψ 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 ψ . 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 ∂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 ∂U must equal to the minus mean curvature (the sum of all principal curvatures) in the direction of ν . 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 ∂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 χ . We say that U is (D^{-1}) -calibrable (with signature χ) 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 χ).

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 χ . 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 λ 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 λ 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 χ 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 λ 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_{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 λ 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 χ . 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 χ . Then there exists a calibration \bar{Z} for (U, χ) of form $\bar{Z}(x) = z(|x|) \frac{x}{|x|}$.*

Proof. Let Z be any calibration for (U, χ) . Let μ_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, χ) . 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 ∂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 λ, 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 ∂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 ∂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

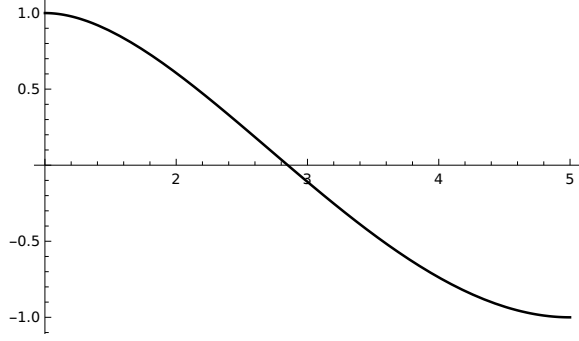
$$\begin{aligned} c_0 &= \frac{-(R_1 - R_0)R_1^n R_0^n ((n-1)(n-2)(R_1 + R_0)^2 - 2R_0 R_1) + 2R_1^3 R_0^{2n} - 2R_0^3 R_1^{2n}}{R_1 R_0 (R_1^n R_0^n (n^2 (R_1^2 - R_0^2)^2 + 8R_1^2 R_0^2) - 4R_1^2 R_0^{2n+2} - 4R_0^2 R_1^{2n+2})}, \\ c_1 &= \frac{(R_1 + R_0) (R_1^n (2(n-1)R_1^2 + (n+2)R_1 R_0 - (n+2)R_0^2)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

This can be rewritten in a form emphasizing homogeneity:

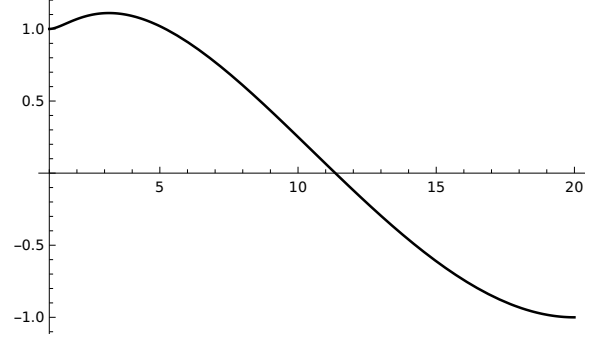
$$\begin{aligned} c_0 &= \frac{-(Q-1)Q^n ((n-1)(n-2)(Q+1)^2 - 2Q) + 2Q^3 - 2Q^{2n}}{Q^{n-2} (n^2 (Q^2 - 1)^2 + 8Q^2) - 4 - 4Q^{2n}} R_1^{-3}, \\ c_1 &= \frac{(Q+1) (Q^n (2(n-1)Q^2 + (n+2)Q - (n+2)) + ((n+2)Q^2 - (n+2)Q - 2(n-1)))}{2(n^2 - 4)Q^n + 4 + 4Q^{2n} - n^2 Q^{n+2} - n^2 Q^{n-2}} R_1^{n-3}, \\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

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 ∂B_{R_0} and $\chi = -1$ on ∂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

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

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

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

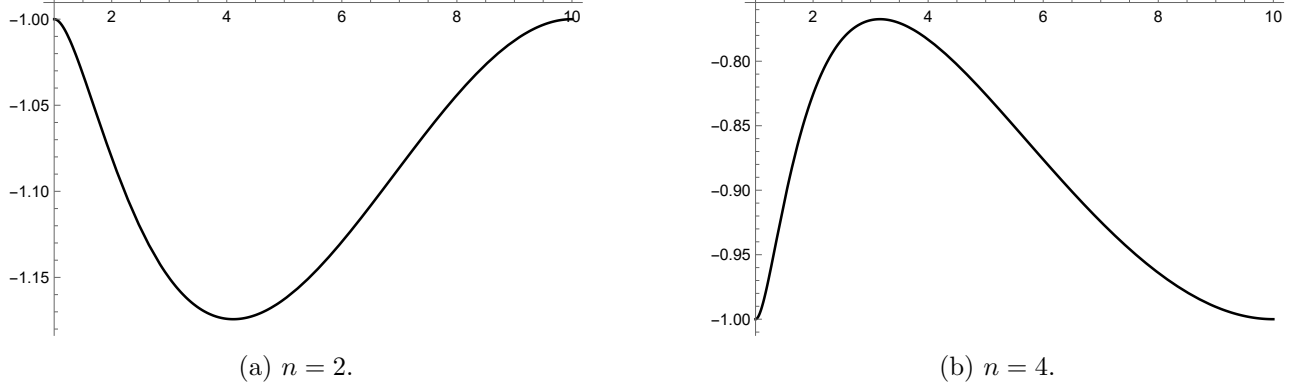


Figure 2: Plots of z for an annulus with non-constant signature for two values of n , with $R_0 = 1$, $R_1 = 10$.

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

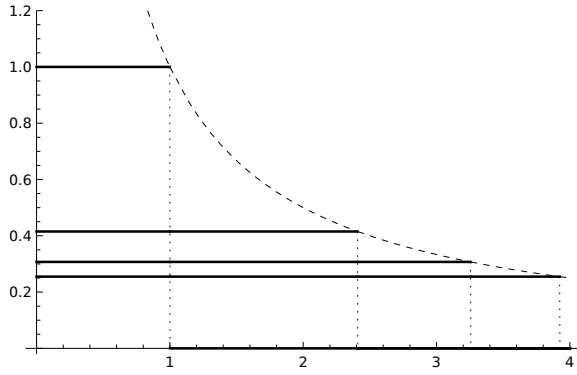
$$u_t(t, x) = -\frac{(n-1)(n-3)}{|x|^3} \quad \text{in } \mathcal{D}'(\mathbb{R}^n \setminus \overline{B_{R(t)}}). \quad (7.11)$$

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

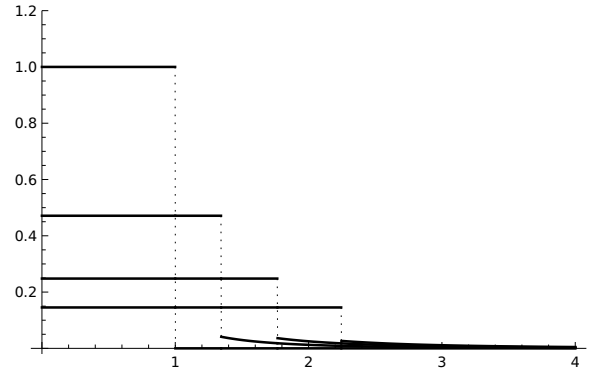
$$u(t, x) = a(t)\mathbf{1}_{B_{R(t)}} + \frac{t}{|x|^3}\mathbf{1}_{\mathbb{R}^2 \setminus B_{R(t)}}. \quad (7.12)$$

We have:

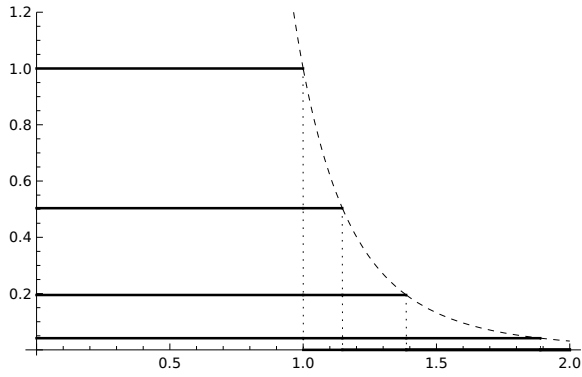
$$\begin{aligned} \operatorname{div} Z_{in}(x) &= 2\frac{|x|^2}{R^3} - 3\frac{1}{R}, & \operatorname{div} Z_{out}(x) &= -\frac{1}{|x|}, \\ \nabla \operatorname{div} Z_{in}(x) &= 4\frac{x}{R^3}, & \nabla \operatorname{div} Z_{out}(x) &= \frac{x}{|x|^3}, \\ u_t &= -\Delta \operatorname{div} Z = -\frac{8}{R^3}\mathcal{L}^2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



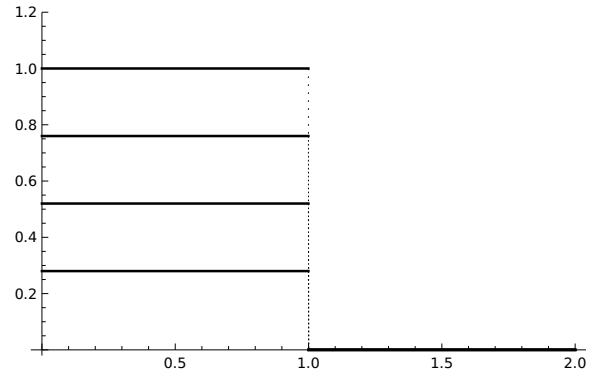
(a) Case $n = 1$. Solid lines: plots of $u(t, \cdot)$ for $t = 0, t = 0.8, t = 1.6, t = 2.4$.



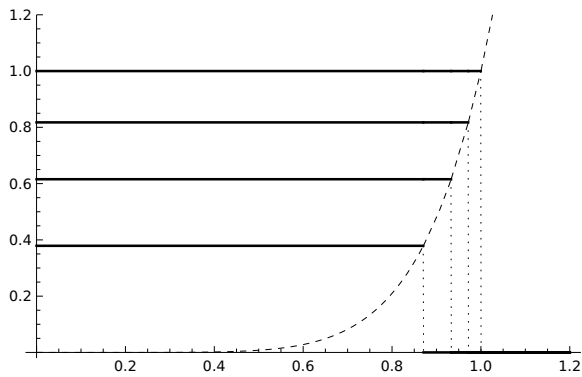
(b) Case $n = 2$. Solid lines: plots of $u(t, \cdot)$ for $t = 0, t = 0.1, t = 0.2, t = 0.3$.



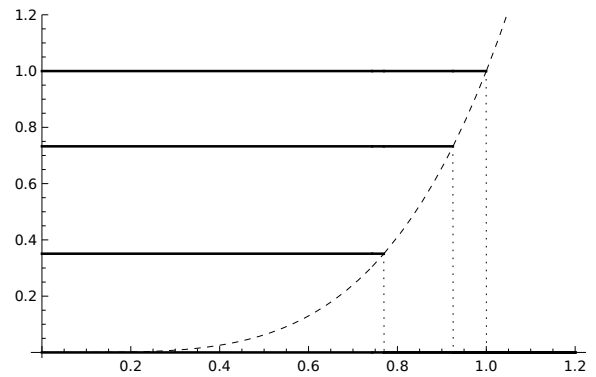
(c) Case $n = 3$. Solid lines: plots of $u(t, \cdot)$ for $t = 0, t = 0.04, t = 0.08, t = 0.12$.



(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:

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

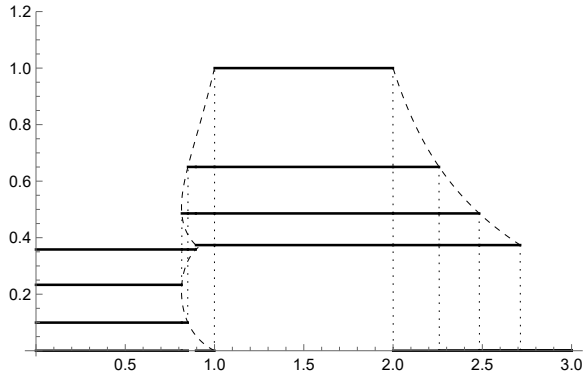
which coincides with (7.5). Repeating the calculations following (7.5), we obtain a solution of form (7.2) satisfying (7.8), i. e.

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

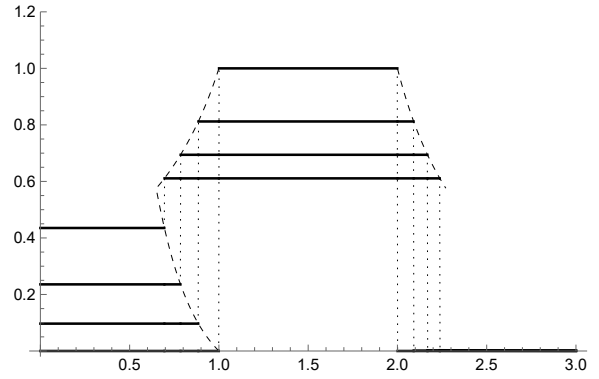
Note that now, as opposed to the case $n \geq 3$, the coefficient multiplying t is positive. Like in $n = 2$, the extinction time is infinite and we have

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

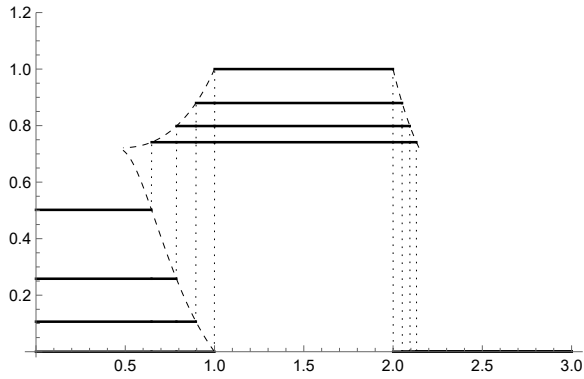
This concludes the proof of Theorem 4.



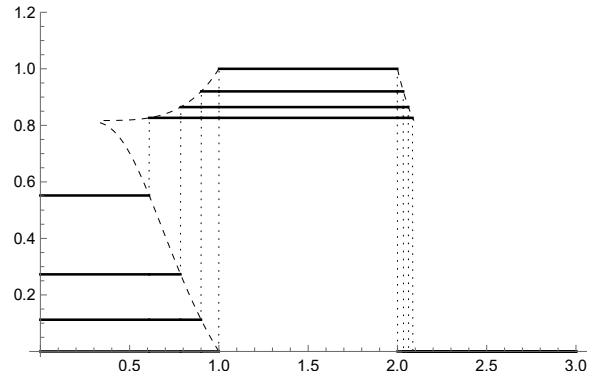
(a) Case $n = 1$. Solid lines: plots of $u(t, \cdot)$ for $t = 0, t = 0.025, t = 0.5, t = 0.075$.



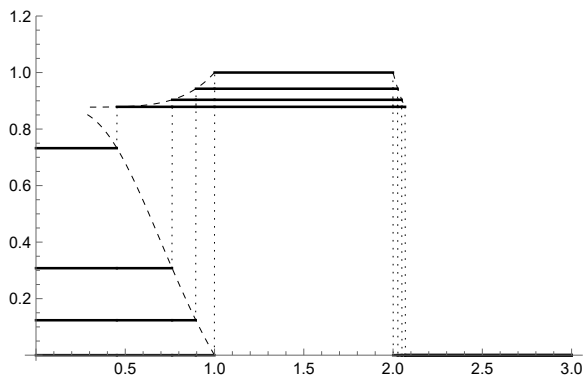
(b) Case $n = 2$. Solid lines: plots of $u(t, \cdot)$ for $t = 0, t = 0.01, t = 0.02, t = 0.03$.



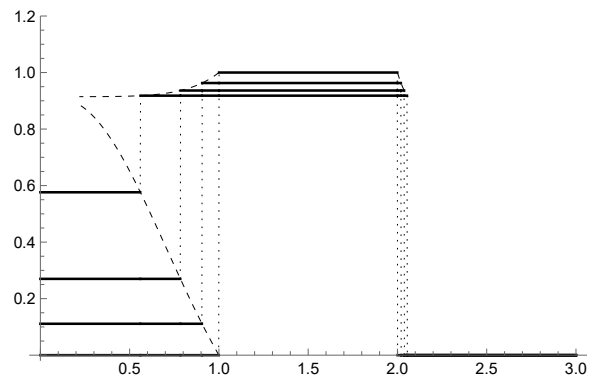
(c) Case $n = 3$. Solid lines: plots of $u(t, \cdot)$ for $t = 0, t = 0.006, t = 0.012, t = 0.018$.



(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 λ^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 :

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

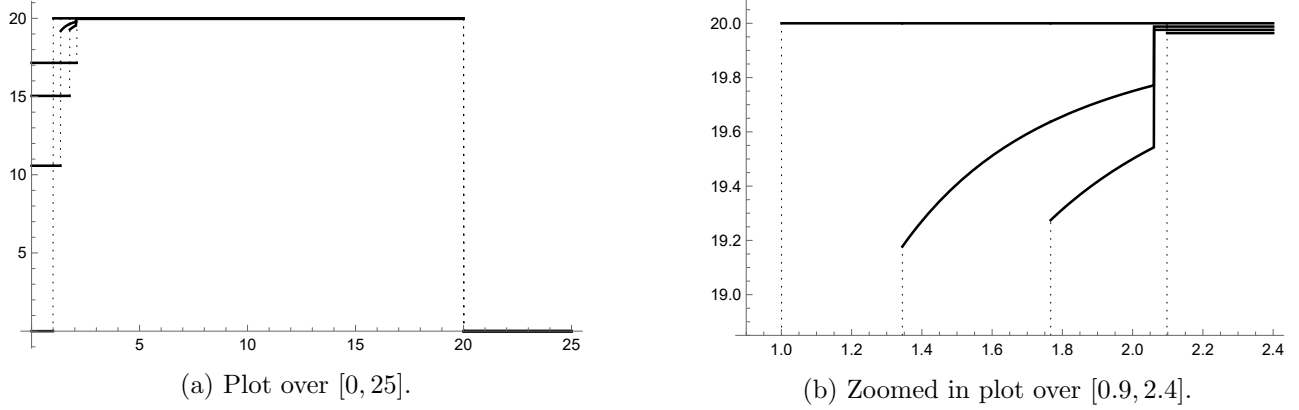


Figure 5: Plots of the solution $u(t, x)$ emanating from the characteristic function of annulus $A_{R_1}^{R_2}$ with $R_0 = 1$, $R_1 = 20$ in $n = 2$ as a function of $|x|$ for $t = 0$, $t = 2$, $t = 4$, $t = 6$.

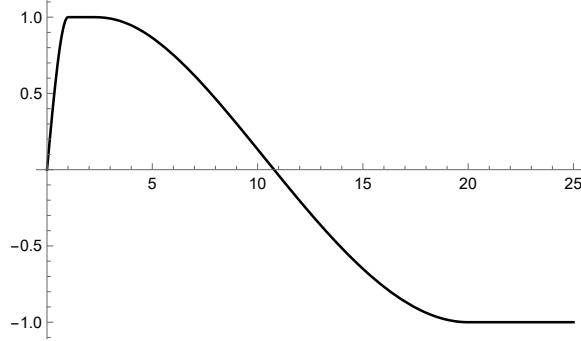


Figure 6: Plot of the Cahn-Hoffman vector field $Z(t, x)$ for the characteristic function of annulus $A_{R_1}^{R_2}$ with $R_0 = 1$, $R_1 = 20$ in $n = 2$ as a function of $|x|$.

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 λ^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, λ^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.

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