

Reconstruction of material losses by perimeter penalization and phase-field methods

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Abstract

We treat the inverse problem of determining material losses, such as cavities, in a conducting body, by performing electrostatic measurements at the boundary. We develop a numerical approach, based on variational methods, to reconstruct the unknown material loss by a single boundary measurement of current and voltage type.

The method is based on the use of phase-field functions to model the material losses and on a perimeter-like penalization to regularize the otherwise ill-posed problem. We justify the proposed approach by a convergence result, as the error on the measurement goes to zero.

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1 Introduction

In many inverse or optimal shape problems arising in the applications, the aim is to reconstruct the shape of an object, usually represented by an unknown open set, satisfying certain requirements. If we restrict ourselves to a variational formulation, for the sake of simplicity, we look for the shape minimizing a given functional F among all the admissible shapes. The shape is often modeled as a binary function, that is the open set is described through its characteristic function.

Two of the main issues for a satisfactory numerical resolution of this kind of problems are the following. First of all, and especially for inverse problems, the problem may be ill-posed, that is stability is missing or, in other words, F is not continuous. Second, numerically handling shapes or sets is not an easy task from the implementation point of view. The first issue is usually tackled by a regularization method, namely by adding to the functional a term penalizing the binary function with respect to some BV -related norm. For most applications, this should be enough for ensuring a regularization without being

restricting on the class of admissible unknowns. Often the BV -related norm is simply a perimeter-like penalization. About the handling of shapes or sets in computations, in many cases this is performed by associating to the open set a smooth function describing it. For example, one way of doing it is to replace the characteristic function of an open set D with a smooth function, referred to as a *phase-field function*, which is close to 0 outside D , close to 1 inside D , and has a quick transition from 0 to 1 across the boundary of D . Another way is the so-called level-set method, where D is identified with the sublevel set $\{\psi < 0\}$ of a smooth function ψ .

We are interested in using perimeter-like regularizations and phase-field functions for solving inverse or optimal shape problems, in particular those that are not well-posed. We aim to prove in a rigorous way that this kind of approach provides a good approximation of the original problem, allowing us at the same time to tame the ill-posedness and to have a formulation amenable to be easily implemented. A cornerstone of this method is the approximation, in the sense of Γ -convergence, of the perimeter functional by functionals defined on phase-field functions, due to Modica and Mortola, [16]. Since [15] such a result has found innumerable applications. In fact, whenever the functional F is continuous in a suitable way, the invariance of Γ -convergence by continuous perturbations permits to obtain an analogous Γ -convergence result if we add to F the perimeter penalization. Whenever the problem is ill-posed, that is F is not continuous, a corresponding convergence result is not straightforward any more. Since we believe that the method is valuable also in the ill-posed case, it would be important to justify it in a rigorous way, in general through a convergence result inspired by Γ -convergence techniques, for various interesting applications.

In this paper we perform such an analysis for the following inverse problem, arising from non-destructive evaluation. We aim to determine perfectly insulating defects in a homogeneous and isotropic conducting body by performing electrostatic measurements of voltage and current type at the boundary. The conducting body is contained in Ω , a bounded domain of \mathbb{R}^N , $N \geq 2$. The defects may have different geometrical properties, for instance we may have at the same time *cracks* (either interior or surface breaking), or *material losses* (either interior, that is cavities, or at the boundary). We denote with K the union of the boundaries of these defects, whereas $\tilde{\gamma}$ is a part of the boundary of Ω which is accessible, known and disjoint from K . If a current density $f \in L^2(\tilde{\gamma})$, with zero mean, is applied on $\tilde{\gamma}$, then the electrostatic potential $u = u(f, K)$ is the solution to the following Neumann boundary value problem

$$(1.1) \quad \begin{cases} \Delta u = 0 & \text{in } \Omega \setminus K, \\ \nabla u \cdot \nu = f & \text{on } \tilde{\gamma}, \\ \nabla u \cdot \nu = 0 & \text{on } \partial(\Omega \setminus K) \setminus \tilde{\gamma}. \end{cases}$$

We call G_K the connected component of $\Omega \setminus K$ which is reachable from $\tilde{\gamma}$ and we say that a defect is a *material loss* if G_K is equal to the interior of its closure, that is if no crack-type defect is present.

The value of u , that is the voltage, may be measured on another part of the boundary of Ω , say γ , which we assume to be accessible, known, disjoint from K and belonging to ∂G_K . We call g such a measurement, that is $g = u|_\gamma$. For simplicity, we may also assume that γ coincides with $\tilde{\gamma}$. If the defect K is unknown, we aim to recover its shape and location, that is G_K , by prescribing

one or more current densities f and measuring the corresponding voltage on γ , $g = u|_\gamma$, where u solves (1.1). In mathematical words, we are given one or more pairs of Cauchy data (g, f) on a known part of the boundary and we aim to reconstruct the domain of validity of the elliptic equation.

Here we are interested in the reconstruction only of material losses, that is cavities or material losses at the boundary, and for simplicity we refer to it as the *inverse cavity problem*. It is well-known that, in every dimension, a single boundary measurement is enough to reconstruct a material loss, thus providing uniqueness for the inverse problem, see for instance [21] for a proof with minimal regularity assumptions on the unknown material loss. Stability results have been proved in [1] for the three dimensional case and in [2] for the planar case, where also the instability character of the problem has been explicitly shown.

We notice that u , the electrostatic potential solution to (1.1), is constant on any connected component of $\Omega \setminus K$ different from G_K . The key observation is that its jump set in Ω is essentially contained in K . The uniqueness result recalled before actually allows us to say more, in fact the jump set of u uniquely determines G_K , that is the unknown material loss. Therefore we are interested in the reconstruction of the electrostatic potential u and especially of its (unknown) jump set. This suggests the possibility to set up a reconstruction procedure by solving a free-discontinuity problem related to the function u .

The main difficulty for the reconstruction is due to the ill-posedness of the problem. In fact, since they are measured, the Cauchy data that are available are not exact. Since the problem is severely ill-posed, such an error on the measurements may lead to a much greater error on the reconstructed defect. Furthermore, the inverse problem is nonlinear. In fact, even if the direct problem (1.1) is linear, the dependence of the electrostatic potential u , and of its values on γ , from the defect K is nonlinear. Finally, from a numerical point of view, the fact that the unknown is a set, namely G_K , introduces an additional complication for the implementation.

We propose a variational method to tackle at the same time these difficulties. The idea is to use a perimeter-like penalization to regularize the problem and to replace the unknown set G_K with its characteristic function and, in turn, with a phase-field function, to obtain a formulation that may be implemented numerically. Namely, the regularization we propose is related to the so-called Modica-Mortola functional, an approximation of the perimeter when phase-field functions are used. We might construct a family of functionals, depending on the noise level on the measurements ε , to be minimized with respect to the variable u (the reconstructed potential) and the phase-field variable v . However, to simplify the implementation we would rather have a functional depending on the phase-field variable v only. Thus, we take u depending on v , $u = u(v)$, as a solution to an almost degenerate elliptic problem whose coefficient is given by a slight modification of v , depending on ε . In other words, we replace the direct cavity problem with an elliptic problem in Ω where the coefficient of the equation is close to 1 in G_K , close to 0 outside G_K , with a quick transition across the boundary G_K . The method consists then of minimizing the so-obtained functionals, depending on ε , with respect to the phase-field variable v only. We remark that the reconstructed material loss may be simply computed by a suitable thresholding of the minimizing phase-field and that an approximation of the looked-for electrostatic potential is given by $u = u(v)$ where v is the minimizing phase-field.

The main result of the paper, Theorem 4.2, is that the corresponding minimizers v_ε converge, as $\varepsilon \rightarrow 0^+$, to the characteristic function of G_K , thus identifying the looked-for material loss, and that $u_\varepsilon = u(v_\varepsilon)$ converge to the looked-for potential u . Such a convergence result, whose proofs is obtained by techniques borrowed by Γ -convergence, provides a rigorous justification of the method. About the material loss to be reconstructed, this is assumed to satisfy a Lipschitz type regularity. We finally remark that the method makes use of a single measurement and that is enough to reconstruct the whole unknown material loss K .

If we instead allow the unknown defect not to be a material loss, that is it may include crack-type defects, for simplicity we refer to this problem as the *inverse crack problem*. About uniqueness, stability and reconstruction results on the inverse crack problem, we refer to [7] and the references therein. The main difference between the two cases is that for the determination of cracks one measurement is not enough, however, at least in the planar case, two suitably chosen measurements are sufficient. In [22, 23] a corresponding variational approach for the inverse crack problem has been developed. Again such an approach makes use of a penalization on the $(N - 1)$ -dimensional measure of the defects and of phase-field functions. Namely, it was constructed a family of functionals, again depending on the noise level on the measurements ε , to be minimized with respect to the variable u (the reconstructed potential) and the phase-field variable v . Instead of the perimeter functional and the Modica-Mortola functional, the Mumford-Shah functional [18] and its approximation, in the sense of Γ convergence, due to Ambrosio and Tortorelli [4, 5] were used, respectively. Also in this case a convergence result guaranteed a justification of the method. In Section 5 we recall the results obtained in [23] for the inverse crack problem and we compare with those obtained here for the inverse cavity problem. The main difficulty in the implementation of the method of [23] is that the functional to be minimized depends on two variables, the variable u , which should approximate the electrostatic potential, and the variable v , which is the phase-field variable that should approximate the jump set of the potential and hence the defect. It would be desirable to formulate the problem depending on one variable only, for instance only on the phase-field variable. Unfortunately such a formulation, which is proved here for the material loss case, may not be feasible. In fact, Section 5 is devoted to show that the result in [23] is essentially optimal, through several counterexamples. Moreover, more regularity is needed for the unknown defects of crack-type, namely a regularity assumptions of C^1 type, instead of Lipschitz, have to be imposed. Thus, we show that restricting ourselves to the reconstruction of material losses allows us to gain the following advantages. First, we may lower the a priori assumptions on the unknown defect to ones which are more suited for applications. More importantly, we obtain and justify a formulation which looks more natural and quite simpler to be implemented.

We finally wish to mention that a numerical implementation, based on the results of this paper and on those of [23], may be found in [20]. The corresponding numerical experiments show the validity of these methods also from a practical point of view.

The plan of the paper is the following. After a preliminaries section, Section 2, we describe the setting of the direct and inverse problem in Section 3. We treat the material loss case in in Section 4, where there is the main result of the paper, Theorem 4.2. In Section 5, we recall the results for the inverse

crack problem proved in [23] and we compare the crack and material loss cases and discuss their differences. In particular we show the optimality of the result of [23]. Finally, in Section 6 we deal with the differentiability of the functionals involved. Such differentiability is crucial for developing the algorithm used in [20].

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2 Preliminaries

Throughout the paper the integer $N \geq 2$ will denote the space dimension. We remark that we shall sometimes drop the dependence of any constant upon N , the space dimension. For every $x \in \mathbb{R}^N$, we shall set $x = (x', x_N)$, where $x' \in \mathbb{R}^{N-1}$ and $x_N \in \mathbb{R}$, and, for any $r > 0$, we shall denote by $B_r(x)$ the open ball in \mathbb{R}^N centred at x of radius r . Usually we shall write B_r instead of $B_r(0)$. For any subset $E \subset \mathbb{R}^N$ and any $r > 0$, we denote $B_r(E) = \bigcup_{x \in E} B_r(x)$.

For any non-negative integer k we denote by \mathcal{H}^k the k -dimensional Hausdorff measure. For Borel subsets of \mathbb{R}^N the N -dimensional Hausdorff measure coincides with \mathcal{L}^N , the N -dimensional Lebesgue measure. Furthermore, if $\gamma \subset \mathbb{R}^N$ is a smooth manifold of dimension k , then \mathcal{H}^k restricted to γ coincides with its k -dimensional surface measure. For any Borel $E \subset \mathbb{R}^N$ we let $|E| = \mathcal{L}^N(E)$.

We recall that a bounded open set $\Omega \subset \mathbb{R}^N$ is said to have a *Lipschitz boundary* if for every $x \in \partial\Omega$ there exist a Lipschitz function $\varphi : \mathbb{R}^{N-1} \rightarrow \mathbb{R}$ and a positive constant r such that for any $y \in B_r(x)$ we have, up to a rigid transformation,

$$y \in \Omega \quad \text{if and only if} \quad y_N < \varphi(y').$$

We observe that in this case the boundary of Ω has finite $(N-1)$ -dimensional Hausdorff measure, that is $\mathcal{H}^{N-1}(\partial\Omega) < +\infty$.

We say that a function $\varphi : A \rightarrow B$, A and B being metric spaces, is *bi-Lipschitz* if it is injective and φ and $\varphi^{-1} : \varphi(A) \rightarrow A$ are both Lipschitz functions. If both the Lipschitz constants of φ and φ^{-1} are bounded by $L \geq 1$, then we say that φ is *bi-Lipschitz* with constant L .

We recall some basic notation and properties of functions of bounded variation and sets of finite perimeter. For a more comprehensive treatment of these subjects see, for instance, [3, 11, 13].

Given a bounded open set $\Omega \subset \mathbb{R}^N$, we denote by $BV(\Omega)$ the Banach space of *functions of bounded variation*. We recall that $u \in BV(\Omega)$ if and only if $u \in L^1(\Omega)$ and its distributional derivative Du is a bounded vector measure. We endow $BV(\Omega)$ with the standard norm as follows. Given $u \in BV(\Omega)$, we denote by $|Du|$ the total variation of its distributional derivative and we set $\|u\|_{BV(\Omega)} = \|u\|_{L^1(\Omega)} + |Du|(\Omega)$. We recall that whenever $u \in W^{1,1}(\Omega)$, then $u \in BV(\Omega)$ and $|Du|(\Omega) = \int_\Omega |\nabla u|$, therefore $\|u\|_{BV(\Omega)} = \|u\|_{L^1(\Omega)} + \|\nabla u\|_{L^1(\Omega)} = \|u\|_{W^{1,1}(\Omega)}$.

We say that a sequence of $BV(\Omega)$ functions $\{u_h\}_{h=1}^\infty$ *weakly* converges* in $BV(\Omega)$ to $u \in BV(\Omega)$ if and only if u_h converges to u in $L^1(\Omega)$ and Du_h weakly*

converges to Du in Ω , that is

$$(2.1) \quad \lim_h \int_{\Omega} v dDu_h = \int_{\Omega} v dDu \quad \text{for any } v \in C_0(\Omega).$$

By Proposition 3.13 in [3], we have that if a sequence of $BV(\Omega)$ functions $\{u_h\}_{h=1}^{\infty}$ is bounded in $BV(\Omega)$ and converges to u in $L^1(\Omega)$, then $u \in BV(\Omega)$ and u_h converges to u weakly* in $BV(\Omega)$.

Let Ω be a bounded open set with Lipschitz boundary. A sequence of $BV(\Omega)$ functions $\{u_h\}_{h=1}^{\infty}$ such that $\sup_h \|u_h\|_{BV(\Omega)} < +\infty$ admits a subsequence converging weakly* in $BV(\Omega)$ to a function $u \in BV(\Omega)$, see for instance Theorem 3.23 in [3]. As a corollary, we infer that for any $C > 0$ the set $\{u \in BV(\Omega) : \|u\|_{BV(\Omega)} \leq C\}$ is a compact subset of $L^1(\Omega)$.

Let E be a bounded Borel set contained in $B_R \subset \mathbb{R}^N$. We shall denote by χ_E its characteristic function. We notice that E is compactly contained in B_{R+1} , which we shall denote by $E \Subset B_{R+1}$. We say that E is a *set of finite perimeter* if χ_E belongs to $BV(B_{R+1})$ and we call the number $P(E) = |D\chi_E|(B_{R+1})$ its *perimeter*.

Let us further remark that the intersection of two sets of finite perimeter is still a set of finite perimeter. Moreover, whenever E is open and $\mathcal{H}^{N-1}(\partial E)$ is finite, then E is a set of finite perimeter, see for instance [11, Section 5.11, Theorem 1]. Therefore a bounded open set Ω with Lipschitz boundary is a set of finite perimeter and its perimeter $P(\Omega)$ coincides with $\mathcal{H}^{N-1}(\partial\Omega)$.

For any bounded open set Ω , we define the following *perimeter functional* $P : L^1(\Omega) \rightarrow [0, +\infty]$ such that

$$(2.2) \quad P(u) = \begin{cases} c|Du|(\Omega) & \text{if } u \in BV(\Omega) \text{ and } u \in \{0, 1\} \text{ a.e.,} \\ +\infty & \text{otherwise,} \end{cases}$$

where c is a positive constant to be chosen later. We observe that $P(u) = cP(E)$ if $u = \chi_E$ and E is a set of finite perimeter compactly contained in Ω .

We denote by $SBV(\Omega)$ the space of *special functions of bounded variation*. For any $u \in SBV(\Omega)$, the density of the absolutely continuous part of Du with respect to \mathcal{L}^N will be denoted by ∇u , the *approximate gradient* of u . The singular part, with respect to \mathcal{L}^N , of Du is concentrated on $J(u)$, $J(u)$ being the *approximate discontinuity set* (or *jump set*) of u in Ω . We further say that a function $u \in GSBV(\Omega)$, the space of *generalized functions of bounded variation*, if $u \in L^1(\Omega)$ and for any $T > 0$ its *truncation* $u_T = (-T) \vee (T \wedge u) \in SBV(\Omega)$. Let us recall that the approximate gradient ∇u of $u \in GSBV(\Omega)$ is defined almost everywhere and coincides with ∇u_T almost everywhere on $\{u = u_T\}$, and that $J(u) = \bigcup_{T>0} J(u_T)$.

The special functions of bounded variation have important compactness and semicontinuity properties, see for instance [3, Theorem 4.7 and Theorem 4.8].

We remark that if $u \in BV(\Omega)$ and $u \in \{0, 1\}$ almost everywhere in Ω , then $u \in SBV(\Omega)$ and $P(u) = c|Du|(\Omega) = c\mathcal{H}^{N-1}(J(u))$.

Let us define the so-called *Mumford-Shah functional*, introduced in [18] in the context of image segmentation. Let us fix positive constants b and c . Let $\mathcal{MS} : L^1(\Omega) \rightarrow [0, +\infty]$ be given by

$$(2.3) \quad \mathcal{MS}(u) = b \int_{\Omega} |\nabla u|^2 + c\mathcal{H}^{N-1}(J(u)) \quad \text{if } u \in GSBV(\Omega),$$

whereas $\mathcal{MS}(u) = +\infty$ otherwise.

Let us introduce the following Γ -convergence results concerning the approximation of the perimeter functional and the Mumford-Shah functional by phase-field functionals. For the definition and properties of Γ -convergence we refer to [8]. The perimeter approximation is due to Modica and Mortola, [16], whereas the Mumford-Shah functional approximation is due to Ambrosio and Tortorelli, [4, 5]. We shall follow the notation and proofs contained in [6].

Throughout the paper, for any p , $1 \leq p \leq +\infty$, we shall denote its conjugate exponent by p' , that is $p^{-1} + (p')^{-1} = 1$. Let $W : \mathbb{R} \rightarrow [0, +\infty)$ be a continuous function such that $W(t) = 0$ if and only if $t \in \{0, 1\}$. Let $c_W = \int_0^1 \sqrt{W(t)} dt$. In the definition of the perimeter functional we pick $c = 2c_W$. For instance, we may choose $W(t) = 9t^2(t-1)^2$ for any $t \in \mathbb{R}$, whence $c = 2c_W = 1$.

The following approximation result is due to Modica and Mortola, [16], see also [6, Theorem 4.13].

Theorem 2.1. *Let $\Omega \subset \mathbb{R}^N$ be a bounded open set with Lipschitz boundary. For any $\eta > 0$ we define the functional $P_\eta : L^1(\Omega) \rightarrow [0, +\infty]$ as follows*

$$(2.4) \quad P_\eta(v) = \begin{cases} \frac{1}{\eta} \int_\Omega W(v) + \eta \int_\Omega |\nabla v|^2 & \text{if } v \in W^{1,2}(\Omega, [0, 1]), \\ +\infty & \text{otherwise.} \end{cases}$$

Then we have that, with respect to the metric of $L^1(\Omega)$, P_η Γ -converges to P as $\eta \rightarrow 0^+$.

Here $W^{1,2}(\Omega, [0, 1]) = \{v \in W^{1,2}(\Omega) : 0 \leq v \leq 1 \text{ a.e. in } \Omega\}$. We note that the result does not change if in the definition of P_η we omit the constraint

$$0 \leq v \leq 1 \text{ a.e. in } \Omega.$$

Also the following result, due to Modica, [15], will be useful.

Proposition 2.2. *For any $C > 0$ and any $\eta > 0$, let us define*

$$A_C = \{v \in L^1(\Omega) : 0 \leq v \leq 1 \text{ a.e. and } P_\eta(v) \leq C\}.$$

Then A_C is precompact in $L^1(\Omega)$.

Remark 2.3. With the same proof, we can show the following. Let us consider any family $\{v_\eta\}_{0 < \eta \leq \eta_0}$ such that, for some positive constant C and for any η , $0 < \eta \leq \eta_0$, we have $0 \leq v_\eta \leq 1$ almost everywhere and $P_\eta(v_\eta) \leq C$. Then $\{v_\eta\}_{0 < \eta \leq \eta_0}$ is precompact in $L^1(\Omega)$.

Let us fix q , $1 < q < +\infty$. Let $V : \mathbb{R} \rightarrow [0, +\infty)$ be a continuous function such that $V(t) = 0$ if and only if $t = 1$ and let $c_V = \int_0^1 \sqrt{V(t)} dt$. Let $\psi : \mathbb{R} \rightarrow \mathbb{R}$ be a continuous non-decreasing function such that $\psi(0) = 0$, $\psi(1) = 1$ and $\psi(t) > 0$ if $t > 0$. For any $\eta > 0$, let us fix $o_\eta = o_\eta(q) \geq 0$ such that $\lim_{\eta \rightarrow 0^+} o_\eta/\eta^{q-1} = 0$. Finally, we define $\psi_\eta = (1 - o_\eta)\psi + o_\eta$. Provided $o_\eta < 1$, we have that ψ_η is a continuous, non-decreasing function such that $\psi_\eta(0) = o_\eta$ and $\psi_\eta(1) = 1$.

For instance, we may choose $V(t) = (t-1)^2/4$ for any $t \in \mathbb{R}$, whence $4c_V = 1$. About ψ , we may take $\psi(t) = t^\gamma$, $\gamma > 0$, if $t \geq 0$, while $\psi(t) = 0$ if

$t < 0$ Alternatively, we may choose $\psi(t) = 0$ if $t < 0$, $\psi(t) = -2t^3 + 3t^2$ for any $t \in [0, 1]$, and $\psi(t) = 1$ for any $t > 1$. We may finally take $o_\eta(q) = \eta^q$.

Then, for any $\eta > 0$, we define the following functional $\mathcal{AT}_\eta^q : L^1(\Omega) \times L^1(\Omega) \rightarrow [0, +\infty]$ by

$$(2.5) \quad \mathcal{AT}_\eta^q(u, v) = b \int_{\Omega} \psi_\eta(v) |\nabla u|^q + \frac{1}{\eta} \int_{\Omega} V(v) + \eta \int_{\Omega} |\nabla v|^2$$

if $u \in W^{1,q}(\Omega)$ and $v \in W^{1,2}(\Omega, [0, 1])$,

whereas $\mathcal{AT}_\eta^q(u, v) = +\infty$ otherwise. We shall refer to \mathcal{AT}_η^q as the Ambrosio-Tortorelli functional.

Let us define the following variant of the Mumford-Shah functional. The main difference is that we allow the exponent q to be different for 2, requiring only that $1 < q < +\infty$. For reasons which will appear evident soon, we also add a formal variable v and we pick $c = 4c_V$. We define the functional $\mathcal{MS}^q : L^1(\Omega) \times L^1(\Omega) \rightarrow [0, +\infty]$ by

$$(2.6) \quad \mathcal{MS}^q(u, v) = b \int_{\Omega} |\nabla u|^q + 4c_V \mathcal{H}^{N-1}(J(u))$$

if $u \in GSBV(\Omega)$ and $v = 1$ a.e. in Ω ,

whereas $\mathcal{MS}^q(u, v) = +\infty$ otherwise.

The Ambrosio-Tortorelli functional approximates the Mumford-Shah functional, in the sense of Γ -convergence. Such an important approximation result is due to Ambrosio and Tortorelli, [4, 5], see also [6].

Theorem 2.4. *With respect to the metric of $L^1(\Omega) \times L^1(\Omega)$, we have that, as $\eta \rightarrow 0^+$, \mathcal{AT}_η^q Γ -converges to \mathcal{MS}^q .*

Let us review some regularity results which will be needed in the sequel. Most of these results are a consequence of a theorem by Meyers, [17], see also [12], and of standard regularity estimates, and we shall omit the proofs. Let $\Omega \subset \mathbb{R}^N$ be a bounded open set with Lipschitz boundary. Let $A = A(x)$, $x \in \Omega$, be an $N \times N$ matrix whose entries are measurable and such that, for some $0 < \lambda < 1$, we have

$$(2.7) \quad \begin{aligned} A(x)\xi \cdot \xi &\geq \lambda|\xi|^2 && \text{for any } \xi \in \mathbb{R}^N \text{ and for a.e. } x \in \Omega, \\ \|A\|_{L^\infty(\Omega)} &\leq \lambda^{-1}. \end{aligned}$$

We remark that for any matrix A , by $\|A\|$ we denote the norm of the matrix as a linear operator.

Let $f \in L^s(\partial\Omega)$, with $s > 1$ if $N = 2$ or $s \geq 2(N-1)/N$ if $N \geq 3$, be such that $\int_{\partial\Omega} f = 0$ and let $F \in L^p(\Omega, \mathbb{R}^N)$, with $p \geq 2$. Let us denote $W_*^{1,2}(\Omega) = \{u \in W^{1,2}(\Omega) : \int_{\Omega} u = 0\}$. Then, there exists a unique $u \in W_*^{1,2}(\Omega)$ such that

$$(2.8) \quad \int_{\Omega} A \nabla u \cdot \nabla v = \int_{\Omega} F \cdot \nabla v + \int_{\partial\Omega} f v \quad \text{for any } v \in W^{1,2}(\Omega).$$

This is the weak formulation of

$$\begin{cases} \operatorname{div}(A \nabla u) = \operatorname{div}(F) & \text{in } \Omega \\ A \nabla u \cdot \nu = f & \text{on } \partial\Omega. \end{cases}$$

The following regularity result holds true.

Proposition 2.5. *Under the previous assumptions, the following regularity properties hold.*

First of all, we have, for a constant C_0 depending on N, λ, p, s and Ω only,

$$(2.9) \quad \|u\|_{W^{1,2}(\Omega)} \leq C_0 (\|F\|_{L^p(\Omega)} + \|f\|_{L^s(\partial\Omega)}).$$

If $p > N$ and $s > N - 1$, there exist a constant $C_1 > 0$ such that

$$(2.10) \quad \|u\|_{L^\infty(\Omega)} \leq C_1 (\|F\|_{L^p(\Omega)} + \|f\|_{L^s(\partial\Omega)}).$$

Here C_1 depends on N, λ, p, s and Ω only.

There exists a constant $Q > 2$, depending on N, λ and Ω only ($Q \rightarrow 2$ if $\lambda \rightarrow 0^+$), such that if p satisfies $2 < p < Q$ and $s \geq (N - 1)p/N$, then

$$(2.11) \quad \|\nabla u\|_{L^p(\Omega)} \leq C_2 (\|F\|_{L^p(\Omega)} + \|f\|_{L^s(\partial\Omega)}).$$

Here C_2 depends on N, λ, p, s and Ω only.

We conclude that if $s > N - 1$, there exists a constant $q(\lambda) > 2$, depending on N, λ, s and Ω only, such that for any $p, 2 \leq p \leq q(\lambda)$, we have

$$(2.12) \quad \|\nabla u\|_{L^p(\Omega)} \leq C_3 (\|F\|_{L^p(\Omega)} + \|f\|_{L^s(\partial\Omega)}),$$

in particular, if $p = q(\lambda)$

$$(2.13) \quad \|\nabla u\|_{L^{q(\lambda)}(\Omega)} \leq C_4 (\|F\|_{L^{q(\lambda)}(\Omega)} + \|f\|_{L^s(\partial\Omega)}).$$

Here C_3 depends on N, λ, p, s and Ω only, whereas C_4 depends on N, λ, s and Ω only.

Remark 2.6. Let us observe that Q and $q(\lambda)$ converges to 2 as $\lambda \rightarrow 0^+$, whereas all the constants C_0 – C_4 might tend to $+\infty$ as $\lambda \rightarrow 0^+$. Let us also remark that the same kinds of estimates hold true if we replace $W_*^{1,2}(\Omega)$ with, for instance,

$$W_E^{1,2}(\Omega) = \left\{ u \in W^{1,2}(\Omega) : \int_E u = 0 \right\}$$

where E is a Borel subset of $\partial\Omega$ with non-empty interior, clearly with respect to the induced topology of $\partial\Omega$. In this case, the constants C_0 – C_4 might depend on E as well.

We conclude this section with the following lemma, in which we state a Caccioppoli inequality.

Lemma 2.7. *Let us assume that $A = A(x)$, $x \in B_{2R}$, is a symmetric $N \times N$ matrix whose entries belong to $L^\infty(B_{2R})$. We also assume that, for some constants $0 < \lambda < \Lambda$, we have*

$$\lambda|\xi|^2 \leq A(x)\xi \cdot \xi \leq \Lambda|\xi|^2 \quad \text{for every } \xi \in \mathbb{R}^N \text{ and for a.e. } x \in B_{2R}.$$

Let $w \in L^\infty(B_{2R})$ be a weight satisfying $0 < \varepsilon \leq w \leq 1$ almost everywhere in B_{2R} . If u solves in a weak sense

$$\operatorname{div}(wA\nabla u) = 0 \quad \text{in } B_{2R},$$

then

$$(2.14) \quad \int_{B_R} w|\nabla u|^2 \leq \frac{C}{R^2} \int_{B_{2R}} wu^2$$

where C depends on λ and Λ only.

PROOF. In order to prove (2.14) it is enough to take a cutoff function χ such that $\chi \in C_0^\infty(B_{2R})$, $0 \leq \chi \leq 1$ on B_{2R} , and $\chi \equiv 1$ on B_R . We may also assume that for an absolute constant C we have $|\nabla\chi| \leq C/R$ on B_{2R} . Then we use the test function $u\chi^2$ in the equation and with simple computations we obtain (2.14). \square

3 The direct problem and setting of the inverse problem

Let Ω , Ω_1 and $\tilde{\Omega}_1$ be three bounded domains contained in \mathbb{R}^N , $N \geq 2$, with Lipschitz boundaries such that $\Omega_1 \subset \tilde{\Omega}_1 \subset \Omega$ and the following properties are satisfied. First, $\Omega \setminus \tilde{\Omega}_1$ is not empty and $\text{dist}(\overline{\Omega_1}, \partial\tilde{\Omega}_1 \cap \Omega) > 0$. Then, there exist γ and $\tilde{\gamma}$, closed subsets of $\partial\Omega$, which are contained in the interior of $\partial\Omega \cap \partial\Omega_1$ and whose interiors, with respect to the induced topology of $\partial\Omega$, are not empty.

We assume that Ω , Ω_1 , $\tilde{\Omega}_1$, γ and $\tilde{\gamma}$ are fixed throughout the paper.

Let K_0 be an *admissible defect*, that is K_0 is a non-empty compact set contained in $\overline{\Omega}$ such that $\text{dist}(K_0, \overline{\Omega_1}) > 0$. We denote with G_{K_0} the connected component of $\Omega \setminus K_0$ such that $\tilde{\Omega}_1 \subset G_{K_0}$. We observe that $\gamma \cup \tilde{\gamma} \subset \partial G_{K_0}$. We remark that if $K_0 \subset \partial\Omega$ then $G_{K_0} = \Omega$, that is no defect is present in the conductor.

We say that an admissible defect K_0 is a *material loss defect*, or *material loss* for short, if G_{K_0} is equal to the interior of its own closure (that is no crack-type defect is allowed).

Remark 3.1. If the defect K_0 to be reconstructed is compactly contained in Ω and $\partial\Omega$ is connected, then we may take γ and $\tilde{\gamma}$ equal to $\partial\Omega$. Furthermore, if $\partial\Omega$ is regular enough and we a priori know that $\text{dist}(K_0, \partial\Omega) > \delta$, for some $\delta > 0$ small enough, then we may choose $\Omega_1 = \{x \in \Omega : \text{dist}(x, \partial\Omega) < \delta/2\}$ and $\tilde{\Omega}_1 = \{x \in \Omega : \text{dist}(x, \partial\Omega) < 3\delta/4\}$.

Let us fix a number s , $s > N - 1$, which shall be kept fixed throughout the paper. Let us prescribe $f_0 \in L^s(\partial\Omega)$ such that $\int_{\partial\Omega} f_0 = 0$, $f_0 \not\equiv 0$ and $\text{supp}(f_0) \subset \tilde{\gamma}$.

Let the electrostatic potential $u_0 = u(f_0, K_0)$ be the weak solution to the following Neumann boundary value problem

$$(3.1) \quad \begin{cases} \Delta u_0 = 0 & \text{in } \Omega \setminus K_0, \\ \nabla u_0 \cdot \nu = f_0 & \text{on } \tilde{\gamma}, \\ \nabla u_0 \cdot \nu = 0 & \text{on } \partial(\Omega \setminus K_0) \setminus \tilde{\gamma}, \end{cases}$$

with the normalization conditions

$$(3.2) \quad \int_{\gamma} u_0 = 0,$$

and

$$(3.3) \quad u_0 = 0 \quad \text{almost everywhere in } \Omega \setminus G_{K_0}.$$

Let us recall that our measured additional information is $g_0 \in L^2(\gamma)$ where $g_0 = u_0|_{\gamma}$. By (3.2), we have $\int_{\gamma} g_0 = 0$.

We observe that there exists a unique solution u_0 to (3.1)-(3.2)-(3.3) and that it satisfies the following regularity properties, see [22] for further details.

There exists a constant $C_1 > 0$, depending on s , Ω , Ω_1 , $\tilde{\Omega}_1$, γ and $\tilde{\gamma}$ only, such that

$$(3.4) \quad \|\nabla u_0\|_{L^2(\Omega \setminus K_0)} \leq C_1 \|f_0\|_{L^s(\tilde{\gamma})},$$

$$(3.5) \quad \|u_0\|_{L^\infty(\Omega)} \leq C_1 \|f_0\|_{L^s(\tilde{\gamma})}.$$

The estimate (3.5) guarantees that u_0 belongs to $W^{1,2}(\Omega \setminus K_0)$. Furthermore, under the additional assumption that $\mathcal{H}^{N-1}(K_0) < +\infty$, or equivalently that $\mathcal{H}^{N-1}(\partial G_{K_0}) < +\infty$, we have that u_0 belongs to $SBV(\Omega)$, its approximate discontinuity set $J(u_0)$ satisfies $\mathcal{H}^{N-1}(J(u_0) \setminus \partial G_{K_0}) = 0$ and, finally, ∇u_0 , the weak derivative of u_0 in $\Omega \setminus K_0$, coincides almost everywhere in Ω with the approximate gradient of u_0 , see for instance [3, Proposition 4.4].

For any r , $1 < r < +\infty$, and any Borel set $E \subset \partial\Omega$ whose interior, in the induced topology, is not empty, we define

$$W_E^{1,r}(\Omega) = \left\{ u \in W^{1,r}(\Omega) : \int_E u = 0 \right\}.$$

We observe that, by a generalized Poincaré inequality, on $W_E^{1,r}(\Omega)$ the usual $W^{1,r}(\Omega)$ norm and the norm $\|u\|_{W_E^{1,r}(\Omega)} = \|\nabla u\|_{L^r(\Omega)}$ are equivalent. Therefore, we shall set this second one as the natural norm of $W_E^{1,r}(\Omega)$.

Let us consider a weight w in Ω satisfying the following properties. We assume that $w \in L^\infty(\Omega)$ and that $w \geq \varepsilon$ almost everywhere in Ω , for some $\varepsilon > 0$.

For any such weight w , and any $u_1, u_2 \in W^{1,2}(\Omega)$, we define the bilinear form

$$\langle u_1, u_2 \rangle_w = \int_\Omega w \nabla u_1 \cdot \nabla u_2$$

and we denote the seminorm

$$|u_1|_w = \langle u_1, u_1 \rangle_w^{1/2} = \left(\int_\Omega w |\nabla u_1|^2 \right)^{1/2}.$$

We denote, for any $u_1 \in W^{1,2}(\Omega)$,

$$\|u_1\|_w = \left(\|u_1\|_{L^2(\Omega)}^2 + |u_1|_w^2 \right)^{1/2}.$$

We have that $\|\cdot\|_w$ is an equivalent norm for $W^{1,2}(\Omega)$, and $\langle \cdot, \cdot \rangle_w$ is a scalar product on $W_E^{1,2}(\Omega)$ whose corresponding norm, $|\cdot|_w$, is an equivalent norm for $W_E^{1,2}(\Omega)$, for any set E as before.

For any such weight w and any $f \in L^s(\partial\Omega)$ such that $\int_{\partial\Omega} f = 0$ and $\text{supp}(f) \subset \tilde{\gamma}$, let $u = u(w)$ be the solution to the following Neumann type boundary value problem

$$(3.6) \quad \begin{cases} \text{div}(w \nabla u) = 0 & \text{in } \Omega \\ w \nabla u \cdot \nu = f & \text{on } \partial\Omega \\ \int_\gamma u = 0. \end{cases}$$

The weak formulation of (3.6) is the following. We look for a function $u \in W_{\tilde{\gamma}}^{1,2}(\Omega)$ such that

$$\int_{\Omega} w \nabla u \cdot \nabla u_1 = \int_{\tilde{\gamma}} f u_1 \quad \text{for any } u_1 \in W^{1,2}(\Omega).$$

Obviously we have existence and uniqueness of such a solution. Furthermore, the following regularity result holds for u .

Proposition 3.2. *Under the previous notation and assumptions, let u solve (3.6) for some weight w . We assume that*

$$\|w\|_{L^\infty(\Omega)} \leq A$$

and

$$w(x) = 1 \quad \text{for a.e. } x \in \tilde{\Omega}_1.$$

Then there exists a constant C_2 , depending on s , Ω , Ω_1 , $\tilde{\Omega}_1$, γ , $\tilde{\gamma}$ and A only, such that

$$(3.7) \quad |u|_w \leq C_2 \|f\|_{L^s(\tilde{\gamma})}$$

$$(3.8) \quad \|u\|_{L^\infty(\Omega)} \leq C_2 \|f\|_{L^s(\tilde{\gamma})}.$$

We notice that the constant C_2 does not depend on w or on ε .

PROOF. We sketch the proof of this proposition. Inequality (3.7) follows from an application of Poincaré inequality in Ω_1 . The L^∞ bound (3.8) is a consequence of the maximum principle and may be proved following the same arguments used to prove (3.5), see [22] for details. \square

Let us fix the notation for our inverse problem. Let K_0 be the unknown defect, which for the time being we assume to be just an admissible defect.

We assume that f_0 belongs to $L^s(\partial\Omega)$ and satisfies $\text{supp}(f_0) \subset \tilde{\gamma}$ and $\int_{\partial\Omega} f_0 = 0$. We recall that s is a fixed constant such that $s > N - 1$.

The unknown electrostatic potential is $u_0 = u(K_0, f_0)$, solution to (3.1)-(3.2)-(3.3), and the additional measured data is $g_0 = u_0|_\gamma$. We observe that $g_0 \in L^2(\gamma)$ and $\int_\gamma g_0 = 0$.

Let us fix a noise level ε , $0 < \varepsilon \leq 1$, then the noisy Cauchy data are given by f_ε and g_ε . Here f_ε belongs to $L^s(\partial\Omega)$ and satisfies $\text{supp}(f_\varepsilon) \subset \tilde{\gamma}$ and $\int_{\partial\Omega} f_\varepsilon = 0$, whereas g_ε belongs to $L^2(\gamma)$ and satisfies $\int_\gamma g_\varepsilon = 0$. We assume that

$$(3.9) \quad \|f_0 - f_\varepsilon\|_{L^s(\tilde{\gamma})} \leq \varepsilon \quad \text{and} \quad \|g_0 - g_\varepsilon\|_{L^2(\gamma)} \leq \varepsilon.$$

For any $0 < \varepsilon \leq 1$, let $\eta = \eta(\varepsilon) > 0$ and $a_\varepsilon > 0$ be such that $\lim_{\varepsilon \rightarrow 0^+} \eta(\varepsilon) = 0$ and $\lim_{\varepsilon \rightarrow 0^+} a_\varepsilon = 0$. Further assumptions on $\eta(\varepsilon)$ and a_ε will be imposed later.

Let us fix a constant c_1 , $0 < c_1 < 1$. We recall that $\psi : \mathbb{R} \rightarrow \mathbb{R}$ is a continuous, non-decreasing function such that $\psi(0) = 0$, $\psi(1) = 1$, and $\psi(t) > 0$ if $t > 0$. In particular $\psi(c_1) > 0$. Provided $o_\eta \leq 1/2$, we have that ψ_η is a continuous, non-decreasing function such that $\psi_\eta(0) = o_\eta$ and $\psi_\eta(1) = 1$. Furthermore, $\psi_\eta(c_1) \geq \psi(c_1)/2 > 0$.

In the sequel we shall always assume that

$$0 < o_\eta \leq 1/2 \quad \text{for any } \eta > 0.$$

Without loss of generality, we also assume that ψ , W , and V are bounded all over \mathbb{R} , for instance by a constant A . For any $\eta > 0$, again without loss of generality, we assume that ψ is such that $\psi_\eta \geq o_\eta/2$ all over \mathbb{R} .

To any function $\tilde{v} \in L^1(\Omega)$ we associate the function $v = 1 - \tilde{v}$. We observe that, provided $0 \leq \tilde{v} \leq 1$ almost everywhere in Ω , we also have $0 \leq v \leq 1$ almost everywhere in Ω .

For any $\eta > 0$ and for any $\tilde{v} \in L^1(\Omega)$, let $w_\eta = w_\eta(\tilde{v}) = \psi_\eta(v)$, where $v = 1 - \tilde{v}$. We observe that w_η is such that $\|w_\eta\|_{L^\infty(\Omega)} \leq A + 1/2$ and $w_\eta \geq o_\eta/2$ almost everywhere in Ω . Therefore we define, for any $\tilde{v} \in L^1(\Omega)$ and any ε , $0 < \varepsilon \leq 1$, the function $\tilde{u}_\varepsilon \in W_\gamma^{1,2}(\Omega)$ where $\tilde{u}_\varepsilon = \tilde{u}_\varepsilon(\tilde{v})$ is the solution to the following boundary value problem

$$(3.10) \quad \begin{cases} \operatorname{div}(w_\eta(\tilde{v})\nabla\tilde{u}_\varepsilon) = 0 & \text{in } \Omega \\ w_\eta(\tilde{v})\nabla\tilde{u}_\varepsilon \cdot \nu = f_\varepsilon & \text{on } \partial\Omega, \end{cases}$$

where as usual $\eta = \eta(\varepsilon)$.

We finally fix positive constants a_1 , a_2 , \tilde{q} , $\tilde{\beta}$, and c_2 , $0 < c_1 < c_2 \leq 1$. We also define the following space $W(\Omega) = \{\tilde{v} \in W^{1,2}(\Omega) : \tilde{v} = 0 \text{ a.e. in } \tilde{\Omega}_1\}$. To any $\tilde{v} \in W(\Omega)$ we associate the function $v = 1 - \tilde{v}$. We remark that $v \in W^{1,2}(\Omega)$ and $v = 1$ almost everywhere in $\tilde{\Omega}_1$. All these constants and the notation will be kept fixed throughout the paper.

4 Determination of material losses

In this section, the main of the paper, we shall consider the problem of determining material losses. We begin by defining suitable classes of material losses.

Definition 4.1. Let us fix a positive constant δ . We say that \mathcal{B} is an admissible class of material losses if the following holds. First, any $K \in \mathcal{B}$ is an admissible defect such that $\operatorname{dist}(K, \tilde{\Omega}_1) \geq \delta$ and G_K is a domain with Lipschitz boundary. Second, we assume that, for some constant C , we have $\mathcal{H}^{N-1}(\partial G_K) \leq C$ for any $K \in \mathcal{B}$. Finally, we assume that the set $\{\overline{G_K} : K \in \mathcal{B}\}$ is compact with respect to the Hausdorff distance.

In the remaining part of this section, let us fix \mathcal{B} , an admissible class of material losses in the sense of Definition 4.1. We assume that the unknown defect K_0 belongs to \mathcal{B} . We observe that, as in Proposition 2.5, we have there exist a constant $q > 2$ and a constant $C > 0$, not depending on f_0 , such that $\nabla u_0 \in L^q(\Omega, \mathbb{R}^N)$, in particular

$$\|\nabla u_0\|_{L^q(\Omega)} \leq C \|f_0\|_{L^s(\tilde{\gamma})}.$$

Here the constants q and C depend also on s and on K_0 . In the sequel of the section, we shall fix $q > 2$ as such a constant, which depends on K_0 , among other things. We define $q_1 = (q-2)/(2q)$ and we observe that $0 < q_1 < 1/2$. We also define the following set. For any positive constant a , we say that $v \in H(a)$ if $v \in W^{1,2}(\Omega, [0, 1])$, $v = 1$ almost everywhere in $\tilde{\Omega}_1$ and there exists $K \in \mathcal{B}$ such that $v \geq c_2$ almost everywhere in $\Omega \setminus \overline{B_a}(\Omega \setminus \overline{G_K})$ and $v \leq c_1$ almost everywhere in $\Omega \setminus \overline{B_a}(\overline{G_K})$. We observe that, by the compactness of the class \mathcal{B} with respect to the Hausdorff distance, such a set $H(a)$ is closed with respect to the weak $W^{1,2}(\Omega)$ convergence.

For any ε , $0 < \varepsilon \leq 1$, we define $\tilde{\mathcal{G}}_\varepsilon : W(\Omega) \rightarrow \mathbb{R}$ as follows. For any $\tilde{v} \in W(\Omega)$, recalling that $v = 1 - \tilde{v}$, we set

$$(4.1) \quad \tilde{\mathcal{G}}_\varepsilon(\tilde{v}) = \frac{a_2}{\varepsilon^{\tilde{\beta}}} \int_\gamma |\tilde{u}_\varepsilon - g_\varepsilon|^2 + b \int_\Omega w_\eta(\tilde{v}) |\nabla \tilde{u}_\varepsilon|^2 + \frac{1}{\eta} \int_\Omega W(v) + \eta \int_\Omega |\nabla v|^2.$$

Here $\eta = \eta(\varepsilon)$, $o_\eta = o_\eta(2)$, $w_\eta = w_{\eta(\varepsilon)}(\tilde{v}) = \psi_{\eta(\varepsilon)}(v)$ and $\tilde{u}_\varepsilon = \tilde{u}_\varepsilon(\tilde{v})$ is the solution to (3.10). Here and in the sequel of this section, we may also set the constant $b = 0$, that is we may drop the second term of the functional.

Then, for any ε , $0 < \varepsilon \leq 1$, we define $\mathcal{G}_\varepsilon : L^1(\Omega) \rightarrow [0, +\infty]$ as follows. For any $\tilde{v} \in L^1(\Omega)$ we set

$$(4.2) \quad \mathcal{G}_\varepsilon(\tilde{v}) = \tilde{\mathcal{G}}_\varepsilon(\tilde{v}) \quad \text{if } \tilde{v} \in W(\Omega) \text{ and } v = (1 - \tilde{v}) \in H(a_\varepsilon),$$

whereas $\mathcal{G}_\varepsilon(\tilde{v}) = +\infty$ otherwise.

Theorem 4.2. *Besides the previous notation and assumptions, let us further assume that the following constants satisfy $0 < \tilde{\beta} \leq \tilde{q} \leq 2$, and that*

$$\limsup_{\varepsilon \rightarrow 0^+} \frac{\eta(\varepsilon)^{2q_1}}{\varepsilon^{\tilde{q}}} < +\infty,$$

and, finally, that $a_\varepsilon \geq 2\eta(\varepsilon)$.

Let $u_0 = u(f_0, K_0)$. For any ε , $0 < \varepsilon \leq 1$, let

$$m_\varepsilon = \inf\{\mathcal{G}_\varepsilon(\tilde{v}) : \tilde{v} \in L^1(\Omega)\}.$$

Then we have that, for some constant C , $m_\varepsilon \leq C$ for any ε , $0 < \varepsilon \leq 1$. Furthermore, if $\tilde{v}_\varepsilon \in L^1(\Omega)$ is such that

$$\mathcal{G}_\varepsilon(\tilde{v}_\varepsilon) \leq C \quad \text{for any } \varepsilon, \quad 0 < \varepsilon \leq 1,$$

the following holds. For any ε , let $v_\varepsilon = 1 - \tilde{v}_\varepsilon$ and $\tilde{u}_\varepsilon = \tilde{u}_\varepsilon(\tilde{v}_\varepsilon)$. Then we have that $\psi_{\eta(\varepsilon)}(v_\varepsilon)\tilde{u}_\varepsilon \rightarrow u_0$ strongly in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$, and $\psi_{\eta(\varepsilon)}(v_\varepsilon)\nabla \tilde{u}_\varepsilon$ converges to ∇u_0 weakly in $L^2(\Omega)$.

Furthermore, for any constant c , $c_1 < c < c_2$, the sets $\overline{\{v_\varepsilon > c\}}$ converge, as $\varepsilon \rightarrow 0^+$, to $\overline{G_{K_0}}$ in the Hausdorff distance.

Remark 4.3. We remark that the theorem in particular hold for a family $\tilde{v}_\varepsilon \in L^1(\Omega)$ of minimizers or quasi-minimizers, that is satisfying

$$\lim_{\varepsilon \rightarrow 0^+} (\mathcal{G}_\varepsilon(\tilde{v}_\varepsilon) - m_\varepsilon) = 0.$$

PROOF. By Proposition 3.2, we infer that there exists a constant C such that for any ε , $0 < \varepsilon \leq 1$, and for any $\tilde{v} \in W(\Omega)$, we have

$$\int_\Omega w_\eta(\tilde{v}) |\nabla \tilde{u}_\varepsilon(\tilde{v})|^2 \leq C \quad \text{and} \quad \|\tilde{u}_\varepsilon(\tilde{v})\|_{L^\infty(\Omega)} \leq C.$$

By a construction pretty similar to the one used in Proposition 4.5 in [23], we may construct $\tilde{v}_\varepsilon \in W(\Omega)$ and $u_\varepsilon \in W^{1,2}(\Omega)$, for any ε , $0 < \varepsilon \leq 1$, such that the following properties hold. For any ε , $0 < \varepsilon \leq 1$, first $v_\varepsilon = (1 - \tilde{v}_\varepsilon) \in H(a_\varepsilon)$ and

$$\frac{1}{\eta} \int_\Omega W(v_\varepsilon) + \eta \int_\Omega |\nabla v_\varepsilon|^2 \leq C.$$

Second, on γ we have $u_\varepsilon|_\gamma = u_0|_\gamma = g_0$. Finally,

$$\int_{\Omega} w_\eta(\tilde{v}_\varepsilon) |\nabla(u_\varepsilon - \tilde{u}_\varepsilon(\tilde{v}_\varepsilon))|^2 \leq C\varepsilon^{\bar{q}}.$$

By Poincaré inequality in Ω_1 , we conclude that

$$(4.3) \quad \int_{\gamma} |\tilde{u}_\varepsilon(\tilde{v}_\varepsilon) - g_\varepsilon|^2 \leq 2 \left(\int_{\gamma} |\tilde{u}_\varepsilon(\tilde{v}_\varepsilon) - u_\varepsilon|^2 + \int_{\gamma} |g_0 - g_\varepsilon|^2 \right) \leq C(\varepsilon^{\bar{q}} + \varepsilon^2).$$

We immediately conclude that for a constant C we have $m_\varepsilon \leq C$ for any ε , $0 < \varepsilon \leq 1$.

For any $n \in \mathbb{N}$, let $\varepsilon_n > 0$ be such that $\lim_n \varepsilon_n = 0$ and let $\mathcal{G}_n = \mathcal{G}_{\varepsilon_n}$. Let \tilde{v}_n be such that $\mathcal{G}_n(\tilde{v}_n) \leq C$, for any $n \in \mathbb{N}$. Then by Remark 2.3, we obtain that, up to a subsequence, v_n converges in $L^1(\Omega)$, and actually in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$, and almost everywhere in Ω , to a function v . Such a function v is such that $P(v)$ is finite. Furthermore, by the definition of $H(a)$ and the compactness properties of \mathcal{B} , we may also assume that there exists $K \in \mathcal{B}$ such that $v = 1$ almost everywhere in G_K and $v = 0$ almost everywhere in $\Omega \setminus G_K$. In other words, $v = \chi_{G_K}$.

Let us call $w_n = w_{\eta(\varepsilon_n)}(\tilde{v}_n)$ and $\tilde{u}_n = \tilde{u}_{\varepsilon_n}(\tilde{v}_n)$. Let us notice that $\sqrt{w_n} \nabla \tilde{u}_n$ is uniformly bounded in $L^2(\Omega, \mathbb{R}^N)$, therefore, up to a subsequence, $\sqrt{w_n} \nabla \tilde{u}_n$ converges to $V \in L^2(\Omega, \mathbb{R}^N)$ weakly in $L^2(\Omega, \mathbb{R}^N)$.

Since $v_n \rightarrow \chi_{G_K}$ almost everywhere in Ω , we conclude that also w_n and $\sqrt{w_n}$ converge to χ_{G_K} almost everywhere in Ω and in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$.

For any $B_{2R}(y) \subset (\Omega \setminus \overline{G_K})$, we have that w_n converges to zero almost everywhere in $B_{2R}(y)$. By the uniform L^∞ bound on \tilde{u}_n and by the dominated convergence theorem, we conclude that $\int_{B_{2R}(y)} w_n \tilde{u}_n^2 \rightarrow 0$ as $n \rightarrow \infty$. By the Caccioppoli inequality described in Lemma 2.7, (2.14), we conclude that $\sqrt{w_n} \nabla \tilde{u}_n$ converges to 0 strongly in $L^2(B_{2R}(y), \mathbb{R}^N)$, consequently $V = 0$ almost everywhere in $\Omega \setminus \overline{G_K}$. We conclude that $w_n \nabla \tilde{u}_n$ weakly converges to V in $L^2(\Omega, \mathbb{R}^N)$ as well. On the other hand, again up to subsequences and by using the property of $H(a)$, we may follow the arguments of the proof of Proposition 4.3 in [23] in order to find a function \tilde{u} with the following properties. First, $\sqrt{w_n} \tilde{u}_n$ converges to \tilde{u} almost everywhere in Ω and consequently in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$. Second, by the same reasoning above, we conclude that $\tilde{u} = 0$ almost everywhere in $\Omega \setminus \overline{G_K}$ and that also $w_n \tilde{u}_n$ converges to \tilde{u} almost everywhere in Ω and in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$. Then, we have that $\tilde{u} \in W^{1,2}(G_K)$, \tilde{u} is harmonic in G_K and $\nabla \tilde{u} = V$ in G_K . We also have that on γ and on $\tilde{\gamma}$, \tilde{u}_n converges to \tilde{u} strongly in $L^p(\gamma \cap \tilde{\gamma})$ for any p , $1 \leq p < +\infty$. As a consequence, $\tilde{u} = g_0$ on γ .

Let us take any function $\varphi \in W^{1,2}(G_K)$. Since G_K is a domain with Lipschitz boundary, therefore it is an extension domain, we can find a function $\tilde{\varphi} \in W^{1,2}(\Omega)$ such that $\tilde{\varphi} = \varphi$ on G_K . We conclude that for any $n \in \mathbb{N}$ we have

$$\int_{\Omega} w_n \nabla \tilde{u}_n \cdot \nabla \tilde{\varphi} = \int_{\tilde{\gamma}} f_{\varepsilon_n} \tilde{\varphi}.$$

Since, as $n \rightarrow \infty$,

$$\int_{\Omega} w_n \nabla \tilde{u}_n \cdot \nabla \tilde{\varphi} \rightarrow \int_{G_K} \nabla \tilde{u} \cdot \nabla \tilde{\varphi}$$

and

$$\int_{\tilde{\gamma}} f_{\varepsilon_n} \tilde{\varphi} \rightarrow \int_{\tilde{\gamma}} f_0 \tilde{\varphi},$$

we conclude that

$$\int_{G_K} \nabla \tilde{u} \cdot \nabla \varphi = \int_{\tilde{\gamma}} f_0 \varphi \quad \text{for any } \varphi \in W^{1,2}(G_K).$$

Then \tilde{u} solves (3.1)-(3.2)-(3.3) with K_0 replaced by K . Then, since $\tilde{u} = u_0$ on γ , we conclude by using Theorems 3.3 and 3.6 in [21] that $\tilde{u} = u_0$ almost everywhere in Ω and that $G_K = G_{K_0}$. The rest of the proof easily follows. \square

We conclude this section with the following existence results, which may be easily proved by the direct method.

Proposition 4.4. *The following problems admit a solution.*

- (i) $\min \tilde{\mathcal{G}}_\varepsilon$ on $W(\Omega)$, with constraint $0 \leq \tilde{v} \leq 1$.
- (ii) $\min \tilde{\mathcal{G}}_\varepsilon$ on $W(\Omega)$, with constraints $0 \leq \tilde{v} \leq 1$ and $v \in H(a_\varepsilon)$ (that is there exists the minimum of \mathcal{G}_ε over $L^1(\Omega)$).

5 The crack case

In this section we shall deal with the determination of general defects, in particular of cracks. We begin by recalling results proved in [23]. We include them here for the convenience of the reader and to compare them with the new results devoted to the determination of material losses, in particular of cavities, which we treated in Section 4. For what concerns the classes of admissible defects we shall use in this section, let us begin with the following definition. We limit ourselves to the two or three-dimensional case, however it is not difficult to see how these definitions can be generalized to higher dimensions.

If $N = 2$, fixed a positive constant $L \geq 1$, we say that Γ is an *L-Lipschitz*, or *L-C^{0,1}*, arc if, up to a rigid transformation, $\Gamma = \{(x, y) \in \mathbb{R}^2 : -a/2 \leq x \leq a/2, y = \varphi_1(x)\}$, where $L^{-1} \leq a \leq L$ and $\varphi_1 : \mathbb{R} \rightarrow \mathbb{R}$ is a Lipschitz map with Lipschitz constant bounded by L and such that $\varphi_1(0) = 0$. For any α , $0 \leq \alpha \leq 1$, we say that Γ is an *L-C^{1,\alpha}* arc if φ_1 is $C^{1,\alpha}$ and its $C^{1,\alpha}$ norm is bounded by L . The points $(a/2, \varphi_1(a/2))$ and $(-a/2, \varphi_1(-a/2))$ will be called the *vertices* or *endpoints* of the arc Γ .

Let us consider now the case $N = 3$. Let T be the closed equilateral triangle which is contained in the plane $\Pi = \{(x, y, z) \in \mathbb{R}^3 : z = 0\}$ with vertices $V_1 = (0, 1, 0)$, $V_2 = (-\sqrt{3}/2, -1/2, 0)$ and $V_3 = (\sqrt{3}/2, -1/2, 0)$ and $T' \subset \mathbb{R}^2$ be its projection on the plane Π . Fixed a positive constant $L \geq 1$, we call an *L-Lipschitz*, or *L-C^{0,1}*, *generalized triangle* a set Γ such that, up to a rigid transformation, $\Gamma = \{(x, y, z) \in \mathbb{R}^3 : (x, y) \in \varphi(T'), z = \varphi_1(x, y)\}$, where $\varphi : \mathbb{R}^2 \rightarrow \mathbb{R}^2$ is a bi-Lipschitz function with constant L such that $\varphi(0) = 0$ and $\varphi_1 : \mathbb{R}^2 \rightarrow \mathbb{R}$ is a Lipschitz map with Lipschitz constant bounded by L and such that $\varphi_1(0) = 0$. For any α , $0 \leq \alpha \leq 1$, we say that Γ is an *L-C^{1,\alpha}* *generalized triangle* if φ_1 is $C^{1,\alpha}$ and its $C^{1,\alpha}$ norm is bounded by L .

In both cases, the image through φ of any vertex or side of T' will be called a generalized vertex or generalized side of $\varphi(T')$, respectively. The image on

the graph of φ_1 of one of the generalized vertices of $\varphi(T')$ will be called a *generalized vertex* of Γ , whereas the image of one of the generalized sides of $\varphi(T')$ will be called a *generalized side* of Γ . We also remark that there exists a constant $L_1 > 0$, depending on L only, such that we can find $\varphi_2 : \mathbb{R}^3 \rightarrow \mathbb{R}^3$, a bi-Lipschitz function with constant L_1 , such that $\Gamma = \varphi_2(T)$.

Definition 5.1. Let us assume that $\Omega \subset B_R \subset \mathbb{R}^N$, with $R \geq 1$ and $N = 2, 3$. For any positive constants $L \geq 1$, δ and c , $c < 1$, any $k = 0, 1$ and α , $0 \leq \alpha \leq 1$, such that $k + \alpha \geq 1$, we define $\mathcal{B}(N, (k, \alpha), L, \delta, c)$ in the following way. We say that $A \in \mathcal{B}(N, (k, \alpha), L, \delta, c)$ if and only if $A \subset \overline{B_{2R}}$, there exists a positive integer n , depending on A , such that $A = \bigcup_{i=1}^n \Gamma_i$, Γ_i an L - $C^{k, \alpha}$ arc (if $N = 2$) or generalized triangle (if $N = 3$) for any $i = 1, \dots, n$, such that the following conditions are satisfied:

- i) for any $i, j \in \{1, \dots, n\}$ with $i \neq j$, we have that either $\Gamma_i \cap \Gamma_j$ is not empty or $\text{dist}(\Gamma_i, \Gamma_j) \geq \delta$;
- ii) for any $i, j \in \{1, \dots, n\}$ with $i \neq j$, if $\Gamma_i \cap \Gamma_j$ is not empty then $\Gamma_i \cap \Gamma_j$ is a common endpoint V if $N = 2$ and either a common generalized vertex V or a common generalized side γ if $N = 3$. Furthermore, in such a case, for any $x \in \Gamma_i$ we have $\text{dist}(x, \Gamma_j) \geq c|x - V|$ or $\text{dist}(x, \Gamma_j) \geq c \text{dist}(x, \gamma)$, respectively.

Let us remark that there exists an integer M , depending on N, R, L, δ and c only, such that for any $A \in \mathcal{B}(N, (k, \alpha), L, \delta, c)$ we have that $n \leq M$.

More importantly, we have that any of the classes \mathcal{B} described in Definition 5.1 is non-empty, is composed of non-empty compact sets and it is compact with respect to the Hausdorff distance. Finally, if A belongs to any of these classes, then $\mathcal{H}^{N-1}(A)$ is bounded by a constant depending on the class only.

For the time being, let us fix \mathcal{B} as one of the classes of Definition 5.1. We call the constant k, α, L, δ and c the *a priori data* related to \mathcal{B} . For any such class \mathcal{B} we call \mathcal{B}' the class of admissible defects K such that $\text{dist}(K, \overline{\Omega_1}) \geq \delta$, $\mathcal{H}^{N-2}(K \cap \partial\Omega) < +\infty$ and there exists $A \in \mathcal{B}$ such that $K \subset A$ and $\mathcal{H}^{N-2}(K \cap A \setminus \overline{K}) < +\infty$.

Moreover, we say that $K \in \mathcal{B}'$ satisfies Assumption A if the following holds.

Assumption A. We assume that, for any $x_0 \in K \cap \Omega$, there exists $r > 0$, depending on x_0 , such that for any U connected component of $(\Omega \setminus K) \cap B_r(x_0)$ we can find $r_1 > 0$, an open set U_1 , such that $U \cap B_{r_1}(x_0) \subset U_1 \subset U$, and a bijective map $T : U_1 \rightarrow (-1, 1)^N$ such that the following properties hold. The maps T and T^{-1} are locally Lipschitz and there exists a constant C such that $\|DT\|$ and $\|DT^{-1}\|$ are bounded by C almost everywhere. By the regularity of $Q = (-1, 1)^N$, T^{-1} can be actually extended up to the boundary and we have that $T^{-1} : [-1, 1]^N \rightarrow \mathbb{R}^N$ is a Lipschitz map with Lipschitz constant bounded by C . Furthermore, if we set $\Gamma = [-1, 1]^{N-1} \times \{1\}$, we require that $T^{-1}(\Gamma) = \partial U_1 \cap K_0$, $T^{-1}(0, \dots, 0, 1) = x_0$ and $T^{-1}(y) \in \Omega \setminus K$ for any $y \in [-1, 1]^N \setminus \Gamma$.

We assume that, for any $x_0 \in K \cap \partial\Omega$, there exists $r > 0$, depending on x_0 , such that for any U connected component of $(\Omega \setminus K) \cap B_r(x_0)$ we can find $r_1 > 0$, an open set U_1 , such that $U \cap B_{r_1}(x_0) \subset U_1 \subset U$, and a bijective map $T : U_1 \rightarrow (0, 1) \times (-1, 1)^{N-1}$ such that the following properties hold. The maps T and T^{-1} are locally Lipschitz and there exists a constant C such that

$\|DT\|$ and $\|DT^{-1}\|$ are bounded by C almost everywhere. By the regularity of $Q_1 = (0, 1) \times (-1, 1)^{N-1}$, T^{-1} can be actually extended up to the boundary and we have that $T^{-1} : \overline{Q_1} \rightarrow \mathbb{R}^N$ is a Lipschitz map with Lipschitz constant bounded by C . Furthermore, if we set $\Gamma_1 = [0, 1] \times [-1, 1]^{N-2} \times \{1\}$ and $\Gamma_2 = \{0\} \times [-1, 1]^{N-1}$, we require that $T^{-1}(\Gamma_1) = \partial U_1 \cap K$, $T^{-1}(\Gamma_2) = \partial U_1 \cap \partial \Omega$, $T^{-1}(0, \dots, 0, 1) = x_0$ and $T^{-1}(y) \in \Omega \setminus K$ for any $y \in \overline{Q_1} \setminus (\Gamma_1 \cup \Gamma_2)$.

In the sequel we shall fix positive constants $L \geq 1$, δ and c , $c < 1$, and α , $0 \leq \alpha \leq 1$. We also assume that $\Omega \subset B_R$, for some fixed constant $R \geq 1$. Let $\mathcal{B} = \mathcal{B}(N, (1, \alpha), L, \delta, c)$. We assume that the unknown defect K_0 belongs to \mathcal{B}' and that it satisfies Assumption A. We recall that examples of defects satisfying Assumption A are described in [22, 23].

The next proposition states that the gradient of u_0 satisfies a higher integrability property.

Proposition 5.2. *Under the previous assumptions, there exist a constant $q > 2$ and a constant $C > 0$, which do not depend on f_0 , such that $\nabla u_0 \in L^q(\Omega, \mathbb{R}^N)$, in particular*

$$\|\nabla u_0\|_{L^q(\Omega)} \leq C \|f_0\|_{L^s(\tilde{\gamma})}.$$

PROOF. See the proof of Proposition 4.5 in [22]. \square

We remark that the constants q and C in Proposition 5.2 depend also on s and on K_0 .

For any $a > 0$, we call $H_1(a)$ the set of functions $v \in W^{1,2}(\Omega, [0, 1])$ such that $v = 1$ almost everywhere in $\tilde{\Omega}_1$ and for some $A \in \mathcal{B}$ we have $v \geq c_1$ almost everywhere in $\Omega \setminus \overline{B}_a(A)$, where again c_1 is a constant such that $0 < c_1 < 1$.

For any $0 < \varepsilon \leq 1$ and any $q \geq 2$, let us define $\tilde{\mathcal{F}}_\varepsilon^q : W_\gamma^{1,q}(\Omega) \times W(\Omega) \rightarrow \mathbb{R}$ as follows. For any $(u, \tilde{v}) \in W_\gamma^{1,q}(\Omega) \times W(\Omega)$, recalling that $v = 1 - \tilde{v}$, we set

$$(5.1) \quad \tilde{\mathcal{F}}_\varepsilon^q(u, \tilde{v}) = \frac{a_1}{\varepsilon^{\tilde{q}}} |u - \tilde{u}_\varepsilon|_{w_\eta}^2 + \frac{a_2}{\varepsilon^{\tilde{\beta}}} \int_\gamma |u - g_\varepsilon|^2 + b \int_\Omega \psi_\eta(v) |\nabla u|^q + \frac{1}{\eta} \int_\Omega V(v) + \eta \int_\Omega |\nabla v|^2.$$

Here $\eta = \eta(\varepsilon)$, $o_\eta = o_\eta(q)$, $w_\eta = w_{\eta(\varepsilon)}(\tilde{v}) = \psi_{\eta(\varepsilon)}(v)$ and $\tilde{u}_\varepsilon = \tilde{u}_\varepsilon(\tilde{v})$ is the solution to (3.10). We also recall that

$$|u - \tilde{u}_\varepsilon|_{w_\eta}^2 = \int_\Omega \psi_{\eta(\varepsilon)}(v) |\nabla(u - \tilde{u}_\varepsilon)|^2 = \int_\Omega \psi_{\eta(\varepsilon)}(v) |\nabla u|^2 - 2 \int_{\tilde{\gamma}} f_\varepsilon u + \int_{\tilde{\gamma}} f_\varepsilon \tilde{u}_\varepsilon.$$

Then, for any $0 < \varepsilon \leq 1$ and any $q \geq 2$, we define $\mathcal{F}_\varepsilon^q$ as the following functional on $L^1(\Omega) \times L^1(\Omega)$. For any $(u, \tilde{v}) \in L^1(\Omega) \times L^1(\Omega)$ we set

$$(5.2) \quad \mathcal{F}_\varepsilon^q(u, \tilde{v}) = \tilde{\mathcal{F}}_\varepsilon^q(u, \tilde{v}) \quad \text{if } (u, \tilde{v}) \in W_\gamma^{1,q}(\Omega) \times W(\Omega) \text{ and } v = (1 - \tilde{v}) \in H_1(a_\varepsilon),$$

whereas $\mathcal{F}_\varepsilon^q(u, \tilde{v}) = +\infty$ otherwise.

Now we shall fix the constant $q > 2$ as the one defined in Proposition 5.2, which depends on K_0 , among other things. Again we set $q_1 = (q - 2)/(2q)$ and we observe that $0 < q_1 < 1/2$. The following convergence result is the main result of [23].

Theorem 5.3. *Besides the previous notation and assumptions, let us further assume that the following constants satisfy $0 < \tilde{q} \leq 2$, $0 < \tilde{\beta} \leq 2$, and that*

$$\limsup_{\varepsilon \rightarrow 0^+} \frac{\eta(\varepsilon)^{2q_1}}{\varepsilon^{\tilde{q}}} < +\infty,$$

and, finally, that $a_\varepsilon \geq 2\eta(\varepsilon)$.

Let $u_0 = u(f_0, K_0)$. Then there exists a constant E_0 , E_0 depending on s , Ω , Ω_1 , $\tilde{\Omega}_1$, γ and $\tilde{\gamma}$ only, such that for any E , $E_0 \leq E < +\infty$, the following holds.

For any $0 < \varepsilon \leq 1$, let

$$m_\varepsilon = \inf\{\mathcal{F}_\varepsilon^q(u, \tilde{v}) : (u, \tilde{v}) \in L^1(\Omega) \times L^1(\Omega) \text{ and } \|u\|_{L^\infty(\Omega)} \leq E\}.$$

Then we have that, for some constant C , $m_\varepsilon \leq C$ for any $0 < \varepsilon \leq 1$.

For any $n \in \mathbb{N}$, let $\varepsilon_n > 0$ be such that $\lim_n \varepsilon_n = 0$ and let $(u_n, \tilde{v}_n) \in L^1(\Omega) \times L^1(\Omega)$ be such that $\|u_n\|_{L^\infty(\Omega)} \leq E$ and

$$\mathcal{F}_{\varepsilon_n}^q(u_n, \tilde{v}_n) \leq C \quad \text{for any } n \in \mathbb{N}.$$

Then, up to a subsequence, $u_n \rightarrow u$ strongly in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$, and $\psi_{\eta(\varepsilon_n)}(v_n) \nabla u_n \rightarrow \nabla u$ strongly in $L^p(\Omega)$ for any $2 \leq p < q$, where $u = u_0$ almost everywhere in G_{K_0} and $\nabla u = \nabla u_0$ almost everywhere in Ω .

Furthermore, there exist compact sets $\tilde{A} \subset \bar{\Omega}$ and $A \in \mathcal{B}$, such that $\tilde{A} \subset A$ and $\mathcal{H}^{N-1}(J(u) \setminus \tilde{A}) = 0$, satisfying the following property. For any constant c , $0 < c \leq c_1$, the sets $\{v_n < c\}$ converge, as $n \rightarrow \infty$, to \tilde{A} in the Hausdorff distance.

An analogous to Proposition 4.4 holds true, again easily proved by the direct method.

Proposition 5.4. *Let E_0 be as in Theorem 5.3. Then for any p , $2 \leq p \leq q$, and any E , $E_0 \leq E \leq +\infty$, the following problems admit a solution.*

- (i) $\min \tilde{\mathcal{F}}_\varepsilon^p$ on $W_\gamma^{1,p}(\Omega) \times W(\Omega)$, with constraints $0 \leq \tilde{v} \leq 1$ and $\|u\|_{L^\infty(\Omega)} \leq E$.
- (ii) $\min \tilde{\mathcal{F}}_\varepsilon^p$ on $W_\gamma^{1,p}(\Omega) \times W(\Omega)$, with constraints $0 \leq \tilde{v} \leq 1$, $v \in H_1(a_\varepsilon)$ and $\|u\|_{L^\infty(\Omega)} \leq E$ (that is there exists the minimum of $\mathcal{F}_\varepsilon^p$ over $L^1(\Omega) \times L^1(\Omega)$ with the same L^∞ bound on u).

Let us now consider the main differences between the cracks and material losses cases. Our aim is to show the optimality of Theorem 5.3, by showing that a reduction to a functional depending on the phase-variable only, with similar convergence properties, may not be feasible. As we have shown in the previous section such a reduction is instead possible in the material loss case.

By Proposition 2.5 and Proposition 3.2, we infer that there exists a constant C such that for any ε , $0 < \varepsilon \leq 1$, and for any $\tilde{v} \in W(\Omega)$, we have

$$\int_\Omega w_\eta(\tilde{v}) |\nabla \tilde{u}_\varepsilon(\tilde{v})|^2 \leq C \quad \text{and} \quad \|\tilde{u}_\varepsilon(\tilde{v})\|_{L^\infty(\Omega)} \leq C.$$

Furthermore, there exists $q(\varepsilon) > 2$, depending on N , Ω , s and ε only, such that $\tilde{u}_\varepsilon(\tilde{v})$ belongs to $W^{1,q(\varepsilon)}(\Omega)$. We can also find a constant C_1 , depending on N , Ω , γ , s , $\|f_0\|_{L^s(\partial\Omega)}$ and ε only, such that for any $\tilde{v} \in W(\Omega)$

$$\|\nabla \tilde{u}_\varepsilon(\tilde{v})\|_{L^{q(\varepsilon)}(\Omega)} \leq C_1.$$

We remark that the dependence of $q(\varepsilon)$ on ε is through $o_{\eta(\varepsilon)}$ and that, unfortunately, it might happen that $q(\varepsilon) \rightarrow 2^+$ and $C_1 \rightarrow +\infty$ as $\varepsilon \rightarrow 0^+$.

Let us consider the following operator. For any ε , $0 < \varepsilon \leq 1$, we define $\mathcal{H}_\varepsilon : W(\Omega) \rightarrow W_\gamma^{1,2}(\Omega)$ as follows

$$\mathcal{H}_\varepsilon(\tilde{v}) = \tilde{u}_\varepsilon(\tilde{v}) \quad \text{for any } \tilde{v} \in W(\Omega).$$

We recall that for any r , $1 < r < +\infty$, we endow $W_\gamma^{1,r}(\Omega)$ with the norm $\|u\|_{W_\gamma^{1,r}(\Omega)} = \|\nabla u\|_{L^r(\Omega)}$ for any $u \in W_\gamma^{1,r}(\Omega)$. We observe that \mathcal{H}_ε is continuous with respect to the weak- $W^{1,2}(\Omega)$ convergence in $W(\Omega)$ and strong convergence in $W_\gamma^{1,2}(\Omega)$.

We obtain that for any q , $2 \leq q \leq q(\varepsilon)$, we have that $\mathcal{H}_\varepsilon : W(\Omega) \rightarrow W_\gamma^{1,q}(\Omega)$ and that for any q , $2 \leq q < q(\varepsilon)$, \mathcal{H}_ε is continuous again with respect to the weak- $W^{1,2}(\Omega)$ convergence in $W(\Omega)$ and strong convergence in $W_\gamma^{1,q}(\Omega)$.

Then for any $q \geq 2$, let us define $\hat{\mathcal{F}}_\varepsilon^q : W(\Omega) \rightarrow [0, +\infty]$ as follows. For any $\tilde{v} \in W(\Omega)$ we set

$$(5.3) \quad \hat{\mathcal{F}}_\varepsilon^q(\tilde{v}) = \tilde{\mathcal{F}}_\varepsilon^q(\mathcal{H}_\varepsilon(\tilde{v}), \tilde{v}) = \frac{a_2}{\varepsilon^{\tilde{\beta}}} \int_\gamma |\tilde{u}_\varepsilon(\tilde{v}) - g_\varepsilon|^2 + b \int_\Omega \psi_\eta(v) |\nabla \tilde{u}_\varepsilon(\tilde{v})|^q + \frac{1}{\eta} \int_\Omega V(v) + \eta \int_\Omega |\nabla v|^2.$$

Let us notice that for any $q \geq 2$, we have that there exists $\min \hat{\mathcal{F}}_\varepsilon^q$ on $W(\Omega)$ with the constraint $0 \leq \tilde{v} \leq 1$, and with the constraints $0 \leq \tilde{v} \leq 1$ and $v \in H_1(a_\varepsilon)$ as well.

We investigate whether, for some $q \geq 2$, we may have convergence properties for $\hat{\mathcal{F}}_\varepsilon^q$ as we have for $\tilde{\mathcal{F}}_\varepsilon^q$. We observe that $\hat{\mathcal{G}}_\varepsilon$ is equal to $\hat{\mathcal{F}}_\varepsilon^2$ but to replace the single-well potential V with the double-well potential W . It would be desirable to have a convergence result for $\hat{\mathcal{F}}_\varepsilon^2$, or at least for $\hat{\mathcal{F}}_\varepsilon^q$ with some $q > 2$, as we have for $\hat{\mathcal{G}}_\varepsilon$, Theorem 4.2. By counterexamples we show that difficulties arise in both cases. We begin with the case $q = 2$ and then we deal with the case $q > 2$.

By the construction used in [23, Proposition 4.5] the next proposition immediately follows.

Proposition 5.5. *Under the assumptions of Theorem 5.3 and if $0 < \tilde{\beta} \leq \tilde{q} \leq 2$, we can find \tilde{v}_ε for any ε , $0 < \varepsilon \leq 1$, such that the following holds. For any ε , $0 < \varepsilon \leq 1$, we have, first, that*

$$\hat{\mathcal{F}}_\varepsilon^2(\tilde{v}_\varepsilon) \leq C.$$

Second, $\{v_\varepsilon < 1/2\} = \{x \in \Omega : \text{dist}(x, K_0) < \xi_\eta + \eta/2\}$ where $\xi_\eta = \sqrt{\eta o_\eta}$. Finally, for any $n \in \mathbb{N}$, let $\varepsilon_n > 0$ be such that $\lim_n \varepsilon_n = 0$ and let $\tilde{v}_n = \tilde{v}_{\varepsilon_n}$ and $\tilde{u}_n = \tilde{u}_{\varepsilon_n}(\tilde{v}_n)$. Then, up to a subsequence, $\tilde{u}_n \rightarrow u$ strongly in $L^p(\Omega)$ for any p , $1 \leq p < +\infty$, and $\psi_{\eta(\varepsilon_n)}(v_n) \nabla \tilde{u}_n \rightarrow \nabla u$ strongly in $L^2(\Omega)$, where $u = u_0$ almost everywhere in G_{K_0} and $\nabla u = \nabla u_0$ almost everywhere in Ω .

In terms of Γ -convergence, we have obtained a kind of Γ -limsup inequality. What is missing is the corresponding Γ -liminf inequality, because taking $q = 2$ does not guarantee enough compactness. In fact the solutions to the corresponding weighted elliptic problems may converge to a function which is not a solution to a material loss direct problem, as we shall show in Example 5.6 where we use the instability of the Neumann problem with respect to boundary variations.

In any case, trying to solve the inverse problem by minimizing $\hat{\mathcal{J}}_\varepsilon^2$ on $W(\Omega)$ with the constraints $0 \leq \tilde{v} \leq 1$ and $v \in H_1(a_\varepsilon)$, might be a good strategy. We recall that in this case the assumption $0 < \tilde{\beta} \leq \tilde{q} \leq 2$ should be adopted. In fact, minimizing $\hat{\mathcal{J}}_\varepsilon^2$ is numerically simpler than minimizing $\tilde{\mathcal{J}}_\varepsilon^q$ and still leads to good numerical reconstructions. In fact this method is adopted in [20] and the numerical simulations presented there show its efficacy.

Example 5.6. Let us consider the following example. Let D be a smooth bounded domain of \mathbb{R}^{N-1} , $N \geq 2$, and let $\lambda^2 > 0$ be a Neumann eigenvalue for $-\Delta$ on D and let f be a corresponding eigenfunction, that is

$$\begin{cases} -\Delta f = \lambda^2 f & \text{in } D \\ \nabla f \cdot \nu = 0 & \text{on } \partial D. \end{cases}$$

We notice that $\int_D f = 0$ and we may normalize f in such a way that $\int_D |f|^2 = 1$.

For some constant $T > 2$, to be fixed later, let $\Omega = D \times (0, T)$ and let $G_{K_0} = D \times (0, 2)$, that is $K_0 = \bar{D} \times \{2\}$. Let $\gamma = \tilde{\gamma} = \bar{D} \times \{0\}$. Then let u_0 be a solution to

$$\begin{cases} \Delta u_0 = 0 & \text{in } \Omega \\ \nabla u_0 \cdot \nu = f & \text{on } D \times \{0\} \\ \nabla u_0 \cdot \nu = 0 & \text{on } K_0 \\ \nabla u_0 \cdot \nu = 0 & \text{on } \partial D \times (0, T). \end{cases}$$

We normalize u_0 in such a way that $\int_\gamma u_0 = 0$ and, by separation of variables, we have that

$$u_0(x, y) = \frac{f(x)}{\lambda} \left[\frac{\cosh(2\lambda)}{\sinh(2\lambda)} \cosh(\lambda y) - \sinh(\lambda y) \right], \quad x \in D, \quad y \in (0, 2),$$

whereas u_0 may be chosen identically equal to 0 in $D \times (2, T)$.

By a simple computation, again by separation of variables, we may find $T > 2$ and $\mu > 0$ and two functions u^- and u^+ such that the following conditions hold. First, $u^- = u_0$ in $D \times (0, 1)$ and u^+ solves

$$\begin{cases} \Delta u^+ = 0 & \text{in } D \times (1, T) \\ \nabla u^+ \cdot \nu = 0 & \text{on } D \times \{T\} \\ \nabla u^+ \cdot \nu = 0 & \text{on } \partial D \times (1, T). \end{cases}$$

Second, the following transmission condition holds true on $D \times \{1\}$

$$u_y^-(x, 1) = u_y^+(x, 1) = \mu(u^+(x, 1) - u^-(x, 1)), \quad x \in D.$$

By following [19], we may then construct a Neumann sieve $K_\delta \subset K_0$, $\delta > 0$, such that if u_δ solves

$$\begin{cases} \Delta u_\delta = 0 & \text{in } \Omega \setminus K_\delta \\ \nabla u_\delta \cdot \nu = f & \text{on } \gamma \\ \nabla u_\delta \cdot \nu = 0 & \text{on } \partial(\Omega \setminus K_\delta) \setminus \gamma \\ \int_\gamma u_\delta = 0, \end{cases}$$

the following holds. We have that, as $\delta \rightarrow 0^+$, u_δ converges to $u^- = u_0$ weakly in $H^1(D \times (0, 1))$, and strongly in $L^2(D \times (0, 1))$, and u_δ converges to u^+ weakly in $H^1(D \times (1, T))$, and strongly in $L^2(D \times (1, T))$. Therefore, the Cauchy data of u_δ on γ converges, for instance in $L^2(\gamma)$, to the Cauchy data of u_0 on γ .

By using Proposition 5.5 to approximate K_δ and u_δ , for any $n \in \mathbb{N}$ we can find $\varepsilon_n > 0$, $\eta_n > 0$ and \tilde{v}_n such that

$$\hat{\mathcal{F}}_{\varepsilon_n}^2(\tilde{v}_n) \leq C \quad \text{for any } n \in \mathbb{N},$$

and that, as $n \rightarrow \infty$, the following holds. First, $\varepsilon_n \rightarrow 0^+$ and $\eta_n \rightarrow 0^+$. Second, if $\tilde{u}_n = \tilde{u}_{\varepsilon_n}(\tilde{v}_n)$, $n \in \mathbb{N}$, then we have that \tilde{u}_n converges to $u^- = u_0$ strongly in $L^2(D \times (0, 1))$ and \tilde{u}_n converges to u^+ strongly in $L^2(D \times (1, T))$. Furthermore, $\nabla \tilde{u}_n \cdot \nu|_\gamma = \nabla u_0 \cdot \nu|_\gamma$ for any $n \in \mathbb{N}$ and

$$\|\tilde{u}_n - u_0\|_{L^2(\gamma)} \rightarrow 0 \quad \text{as } n \rightarrow \infty.$$

Therefore, even if the Cauchy data of u_0 on γ are well approximated by those of $\tilde{u}_{\varepsilon_n}(\tilde{v}_n)$, we have that $v_n = 1 - \tilde{v}_n$ is small in a region close to the corresponding Neumann sieve which is far away from the actual location of the looked-for defect K_0 . This example shows also the difficulty in proving a convergence result without imposing any further condition on the region where v is small.

On the other hand, one might try to minimize $\hat{\mathcal{F}}_\varepsilon^q$ on $W(\Omega)$ for some $q > 2$. If we take $q > 2$, then compactness and convergence would follow as a simple consequence of Theorem 5.3, but we may not guarantee that we can find a sequence of phase-field functions \tilde{v}_n such that $\hat{\mathcal{F}}_{\varepsilon_n}^q(\tilde{v}_n)$ is uniformly bounded.

Again we use the constraints $0 \leq \tilde{v} \leq 1$ and $v \in H_1(a_\varepsilon)$. If one would be able to find \tilde{v}_ε , $0 < \varepsilon \leq 1$, such that $\hat{\mathcal{F}}_\varepsilon^q(\tilde{v}_\varepsilon) \leq C$ for any $0 < \varepsilon \leq 1$ for some constant C , then by Proposition 4.3 in [23], we would obtain the results of Theorem 5.3, replacing $\tilde{\mathcal{F}}_\varepsilon^q$ with $\hat{\mathcal{F}}_\varepsilon^q$, even allowing E to be equal to $+\infty$.

We believe that constructing such functions \tilde{v}_ε for some $q > 2$ is a difficult task and that minimizing $\hat{\mathcal{F}}_\varepsilon^q$ for some $q > 2$ might lead to a not correct reconstruction. In Proposition 5.7 below we show the difficulty of obtaining such a uniform bound.

In order to have higher integrability of the gradient of $\tilde{u}_\varepsilon(\tilde{v})$, we need to guarantee that $w_\eta(\tilde{v}) = \psi_\eta(v)$ is a weight satisfying certain properties, for instance those described by Stredulinsky in [25]. An important class of weights for which these properties are satisfied is the so-called Muckenhoupt class A_2 .

We recall that w , a non-negative measurable function over \mathbb{R}^N , is a *weight* if $0 < w < +\infty$ almost everywhere and w is locally integrable. We say that a weight w belongs to the *Muckenhoupt class* A_2 if there exists a constant C such that for any ball $B \subset \mathbb{R}^N$ we have

$$(5.4) \quad \left(\frac{1}{|B|} \int_B w \right) \left(\frac{1}{|B|} \int_B w^{-1} \right) \leq C.$$

The best constant C for which (5.4) holds is usually referred to as the A_2 -constant of w . We observe that the A_2 -constant of w is always greater than or equal to 1. For more details about the Muckenhoupt weights and weighted elliptic equations, we refer for instance to [14].

Therefore, a reasonable assumption is to take $w = w_\eta(\tilde{v})$ belonging to the Muckenhoupt class A_2 and such that its A_2 -constant is bounded by C , for some fixed C . Without loss of generality we can assume that $0 \leq w \leq 1$ almost

everywhere and that $w = 1$ outside a given ball B_{2R} . Consequently, we infer that there exists a constant C_1 , depending on C and R only, such that

$$\int_{B_R} w^{-1} \leq C_1.$$

Proposition 5.7. *Let us fix $q > 2$. Let $\varepsilon_n, n \in \mathbb{N}$, be a sequence of positive numbers such that $\lim_n \varepsilon_n = 0$. For any $n \in \mathbb{N}$, let $\hat{\mathcal{F}}_n^q = \hat{\mathcal{F}}_{\varepsilon_n}^q$ and let $\tilde{v}_n \in W(\Omega)$ be such that the following holds. For any $n \in \mathbb{N}$, we assume that $0 \leq \tilde{v}_n \leq 1$ and we set $\eta_n = \eta(\varepsilon_n)$, $v_n = (1 - \tilde{v}_n) \in H_1(a_{\varepsilon_n})$, $w_n = w_\eta(\tilde{v}_n)$ and $\tilde{u}_n = \tilde{u}_{\varepsilon_n}(\tilde{v}_n)$. For any $n \in \mathbb{N}$, we assume*

$$(5.5) \quad \int_{\Omega} w_n^{-1} \leq C_1$$

and

$$\int_{\Omega} w_n |\nabla \tilde{u}_n|^q + \frac{1}{\eta_n} \int_{\Omega} V(v_n) + \eta_n \int_{\Omega} |\nabla v_n|^2 \leq C.$$

Let us consider \tilde{u} as the solution to

$$(5.6) \quad \begin{cases} \Delta \tilde{u} = 0 & \text{in } \Omega \\ \nabla \tilde{u} \cdot \nu = f_0 & \text{on } \partial\Omega. \end{cases}$$

We assume that K_0 satisfies the assumption of Theorem 5.3 and that $\tilde{u} \neq u_0$, in particular that $\tilde{u}|_{\gamma} \neq u_0|_{\gamma} = g_0$. We also assume, for the time being, that Ω and f_0 are regular enough to guarantee that $\tilde{u} \in L^\infty(\Omega)$ and $\nabla \tilde{u} \in L^\infty(\Omega, \mathbb{R}^N)$. We may also assume that actually \tilde{v}_n provides a good approximation of K_0 , namely that $\{v_n < 1\} \subset B_{a_n}(K_0)$ where $a_n = a_{\varepsilon_n}$.

Then we have that, as $n \rightarrow \infty$, $w_n \nabla \tilde{u}_n$ converges to $\nabla \tilde{u}$ weakly in $L^2(\Omega)$. Consequently, as $n \rightarrow \infty$, we also have that $\int_{\gamma} |\tilde{u}_n - g_{\varepsilon_n}|^2 \rightarrow \int_{\gamma} |\tilde{u} - g_0|^2 \neq 0$ and $\hat{\mathcal{F}}_{\varepsilon_n}^q(\tilde{v}_n) \rightarrow +\infty$.

PROOF. Let us compute

$$\int_{\Omega} w_n |\nabla \tilde{u}_n - \nabla \tilde{u}|^2 = \int_{\partial\Omega} (f_{\varepsilon_n} - f_0)(\tilde{u}_n - \tilde{u}) + \int_{\Omega} (1 - w_n) \nabla \tilde{u} \cdot \nabla (\tilde{u}_n - \tilde{u}).$$

By the uniform bound on \tilde{u}_n , we easily obtain that

$$\int_{\partial\Omega} (f_{\varepsilon_n} - f_0)(\tilde{u}_n - \tilde{u}) \rightarrow 0 \quad \text{as } n \rightarrow \infty.$$

Let us evaluate the other term. We have

$$\int_{\Omega} (1 - w_n) \nabla \tilde{u} \cdot \nabla (\tilde{u}_n - \tilde{u}) = \int_{\Omega} \frac{(1 - w_n)}{\sqrt[q]{w_n}} \sqrt[q]{w_n} \nabla \tilde{u} \cdot \nabla (\tilde{u}_n - \tilde{u}).$$

We apply Hölder inequality with coefficients q, p and r such that $q^{-1} + p^{-1} + r^{-1} = 1$ and we obtain

$$\begin{aligned} \int_{\Omega} (1 - w_n) \nabla \tilde{u} \cdot \nabla (\tilde{u}_n - \tilde{u}) &\leq \\ &\left(\int_{\Omega} w_n^{-r/q} \right)^{1/r} \left(\int_{\Omega} w_n |\nabla (\tilde{u}_n - \tilde{u})|^q \right)^{1/q} \left(\int_{\Omega} |1 - w_n|^p |\nabla \tilde{u}|^p \right)^{1/p}. \end{aligned}$$

We use our assumptions to infer

$$\int_{\Omega} (1 - w_n) \nabla \tilde{u} \cdot \nabla (\tilde{u}_n - \tilde{u}) \leq C \left(\int_{\Omega} w_n^{-r/q} \right)^{1/r} \left(\int_{\Omega} |1 - w_n|^p \right)^{1/p},$$

with C independent of n . We may choose r such that $0 < r/q < 1$, therefore, since $0 \leq w_n \leq 1$, we have $w_n^{-r/q} \leq w_n^{-1}$. Hence, by (5.5), we conclude that $\int_{\Omega} w_n |\nabla \tilde{u}_n - \nabla \tilde{u}|^2$ goes to zero as $n \rightarrow \infty$. We obtain that $\int_{\gamma} |\tilde{u}_n - \tilde{u}|^2$ goes to zero as well. We then apply Theorem 4.4 in [23] and the proof is concluded. \square

Therefore, even if v_n is a good phase-field approximation of K_0 , \tilde{u}_n is not a good approximation of u_0 . In order to have that \tilde{u}_n approximates u_0 , we need to require that v_n is very small close to K_0 , in such a way that violates (5.5). In turn, this might suggest the fact that higher integrability and the correct approximation might in some sense oppose each other.

Let us conclude by observing that (5.5) is a kind of minimal condition to have $\int_{\Omega} w_n |\nabla \tilde{u}_n|^q$ uniformly bounded. We wish to point out that potential theory for weights whose inverse is not integrable has been developed, see for instance [10] for the case of weights $w = \omega^{1-p/N}$, $1 < p < N$, where ω is a so-called strong A_{∞} -weight. Strong A_{∞} -weights have been introduced in [9]. Following [24], an important example of strong A_{∞} -weights is given by

$$\omega(x) = \min\{1, \text{dist}(x, A)^s\}, \quad x \in \mathbb{R}^N,$$

where $s > 0$ and A is a suitable compact set. In [24, Proposition 4.4] it is shown that ω is a strong A_{∞} -weight for any $s > 0$ provided A is uniformly disconnected. On the other hand, no strong A_{∞} -weight may vanish on a rectifiable curve, therefore this class of weights seems to be not apt to approximate hypersurfaces as we require in our application.

6 Differentiability of the functionals

In this last section, we investigate the differentiability properties of $\tilde{\mathcal{F}}_{\varepsilon}^q$, $\tilde{\mathcal{G}}_{\varepsilon}$ and $\tilde{\mathcal{J}}_{\varepsilon}^q$, for a fixed ε , $0 < \varepsilon \leq 1$, and any $q \geq 2$. For this purpose, we further assume that the functions ψ , V and W are actually of class C^1 and such that their derivatives are bounded and uniformly continuous all over \mathbb{R} .

We define the following spaces. For any p , $2 \leq p \leq +\infty$, let us call $L_p(\Omega) = \{\tilde{v} \in L^p(\Omega) : \tilde{v} = 0 \text{ a.e. in } \tilde{\Omega}_1\}$ and $W_p(\Omega) = W^{1,2}(\Omega) \cap L_p(\Omega)$, with norm $\|\tilde{v}\|_{L_p(\Omega)} = \|\tilde{v}\|_{L^p(\Omega)}$ and $\|\tilde{v}\|_{W_p(\Omega)} = \|\tilde{v}\|_{L^p(\Omega)} + \|\nabla \tilde{v}\|_{L^2(\Omega)}$. To any $\tilde{v} \in L^2(\Omega)$ we as usual associate the function $v = 1 - \tilde{v}$. If \tilde{v} belongs either to $L_p(\Omega)$ or to $W_p(\Omega)$, then $v \in L^p(\Omega)$, $v = 1$ almost everywhere in $\tilde{\Omega}_1$, and, provided $0 \leq \tilde{v} \leq 1$ almost everywhere in Ω , we also have $0 \leq v \leq 1$ almost everywhere in Ω . We observe that $W_2(\Omega) = W(\Omega)$ as previously defined. We also recall that $W_{\gamma}^{1,q}(\Omega)$ is equipped with the norm $\|u\|_{W_{\gamma}^{1,q}(\Omega)} = \|\nabla u\|_{L^q(\Omega)}$ for any $u \in W_{\gamma}^{1,q}(\Omega)$.

We recall that for any ε , $0 < \varepsilon \leq 1$, we define $\mathcal{H}_{\varepsilon} : L_2(\Omega) \rightarrow W_{\gamma}^{1,2}(\Omega)$ as follows

$$\mathcal{H}_{\varepsilon}(\tilde{v}) = \tilde{u}_{\varepsilon}(\tilde{v}) \quad \text{for any } \tilde{v} \in L_2(\Omega).$$

It can be shown that for any $\tilde{v}_0 \in L_2(\Omega)$ such an operator $\mathcal{H}_{\varepsilon}$ is differentiable in \tilde{v}_0 with respect to the $L^{\infty}(\Omega)$ norm. Let $D\mathcal{H}_{\varepsilon}(\tilde{v}_0) : L^{\infty}(\Omega) \rightarrow W_{\gamma}^{1,2}(\Omega)$ be

the differential in \tilde{v}_0 . Then for any \tilde{v} in $L_\infty(\Omega)$ we have

$$D\mathcal{H}_\varepsilon(\tilde{v}_0)[\tilde{v}] = U_\varepsilon(\tilde{v}_0, \tilde{v})$$

where $U_\varepsilon = U_\varepsilon(\tilde{v}_0, \tilde{v}) \in W_\gamma^{1,2}(\Omega)$ solves the following problem

$$(6.1) \quad \begin{cases} \operatorname{div}(\psi_\eta(v_0)\nabla U_\varepsilon) = \operatorname{div}(\psi'_\eta(v_0)\tilde{v}\nabla(\mathcal{H}_\varepsilon(\tilde{v}_0))) & \text{in } \Omega, \\ \psi_\eta(v_0)\nabla U_\varepsilon \cdot \nu = 0 & \text{on } \partial\Omega. \end{cases}$$

Here, obviously, $v_0 = 1 - \tilde{v}_0$.

We recall that for any vector valued function $G \in L^2(\Omega, \mathbb{R}^N)$, $\operatorname{div}(G)$ defines a functional on $W^{1,2}(\Omega)$ in the following way

$$\operatorname{div}(G)[\phi] = - \int_\Omega G \cdot \nabla \phi \quad \text{for any } \phi \in W^{1,2}(\Omega).$$

Therefore, the weak formulation of (6.1) is looking for a function $U_\varepsilon \in W_\gamma^{1,2}(\Omega)$ such that

$$\int_\Omega \psi_\eta(v_0)\nabla U_\varepsilon \cdot \nabla \varphi = \int_\Omega \psi'_\eta(v_0)\tilde{v}\nabla(\mathcal{H}_\varepsilon(\tilde{v}_0)) \cdot \nabla \varphi \quad \text{for any } \varphi \in W^{1,2}(\Omega).$$

Here, and analogously in the sequel, the differentiability has to be understood in the following sense. For any \tilde{v} in $L_\infty(\Omega)$

$$\mathcal{H}_\varepsilon(\tilde{v}_0 + \tilde{v}) = \mathcal{H}_\varepsilon(\tilde{v}_0) + D\mathcal{H}_\varepsilon(\tilde{v}_0)[\tilde{v}] + R(\tilde{v})$$

where

$$\lim_{\|\tilde{v}\|_{L^\infty(\Omega)} \rightarrow 0} \frac{\|R(\tilde{v})\|_{W_\gamma^{1,2}(\Omega)}}{\|\tilde{v}\|_{L^\infty(\Omega)}} = 0.$$

For any $q \geq 2$, let us consider the functional $\tilde{\mathcal{F}}_\varepsilon^q : W_\gamma^{1,q}(\Omega) \times W(\Omega) \rightarrow \mathbb{R}$. For any $(u_0, \tilde{v}_0) \in W_\gamma^{1,q}(\Omega) \times W(\Omega)$, $\tilde{\mathcal{F}}_\varepsilon^q$ is differentiable in (u_0, \tilde{v}_0) , with respect to the $W_\gamma^{1,q}(\Omega) \times W_\infty(\Omega)$ norm. Let $D\tilde{\mathcal{F}}_\varepsilon^q(u_0, \tilde{v}_0) : W_\gamma^{1,q}(\Omega) \times W_\infty(\Omega) \rightarrow \mathbb{R}$ be the differential in (u_0, \tilde{v}_0) . Then, for any $(u, \tilde{v}) \in W_\gamma^{1,q}(\Omega) \times W_\infty(\Omega)$, we have

$$(6.2) \quad D\tilde{\mathcal{F}}_\varepsilon^q(u_0, \tilde{v}_0)[(u, \tilde{v})] = \begin{aligned} & \frac{a_1}{\varepsilon^{\frac{q}{2}}} \int_\Omega (2\psi_\eta(v_0)\nabla u_0 \cdot \nabla u - \psi'_\eta(v_0)|\nabla u_0|^2 \tilde{v}) + \\ & \frac{a_1}{\varepsilon^{\frac{q}{2}}} \int_{\tilde{\gamma}} (f_\varepsilon U_\varepsilon(\tilde{v}_0, \tilde{v}) - 2f_\varepsilon u) + \frac{2a_2}{\varepsilon^{\frac{2}{\beta}}} \int_\gamma (u_0 - g_\varepsilon)u + \\ & b \int_\Omega (q\psi_\eta(v_0)|\nabla u_0|^{q-2} \nabla u_0 \cdot \nabla u - \psi'_\eta(v_0)|\nabla u_0|^q \tilde{v}) + \\ & \frac{1}{\eta} \int_\Omega (-V'(v_0)\tilde{v}) + 2\eta \int_\Omega \nabla \tilde{v}_0 \cdot \nabla \tilde{v}. \end{aligned}$$

With the same computation, we infer that the functionals $\hat{\mathcal{F}}_\varepsilon^2 : W(\Omega) \rightarrow \mathbb{R}$ and $\hat{\mathcal{G}}_\varepsilon : W(\Omega) \rightarrow \mathbb{R}$ are differentiable in \tilde{v}_0 for any $\tilde{v}_0 \in W(\Omega)$, with respect to the $W_\infty(\Omega)$ norm. Let $D\hat{\mathcal{F}}_\varepsilon^2(\tilde{v}_0) : W_\infty(\Omega) \rightarrow \mathbb{R}$ and $D\hat{\mathcal{G}}_\varepsilon(\tilde{v}_0) : W_\infty(\Omega) \rightarrow \mathbb{R}$ be

the differentials in (u_0, \tilde{v}_0) . Then, for any $\tilde{v} \in W_\infty(\Omega)$, we have

$$(6.3) \quad D\hat{\mathcal{F}}_\varepsilon^2(\tilde{v}_0)[\tilde{v}] = \frac{2a_2}{\varepsilon^{\tilde{\beta}}} \int_\gamma (\mathcal{H}_\varepsilon(\tilde{v}_0) - g_\varepsilon) U_\varepsilon(\tilde{v}_0, \tilde{v}) + \\ b \int_\Omega (2\psi_\eta(v_0) \nabla \mathcal{H}_\varepsilon(\tilde{v}_0) \cdot \nabla U_\varepsilon(\tilde{v}_0, \tilde{v}) - \psi'_\eta(v_0) |\nabla \mathcal{H}_\varepsilon(\tilde{v}_0)|^2 \tilde{v}) + \\ \frac{1}{\eta} \int_\Omega (-V'(v_0) \tilde{v}) + 2\eta \int_\Omega \nabla \tilde{v}_0 \cdot \nabla \tilde{v}$$

and

$$(6.4) \quad D\tilde{\mathcal{G}}_\varepsilon(\tilde{v}_0)[\tilde{v}] = \frac{2a_2}{\varepsilon^{\tilde{\beta}}} \int_\gamma (\mathcal{H}_\varepsilon(\tilde{v}_0) - g_\varepsilon) U_\varepsilon(\tilde{v}_0, \tilde{v}) + \\ b \int_\Omega (2\psi_\eta(v_0) \nabla \mathcal{H}_\varepsilon(\tilde{v}_0) \cdot \nabla U_\varepsilon(\tilde{v}_0, \tilde{v}) - \psi'_\eta(v_0) |\nabla \mathcal{H}_\varepsilon(\tilde{v}_0)|^2 \tilde{v}) + \\ \frac{1}{\eta} \int_\Omega (-W'(v_0) \tilde{v}) + 2\eta \int_\Omega \nabla \tilde{v}_0 \cdot \nabla \tilde{v}.$$

It might be useful to have differentiability properties with respect to the $W_p(\Omega)$ norm, with p finite. In fact in this case $W_p(\Omega)$ is a strictly convex real reflexive Banach space and this is useful when we need to apply a gradient method in a numerical implementation, see [20] for details on the use of this information. In order to obtain such differentiability, let us now assume that ψ' , V' and W' are Hölder continuous for some exponent α , $0 < \alpha \leq 1$, all over \mathbb{R} .

We recall again that, by Proposition 2.5 and Proposition 3.2, there exists a constant C such that for any ε , $0 < \varepsilon \leq 1$, and for any $\tilde{v} \in L_2(\Omega)$, we have

$$\int_\Omega w_\eta(\tilde{v}) |\nabla \tilde{u}_\varepsilon(\tilde{v})|^2 \leq C \quad \text{and} \quad \|\tilde{u}_\varepsilon(\tilde{v})\|_{L^\infty(\Omega)} \leq C.$$

Furthermore, there exists $q(\varepsilon) > 2$, depending on N , Ω , s and ε only, such that $\tilde{u}_\varepsilon(\tilde{v})$ belongs to $W^{1,q(\varepsilon)}(\Omega)$. We can also find a constant C_1 , depending on N , Ω , γ , s , $\|f_0\|_{L^s(\partial\Omega)}$ and ε only, such that for any $\tilde{v} \in L_2(\Omega)$

$$\|\nabla \tilde{u}_\varepsilon(\tilde{v})\|_{L^{q(\varepsilon)}(\Omega)} \leq C_1.$$

We remark that the dependence of $q(\varepsilon)$ on ε is through $o_{\eta(\varepsilon)}$ and that, unfortunately, it might happen that $q(\varepsilon) \rightarrow 2^+$ and $C_1 \rightarrow +\infty$ as $\varepsilon \rightarrow 0^+$.

However, we may conclude that $\mathcal{H}_\varepsilon : L_2(\Omega) \rightarrow W_\gamma^{1,q(\varepsilon)}(\Omega)$ and its image is bounded in $W_\gamma^{1,q(\varepsilon)}(\Omega)$. Furthermore, again by Proposition 2.5, we infer that for any $\tilde{v}_0 \in L_2(\Omega)$ we may define as before $D\mathcal{H}_\varepsilon(\tilde{v}_0)$ and prove that $D\mathcal{H}_\varepsilon(\tilde{v}_0) : L_{q(\varepsilon)(q(\varepsilon)+2)/(q(\varepsilon)-2)}(\Omega) \rightarrow W_\gamma^{1,(q(\varepsilon)+2)/2}(\Omega)$ is a bounded linear operator.

Let $p(\varepsilon) = q(\varepsilon) \frac{q(\varepsilon)+2}{q(\varepsilon)-2}$. Then, straightforward but lengthy computations allow us to show that for any $\tilde{v}_0 \in L_2(\Omega)$, \mathcal{H}_ε is differentiable in \tilde{v}_0 with respect to the $L^{p(\varepsilon)}(\Omega)$ and the $W_\gamma^{1,2}(\Omega)$ norms. The differential is still given by (6.1). We immediately infer that for any p , $p \geq p(\varepsilon)$, $\tilde{\mathcal{G}}_\varepsilon$ is differentiable in \tilde{v}_0 , for any $\tilde{v}_0 \in W(\Omega)$, with respect to the $W_p(\Omega)$ norm, with the differential given by (6.4).

By an interpolation inequality, we may find $q_1(\varepsilon)$, $2 < q_1(\varepsilon) < (q(\varepsilon) + 2)/2$, depending on $q(\varepsilon)$ and α only, such that for any q , $2 \leq q \leq q_1(\varepsilon)$, and any

p , $p \geq p(\varepsilon)$, we have that, for any $\tilde{v}_0 \in L_2(\Omega)$, \mathcal{H}_ε is differentiable in \tilde{v}_0 with respect to the $L^p(\Omega)$ and the $W_\gamma^{1,q}(\Omega)$ norms. Obviously the differential is still given by (6.1).

We conclude that for such q and p , and any $\tilde{v} \in W(\Omega)$, we have that $\hat{\mathcal{F}}_\varepsilon^q$ is differentiable in \tilde{v}_0 , for any $\tilde{v}_0 \in W(\Omega)$, with respect to the $W_p(\Omega)$ norm. Its differential is given by the following formula. For any $\tilde{v} \in W_p(\Omega)$ we have

$$(6.5) \quad D\hat{\mathcal{F}}_\varepsilon^q(\tilde{v}_0)[\tilde{v}] = \frac{2a_2}{\varepsilon^{\tilde{\beta}}} \int_\gamma (\mathcal{H}_\varepsilon(\tilde{v}_0) - g_\varepsilon) U_\varepsilon(\tilde{v}_0, \tilde{v}) + \\ b \int_\Omega (q\psi_\eta(v_0) |\nabla \mathcal{H}_\varepsilon(\tilde{v}_0)|^{q-2} \nabla \mathcal{H}_\varepsilon(\tilde{v}_0) \cdot \nabla U_\varepsilon(\tilde{v}_0, \tilde{v}) - \psi'_\eta(v_0) |\nabla \mathcal{H}_\varepsilon(\tilde{v}_0)|^q \tilde{v}) + \\ \frac{1}{\eta} \int_\Omega (-V'(v_0)\tilde{v}) + 2\eta \int_\Omega \nabla \tilde{v}_0 \cdot \nabla \tilde{v}.$$

An important final remark is the following. If $N = 2$, then we may actually choose $p(\varepsilon) = 2$, and we observe that $W_2(\Omega)$ is a Hilbert space, with the scalar product $\int_\Omega \nabla \tilde{v}_1 \cdot \nabla \tilde{v}_2$ for any $\tilde{v}_1, \tilde{v}_2 \in W_2(\Omega)$. If $N > 2$, then it might happen that $p(\varepsilon) > 2$ and therefore $W_{p(\varepsilon)}(\Omega)$ has not a Hilbert space structure anymore. However, since $p(\varepsilon)$ is finite, $W_{p(\varepsilon)}(\Omega)$ is still a strictly convex real reflexive Banach space.

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