Hölder Continuity and Optimal Control for Nonsmooth Elliptic Problems

R. Haller-Dintelmann · C. Meyer · J. Rehberg · A. Schiela

Published online: 30 May 2009 © Springer Science+Business Media, LLC 2009

Abstract The well known De Giorgi result on Hölder continuity for solutions of the Dirichlet problem is re-established for mixed boundary value problems, provided that the underlying domain is a Lipschitz domain and the border between the Dirichlet and the Neumann boundary part satisfies a very general geometric condition. Implications of this result for optimal control theory are presented.

Keywords Elliptic problems \cdot Mixed boundary value problems \cdot Hölder continuity \cdot Optimal control

1 Introduction

In the last decades it has been anticipated in applied analysis that many elliptic problems originating from science, engineering, and technology possess nonsmooth data. This means that they often live on nonsmooth domains, the coefficients are nonsmooth and, thirdly, they often exhibit mixed boundary conditions, see [3, 43] and the references cited therein, see also [21, 46]. In this paper we prove Hölder continuity for the solution u of

R. Haller-Dintelmann

C. Meyer (🖂) · J. Rehberg

J. Rehberg e-mail: rehberg@wias-berlin.de

A. Schiela Konrad-Zuse-Zentrum für Informationstechnik Berlin, Takustr. 7, 14195 Berlin, Germany e-mail: schiela@zib.de

Technische Universität Darmstadt, Schlossgartenstr. 7, 64289 Darmstadt, Germany e-mail: haller@mathematik.tu-darmstadt.de

Weierstrass Institute for Applied Analysis and Stochastics, Mohrenstr. 39, 10117 Berlin, Germany e-mail: cmeyer@gsc.tu-darmstadt.de

$$-\nabla \cdot \rho \nabla u + u = f + \sum_{j=1}^{d} \frac{\partial f_j}{\partial x_j}, \quad f \in L^{q/2}, f_j \in L^q,$$
(1.1)

where q is larger than the space dimension, and mixed boundary conditions are incorporated, see Theorem 3.3. The result is to be seen in the tradition of Stampacchia's paper [51] (see also [40, 41]), where Hölder continuity already was achieved for mixed boundary value problems, but under rather technical conditions—difficult to verify in applications. Here, we generalize the Stampacchia result in space dimensions d = 2, 3 and 4 to Lipschitz domains, provided that the Dirichlet boundary part satisfies a very general compatibility condition—purely topological in nature and easy to check at least for d = 2, 3, see Theorems 5.2 and 5.4 below.

Note that the admissible distributional right-hand sides in (1.1) allow for jumps in the conormal derivative of solutions across internal interfaces. This means, e.g. in electrostatics, that the jump in the normal component of the displacement $v_+ \cdot \varepsilon \nabla \varphi - v_- \cdot \varepsilon \nabla \varphi$ across a prescribed interface equals the surface charge density on the interface, and this surface charge density is represented by a distribution on the underlying domain Ω .

Divergence type operators as in (1.1) are of fundamental significance in many application areas. This is the case not only in mechanics (see [38, Chaps. IV, V]), thermodynamics (see [50]) and electrodynamics (see [49]) of heterogeneous media, but also in mining, multiphase flow, mathematical biology (see [6, 20]) and semiconductor device simulation (see [21, 24, 46]), in particular quantum electronics (see [5, 36, 39, 56, 57]).

The non-homogeneous coefficient function ρ represents varying material properties as the context requires. It may be thermal conductivity in a heat equation (see [50, Sect. 21]) or dielectric permittivity in a Poisson equation, or diffusivity in a transport equation (see for instance [46, Sect. 2.2] for carrier continuity equations) or effective electron mass in a Schrödinger equation (see [36]).

Continuity of solutions to (1.1) plays an important role for the discussion of stateconstrained optimal control problems (see for instance [8]). Hence it is rather natural to use the above mentioned result for the discussion of semilinear elliptic control problems with pointwise inequality constraints on the state, which is done here in a very general setting. Such problems have been discussed by numerous authors before (see for instance [1, 8, 11] and the references therein). Concerning second-order sufficient optimality conditions, some progress has recently been made in a contribution of Casas et al. [12]. Here we show that, based on the regularity results of Theorem 3.3, the analysis, developed in [12], is also applicable to problems with mixed boundary conditions for semilinear elliptic state-constrained optimal control problems represents the genuine contribution of this paper from the viewpoint of optimal control theory.

The outline of the paper is as follows: first we introduce some notation. In Sect. 3 we formulate our regularity result, which is proved in Sect. 4. In Sect. 5 we give an alternative characterization for Gröger's regular sets, which represent the geometric setting for the domains under consideration and the associated Dirichlet boundary parts, in the 2d and 3d case. Finally, the relevance of the Hölder property for the

discussion of semilinear elliptic optimal control problems with pointwise state constraints is pointed out in Sect. 6.

2 Notation

Throughout this paper, $\Omega \subset \mathbb{R}^d$ always denotes a bounded Lipschitz domain (see [29, Chap. 1.2] for the definition) and $\Gamma \subset \partial \Omega$ is an open part of its boundary. In particular, we often use the cube $K := \{x \in \mathbb{R}^d : -1 < x_j < 1 \text{ for } 1 \le j \le d\}$, the half cube $K_- := \{x \in K : x_d < 0\}$, its upper plate $\Sigma := \{x \in K : x_d = 0\}$ and, lastly, the half of this, $\Sigma_0 := \{x \in \Sigma : x_{d-1} < 0\}$. The symbol $C^{\alpha}(\Omega)$ stands for the usual Hölder space on Ω , see [35] or [53]. $W^{1,p}(\Omega)$ denotes the Sobolev space on Ω consisting of those $L^p(\Omega)$ functions whose first order distributional derivatives also belong to $L^p(\Omega)$ (see [29] or [42]). We use the symbol $W_{\Gamma}^{1,p}(\Omega)$ for the closure of

$$\left\{ v|_{\Omega} : v \in C^{\infty}(\mathbb{R}^d), \text{ supp } v \cap (\partial \Omega \setminus \Gamma) = \emptyset \right\}$$
(2.1)

in $W^{1,p}(\Omega)$. Note that Ω enjoys the extension property for $W^{1,p}(\Omega)$ in view of being a bounded Lipschitz domain, see [23, Theorem 3.10] or [42, Chap. 1.1.16]. Thus, in case of $\Gamma = \partial \Omega$ the space $W_{\Gamma}^{1,p}(\Omega)$ is identical with the usual Sobolev space $W^{1,p}(\Omega)$. If $\Gamma = \emptyset$ we write as usual $W_0^{1,p}(\Omega)$ instead of $W_{\emptyset}^{1,p}(\Omega)$. $W_{\Gamma}^{-1,p'}(\Omega)$ denotes the dual to $W_{\Gamma}^{1,p}(\Omega)$ and $W^{-1,p'}(\Omega)$ denotes the dual to $W_0^{1,p}(\Omega)$, when $\frac{1}{p} + \frac{1}{p'} = 1$ holds. If Ω is understood, then we sometimes abbreviate $W_{\Gamma}^{\pm 1,p}$, $W_0^{1,p}$ and $W^{-1,p}$, respectively. Please notice that all functional spaces under consideration are regarded as complex ones. By $\langle \cdot, \cdot \rangle_X$ we indicate the duality between a Banach space X and its dual. Finally, γ denotes a generic constant not always of the same numerical value.

3 The Regularity Result

Definition 3.1 Let $\Lambda \subset \mathbb{R}^d$ be a bounded domain and Υ a (relatively) open part of its boundary $\partial \Lambda$. Then we call $\Lambda \cup \Upsilon$ *regular* (in the sense of Gröger [30]), if for every $x \in \partial \Lambda$ there are two open sets $\mathcal{U}_x, \mathcal{V}_x \subset \mathbb{R}^d$ and a bi-Lipschitz transform Ψ_x from \mathcal{U}_x onto \mathcal{V}_x , such that $x \in \mathcal{U}_x, \Psi_x(x) = 0$ and $\Psi_x(\mathcal{U}_x \cap (\Lambda \cup \Upsilon))$ either coincides with K_- or with $K_- \cup \Sigma$ or with $K_- \cup \Sigma_0$.

Assumption 3.2 Let ρ be a Lebesgue measurable, essentially bounded function on Ω , taking its values in the set of real $d \times d$ matrices, that additionally satisfies the usual (strong) ellipticity condition

$$\mathbf{y} \cdot \boldsymbol{\rho}(\mathbf{x}) \, \mathbf{y} \ge \iota |\mathbf{y}|^2, \quad \mathbf{y} \in \mathbb{R}^d,$$
(3.1)

for almost all $x \in \Omega$ and some $\iota > 0$.

Given a coefficient function ρ , satisfying this assumption, we define the operator $-\nabla \cdot \rho \nabla + 1: W_{\Gamma}^{1,2}(\Omega) \to W_{\Gamma}^{-1,2}(\Omega)$ by

$$\langle -\nabla \cdot \rho \nabla v + v, w \rangle_{W_{\Gamma}^{-1,2}} := \int_{\Omega} (\rho \nabla v \cdot \nabla w + vw) \, d\mathbf{x} \,, \quad v, w \in W_{\Gamma}^{1,2}(\Omega). \tag{3.2}$$

Then our first main result reads as follows.

Theorem 3.3 Suppose $2 \le d \le 4$ and q > d. Suppose further that $\Omega \cup \Gamma$ is regular and that Assumption 3.2 is satisfied. Then there is an $\alpha > 0$, such that $(-\nabla \cdot \rho \nabla + 1)^{-1} \max W_{\Gamma}^{-1,q}(\Omega)$ continuously into $C^{\alpha}(\Omega)$.

Remark 3.4

- (a) The cases $\Gamma = \emptyset$ (Dirichlet boundary condition) and $\Gamma = \partial \Omega$ (Neumann boundary condition) are explicitly allowed.
- (b) It is not hard to see that the right-hand side of (1.1) defines an element of $W_{\Gamma}^{-1,q}(\Omega)$. Conversely, any element of $W_{\Gamma}^{-1,q}(\Omega)$ may be represented this way, see [58, Chap. 4.3].

Corollary 3.5 The result of Theorem 3.3 carries over to problems with Robin boundary conditions on Γ , if the function \varkappa , defining the boundary term in the bilinear form (see (4.18) below) is from $L^{\infty}(\Gamma, \sigma)$ (σ being the induced boundary measure on Γ , cf. [32, Sect. 3]).

Corollary 3.6 Let a non-negative function $V \in L^{\infty}(\Omega)$ be given. Moreover, assume that, if $\operatorname{meas}(\partial \Omega \setminus \Gamma) = 0$, then there is a subset Ω_+ of Ω of positive measure, where V is strictly positive. Then, similarly to Corollary 3.5, Theorem 3.3 also applies to problems of the form $-\nabla \cdot \rho \nabla v + Vv = f$, $f \in W_{\Gamma}^{-1,q}(\Omega)$.

Corollary 3.7 Let D_q denote the domain of the maximal restriction of $-\nabla \cdot \rho \nabla + 1$ to the space $W_{\Gamma}^{-1,q}(\Omega)$. Then, under the suppositions of Theorem 3.3, even the complex interpolation space $[D_q, W_{\Gamma}^{-1,q}]_{\tau}$ continuously embeds into a Hölder space $C^{\beta}(\Omega)$, if τ and β are sufficiently close to 0.

Remark 3.8 Corollary 3.7 may be of use for the treatment of parabolic equations, see [2, 45, 48].

4 Proof of the Regularity Result

Let us start by commenting on the philosophy of the proof: the problem will be localized by means of a suitably chosen partition of unity, afterwards transformed by bi-Lipschitz mappings and, if necessary, by reflection. In any case one ends up with a Dirichlet problem on either a ball, the half cube K_{-} or the cube K. Then a well known regularity result (see Proposition 4.4) may be applied.

In order to perform this procedure we first quote two results from the literature and afterwards establish some auxiliary results, which will justify the required technical steps.

4.1 Known Results

Proposition 4.1 ([31], see also [30]) Suppose that ρ satisfies Assumption 3.2. If $\Omega \cup \Gamma$ is regular, then there is a $q_0 > 0$ such that for all $q \in [2 - q_0, 2 + q_0]$ the operator $-\nabla \cdot \rho \nabla + 1$ provides a topological isomorphism between $W_{\Gamma}^{1,q}(\Omega)$ and $W_{\Gamma}^{-1,q}(\Omega)$.

Remark 4.2 In the quoted papers the assertion of Proposition 4.1 is shown for q from an interval $[2, 2+q_0]$. The case $q \in [2-q_0, 2]$ is achieved by considering the adjoint operator.

Remark 4.3 It is clear by Sobolev embedding that in the two dimensional case the assertion of Theorem 3.3 already follows from Proposition 4.1.

Proposition 4.4 (See [37, Chap. III.14], [35, Theorem C.2], [13, Chap. 4], see also [15]) Let Λ be a ball or a cuboid and let the coefficient function ω (mutatis mutandis) satisfy Assumption 3.2. If q > d, then there is an $\alpha > 0$, such that

$$(-\nabla \cdot \omega \nabla)^{-1} : W^{-1,q}(\Lambda) \to C^{\alpha}(\Lambda)$$
 (4.1)

is continuous.

Remark 4.5 Usually, Proposition 4.4 is proved only for real spaces, but it is straightforward to extend this to the complex case: one considers for any element T of the complex Sobolev space $W_{\Gamma}^{-1,q}$ the linear forms $\frac{T+T^*}{2}$ and $\frac{T-T^*}{2i}$, where T^* is defined by $T^*\psi := \overline{T\psi}$. Obviously, both take real values when applied to real functions and satisfy $\frac{T+T^*}{2} + i\frac{T-T^*}{2i} = T$. Thus, one may use the result for the corresponding real spaces.

4.2 Auxiliary Results

Lemma 4.6 Let $\Omega \cup \Gamma$ be regular and let $\mathcal{U} \subset \mathbb{R}^d$ be open, such that $\Omega_{\bullet} := \Omega \cap \mathcal{U}$ is also a Lipschitz domain. Furthermore, we put $\Gamma_{\bullet} := \Gamma \cap \mathcal{U}$ and fix an arbitrary function $\eta \in C_0^{\infty}(\mathbb{R}^d)$ with supp $\eta \subset \mathcal{U}$. Then for any $q \in [1, \infty[$ we have the following assertions:

- (i) If $v \in W_{\Gamma}^{1,q}(\Omega)$, then $\eta v|_{\Omega_{\bullet}} \in W_{\Gamma_{\bullet}}^{1,q}(\Omega_{\bullet})$.
- (ii) Let for any $v \in L^1(\Omega_{\bullet})$ the symbol \tilde{v} indicate the extension of v to Ω by zero. Then the mapping

$$W^{1,q}_{\Gamma_{\bullet}}(\Omega_{\bullet}) \ni v \mapsto \widetilde{\eta v}$$

has its image in $W^{1,q}_{\Gamma}(\Omega)$ and is continuous.

Proof For the proof of both points we will employ the following well known set inclusion (cf. [16, Chap. 3.8]):

$$(\partial \Omega \cap \mathcal{U}) \cup (\Omega \cap \partial \mathcal{U}) \subset \partial \Omega_{\bullet} \subset (\partial \Omega \cap \mathcal{U}) \cup (\Omega \cap \partial \mathcal{U}).$$
(4.2)

(i) First one observes that the multiplication with η and the restriction is a continuous mapping from $W_{\Gamma}^{1,q}(\Omega)$ into $W^{1,q}(\Omega_{\bullet})$. Thus, it suffices to show the assertion only for elements of the dense subset

$$\{v|_{\Omega}: v \in C^{\infty}(\mathbb{R}^d), \text{ supp } v \cap (\partial \Omega \setminus \Gamma) = \emptyset\},\$$

what we will do now. One has by (4.2)

$$\operatorname{supp}(\eta v) \cap (\partial \Omega_{\bullet} \setminus \Gamma_{\bullet}) \subset \operatorname{supp} \eta \cap \operatorname{supp} v \cap \left[\left((\partial \Omega \cap \mathcal{U}) \cup (\overline{\Omega} \cap \partial \mathcal{U}) \right) \setminus \left(\Gamma \cap \mathcal{U} \right) \right].$$

Since
$$(\overline{\Omega} \cap \partial \mathcal{U}) \cap (\Gamma \cap \mathcal{U}) = \emptyset$$
, we see

$$\left((\partial \Omega \cap \mathcal{U}) \cup (\overline{\Omega} \cap \partial \mathcal{U}) \right) \setminus \left(\Gamma \cap \mathcal{U} \right) = \left((\partial \Omega \cap \mathcal{U}) \setminus (\Gamma \cap \mathcal{U}) \right) \cup \left((\overline{\Omega} \cap \partial \mathcal{U}) \setminus (\Gamma \cap \mathcal{U}) \right)$$
$$= \left((\partial \Omega \setminus \Gamma) \cap \mathcal{U} \right) \cup (\overline{\Omega} \cap \partial \mathcal{U}).$$

This, together with supp $\eta \subset \mathcal{U}$ yields

 $\operatorname{supp}(\eta v) \cap (\partial \Omega_{\bullet} \setminus \Gamma_{\bullet}) \subset \operatorname{supp} \eta \cap \operatorname{supp} v \cap ((\partial \Omega \setminus \Gamma) \cap \mathcal{U}) = \emptyset.$

(ii) Let $v \in C_0^{\infty}(\mathbb{R}^d)$ with supp $v \cap (\partial \Omega_{\bullet} \setminus \Gamma_{\bullet}) = \emptyset$. Since by the right-hand side of (4.2) we have

$$\partial \Omega_{\bullet} \setminus \Gamma_{\bullet} \supseteq (\partial \Omega \cap \mathcal{U}) \setminus \Gamma_{\bullet} = \mathcal{U} \cap (\partial \Omega \setminus \Gamma),$$

it follows supp $v \cap (\mathcal{U} \cap (\partial \Omega \setminus \Gamma)) = \emptyset$. Combining this with supp $\eta \subset \mathcal{U}$, we obtain

$$\operatorname{supp}(\eta v) \cap (\partial \Omega \setminus \Gamma) = \operatorname{supp}(\eta v) \cap (\mathcal{U} \cap (\partial \Omega \setminus \Gamma) = \emptyset,$$

so $\eta v|_{\Omega} \in W_{\Gamma}^{1,q}(\Omega)$. Furthermore, it is not hard to see that $\|\eta v\|_{W^{1,q}(\Omega)} \leq \gamma_{\eta} \|v\|_{W^{1,q}(\Omega_{\bullet})}$, where the constant γ_{η} is independent from v. Thus, the assertion follows, since $\{v|_{\Omega_{\bullet}} : v \in C_{0}^{\infty}(\mathbb{R}^{d}), \operatorname{supp}(v) \cap (\partial \Omega_{\bullet} \setminus \Gamma_{\bullet}) = \emptyset$ is dense in $W_{\Gamma_{\bullet}}^{1,q}(\Omega_{\bullet})$ and $W_{\Gamma}^{1,q}(\Omega)$ is closed in $W^{1,q}(\Omega)$.

Lemma 4.7 Let Ω , Γ , \mathcal{U} , η , Ω_{\bullet} and Γ_{\bullet} be as in the foregoing lemma. Denote by ρ_{\bullet} the restriction of the coefficient function ρ to Ω_{\bullet} and let the operator $-\nabla \cdot \rho_{\bullet} \nabla$: $W_{\Gamma_{\bullet}}^{1,2}(\Omega_{\bullet}) \rightarrow W_{\Gamma_{\bullet}}^{-1,2}(\Omega_{\bullet})$ be defined analogously to (3.2). Assume $v \in W_{\Gamma}^{1,2}(\Omega)$ to be the solution of

$$-\nabla \cdot \rho \nabla v + v = f \in W_{\Gamma}^{-1,2}(\Omega).$$
(4.3)

Then the following holds true:

(i) For all $q \in [1, \infty[$ the linear form

$$f_{\bullet}: w \mapsto \langle f, \widetilde{\eta w} \rangle_{W_{\mathfrak{p}}^{-1,q}(\Omega)},$$

where $\widetilde{\eta w}$ again means the extension by zero to the whole Ω , is well defined and continuous on $W_{\Gamma_{\bullet}}^{1,q'}(\Omega_{\bullet})$, whenever $f \in W_{\Gamma}^{-1,q}(\Omega)$.

(ii) If we denote the linear form

$$W^{1,2}_{\Gamma_{\bullet}}(\Omega_{\bullet}) \ni w \mapsto \int_{\Omega_{\bullet}} v \rho_{\bullet} \nabla \eta \cdot \nabla w \, d\mathbf{x},$$

by T_v , then $u := \eta v|_{\Omega_{\bullet}}$ satisfies

$$-\nabla \cdot \rho_{\bullet} \nabla u = -\eta v|_{\Omega_{\bullet}} - \rho_{\bullet} \nabla v|_{\Omega_{\bullet}} \cdot \nabla \eta|_{\Omega_{\bullet}} + T_v + f_{\bullet} =: f^{\bullet}.$$
(4.4)

(iii) Assume now $2 \le d \le 4$. If $f \in W_{\Gamma}^{-1,q}(\Omega)$ for $a \ q > d$, then there is $a \ p > d$ such that $f^{\bullet} \in W_{\Gamma_{\bullet}}^{-1,p}(\Omega_{\bullet})$. Moreover, the mapping $W_{\Gamma}^{-1,q}(\Omega) \ni f \mapsto f^{\bullet} \in W_{\Gamma_{\bullet}}^{-1,p}(\Omega_{\bullet})$ is continuous.

Proof (i) The mapping $f \mapsto f_{\bullet}$ is the adjoint to $v \mapsto \tilde{\eta v}$, which maps by the preceding lemma $W_{\Gamma_{\bullet}}^{1,q'}(\Omega_{\bullet})$ continuously into $W_{\Gamma}^{1,q'}(\Omega)$.

(ii) For every $w \in W^{1,2}_{\Gamma_{\bullet}}(\Omega_{\bullet})$ we have

$$\begin{split} \langle -\nabla \cdot \rho_{\bullet} \nabla u, w \rangle_{W_{\Gamma_{\bullet}}^{-1,2}(\Omega_{\bullet})} &+ \int_{\Omega_{\bullet}} \eta v w \, dx \\ &= \int_{\Omega_{\bullet}} \rho_{\bullet} \nabla (\eta v) \cdot \nabla w \, dx + \int_{\Omega_{\bullet}} \eta v w \, dx \\ &= \int_{\Omega_{\bullet}} v \rho_{\bullet} \nabla \eta \cdot \nabla w \, dx + \int_{\Omega_{\bullet}} \eta \rho_{\bullet} \nabla v \cdot \nabla w \, dx + \int_{\Omega} v \, \tilde{\eta w} \, dx \\ &= \int_{\Omega_{\bullet}} v \rho_{\bullet} \nabla \eta \cdot \nabla w \, dx + \int_{\Omega_{\bullet}} \rho_{\bullet} \nabla v \cdot \nabla (\eta w) \, dx - \int_{\Omega_{\bullet}} w \rho_{\bullet} \nabla v \cdot \nabla \eta \, dx \\ &+ \int_{\Omega} v \, \tilde{\eta w} \, dx \end{split}$$

and by (4.3) we see

$$\begin{split} \int_{\Omega_{\bullet}} \rho_{\bullet} \nabla v \cdot \nabla(\eta w) \, d\mathbf{x} + \int_{\Omega} v \, \widetilde{\eta w} \, d\mathbf{x} &= \int_{\Omega} \rho \nabla v \cdot \nabla(\widetilde{\eta w}) \, d\mathbf{x} + \int_{\Omega} v \, \widetilde{\eta w} \, d\mathbf{x} \\ &= \langle f, \widetilde{\eta w} \rangle_{W_{\Gamma}^{-1,2}(\Omega)}. \end{split}$$

Applying the definition of T_v and f_{\bullet} and afterwards subtracting $\int_{\Omega_{\bullet}} \eta v w \, dx$ from both sides yields the assertion.

(iii) We regard the terms in (4.4) from left to right. For the first summand the assertion is obvious.

According to Proposition 4.1 there is an $\epsilon > 0$ such that the solution v of $-\nabla \cdot \rho \nabla v + v = f$ is from $W_{\Gamma}^{1,2+\epsilon}(\Omega)$, what implies $\nabla v \in L^{2+\epsilon}(\Omega)$. Furthermore, $|\nabla \eta| \in L^{\infty}(\Omega_{\bullet})$ and $|\rho_{\bullet}\nabla v| \in L^{2+\epsilon}(\Omega_{\bullet})$. Consequently we have $\rho_{\bullet}\nabla v \cdot \nabla \eta \in L^{2+\epsilon}(\Omega_{\bullet})$. Now, whenever $1/p \ge (d-2-\epsilon)/(d(2+\epsilon))$ we have the embedding $L^{2+\epsilon}(\Omega_{\bullet}) \hookrightarrow W_{\Gamma_{\bullet}}^{-1,p}(\Omega_{\bullet})$. Since we restricted the dimension to $2 \le d \le 4$, there is always a $p > 4 \ge d$ satisfying that condition, so the second term is also fine. On the other hand, we have $v \in W^{1,2+\epsilon}_{\Gamma}(\Omega) \hookrightarrow L^{4+\delta}(\Omega)$ for a $\delta = \delta(d) > 0$. Thus, concerning T_v , we can estimate

$$|\langle T_{v}, w \rangle_{W_{\Gamma_{\bullet}}^{-1,4+\delta}(\Omega_{\bullet})}| \leq \|v\|_{L^{4+\delta}(\Omega_{\bullet})} \|\rho\|_{L^{\infty}(\Omega;\mathbb{C}^{d\times d})} \|\nabla\eta\|_{L^{\infty}(\Omega_{\bullet})} \|w\|_{W_{\Gamma_{\bullet}}^{1,(4+\delta)'}(\Omega_{\bullet})}$$

The claim on f_{\bullet} follows from (i), while the proof of the last assertion is implicitly contained in the above considerations.

Remark 4.8 It is the lack of integrability of the gradient of v (see the counterexample in [18, Chap. 4]) together with the quality of the needed Sobolev embeddings, which prevents the applicability of this localization procedure to higher dimensions and thus limits our central result to the dimensions up to 4.

The reader may wonder why we start with the operator $-\nabla \cdot \rho \nabla + 1$ and consider the operator $-\nabla \cdot \rho_{\bullet} \nabla$ after the localization. The reason for this is the following: in order to include the pure Neumann case one should consider the operator $-\nabla \cdot \rho \nabla + 1$. On the other hand, the resulting localized operators $-\nabla \cdot \rho_{\bullet} \nabla$ exhibit in any case a nontrivial Dirichlet boundary part (making them invertible) and the subsequent transformation techniques are technically simpler for the pure operator $-\nabla \cdot \rho_{\bullet} \nabla$.

Proposition 4.9 Let $\Lambda \subset \mathbb{R}^d$ be a bounded Lipschitz domain and Υ be an open subset of its boundary. Assume that ϕ is a mapping from a neighborhood of Λ into \mathbb{R}^d that is bi-Lipschitz. Let us denote $\phi(\Lambda) =: \Lambda_{\Delta}$ and $\phi(\Upsilon) =: \Upsilon_{\Delta}$. Then the following is true:

(i) For any $p \in [1, \infty[$, the mapping ϕ induces a linear, topological isomorphism

$$\Phi_p: W^{1,p}_{\Upsilon_{\Delta}}(\Lambda_{\Delta}) \to W^{1,p}_{\Upsilon}(\Lambda),$$

which is given by $(\Phi_p f)(\mathbf{x}) = f(\phi(\mathbf{x})) = (f \circ \phi)(\mathbf{x}).$

- (ii) $\Phi_{p'}^*$ is a linear, topological isomorphism between $W_{\Upsilon}^{-1,p}(\Lambda)$ and $W_{\Upsilon_{\Lambda}}^{-1,p}(\Lambda_{\Delta})$.
- (iii) If ω is a bounded, measurable function on Λ , taking its values in the set of $d \times d$ matrices, then

$$\Phi_{p'}^* \nabla \cdot \omega \nabla \Phi_p = \nabla \cdot \omega_\Delta \nabla \tag{4.5}$$

with

$$\omega_{\Delta}(\mathbf{y}) := \frac{1}{|\det(D\phi)(\phi^{-1}(\mathbf{y}))|} (D\phi)(\phi^{-1}(\mathbf{y})) \,\omega(\phi^{-1}(\mathbf{y})) \left(D\phi\right)^{T}(\phi^{-1}(\mathbf{y})) \tag{4.6}$$

for almost all $y \in \Lambda_{\Delta}$. Here, $D\phi$ denotes the Jacobian of ϕ and det $(D\phi)$ the corresponding determinant.

(iv) If ω satisfies Assumption 3.2, then ω_{Δ} also does.

Proof The proof of (i) is contained in [28, Theorem 2.10)]. Assertion (ii) follows from (i) by duality, while (iii) is well known, see [33] for an explicit verification or [4, Chap. 0.8]. Finally, (iv) is implied by (4.6) and the fact that for a bi-Lipschitz ϕ the Jacobian $D\phi$ and its inverse $(D\phi)^{-1}$ are essentially bounded (see [19, Chap. 3.1]).



Fig. 1 $K_- \cup \Sigma_0$ and $\chi_1(K_- \cup \Sigma_0)$

The next lemma makes clear that within the class of bi-Lipschitz transformations one only needs the two local model sets K_- and $K_- \cup \Sigma$, if one dispenses with the condition $\Psi_x(x) = 0$ (cf. Definition 3.1):

Lemma 4.10 There is a bi-Lipschitz mapping $\Psi : \mathbb{R}^d \to \mathbb{R}^d$ which maps $K_- \cup \Sigma_0$ onto $K_- \cup \Sigma$.

Proof Let us first consider the case d = 2. We define on the lower halfspace $\{(x, y) : y \le 0\}$

$$\chi_{1}(x, y) := \begin{cases} (x - y/2, y/2), & \text{if } x \le 0, y \ge x, \\ (x/2, -x/2 + y), & \text{if } x \le 0, y < x, \\ (x/2, x/2 + y), & \text{if } x > 0, y < -x, \\ (x + y/2, y/2), & \text{if } x > 0, y \ge -x. \end{cases}$$
(4.7)

Observing that χ_1 acts as the identity on the *x*-axis, we may define χ_1 on the upper half space $\{(x, y) : y > 0\}$ also as the identity and thus obtain a globally bi-Lipschitz transformation χ_1 from \mathbb{R}^2 onto itself that transforms $K_- \cup \Sigma_0$ onto the triangle shown in Fig. 1.

Next, we define the bi-Lipschitz mapping $\chi_2 : \mathbb{R}^2 \to \mathbb{R}^2$ by

$$\chi_2(x, y) := \begin{cases} (x, x+2y+1), & \text{if } x \le 0, \\ (x, -x+2y+1), & \text{if } x > 0, \end{cases}$$
(4.8)

in order to get the geometric arrangement in Fig. 2.

If ϑ is the (clockwise) rotation by $\pi/4$, we thus achieved that $\chi := \vartheta \chi_2 \chi_1 : \mathbb{R}^2 \to \mathbb{R}^2$ is bi-Lipschitz and satisfies

$$\chi(K_{-} \cup \Sigma_{0}) = \left\{ (x, y) : -\frac{1}{\sqrt{2}} < x < \frac{1}{\sqrt{2}}, -\frac{1}{\sqrt{2}} < y \le \frac{1}{\sqrt{2}} \right\}.$$

Let $\varsigma : \mathbb{R}^2 \to \mathbb{R}^2$ be the affine mapping $(x, y) \mapsto (\sqrt{2}x, \frac{1}{\sqrt{2}}y - \frac{1}{2})$. Then $\Psi_2 := \varsigma \chi$ maps $K_- \cup \Sigma_0$ bi-Lipschitzian onto $K_- \cup \Sigma$ in the 2*d* case.

If d > 2, one simply puts $\Psi(x_1, ..., x_d) := (x_1, ..., x_{d-2}, \Psi_2(x_{d-1}, x_d))$.

Springer

Fig. 2 $\chi_2(\chi_1(K_- \cup \Sigma_0))$



Proposition 4.11 Let for any $\mathbf{x} = (x_1, \ldots, x_d) \in \mathbb{R}^d$ the symbol \mathbf{x}_- denote the element $(x_1, \ldots, x_{d-1}, -x_d)$. Further, for a $d \times d$ matrix ϱ , we define the matrix ϱ^- by

$$\varrho_{j,k}^{-} := \begin{cases} \varrho_{j,k}, & \text{if } j, k < d, \\ -\varrho_{j,k}, & \text{if } j = d \text{ and } k \neq d \text{ or } k = d \text{ and } j \neq d, \\ \varrho_{j,k}, & \text{if } j = k = d. \end{cases}$$
(4.9)

Let ω be a bounded, measurable function on K_{-} taking its values in the set of real, symmetric $d \times d$ matrices. We define the matrix valued function $\hat{\omega}$ on K by

$$\hat{\omega}(\mathbf{x}) := \begin{cases} \omega(\mathbf{x}), & \text{if } \mathbf{x} \in K_{-}, \\ \left(\omega(\mathbf{x}_{-})\right)^{-}, & \text{if } \mathbf{x}_{-} \in K_{-}, \\ 0, & \text{if } \mathbf{x} \in \Sigma. \end{cases}$$
(4.10)

Then we have the following assertions for every $p \in]1, \infty[$:

(i) If $\psi \in W_{\Sigma}^{1,p}(K_{-})$ satisfies $-\nabla \cdot \omega \nabla \psi = f \in W_{\Sigma}^{-1,p}(K_{-})$, then $-\nabla \cdot \hat{\omega} \nabla \hat{\psi} = \hat{f} \in W^{-1,p}(K)$ holds for $\hat{\psi}$ with

$$\hat{\psi}(\mathbf{x}) = \begin{cases} \psi(\mathbf{x}), & \text{if } \mathbf{x} \in K_{-}, \\ \psi(\mathbf{x}_{-}), & \text{if } \mathbf{x}_{-} \in K_{-}, \end{cases}$$

and \hat{f} defined by $\langle \hat{f}, \varphi \rangle_{W^{-1,p}(K)} := \langle f, \varphi |_{K_{-}} + \varphi_{-} |_{K_{-}} \rangle_{W_{\Sigma}^{-1,p}(K_{-})}$. Here, the function φ_{-} is defined by $\varphi_{-}(\mathbf{x}) := \varphi(\mathbf{x}_{-})$.

(ii) The mapping $W_{\Sigma}^{-1,p}(K_{-}) \ni f \mapsto \hat{f} \in W^{-1,p}(K)$ is continuous.

Proof (i) It is known that $\hat{\psi}$ belongs to $W_0^{1,p}(K)$, see [23, Lemma 3.4]. Thus, it remains to show $-\nabla \cdot \hat{\omega} \nabla \hat{\psi} = \hat{f}$ as an equation in $W^{-1,p}(K)$. Since every test function from $W_0^{1,p}(K)$ may be split up into a symmetric and an anti-symmetric part, we

may look at these two cases separately. For every anti-symmetric $\varphi \in W_0^{1,p}(K)$. i.e. $\varphi(\mathbf{x}_-) = -\varphi(\mathbf{x})$, the symmetry of $\hat{\omega}, \hat{\psi}$ and \hat{f} yields immediately

$$\langle -\nabla \cdot \hat{\omega} \nabla \hat{\psi}, \varphi \rangle_{W_0^{-1,p}(K)} = 0 = \langle \hat{f}, \varphi \rangle_{W_0^{-1,p}(K)}$$

In the case of symmetric functions φ the assertion is obtained by the definitions of $\hat{\psi}, \hat{f}, -\nabla \cdot \omega \nabla, -\nabla \cdot \hat{\omega} \nabla$ and straightforward calculations, based on Proposition 4.9 when applied to the transformation $x \mapsto x_-$.

(ii) The operator $f \mapsto \hat{f}$ is the adjoint to $\varphi \mapsto (\varphi|_{K_{-}} + \varphi_{-}|_{K_{-}})$.

4.3 Core of the Proof

By Definition 3.1 and Lemma 4.10, for every $x \in \partial \Omega$ there is an open neighborhood \mathcal{U}_x of x, an open set \mathcal{W}_x and a bi-Lipschitz mapping Ψ_x from \mathcal{U}_x onto \mathcal{W}_x such that $\Psi_x((\Omega \cup \Gamma) \cap \mathcal{U}_x)$ equals either K_- or $K_- \cup \Sigma$. Take for every point $x \in \Omega$ a ball $\mathcal{B}_x \subset \Omega$ centered at x. Obviously, the system $\{\mathcal{U}_x\}_{x\in\partial\Omega} \cup \{\mathcal{B}_x\}_{x\in\Omega}$ forms an open covering of $\overline{\Omega}$. Take a finite subcovering $\mathcal{U}_{x_1}, \ldots, \mathcal{U}_{x_k}, \mathcal{B}_{x_1}, \ldots, \mathcal{B}_{x_l}$ and choose a partition of unity $\eta_1, \ldots, \eta_k, \zeta_1, \ldots, \zeta_l$ over $\overline{\Omega}$, which is subordinated to this subcovering.

Assume now $f \in W_{\Gamma}^{-1,q}(\Omega)$ with q > d and that v is the solution of $-\nabla \cdot \rho \nabla v + v = f$. Then, according to Proposition 4.7 with $\Gamma_{\bullet} = \emptyset$ and $\Omega_{\bullet} = B_{x_j}$, every function $\zeta_j v|_{\mathcal{B}_{x_j}}$ satisfies an equation $-\nabla \cdot \rho \nabla (\zeta_j v|_{\mathcal{B}_{x_j}}) = g_j$, where $g_j \in W^{-1,p}(\mathcal{B}_{x_j})$ with p > d and, additionally,

$$\|g_{j}\|_{W^{-1,p}(\mathcal{B}_{x_{j}})} \le \gamma \|f\|_{W^{-1,q}_{p}(\Omega)}$$
(4.11)

with γ independent from f. Hence, by Proposition 4.4, we have $\zeta_j v|_{\mathcal{B}_{x_j}} \in C^{\alpha}(\mathcal{B}_{x_j})$ for an $\alpha = \alpha(j) > 0$ and, moreover,

$$\|\zeta_{j}v\|_{\mathcal{B}_{x_{j}}}\|_{C^{\alpha}(\mathcal{B}_{x_{j}})} \leq \gamma \|g_{j}\|_{W^{-1,p}(\mathcal{B}_{x_{j}})}.$$
(4.12)

Clearly, (4.12) together with (4.11) implies

$$\|\zeta_{j}v\|_{C^{\alpha}(\Omega)} = \|\zeta_{j}v|_{\mathcal{B}_{x_{j}}}\|_{C^{\alpha}(\mathcal{B}_{x_{j}})} \le \gamma \|g_{j}\|_{W^{-1,p}(\mathcal{B}_{x_{j}})} \le \gamma \|f\|_{W^{-1,q}_{\Gamma}(\Omega)}.$$
 (4.13)

Let us now consider the functions $\eta_j v$ for fixed j: putting $\Omega_j := \Omega \cap \mathcal{U}_{x_j}$ and $\Gamma_j := \Gamma \cap \mathcal{U}_{x_j}$ we obtain by Lemma 4.6(i) that each $\eta_j v|_{\Omega_j}$ belongs to $W_{\Gamma_j}^{1,2}(\Omega_j)$. Furthermore, Proposition 4.7 shows that $\eta_j v|_{\Omega_j}$ satisfies an equation $-\nabla \cdot \rho \nabla(\eta_j v|_{\Omega_j}) = f_j$, where $f_j \in W_{\Gamma_j}^{-1,p}(\Omega_j)$ with p > d and, additionally,

$$\|f_{j}\|_{W_{\Gamma_{j}}^{-1,p}(\Omega_{j})} \le \gamma \|f\|_{W_{\Gamma}^{-1,q}(\Omega)}$$
(4.14)

with γ independent from f. Next we consider the 'transformed' function (cf. Proposition 4.9 with $\phi = \Psi_{x_j}^{-1}$) $\psi_j := \Phi_p(\eta_j v|_{\Omega_j}) = (\eta_j v|_{\Omega_j}) \circ \Psi_{x_j}^{-1}$ on K_- , from now on distinguishing the cases

$$\Psi_{\mathbf{x}_{j}}((\Omega \cup \Gamma) \cap \mathcal{U}_{\mathbf{x}_{j}}) = K_{-} \tag{4.15}$$

and

$$\Psi_{\mathbf{x}_i}\big((\Omega \cup \Gamma) \cap \mathcal{U}_{\mathbf{x}_i}\big) = K_- \cup \Sigma. \tag{4.16}$$

If (4.15) is true, then $\Gamma_j := \Gamma \cap \mathcal{U}_{\mathbf{x}_j} \subset \partial(\Omega \cap \mathcal{U}_{\mathbf{x}_j})$ must be empty, since K_- has only inner points and inner points pass to inner points and boundary points to boundary points under a bi-Lipschitz transformation. In particular, this means $\eta_j v|_{\Omega_j} \in W_0^{1,2}(\Omega_j)$. By Proposition 4.9 the function ψ_j belongs to $W_0^{1,2}(K_-)$ and satisfies an equation $-\nabla \cdot \omega \nabla \psi_j = h_j$ with $h_j = (\Phi_{p'}^*)^{-1} f_j \in W^{-1,p}(K_-)$ for the same p > d as above. Thanks to Proposition 4.9(iv), the coefficient function ω again satisfies Assumption 3.2. Thus, by Proposition 4.4, $\psi_j \in C^{\alpha}(K_-)$, where α depends on j, and

$$\|\psi_j\|_{C^{\alpha}(K_-)} \leq \gamma \|h_j\|_{W^{-1,p}(K_-)}$$

with α and γ independent from h_j . Transforming back, this gives $\eta_j v|_{\Omega_j} \in C^{\alpha}(\Omega_j)$ with

$$\|\eta_{j}v\|_{\Omega_{j}}\|_{C^{\alpha}(\Omega_{j})} \leq \gamma \|f_{j}\|_{W^{-1,p}(\Omega_{j})} \leq \gamma \|f\|_{W^{-1,q}(\Omega)},$$
(4.17)

where the last inequality is just (4.14). As the support of $\eta_j v$ has a positive distance to $\Omega \setminus \Omega_j$, the function $\eta_j v$ is from $C^{\alpha}(\Omega)$ with the norm equality $\|\eta_j v\|_{C^{\alpha}(\Omega)} =$ $\|\eta_j v\|_{\Omega_j}\|_{C^{\alpha}(\Omega_j)}$. This, together with (4.17) gives the desired estimate, where γ is independent from f.

Let us now consider the case (4.16). Analogously as before Proposition 4.9 yields that ψ_j belongs to $W_{\Sigma}^{1,2}(K_-)$ and satisfies an equation $-\nabla \cdot \omega \nabla \psi_j = h_j$ with $h_j \in W_{\Sigma}^{-1,p}(K_-)$ and p > d, where ω again satisfies Assumption 3.2. Now, we apply the reflection principle from Proposition 4.11. This leads to a homogeneous Dirichlet problem $-\nabla \cdot \hat{\omega} \nabla \hat{\psi}_j = \hat{h}_j$, where $\hat{h}_j \in W^{-1,p}(K)$. But then Proposition 4.4 gives $\hat{\psi}_j \in C^{\alpha}(K)$ and $\|\hat{\psi}_j\|_{C^{\alpha}(K)} \leq \gamma \|\hat{h}_j\|_{W^{-1,p}(K)}$ (α depending on j). Clearly, this, together with Proposition 4.11(ii), implies $\|\psi_j\|_{C^{\alpha}(K_-)} \leq \gamma \|h_j\|_{W_{\Sigma}^{-1,p}(K_-)}$ and, consequently,

$$\|\eta_j v|_{\Omega_j}\|_{C^{\alpha}(\Omega_j)} \leq \gamma \|f_j\|_{W^{-1,p}_{\Gamma_j}(\Omega_j)} \leq \gamma \|f\|_{W^{-1,q}_{\Gamma}(\Omega)}.$$

Thus, we get $\eta_j v \in C^{\alpha}(\Omega)$ and $\|\eta_j v\|_{C^{\alpha}(\Omega)} \leq \gamma \|f\|_{W_{\Gamma}^{-1,q}(\Omega)}$ with γ independent from *f* as in the previous case. Passing to the minimal $\alpha(j)$, this finishes the proof of Theorem 3.3.

Proof of Corollary 3.5 As is well known ([14, Chap. 1.2], [22, Chap. II.2]), in case of a Robin boundary condition the operator is defined via the bilinear form

$$W_{\Gamma}^{1,2}(\Omega) \times W_{\Gamma}^{1,2}(\Omega) \ni (v,w) \mapsto \int_{\Omega} \rho \nabla v \cdot \nabla w \, d\mathbf{x} + \int_{\Omega} v w \, d\mathbf{x} + \int_{\Gamma} \varkappa v \, w \, d\sigma$$
(4.18)

for some $\varkappa \in L^{\infty}(\Gamma, \sigma)$. It is easy to see that the linear mapping $T : L^{\infty}(\Gamma, \sigma) \to W_{\Gamma}^{-1, p}(\Omega)$, given by

$$\langle T\psi,\varphi\rangle_{W_{\Gamma}^{-1,p}(\Omega)} = \int_{\Gamma} \varkappa\psi\varphi\,d\sigma,$$

D Springer

is well defined and continuous. Denoting the domain of $-\nabla \cdot \rho \nabla + 1$ again by D_q , we have by Theorem 3.3 the compact embedding $D_q \hookrightarrow C^{\alpha}(\Omega) \hookrightarrow L^{\infty}(\Gamma, \sigma)$. Hence, the mapping *T* is relatively compact with respect to $-\nabla \cdot \rho \nabla + 1$ and a classical perturbation theorem [34, Chap. IV.1.3] applies.

Proof of Corollary 3.6 We argue analogously to the proof of Corollary 3.5. Here, $T: L^{\infty}(\Omega) \to W_{\Gamma}^{-1,q}(\Omega)$ is defined by

$$\langle T\psi,\varphi\rangle_{W_{\Gamma}^{-1,q}(\Omega)} = \int_{\Omega} V\psi\varphi\,d\mathbf{x}.$$

Then, due to the compact embedding $D_q \hookrightarrow L^{\infty}(\Omega)$, the same perturbation argument as above yields that the domain of the maximal restriction of $-\nabla \cdot \rho \nabla + V$ to $W_{\Gamma}^{-1,q}(\Omega)$ coincides with D_q . Together with the coercivity of the associated bilinear form that follows from the assumptions on V stated in Corollary 3.6, this gives the assertion.

Proof of Corollary 3.7 Applying a well known re-iteration result on complex interpolation (see [53, Chap. 1.9.3]), we obtain for $\tau \in [0, 1[$

$$[D_q, W_{\Gamma}^{-1,q}(\Omega)]_{\frac{\tau}{2}} \hookrightarrow \left[D_q, [D_q, W_{\Gamma}^{-1,q}(\Omega)]_{\frac{1}{2}}\right]_{\tau}.$$
(4.19)

But the embedding $W_{\Gamma}^{-1,q}(\Omega) \hookrightarrow W_{\Gamma}^{-1,2}(\Omega)$ gives $D_q \hookrightarrow W_{\Gamma}^{1,2}(\Omega)$. Together with $D_q \hookrightarrow C^{\alpha}(\Omega)$, we obtain by (4.19)

$$[D_q, W_{\Gamma}^{-1,q}(\Omega)]_{\frac{\tau}{2}} \hookrightarrow \left[C^{\alpha}(\Omega), [W_{\Gamma}^{1,2}(\Omega), W_{\Gamma}^{-1,2}(\Omega)]_{\frac{1}{2}}\right]_{\tau}.$$

One identifies the interpolation space $[W_{\Gamma}^{1,2}(\Omega), W_{\Gamma}^{-1,2}(\Omega)]_{\frac{1}{2}}$ as the space $L^{2}(\Omega)$ (see [53, Chap. 1.18.10]), what gives

$$[D_q, W_{\Gamma}^{-1,q}(\Omega)]_{\frac{\tau}{2}} \hookrightarrow [C^{\alpha}(\Omega), L^2(\Omega)]_{\tau}.$$
(4.20)

This latter interpolation space is known to embed into another Hölder space $C^{\beta}(\Omega)$, if $\tau > 0$ is chosen sufficiently small (see [27, Chap. 7], see also [54]).

Remark 4.12 If the coefficient matrices are symmetric, the Hölder continuity of the solution for (1.1) may be deduced from the results of [26] and [25] by means of suitable (but nontrivial) embedding theorems. Unfortunately, this is not carried out there.

5 Alternative Characterization for Regular Sets

Gröger's concept of regular sets [30] turned out to be a powerful tool for the treatment of mixed boundary value problems. Not only his regularity result [30], based on this, is exploited in some tens of papers; but the regular sets proved also to be an adequate frame for establishing interpolation results for function spaces which include a trace zero condition on part of the boundary, see [28]. Moreover, it allowed to recover resolvent estimates and thus provided tools for the treatment of parabolic equations, which incorporate mixed boundary conditions, see [27, 31]. All of this shows that the concept of regular sets in itself deserves some investigation in order to simplify things.

In this spirit, the aim of this chapter is to prove that for two and three space dimensions the property of a set $\Lambda \cup \Upsilon$ to be regular in the sense of Gröger (see Definition 3.1) can be characterized by Λ being a Lipschitz domain and a certain topological property of Υ (to be specified in a moment). The point is that the resulting conditions usually can be checked 'by appearance'—in contrast to the original definition. Let us explicitly mention that the underlying class of Lipschitz domains is broad enough to contain e.g. the case of two beams, lying on each other with an angle $\neq \pi$, which together do not form a domain with Lipschitz boundary.

We start with the following observation.

Theorem 5.1 If $\Lambda \cup \Upsilon$ is regular, then Λ is a Lipschitz domain.

Proof Let $\mathbf{x} \in \partial \Lambda$. Then there is, due to the definition, an open neighborhood \mathcal{U} of \mathbf{x} and a bi-Lipschitz mapping $\Psi : \mathcal{U} \to \mathbb{R}^d$, such that $\Psi(\mathbf{x}) = 0$ holds and $\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon))$ equals $K_- \cup \tilde{\Sigma}$, where $\tilde{\Sigma} = \emptyset$ or $\tilde{\Sigma} = \Sigma$ or $\tilde{\Sigma} = \Sigma_0$. This means

$$\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon)) = \Psi((\mathcal{U} \cap \Lambda) \cup (\mathcal{U} \cap \Upsilon)) = K_{-} \cup \tilde{\Sigma}$$
(5.1)

and since inner points pass to inner points and boundary points to boundary points under a bi-Lipschitz transformation, (5.1) implies $\Psi(\mathcal{U} \cap \Lambda) = K_{-}$ in all three cases. Hence, \mathcal{U} may serve as the local chart neighborhood required in the definition of a Lipschitz domain, see [29, Definition 1.2.1.2].

We first deal with the easier case of d = 2.

Theorem 5.2 Let $\Lambda \subset \mathbb{R}^2$ be a bounded Lipschitz domain and $\Upsilon \subset \partial \Lambda$ be an open part of the boundary. Then $\Lambda \cup \Upsilon$ is regular in the sense of Gröger, iff the set $\overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)$ is finite and no connected component of $\partial \Lambda \setminus \Upsilon$ consists of a single point.

Proof In view of Theorem 5.1 and an application of the definition for regular sets it is clear that the condition is necessary. Sufficiency follows from the Lipschitz domain property and the fact that any point $x \in \partial \Lambda$ can only lie in $\Upsilon, \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)$ or in the (relative) interior of $\partial \Lambda \setminus \Upsilon$.

Remark 5.3 It is not hard to see that the given condition is equivalent to the following: Υ is a finite union of open arc pieces from $\partial \Lambda$ and $\partial \Lambda \setminus \Upsilon$ is a finite union of (nondegenerate) closed arc pieces.

Now we come to an intrinsic characterization of regular sets in \mathbb{R}^3 , which we regard as the second essential result of this work.

Theorem 5.4 Let $\Lambda \subset \mathbb{R}^3$ be a bounded Lipschitz domain. Assume Υ to be an open subset of $\partial \Lambda$. Then $\Lambda \cup \Upsilon$ is regular in the sense of Gröger, iff the following two conditions are satisfied:

- (i) $\partial \Lambda \setminus \Upsilon$ is the closure of its interior (within $\partial \Lambda$).
- (ii) For any $\mathbf{x} \in \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)$ there is an open neighborhood \mathcal{N} of \mathbf{x} and a bi-Lipschitz mapping $\kappa : \mathcal{N} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon) \to]-1, 1[.$

Proof According to the definition of regular sets the conditions are necessary.

In order to prove sufficiency, we have to show that for every $x \in \partial \Lambda$ there is an open neighborhood \mathcal{U} of x and a bi-Lipschitz mapping Ψ with $\Psi(x) = 0$, such that $\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon))$ is either K_- or $K_- \cup \Sigma$ or $K_- \cup \Sigma_0$.

We first observe that

$$\partial \Lambda = \Upsilon \cup (\partial \Lambda \setminus \Upsilon)^{\circ} \cup \left[(\partial \Lambda \setminus \Upsilon) \cap \overline{\Upsilon} \right], \tag{5.2}$$

where the closure and the interior are again taken with respect to the topology of $\partial \Lambda$. In the following we will treat these three cases separately.

Let $\mathbf{x} \in \Upsilon$. Since Υ was supposed to be open, there is an open set $\mathcal{U}_1 \subset \mathbb{R}^d$ that contains x and satisfies $\mathcal{U}_1 \cap \partial \Lambda \subset \Upsilon$. Furthermore, as Λ is a Lipschitz domain, there is another open neighborhood $\mathcal{U}_2 \subset \mathbb{R}^d$ of x and a bi-Lipschitz transform Φ from \mathcal{U}_2 onto the (open) cube *K*, such that $\Phi(\mathbf{x}) = 0$, $\Phi(\Lambda \cap \mathcal{U}_2) = K_-$ and $\Phi(\partial \Lambda \cap \mathcal{U}_2) = \Sigma$. Since Φ is in particular a homeomorphism, the set $\Phi(\mathcal{U}_1 \cap \mathcal{U}_2)$ is an open neighborhood of 0 and it is contained in *K*. Thus, it contains a homothety *t K* of *K* for some t > 0. If we define $\mathcal{U} := \Phi^{-1}(tK)$, then $\Phi|_{\mathcal{U}}$ is a bi-Lipschitz mapping from the open neighborhood \mathcal{U} of x onto *t K*, such that $\Phi(\mathcal{U} \cap (\Lambda \cup \Upsilon))$ is the set $t(K_- \cup \Sigma)$. Combining Φ with a homothety, we get a bi-Lipschitz mapping Ψ from \mathcal{U} onto *K* that satisfies $\Psi(\mathbf{x}) = 0$ and $\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon)) = K_- \cup \Sigma$.

Analogously, one proves for the (relatively) inner points $x \in (\partial \Lambda \setminus \Upsilon)^\circ$ the existence of a neighborhood \mathcal{U} and a bi-Lipschitz mapping Ψ onto the open cube K such that $\Psi(x) = 0$ and $\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon))$ is the set K_- .

It remains to consider the points of $(\partial \Lambda \setminus \Upsilon) \cap \overline{\Upsilon}$. Let x be an element of this set. As Λ is a Lipschitz domain, there is an open neighborhood \mathcal{O} of x in \mathbb{R}^3 and a bi-Lipschitz mapping Φ from \mathcal{O} onto the cube $K \subset \mathbb{R}^3$, such that $\Phi(x) = 0$, $\Phi(\Lambda \cap \mathcal{O}) = K_-$ and $\Phi(\partial \Lambda \cap \mathcal{O}) = \Sigma$. Exploiting (ii), we find another open neighborhood \mathcal{N} of x and a bi-Lipschitz mapping κ , such that $\kappa(\overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon) \cap \mathcal{N}) =] - 1, 1[$. Without loss of generality we may assume $\kappa(x) = 0 \in \mathbb{R}$.

Our job is now to combine the good properties of Φ and κ . In order to do so, we first define a smaller neighborhood of x that is contained in $\mathcal{O} \cap \mathcal{N}$. Since $\Phi(\mathcal{O} \cap \mathcal{N})$ is an open neighborhood of $0 \in \mathbb{R}^3$, we find a number $t \in]0, 1[$, such that $tK \subset \Phi(\mathcal{O} \cap \mathcal{N})$ and we set $\mathcal{X} := \Phi^{-1}(tK)$. Clearly, \mathcal{X} then is an open neighborhood of x that is contained in $\mathcal{O} \cap \mathcal{N}$. Additionally, one has

$$\Phi(\Lambda \cap \mathcal{X}) = \Phi(\Lambda) \cap tK = \Phi(\Lambda \cap \mathcal{O}) \cap tK = K_{-} \cap tK = tK_{-}$$

and

$$\Phi(\partial \Lambda \cap \mathcal{X}) = \Phi(\partial \Lambda \cap \mathcal{O}) \cap tK = \Sigma \cap tK =]-t, t[\times]-t, t[\times \{0\}].$$

We define $P : \mathbb{R}^3 \to \mathbb{R}^2$ as the canonical projection onto the first two components and $\hat{\Phi} := P\Phi$ with $\mathcal{X} \cap \partial \Lambda$ as its domain of definition. Note that by the above considerations $\hat{\Phi} : \mathcal{X} \cap \partial \Lambda \to]-t, t[\times]-t, t[$ is a bi-Lipschitz mapping with $\hat{\Phi}(\mathbf{x}) = 0 \in \mathbb{R}^2$. Let $]s_-, s_+[\subset]-1, 1[$ be the maximal interval containing 0, such that $\hat{\Phi}\kappa^{-1}(]s_-, s_+[] \subset]-t, t[\times]-t, t[$ and denote the set $\hat{\Phi}\kappa^{-1}(]s_-, s_+[)$ by \mathcal{C} . It is not hard to see, that \mathcal{C} is the connected component of $\hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon))$, which contains $0 \in \mathbb{R}^2$ within $\hat{\Phi}(\mathcal{X} \cap \partial \Lambda) =]-t, t[\times]-t, t[$. We claim:

$$\hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \setminus \mathcal{C} \quad \text{has a positive distance to } 0 \in \mathbb{R}^2.$$
(5.3)

In fact, the elements of $\hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \setminus \mathcal{C}$ correspond to numbers from the set $[-1, s_-] \cup [s_+, 1[$ with respect to the mapping $\kappa \hat{\Phi}^{-1}$. Since κ is bi-Lipschitz, the image of this set under κ^{-1} has a positive distance to x. From this (5.3) follows from the bi-Lipschitz property of $\hat{\Phi}$.

Let θ be the (bi-Lipschitz) mapping $\hat{\Phi}\kappa^{-1}$: $]s_-, s_+[\rightarrow C$. We will identify $]s_-, s_+[$ by means of the (bi-Lipschitz) embedding $\mathbb{R} \ni x \mapsto (x, 0) \in \mathbb{R}^2$ with the set $]s_-, s_+[\times\{0\}$. Then by a deep lying theorem of Tukia, cf. [55, Theorem B], there exists a bi-Lipschitz extension of θ which maps \mathbb{R}^2 onto itself that we will denote by Θ .

Note that Θ maps $]s_-, s_+[\times\{0\}$ onto C and, in particular, $\Theta(0) = 0 \in \mathbb{R}^2$. As $\hat{\Phi}(\mathcal{X} \cap \partial \Lambda) =]-t, t[\times]-t, t[$ is open in \mathbb{R}^2 and due to (5.3) one finds an $\epsilon \in]0, \min\{-s_-, s_+, t\}]$, such that

$$\left[\hat{\Phi}\left(\mathcal{X}\cap\overline{\Upsilon}\cap(\partial\Lambda\setminus\Upsilon)\right)\setminus\mathcal{C}\right]\cap\Theta\left(\left]-\epsilon,\epsilon\right[\times]-\epsilon,\epsilon\right]=\emptyset$$
(5.4)

and simultaneously

$$\Theta(]-\epsilon,\epsilon[\times]-\epsilon,\epsilon[)\subset]-t,t[\times]-t,t[$$

holds, see Fig. 3.



Fig. 3 $\hat{\Phi}(\mathcal{X} \cap \partial \Lambda)$ and $\Theta^{-1}(\hat{\Phi}(\mathcal{X} \cap \partial \Lambda))$

This ϵ will provide us a suitable neighborhood \mathcal{U} of x to complete our task. In fact, we set $\mathcal{Z} := \Theta(] - \epsilon, \epsilon[\times] - \epsilon, \epsilon[)$ and $\mathcal{U} := \Phi^{-1}(\mathcal{Z} \times] - \epsilon, \epsilon[)$. For the bi-Lipschitz mapping Ψ we define first the mapping ξ on $\mathcal{Z} \times] - \epsilon, \epsilon[$ by

$$\xi(y_1, y_2, y_3) := (\Theta^{-1}(y_1, y_2), y_3)$$

and then set $\Psi := \frac{1}{\epsilon} (\xi \circ \Phi)$ on \mathcal{U} . Note that $\Phi(\mathcal{U})$ is exactly $\mathcal{Z} \times] - \epsilon, \epsilon[$ by construction, which implies that everything is well defined. Furthermore, \mathcal{U} is obviously a neighborhood of x and since ξ and Φ are bi-Lipschitz mappings, Ψ is of the same quality.

Regarding the desired mapping properties of Ψ , we already see

$$\Psi(\mathbf{x}) = \xi(\Phi(\mathbf{x}))/\epsilon = \xi(0)/\epsilon = (\Theta^{-1}(0,0),0)/\epsilon = 0$$

and

$$\epsilon \Psi(\mathcal{U}) = \xi(\mathcal{Z} \times] - \epsilon, \epsilon[) = \Theta^{-1}(\mathcal{Z}) \times] - \epsilon, \epsilon[= \epsilon K.$$

Since $\mathcal{U} \subset \mathcal{O}$, we get

$$\epsilon \Psi(\mathcal{U} \cap \Lambda) = \xi \left((\mathcal{Z} \times] - \epsilon, \epsilon[) \cap \Phi(\Lambda \cap \mathcal{O}) \right) = \epsilon K \cap \xi(K_{-})$$
$$= \epsilon K \cap \left(\Theta^{-1}(] - 1, 1[^{2}) \times] - 1, 0[\right).$$

Observing $\Theta(]-\epsilon, \epsilon[\times]-\epsilon, \epsilon[) \subset]-t, t[\times]-t, t[\subset]-1, 1[^2 \text{ and thus }]-\epsilon, \epsilon[\times]-\epsilon, \epsilon[\subset \Theta^{-1}(]-1, 1[^2), \text{ this yields}$

$$\Psi(\mathcal{U} \cap \Lambda) = K_{-}.\tag{5.5}$$

For the boundary of Λ we get by analogous considerations

$$\Psi(\mathcal{U} \cap \partial \Lambda) = \Sigma, \tag{5.6}$$

so the only thing left to prove is $\Psi(\mathcal{U} \cap \Upsilon) = \Sigma_0$.

First, we focus on the interface $\mathcal{U} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)$ and show that this is mapped to the line $]-1, 1[\times \{0\} \times \{0\}$, i.e. the boundary of Σ_0 in Σ . Then, in a second step, we will show that $\Psi(\mathcal{U} \cap \Upsilon)$ must be exactly one of the half squares Σ_0 or $-\Sigma_0$. We first observe

$$\begin{split} \Theta^{-1} \Big(\mathcal{Z} \cap \hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \Big) &= \Theta^{-1} \Big(\mathcal{Z} \cap \Big[\hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \setminus \mathcal{C} \Big] \Big) \\ & \cup \Theta^{-1} \Big(\mathcal{Z} \cap \hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \cap \mathcal{C} \Big). \end{split}$$

Now, the left part of this union is empty thanks to (5.4). Using $\mathcal{Z} = \Theta(]-\epsilon, \epsilon[\times]-\epsilon, \epsilon[)$ and $\mathcal{C} \subset \hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon))$, we thus obtain

$$\Theta^{-1} \Big(\mathcal{Z} \cap \hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \Big) = (] -\epsilon, \epsilon[\times] -\epsilon, \epsilon[) \cap \Theta^{-1}(\mathcal{C})$$
$$= (] -\epsilon, \epsilon[\times] -\epsilon, \epsilon[) \cap (] -\epsilon, \epsilon[\times\{0\})$$
$$=] -\epsilon, \epsilon[\times\{0\}. \tag{5.7}$$

Deringer

Having in mind that $\Phi(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) = \hat{\Phi}(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) \times \{0\}$ and $\mathcal{U} = \Phi^{-1}(\mathcal{Z} \times] - \epsilon, \epsilon[)$, this implies

$$\begin{split} \epsilon \Psi \big(\mathcal{U} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon) \big) &= \xi \big(\Phi \big(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon) \big) \cap \Phi (U) \big) \\ &= \Theta^{-1} \big(\hat{\Phi} \big(\mathcal{X} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon) \big) \cap \mathcal{Z} \big) \times \{ 0 \} \\ &=] - \epsilon, \epsilon [\times \{ 0 \} \times \{ 0 \} \end{split}$$

and thus

$$\Psi(\mathcal{U} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)) =]-1, 1[\times \{0\} \times \{0\}.$$
(5.8)

Now, we claim:

(*) (5.6) and (5.8) imply, that $\Psi(\mathcal{U} \cap \Upsilon)$ is either]-1, 1[×]-1, 0[× {0} or]-1, 1[×]0, 1[× {0}.

Firstly, (5.6) and (5.8) imply that at least one of the two sets in this claim must contain a point from $\Psi(\mathcal{U} \cap \Upsilon)$. Let in this spirit \mathcal{A} be any of the two sets $]-1, 1[\times]-1, 0[$ $\times \{0\}$ or $]-1, 1[\times]0, 1[\times \{0\}$, which contains at least one point from $\Psi(\mathcal{U} \cap \Upsilon)$. Both the sets $\Psi(\mathcal{U} \cap \Upsilon)$ and $\Psi(\mathcal{U} \cap (\partial \Lambda \setminus \Upsilon)^\circ)$ are open in $\Psi(\mathcal{U} \cap \partial \Lambda)$ and, consequently, the—mutually disjoint—sets $\Psi(\mathcal{U} \cap \Upsilon) \cap \mathcal{A}$ and $\Psi(\mathcal{U} \cap (\partial \Lambda \setminus \Upsilon)^\circ) \cap \mathcal{A}$ are open in \mathcal{A} . Since by (5.8) no points from $\Psi(\mathcal{U} \cap \overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon))$ can lie in \mathcal{A} , we have, according to (5.2) and (5.6), the identity

$$(\Psi(\mathcal{U}\cap\Upsilon)\cap\mathcal{A})\cup(\Psi(\mathcal{U}\cap(\partial\Lambda\setminus\Upsilon)^{\circ})\cap\mathcal{A})=\Psi(\mathcal{U}\cap\partial\Lambda)\cap\mathcal{A}=\mathcal{A}.$$
 (5.9)

 \mathcal{A} is connected, therefore (5.9) can only be true if $\Psi(\mathcal{U} \cap (\partial \Lambda \setminus \Upsilon)^{\circ}) \cap \mathcal{A} = \emptyset$. This means: if any of the two sets $]-1, 1[\times]-1, 0[\times\{0\} \text{ and }]-1, 1[\times]0, 1[\times\{0\} \text{ contains a point from } \Psi(\mathcal{U} \cap \Upsilon)$, then it is a subset of $\Psi(\mathcal{U} \cap \Upsilon)$. But then the other cannot contain a point from $\Psi(\mathcal{U} \cap \Upsilon)$, because in this case it also would be a subset of $\Psi(\mathcal{U} \cap \Upsilon)$, which cannot be true in view of (5.8) and our supposition that $\partial \Lambda \setminus \Upsilon$ is the closure of its interior. This proves the claim (*).

Together with (5.5) this gives $\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon)) = K_- \cup \Sigma_0$ or $\Psi(\mathcal{U} \cap (\Lambda \cup \Upsilon)) = K_- \cup (-\Sigma_0)$. In the first case we have finished the proof, in the second we compose Ψ with a reflection at the *x*-*z*-plane to conclude.

Corollary 5.5 If $\Lambda \subset \mathbb{R}^3$ is a Lipschitzian polyhedron and $\overline{\Upsilon} \cap (\partial \Lambda \setminus \Upsilon)$ is a finite union of line segments, then $\Lambda \cup \Upsilon$ is regular.

Remark 5.6 Theorem 5.4 makes precise an old suggestion of Gröger, see [30, Remark 1]. Unfortunately, the given intrinsic characterization is restricted to the dimensions 2 and 3, because there is no analogue of the Tukia theorem in dimensions above 2. Nevertheless, the by far most important cases concerning applications are covered.

6 Application to Semilinear Elliptic Optimal Control Problems

In the subsequent we will employ the results of the previous sections, in particular Theorem 3.3, to derive necessary and sufficient optimality conditions for the following semilinear elliptic optimal control problem with pointwise state and control constraints and jumping boundary conditions:

$$(P) \begin{cases} \text{minimize} \quad J(y,u) := \int_{\Omega} L(x,y(x)) \, dx + \int_{\Gamma} l(x,y(x),u(x)) \, d\sigma \\ \text{subject to} \quad -\nabla \cdot a \, \nabla y + b(x,y) = f \quad \text{in } \Omega, \\ \nu \cdot a \, y = u \quad \text{on } \Gamma, \\ y = 0 \quad \text{on } \partial\Omega \setminus \Gamma \\ \text{and} \quad u_{\min}(x) \le u(x) \le u_{\max}(x) \quad \text{a.e. on } \Gamma, \\ g(x,y(x)) \le 0 \qquad \text{for all } x \in \overline{\Omega}. \end{cases}$$

As already mentioned in the introduction, necessary and sufficient optimality conditions for semilinear elliptic control problems have been addressed by numerous authors before (cf. for instance [8, 11, 12] and the references therein). In particular, we refer to the recent contribution of Casas et al. [12], where an optimal control problem is analyzed that is very similar to (P), but does not contain mixed boundary conditions. However, as we will see in the following, with the results of Sect. 3 at hand, the analysis of [12] can easily be adapted to (P).

Note that mixed boundary conditions play an important role in various applications. A typical example is the optimal control of an electric potential in a conducting material by adjusting the direct current inducing the potential. In the stationary case, this problem is modelled by the electrostatic equation, an elliptic PDE with homogeneous Neumann boundary conditions at isolated surfaces, homogeneous Dirichlet conditions at the anode and inhomogeneous Neumann boundary conditions at the cathode, where the control enters the system (see for instance [17]). Hence, the arising problem is covered by the general problem (P). We point out that state-constrained optimal control problems with mixed boundary conditions and distributed control can be discussed analogously to the following investigation of (P). However, to keep the discussion concise, we do not consider distributed controls here.

In addition to Assumption 3.2 for the coefficient function a, we require the following conditions to be satisfied by the quantities in (P):

Assumption 6.1 The domain $\Omega \subset \mathbb{R}^d$, $d \leq 4$, is a bounded Lipschitz domain, $\Gamma \subset \partial \Omega$ is an open part of its boundary and $\partial \Omega \setminus \Gamma$ has positive measure. The outward unit normal vector on $\partial \Omega$ is denoted by ν . Moreover, $\Omega \cup \Gamma$ is regular in the sense of Gröger (cf. Theorems 5.2 and 5.4 for the two and three-dimensional case). The coefficient function $a : \mathbb{R} \to \mathbb{R}^{d \times d}$ satisfies the conditions (on ρ) in Assumption 3.2, i.e. in particular, *a* is essentially bounded and fulfills the ellipticity condition (3.1). In addition, the function $b : \Omega \times \mathbb{R} \to \mathbb{R}$ is twice continuously differentiable w.r.t. the second variable and monotone increasing, i.e., $\frac{\partial b}{\partial y}(x, y) \ge 0$ a.e. in Ω . Furthermore, there is an s > d/2 such that

$$b(\cdot, 0) \in L^{s}(\Omega)$$
 and $f \in L^{s}(\Omega)$.

Moreover, for all M > 0 there is a constant $C_{b,M} > 0$, such that

$$\left| \frac{\partial b}{\partial y}(\mathbf{x}, y_1) \right| + \left| \frac{\partial^2 b}{\partial y^2}(\mathbf{x}, y_1) \right| \le C_{b,M} \quad \text{and}$$
$$\left| \frac{\partial^2 b}{\partial y^2}(\mathbf{x}, y_2) - \frac{\partial^2 b}{\partial y^2}(\mathbf{x}, y_1) \right| \le C_{b,M} |y_2 - y_1|$$

for almost all $x \in \Omega$ and all $y_1, y_2 \in \mathbb{R}$ with $|y_1|, |y_2| \leq M$.

Assumption 6.2 The function $l: \Gamma \times \mathbb{R} \times \mathbb{R} \to \mathbb{R}$ is a Carathéodory function of class C^2 w.r.t. the second and third variables. In addition, l is convex w.r.t. the third variable. Moreover, $l(\cdot, 0, 0) \in L^1(\Gamma)$ and for all M > 0 there exist a constant $C_{l,M} > 0$ and a function $\psi_{l,M} \in L^2(\Gamma)$ with

$$\left|\frac{\partial l}{\partial y}(\mathbf{x}, y_1, u_1)\right| + \left|\frac{\partial l}{\partial u}(\mathbf{x}, y_1, u_1)\right| \le \psi_{l,M}(\mathbf{x}), \qquad |D^2_{(y,u)}l(\mathbf{x}, y_1, u_1)| \le C_{l,M}, |D^2_{(y,u)}l(\mathbf{x}, y_2, u_2) - D^2_{(y,u)}l(\mathbf{x}, y_1, u_1)| \le C_{l,M} \left(|y_2 - y_1| + |u_2 - u_1|\right),$$
(6.1)

for almost all $x \in \Gamma$ and all $|y_1|, |y_2|, |u_1|, |u_2| \leq M$. Here, $D_{(y,u)}^2 l$ denotes the Hessian of l w.r.t. (y, u). Furthermore, $L : \Omega \times \mathbb{R} \to \mathbb{R}$ fulfills analogous conditions, i.e. it is of class C^2 w.r.t. the second variable, $L(\cdot, 0, 0) \in L^1(\Omega)$ and for all M > 0 there exist $C_{L,M} > 0$ and $\psi_{L,M} \in L^2(\Omega)$ with

$$\left|\frac{\partial L}{\partial y}(\mathbf{x}, y_1)\right| \le \psi_{L,M}(\mathbf{x}), \qquad \left|\frac{\partial^2}{\partial y^2}L(\mathbf{x}, y_1)\right| \le C_{L,M},$$

$$\left|\frac{\partial^2}{\partial y^2}L(\mathbf{x}, y_2) - \frac{\partial^2}{\partial y^2}L(\mathbf{x}, y_1)\right| \le C_{L,M} |y_2 - y_1|,$$
(6.2)

for almost all $x \in \Omega$ and all $|y_1|, |y_2| \leq M$.

Assumption 6.3 The bounds in the control constraints satisfy $u_{\min}, u_{\max} \in L^{\infty}(\Gamma)$ with $u_{\min}(x) < u_{\max}(x)$ a.e. in Ω . Moreover, $g : \overline{\Omega} \times \mathbb{R} \to \mathbb{R}$ is continuous and twice continuously differentiable w.r.t. the second variable. In addition, $\frac{\partial g}{\partial y}$ and $\frac{\partial^2 g}{\partial y^2}$ are continuous on $\overline{\Omega} \times \mathbb{R}$, and g(x, 0) < 0 is satisfied on $\partial \Omega \setminus \Gamma$.

Note that the last conditions in Assumption 6.3 allow for the existence of a Slater point, which is essential for the derivation of first-order necessary conditions (see Assumption 6.21 below).

6.1 Discussion on the State Equation

We start the discussion of (P) with the analysis of the state equation, i.e.

$$-\nabla \cdot a \,\nabla y + b(\cdot, y) = f \quad \text{in } \Omega,$$

$$\nu \cdot a \, y = u \quad \text{on } \Gamma,$$

$$y = 0 \quad \text{on } \partial \Omega \setminus \Gamma.$$
(6.3)

Definition 6.4 Let $q \in [2, \infty[$ and let $\mathfrak{s}, \mathfrak{r} \in \mathbb{R}$ satisfy $\mathfrak{s} > 1$ and $\mathfrak{r} > 1$, if d = q = 2, and $\mathfrak{s} \ge dq/(d+q)$ and $\mathfrak{r} \ge (d-1)q/d$, otherwise. Moreover, let $\varphi \in L^{\mathfrak{s}}(\Omega)$ and $\psi \in L^{\mathfrak{r}}(\Gamma)$ be given. Then we denote the elements of $W_{\Gamma}^{-1,q}(\Omega)$, associated to φ and ψ , by $\tilde{\varphi}$ and $\tilde{\psi}$, i.e.

$$\langle \tilde{\varphi} , w \rangle_{W_{\Gamma}^{-1,q}(\Omega)} := \int_{\Omega} \varphi \, w \, dx, \qquad \langle \tilde{\psi} , w \rangle_{W_{\Gamma}^{-1,q}(\Omega)} := \int_{\Gamma} \psi \, w \, d\sigma, \quad w \in W_{\Gamma}^{1,q'}(\Omega).$$

$$(6.4)$$

Since embedding and trace theorems guarantee $w \in L^{\mathfrak{s}'}(\Omega)$ and $\tau_{\Gamma} w \in L^{\mathfrak{r}'}(\Gamma)$, if $w \in W_{\Gamma}^{1,q'}(\Omega)$, the integrals in (6.4) are finite and, hence, $\tilde{\varphi}$ and $\tilde{\psi}$ are well defined.

Definition 6.5 Suppose that $f \in L^{\mathfrak{s}}(\Omega)$, $\mathfrak{s} > 2d/(d+2)$, and $u \in L^{\mathfrak{r}}(\Gamma)$, $\mathfrak{r} > (2d-2)/d$. Then a function $y \in W_{\Gamma}^{1,2}(\Omega) \cap L^{\infty}(\Omega)$ is said to be a *solution of* (6.3), if it fulfills the operator equation

$$-\nabla \cdot a \,\nabla y + \tilde{b}(y) = \tilde{f} + \tilde{u} \quad \text{in } W_{\Gamma}^{-1,2}(\Omega), \tag{6.5}$$

where $\tilde{f}, \tilde{u} \in W_{\Gamma}^{-1,2}(\Omega)$ are defined according to Definition 6.4 and $\tilde{b}: L^{\infty}(\Omega) \to W_{\Gamma}^{-1,2}(\Omega)$ is analogously given by

$$\langle \tilde{b}(\mathbf{y}), w \rangle_{W_{\Gamma}^{-1,2}(\Omega)} := \int_{\Omega} b(\mathbf{x}, \mathbf{y}(\mathbf{x})) w(\mathbf{x}) d\mathbf{x}, \quad w \in W_{\Gamma}^{1,2}(\Omega).$$

Note that, due to Assumption 6.1, the Nemyzki operator $\Phi_b(y) := b(\cdot, y(\cdot))$ is continuous from $L^{\infty}(\Omega)$ to $L^{\mathfrak{s}}(\Omega)$, $\mathfrak{s} > 2d/(d+2)$, so \tilde{b} is well defined.

Theorem 6.6 Let $f \in L^s(\Omega)$ with s > d/2 and $u \in L^r(\Gamma)$ with r > d - 1. Then, under Assumption 6.1, there exists a unique solution of (6.3) in the sense of Definition 6.5. Moreover, there is an $\alpha > 0$ such that this solution belongs to $C^{\alpha}(\Omega)$.

Proof The existence of a unique solution in $W^{1,2}_{\Gamma}(\Omega) \cap L^{\infty}(\Omega)$ is standard (cf. for instance [8] or [1] and the references therein). For convenience of the reader, we recall the main arguments. First, one considers a modified nonlinearity given by

$$b_k(\cdot, y) := \begin{cases} b(\cdot, k), & \text{if } y > k, \\ b(\cdot, y), & \text{if } -k \le y \le k, \\ b(\cdot, -k), & \text{if } y < -k, \end{cases}$$

Deringer

with some k > 0. For a nonlinearity of this form, Browder and Minty's theorem for monotone operators immediately implies the existence of a unique solution in $W_{\Gamma}^{1,2}(\Omega)$. Then a classical argument in the spirit of Stampacchia [52] yields

$$\|y\|_{L^{\infty}(\Omega)} \le c_{\infty} \left(\|f\|_{L^{s}(\Omega)} + \|u\|_{L^{r}(\Gamma)} + 1 \right)$$

with a constant c_{∞} independent of f and k. It is easily verified that the mixed boundary conditions do not influence the analysis in [52].

Hence, if we choose $k \ge c_{\infty} (||f||_{L^{s}(\Omega)} + ||u||_{L^{r}(\Gamma)} + 1)$, then the solution of the truncated problem coincides with the one of (6.3). It remains to verify the Hölder continuity of *y*, which follows from Theorem 3.3 together with a classical bootstrapping argument. To see this, rewrite (6.5) as

$$(-\nabla \cdot a \,\nabla + 1)y = g \tag{6.6}$$

with $g := \tilde{f} + \tilde{u} - \tilde{b}(y) + \tilde{y}$ and \tilde{y} according to Definition 6.4. Due to $f, b(\cdot, y) \in L^{s}(\Omega)$, s > d/2, $u \in L^{r}(\Gamma)$, r > d - 1, and $y \in L^{\infty}(\Omega)$, Sobolev embedding theorems give that $g \in W_{\Gamma}^{-1,q}(\Omega)$ for a q > d. Therefore, Theorem 3.3 implies $y \in C^{\alpha}(\Omega)$.

Definition 6.7 For the rest of this section, let s > d/2 and r > d - 1 be fixed, but arbitrary. Moreover, f is a fixed inhomogeneity in $L^s(\Omega)$ (cf. Assumption 6.1). Based on Theorem 6.6, we introduce the *control-to-state operator* $S : L^r(\Gamma) \to W^{1,2}_{\Gamma}(\Omega) \cap C^{\alpha}(\Omega)$, mapping u to the solution of (6.3).

Lemma 6.8 Suppose that there is a sequence $\{u_k\}$ converging weakly to u in $L^r(\Gamma)$. Then $S(u_k) \to S(u)$ in $W_{\Gamma}^{1,2}(\Omega) \cap C^{\alpha}(\Omega)$.

Proof With Theorem 3.3 at hand, the arguments are standard (cf. for instance [12]). Nevertheless, let us recall the basic ideas. In all what follows we use the notation $y_k := S(u_k)$ and y := S(u). The weak convergence of $\{u_k\}$ implies the uniform boundedness of this sequence in $L^r(\Gamma)$ giving in turn that $\{y_k\}$ is uniformly bounded in $C(\overline{\Omega})$. Hence, $\{y_k\}$ and $\{\Phi_b(y_k)\}$, with Φ_b as defined above, converge weakly in $L^s(\Omega)$ with s > d/2, to some z_y and z_{Φ} , respectively. Now define the sequence $\{g_k\}$ in $W_{\Gamma}^{-1,q}(\Omega)$ by $g_k := \tilde{f} + \tilde{u}_k - \tilde{b}(y_k) + \tilde{y}_k$. Due to the compact embedding $L^s(\Omega) \hookrightarrow W_{\Gamma}^{-1,q}(\Omega)$ and the compactness of the trace operator $\tau_{\Gamma} : W_{\Gamma}^{1,q'}(\Omega) \to L^{r'}(\Gamma)$, weak convergences of $\{u_k\}, \{y_k\}, \text{ and } \{\Phi_b(y_k)\}$ imply strong convergence of $\{g_k\}$ in $W_{\Gamma}^{-1,q}(\Omega)$ to $g := \tilde{f} + \tilde{u} - \tilde{z}_{\Phi} + \tilde{z}_y$, where \tilde{z}_{Φ} and \tilde{z}_y again denote the elements in $W_{\Gamma}^{-1,q}(\Omega)$ associated to z_{Φ} and z_y , respectively. Now consider again the auxiliary equation (6.6) with g_k as inhomogeneity. Theorem 3.3 then implies

$$y_k \to \eta := (-\nabla \cdot a \nabla + 1)^{-1} g$$
 in $W^{1,2}_{\Gamma}(\Omega) \cap C^{\alpha}(\Omega)$.

This in particular guarantees $y_k \to \eta$ in $L^{\infty}(\Omega)$ and, hence, $\Phi_b(y_k) \to \Phi_b(\eta)$ in $L^s(\Omega)$, s > d/2, as well as $\tilde{y}_k \to \tilde{\eta}$ and $\tilde{b}(y_k) \to \tilde{b}(\eta)$ in $W_{\Gamma}^{-1,q}(\Omega)$. Consequently, η is the solution of (6.3) associated to u, which implies $y_k \to y$ in $W_{\Gamma}^{1,2}(\Omega) \cap C^{\alpha}(\Omega)$.

Now, we turn to the linearized version of (6.3). Given a $\bar{y} \in L^{\infty}(\Omega)$, the linearized state equation reads as

$$-\nabla \cdot a \,\nabla y + \tilde{b}'(\bar{y})y = \tilde{h} \quad \text{in } W_{\Gamma}^{-1,2}(\Omega), \tag{6.7}$$

where $\tilde{h} \in W_{\Gamma}^{-1,2}(\Omega)$ and $\tilde{b}'(\bar{y}) : W_{\Gamma}^{1,2}(\Omega) \to W_{\Gamma}^{-1,2}(\Omega)$ is defined by

$$\langle \tilde{b}'(\bar{y})y, w \rangle_{W_{\Gamma}^{-1,2}(\Omega)} := \int_{\Omega} \frac{\partial b}{\partial y}(x, \bar{y}(x))y(x)w(x)\,dx, \quad w \in W_{\Gamma}^{1,2}(\Omega).$$

Note that Assumption 6.1 implies $\frac{\partial b}{\partial y}(\mathbf{x}, \bar{y}(\mathbf{x})) \in L^{\infty}(\Omega)$ and that, due to the monotonicity of b, $\frac{\partial b}{\partial y}(\mathbf{x}, \bar{y}(\mathbf{x})) \ge 0$ holds true a.e. in Ω . Hence, an immediate consequence of Corollary 3.6 is the following

Lemma 6.9 Let $\bar{y} \in L^{\infty}(\Omega)$ be given. For every $\tilde{h} \in W_{\Gamma}^{-1,2}(\Omega)$ there is a unique solution $y \in W_{\Gamma}^{1,2}(\Omega)$ of (6.7). Furthermore, if $\tilde{h} \in W_{\Gamma}^{-1,q}(\Omega)$ for some q > d, then $y \in W_{\Gamma}^{1,2}(\Omega) \cap C^{\alpha}(\Omega)$ for some $\alpha > 0$.

In view of Assumption 6.1, the Nemyzki operator $\Phi_b(y) = b(\cdot, y(\cdot))$ clearly is twice continuously Fréchet differentiable in $L^{\infty}(\Omega)$. Thus, together with Lemma 6.9, the implicit function theorem implies the following result (for a detailed proof see for instance [10]).

Theorem 6.10 Under Assumption 6.1 the control-to-state operator *S* is twice continuously Fréchet differentiable from $L^r(\Gamma)$ to $W^{1,2}_{\Gamma}(\Omega) \cap C^{\alpha}(\Omega)$. Its first derivative at $\bar{u} \in L^r(\Gamma)$ in direction $h \in L^r(\Gamma)$ solves

$$-\nabla \cdot a \,\nabla y + \tilde{b}'(\bar{y})y = \tilde{h} \quad in \ W_{\Gamma}^{-1,2}(\Omega), \tag{6.8}$$

where $\bar{y} = S(\bar{u})$ and \tilde{h} denotes the element of $W_{\Gamma}^{-1,2}(\Omega)$ associated to h. Furthermore, $\eta = S''(\bar{u})[h_1, h_2], h_i \in L^r(\Gamma), i = 1, 2$, is the solution of

$$-\nabla \cdot a \,\nabla \eta + \tilde{b}'(\bar{y})\eta = -\tilde{b}''(\bar{y})y_1y_2 \quad in \ W_{\Gamma}^{-1,2}(\Omega) \tag{6.9}$$

with $y_i = S'(\bar{u})h_i$, i = 1, 2, i.e. the solution of (6.8), and

$$\langle \tilde{b}''(\bar{y})y_1y_2, w \rangle_{W_{\Gamma}^{-1,2}(\Omega)} := \int_{\Omega} \frac{\partial^2 b}{\partial y^2}(\mathbf{x}, \bar{y}(\mathbf{x}))y_1(\mathbf{x})y_2(\mathbf{x})w(\mathbf{x})d\mathbf{x}, \quad w \in W_{\Gamma}^{1,2}(\Omega).$$

Note that, due to $h_i \in L^r(\Gamma)$, i = 1, 2, Lemma 6.9 yields $y_i \in L^{\infty}(\Omega)$. In addition, Assumption 6.1 implies $\frac{\partial^2 b}{\partial y^2}(\cdot, \bar{y}(\cdot)) \in L^s(\Omega)$, so $\tilde{b}''(\bar{y})y_1y_2$ is well defined.

6.2 An Adjoint Equation Involving Measures

In the following section, we consider an elliptic PDE with measures as inhomogeneity. Such an equation will arise as adjoint equation in the derivation of first-order conditions for (P), see Sect. 6.3 below. It turns out that the Theorem 3.3 is in particular well suited for the discussion of the adjoint equation, since it immediately yields the existence and uniqueness. The underlying notion of solutions of equations with measures (see Definition 6.13 below) coincides with the one given in [52]. It is to be noted however that the adjoint state, defined in this way, does in general not allow for a distributional characterization in terms of a 'classical' variational formulation (cf. Remark 6.15 below).

Let us consider the following PDE

$$-\nabla \cdot a^{\top} \nabla p + \beta p = \mu_{\Omega} \quad \text{in } \Omega,$$

$$\nu \cdot a^{\top} p = \mu_{\Gamma} \quad \text{on } \Gamma,$$

$$p = 0 \qquad \text{on } \partial\Omega \setminus \Gamma,$$
(6.10)

where Ω and *a* are supposed to fulfill the assumptions of Theorem 3.3. Moreover, β is a fixed, but arbitrary non-negative function in $L^{\infty}(\Omega)$. Furthermore, the inhomogeneity μ is given in $\mathcal{M}(\overline{\Omega})$ which is the space of regular Borel measures that can be identified with the dual of $C(\overline{\Omega})$ by means of the Riesz representation theorem. Moreover, μ_{Ω} and μ_{Γ} denote the restrictions of μ to Ω and Γ , respectively.

Throughout this section, let $q \in \mathbb{R}$ be a fixed number with q > d. As in Corollary 3.7, we denote the domain of the maximal restriction of $-\nabla \cdot a \nabla + 1$ to $W_{\Gamma}^{-1,q}(\Omega)$ by D_q , i.e.

$$D_q := \{ v \in W^{1,2}_{\Gamma}(\Omega) : -\nabla \cdot a \, \nabla v + v \in W^{-1,q}_{\Gamma}(\Omega) \}.$$

Thus, D_q is the subset of $W^{1,2}_{\Gamma}(\Omega)$ the elements v of which admit an continuous extension of the linear form

$$W_{\Gamma}^{1,2}(\Omega) \ni \psi \mapsto \int_{\Omega} a \nabla v \cdot \nabla \psi \, d\mathbf{x} \tag{6.11}$$

to the space $W^{1,q'}_{\Gamma}(\Omega)$.

Remark 6.11 We point out that one can in general not expect that, for $v \in D_q$, the duality $\langle -\nabla \cdot a \nabla v, \psi \rangle$ can be expressed by the right-hand side of (6.11), if ψ is an arbitrary element from $W_{\Gamma}^{1,q'}(\Omega)$.

Lemma 6.12 Endowed with the norm

$$\|v\|_{D_q} := \| -\nabla \cdot a \,\nabla v + v \|_{W_{\Gamma}^{-1,q}(\Omega)}$$
(6.12)

the set D_q is a Banach space. Moreover, it is compactly embedded in $C(\overline{\Omega})$.

Proof By construction $-\nabla \cdot a \nabla + 1 : D_q \to W_{\Gamma}^{-1,q}(\Omega)$ is linear and bijective, and if we equip D_q with the norm $\|\cdot\|_{D_q}$, then $-\nabla \cdot a \nabla + 1$ is an isometry, again by construction. Hence, the completeness of the domain space $(D_q, \|\cdot\|_{D_q})$ follows from the completeness of the image space $W_{\Gamma}^{-1,q}(\Omega)$.

Next we define the operator $A_{\beta}: D_q \to W_{\Gamma}^{-1,q}(\Omega)$ by

$$A_{\beta} y := -\nabla \cdot a \,\nabla y + \beta \, y.$$

Moreover, the embedding operator $E_q : \mathcal{M}(\overline{\Omega}) \to D_q^*$ is given by

$$\langle E_q \ \mu \ , \ v \rangle_{D_q^*} = \int_{\overline{\Omega}} v \, d\mu \quad \forall v \in D_q \subset C(\overline{\Omega}).$$

With these settings at hand, we are now in the position to introduce the notion of solutions to (6.10) fitting to the theory presented in the sections before.

Definition 6.13 A function $p \in W_{\Gamma}^{1,q'}(\Omega)$, q' = q/(q-1) < d/(d-1), is said to be a *solution* of (6.10), if the equation

$$A_{\beta}^{*} p = E_{q} \mu \quad \text{in } D_{q}^{*} \quad \Leftrightarrow \quad \langle A_{\beta} v, p \rangle_{W_{\Gamma}^{-1,q}(\Omega)} = \int_{\Omega} v \, d\mu_{\Omega} + \int_{\Gamma} v \, d\mu_{\Gamma} \quad \forall v \in D_{q}$$

$$\tag{6.13}$$

is satisfied.

Theorem 6.14 *There exists a unique solution of* (6.10) *in the sense of the above Definition* 6.13, *and there holds*

$$\|p\|_{W_{r}^{1,q'}(\Omega)} \le c \, \|\mu\|_{\mathcal{M}(\overline{\Omega})} \tag{6.14}$$

with a constant c > 0 independent of μ .

Proof According to Theorem 3.3 and Corollary 3.6, respectively, $A_{\beta} : D_q \rightarrow W_{\Gamma}^{-1,q}(\Omega)$ is a topological isomorphism. Therefore, $A_{\beta}^* : D_q^* \rightarrow W_{\Gamma}^{1,q'}(\Omega)$ is continuously invertible, too. The continuity of $D_q \hookrightarrow C(\overline{\Omega})$ by Lemma 6.12 finally gives (6.14).

Remark 6.15 Since $W_{\Gamma}^{1,q}(\Omega) \hookrightarrow D_q$ by construction, Definition 6.13 implies that the solution of the operator equation $A_{\beta}^* p = E_q \mu$ in $W_{\Gamma}^{-1,q}(\Omega)$ also solves the following variational formulation

$$\int_{\Omega} (a^{\top} \nabla p \cdot \nabla v + \beta p v) d\mathbf{x} = \int_{\Omega} v d\mu_{\Omega} + \int_{\Gamma} v d\mu_{\Gamma} \quad \forall v \in W_{\Gamma}^{1,q}(\Omega), \quad (6.15)$$

which seems to be a more comfortable equation than the ones in (6.13). However, (6.15) is not equivalent to (6.13), because, in contrast to the operator equations in (6.13), the equation (6.15) does not determine the solution *p* uniquely. Since this is

rather surprising, let us put some more emphasis on this point. Serrin showed in [47] the following fact: Assume that Ω is a ball around 0 in \mathbb{R}^d , $\Gamma = \emptyset$ and r < 2. Then there is a bounded, elliptic coefficient function $a = a_r$ and a non-vanishing function $p \in W_0^{1,r}(\Omega)$, such that (6.15) is satisfied with right-hand side 0. This counterexample shows that the variational formulation (6.15) is not sufficient to describe a solution of (6.10) uniquely, since the space of test functions, i.e. $W_{\Gamma}^{1,q}(\Omega)$, is too poor and has to be extended to D_q . Hence, the operator from $W_{\Gamma}^{1,q'}(\Omega)$ to $W_{\Gamma}^{-1,q'}(\Omega)$, defined by the left-hand side in (6.15), is in general not injective if q < 2 is not contained in an interval around 2 whose radius depends on the data of the problem. Duality then implies that $W_{\Gamma}^{1,q}(\Omega)$ is not always dense in D_q , this equipped with the norm given in (6.12).

It is to be noted however that in the two-dimensional case Proposition 4.1 allows to show the equivalence of (6.15) and (6.13): one uses the fact that the space $\mathcal{M}(\overline{\Omega})$ of finite, regular Borel measures on $\overline{\Omega}$ embeds into $W_{\Gamma}^{-1,r}(\Omega)$ for every r < 2 and takes in this spirit (for our fixed coefficient function *a*) *q* sufficiently close to 2. As a consequence, $D_q = W_{\Gamma}^{1,q}(\Omega)$ is obtained in this case.

Remark 6.16 As shown by Prignet in [44, Sect. 2.1], the notion of solutions in Definition 6.13 coincides with the one given by Stampacchia in [52].

Remark 6.17 Since $W_{\Gamma}^{1,2}(\Omega)$ is dense in $W_{\Gamma}^{1,q'}(\Omega)$ and by continuity of A_{β} , (6.13) is equivalent to

$$\lim_{k \to \infty} \int_{\Omega} (a^{\top} \nabla p_k \cdot \nabla v + \beta \, p_k \, v) \, d\mathbf{x} = \int_{\Omega} v \, d\mu_{\Omega} + \int_{\Gamma} v \, d\mu_{\Gamma} \quad \forall v \in D_q, \quad (6.16)$$

where $\{p_k\} \subset W_{\Gamma}^{1,2}(\Omega)$ is an arbitrary sequence converging to p in $W_{\Gamma}^{1,q'}(\Omega)$.

Remark 6.18 In [1], Alibert and Raymond introduce another notion of solutions to elliptic PDEs with nonsmooth data and measures as inhomogeneity, cf. [1, Theorem 4]. Their definition of solutions is based on (6.15) and a formula of integration by parts to ensure uniqueness. By adapting their arguments to the case with mixed boundary conditions, one can show that this definition of the adjoint equation is equivalent to Definition 6.13 in the sense that they yield the same solution in $W_{\Gamma}^{1,q'}(\Omega)$.

6.3 Necessary and Sufficient Optimality Conditions

With the above results, the analysis of this section is along the lines of [12]. Hence, we shorten the description, if the arguments are analogous to the ones in [12]. First, let us introduce the reduced objective functional $j: L^{\infty}(\Gamma) \to \mathbb{R}$ and the Lagrange function $\mathcal{L}: L^{\infty}(\Gamma) \times \mathcal{M}(\overline{\Omega}) \to \mathbb{R}$ by

$$j(u) := J(S(u), u) = \int_{\Omega} L(\mathbf{x}, S(u)(\mathbf{x})) \, d\mathbf{x} + \int_{\Gamma} l(\mathbf{x}, S(u)(\mathbf{x}), u(\mathbf{x})) \, d\sigma, \quad (6.17)$$

$$\mathcal{L}(u,\mu) := j(u) + \int_{\overline{\Omega}} g(\mathbf{x}, S(u)(\mathbf{x})) d\mu.$$
(6.18)

Deringer

Lemma 6.19 Let $\mu \in \mathcal{M}(\overline{\Omega})$ be arbitrary. Then the Lagrange function is twice continuously Fréchet differentiable w.r.t. u from $L^{\infty}(\Gamma)$ to \mathbb{R} . If $\mu_{\partial\Omega\setminus\Gamma} = 0$, then its first partial derivative at \overline{u} in direction $h \in L^{\infty}(\Gamma)$ is given by

$$\frac{\partial \mathcal{L}}{\partial u}(\bar{u},\mu)h = \int_{\Gamma} \left[\frac{\partial l}{\partial u}(\mathbf{x},\bar{\mathbf{y}}(\mathbf{x}),\bar{u}(\mathbf{x})) + p(\mathbf{x})\right]h(\mathbf{x})\,d\sigma,\tag{6.19}$$

where $\bar{y} = S(\bar{u})$ and $p \in W^{1,q'}_{\Gamma}(\Omega)$ solves

$$-\nabla \cdot a^{\top} \nabla p + \frac{\partial b}{\partial y} (\cdot, \bar{y}) p = \frac{\partial L}{\partial y} (\cdot, \bar{y}) + \frac{\partial g}{\partial y} (\cdot, \bar{y}) \mu_{\Omega} \quad in \ \Omega,$$
$$\nu \cdot a^{\top} p = \frac{\partial l}{\partial y} (\cdot, \bar{y}, \bar{u}) + \frac{\partial g}{\partial y} (\cdot, \bar{y}) \mu_{\Gamma} \quad on \ \Gamma, \qquad (6.20)$$
$$p = 0 \qquad \qquad on \ \partial \Omega \setminus \Gamma,$$

in the sense of Definition 6.13. Moreover, the second derivative of \mathcal{L} at \overline{u} in directions $h_1, h_2 \in L^{\infty}(\Gamma)$ is given by

$$\frac{\partial^{2} \mathcal{L}}{\partial u^{2}}(\bar{u},\mu)h_{1}h_{2} = \int_{\Omega} \left[\frac{\partial^{2} L}{\partial y^{2}}(\mathbf{x},\bar{y},\bar{u})y_{1}y_{2} - p \frac{\partial^{2} b}{\partial y^{2}}(\mathbf{x},\bar{y})y_{1}y_{2} \right] d\mathbf{x}$$
$$+ \int_{\Gamma} \left[\frac{\partial^{2} l}{\partial y^{2}}(\mathbf{x},\bar{y},\bar{u})y_{1}y_{2} + \frac{\partial^{2} l}{\partial y \partial u}(\mathbf{x},\bar{y},\bar{u})(y_{1}h_{2} + y_{2}h_{1}) \right]$$
$$+ \frac{\partial^{2} l}{\partial u^{2}}(\mathbf{x},\bar{y},\bar{u})h_{1}h_{2} d\sigma + \int_{\overline{\Omega}} \frac{\partial^{2} g}{\partial y^{2}}(\mathbf{x},\bar{y},\bar{u})y_{1}y_{2} d\mu \quad (6.21)$$

with $y_i = S'(\bar{u})h_i$, i = 1, 2.

Proof The arguments are standard (cf. [10]). Nevertheless, we shortly describe the derivation of (6.19) to see how the associated theory is influenced by the analysis of the adjoint equation as carried out in Sect. 6.2. The differentiability of \mathcal{L} is an immediate consequence of Assumptions 6.1–6.3 and Theorem 6.10. Concerning the explicit form of $\frac{\partial \mathcal{L}}{\partial u}$, the chain rule yields

$$\frac{\partial \mathcal{L}}{\partial u}(\bar{u},\mu)h = \int_{\Omega} \frac{\partial L}{\partial y}(\mathbf{x},\bar{y},\bar{u})y\,d\mathbf{x} + \int_{\Gamma} \left(\frac{\partial l}{\partial y}(\mathbf{x},\bar{y},\bar{u})y + \frac{\partial l}{\partial u}(\mathbf{x},\bar{y},\bar{u})h\right)d\sigma + \left(\frac{\partial g}{\partial y}(\mathbf{x},\bar{y})y\right)_{\mathcal{M}(\overline{\Omega})}$$
(6.22)

with $\bar{y} = S(\bar{u})$ and $y = \frac{\partial S}{\partial u}(\bar{u})h$, i.e. $y \in D_q$ solves $A_{b'(\bar{y})} y = -\nabla \cdot a \nabla y + \tilde{b}'(\bar{y})y = \tilde{h}$ in $W_{\Gamma}^{-1,q}(\Omega)$. The unique solution of (6.20) in the sense of Definition 6.13 solves

$$\langle A_{b'(\bar{y})}v, p \rangle_{W_{\Gamma}^{-1,q}(\Omega)} = \int_{\Omega} \frac{\partial L}{\partial y}(\mathbf{x}, \bar{y}, \bar{u})v \, d\mathbf{x} + \int_{\Gamma} \frac{\partial l}{\partial y}(\mathbf{x}, \bar{y}, \bar{u})v \, d\sigma$$
$$+ \left\langle \mu, \frac{\partial g}{\partial y}(\mathbf{x}, \bar{y})v \right\rangle_{\mathcal{M}(\overline{\Omega})}$$

Deringer

for all $v \in D_q$. Choosing v = y as test function yields

$$\int_{\Omega} \frac{\partial L}{\partial y}(\mathbf{x}, \bar{y}, \bar{u}) y \, d\mathbf{x} + \int_{\Gamma} \frac{\partial l}{\partial y}(\mathbf{x}, \bar{y}, \bar{u}) y \, d\sigma + \left\langle \mu, \frac{\partial g}{\partial y}(\mathbf{x}, \bar{y}) y \right\rangle_{\mathcal{M}(\overline{\Omega})}$$
$$= \left\langle A_{b'(\bar{y})} y, p \right\rangle_{W_{\Gamma}^{-1,q}(\Omega)} = \left\langle \tilde{h}, p \right\rangle_{W_{\Gamma}^{-1,q}(\Omega)} = \int_{\Gamma} h \, p \, d\sigma.$$

Inserting this into (6.22) gives (6.19). Finally, (6.21) follows from an analogous argument. $\hfill \Box$

Definition 6.20 A function $u \in L^{\infty}(\Gamma)$ is called *feasible* for (P), if it fulfills

$$u_{\min}(\mathbf{x}) \le u(\mathbf{x}) \le u_{\max}(\mathbf{x})$$
 a.e. on Γ ,
 $g(\mathbf{x}, S(u)(\mathbf{x})) \le 0$ for all $\mathbf{x} \in \overline{\Omega}$.

Under Assumptions 6.1–6.3 there is at least one (global) solution of (P), provided that a feasible function exists (see [9, Theorem 8] for the proof).

Let us now turn to necessary optimality conditions for (P). It is well known that certain constraint qualifications are required to discuss pointwise inequality constraints on the state as they occur in (P). Here, we rely on the following *linearized Slater condition*.

Definition 6.21 Let $\bar{u} \in L^{\infty}(\Gamma)$ be feasible for (P). We say that the *linearized Slater condition* is fulfilled at \bar{u} , if there exists a function $\hat{u} \in L^{\infty}(\Gamma)$, such that

$$u_{\min}(\mathbf{x}) \le \hat{u}(\mathbf{x}) \le u_{\max}(\mathbf{x})$$
 a.e. on Γ , (6.23)

$$g(\mathbf{x}, \bar{\mathbf{y}}(\mathbf{x})) + \frac{\partial g}{\partial y}(\mathbf{x}, \bar{\mathbf{y}}(\mathbf{x}))\hat{\mathbf{y}}(\mathbf{x}) < 0 \quad \text{for all } \mathbf{x} \in \overline{\Omega},$$
(6.24)

where $\bar{y} = S(\bar{u})$ and $\hat{y} = S'(\bar{u})(\hat{u} - \bar{u})$.

Note that, due to $\hat{y}|_{\partial\Omega\setminus\Gamma} = 0$, (6.24) yields $g(\mathbf{x}, \bar{y}(\mathbf{x})) = g(\mathbf{x}, 0) < 0$ for all $x \in \partial\Omega \setminus \Gamma$ which is guaranteed by Assumption 6.3. With the existence and regularity results for the state and the adjoint equation obtained before, the theory of first-order necessary conditions for (P) is standard. For the corresponding theorem, we define the *Hamiltonian associated to* (P), denoted by $H : \Gamma \times \mathbb{R} \times \mathbb{R} \to \mathbb{R}$:

$$H(x, y, u, p) := l(x, y, u) + p(u - b(x, y)).$$

The definition of H allows to formulate the first-order necessary conditions in form of Pontryagin's principle. For the corresponding proof we refer to [7].

Theorem 6.22 Suppose that $\bar{u} \in L^{\infty}(\Gamma)$ is a local solution of (P) in the topology of $L^{\infty}(\Gamma)$, i.e., $j(\bar{u}) \leq j(u)$ for all feasible u with $||u - \bar{u}||_{L^{\infty}(\Gamma)} \leq \varepsilon$. Furthermore, denote the state associated to \bar{u} by $\bar{y} \in W_{\Gamma}^{1,2}(\Omega) \cap C^{\alpha}(\Omega)$. Moreover, let Assumptions 6.1–6.3 hold and let the linearized Slater condition be satisfied at \bar{u} . Then there exist a function $p \in W^{1,q'}_{\Gamma}(\Omega)$ and a Borel measure $\mu \in \mathcal{M}(\overline{\Omega})$, such that the adjoint equation (6.20) is fulfilled in the sense of Definition 6.13 and it holds

$$\int_{\overline{\Omega}} (v(\mathbf{x}) - g(\mathbf{x}, \overline{y}(\mathbf{x}))) d\mu(\mathbf{x}) \le 0 \quad \text{for all } v \in C(\overline{\Omega}) \text{ with } v(\mathbf{x}) \le 0 \ \forall \mathbf{x} \in \overline{\Omega}, \quad (6.25)$$

$$H(\mathbf{x}, \bar{y}(\mathbf{x}), \bar{u}(\mathbf{x}), p(\mathbf{x})) = \min_{t \in U_{ad}(\mathbf{x}, \bar{u}(\mathbf{x}))} H(\mathbf{x}, \bar{y}(\mathbf{x}), t, p(\mathbf{x})) \quad a.e. \text{ on } \Gamma,$$
(6.26)

where $U_{ad}(\mathbf{x}, \bar{u}(\mathbf{x})) := [\max\{u_{\min}(\mathbf{x}), \bar{u}(\mathbf{x}) - \varepsilon\}, \min\{u_{\max}(\mathbf{x}), \bar{u}(\mathbf{x}) + \varepsilon\}].$

In [8], it is shown that the Lagrange multiplier associated to the state constraints is concentrated in the Borel set { $x \in \overline{\Omega}$: $g(x, \overline{y}(x)) = 0$ }, such that $\mu_{\partial\Omega\setminus\Gamma} = 0$ since g(x, 0) < 0 on $\partial\Omega\setminus\Gamma$ according to Assumption 6.3. Hence, we obtain homogeneous Dirichlet boundary conditions on $\partial\Omega\setminus\Gamma$ also in the adjoint equation (6.20). The firstorder necessary conditions in Theorem 6.22 can also be formulated in terms of the Lagrangian (see [12] for details).

In all what follows, let \bar{u} again be a fixed local optimum with associated state \bar{y} , adjoint state p, and Lagrange multiplier μ such that (6.25) and (6.26) are fulfilled. For the statement of second-order sufficient conditions accounting for strongly active sets, we have to restrict to the two dimensional case, since the underlying analysis heavily relies on the assumption that $S : L^2(\Omega) \to C(\overline{\Omega})$ (see [12, Sects. 4 and 6.3]). In view of Theorem 6.6, this is not fulfilled in the three and four dimensional case. We start with the definition of the *critical cone associated to* \bar{u} :

$$C(\bar{u}) := \{h \in L^2(\Gamma) : h \text{ satisfies conditions (a), (b), and (c)}\},\$$

where

$$h(\mathbf{x}) \begin{cases} \geq 0, & \text{if } \bar{u}(\mathbf{x}) = u_{\min}(\mathbf{x}), \\ \leq 0, & \text{if } \bar{u}(\mathbf{x}) = u_{\max}(\mathbf{x}), \\ = 0, & \text{if } \frac{\partial H}{\partial u}(x, \bar{y}(\mathbf{x}), \bar{u}(\mathbf{x}), p(\mathbf{x})) \neq 0, \end{cases}$$
(a)

$$\frac{\partial g}{\partial y}(\mathbf{x}, \bar{y}(\mathbf{x}))y(\mathbf{x}) \le 0, \quad \text{if } g(\mathbf{x}, \bar{y}(\mathbf{x})) = 0, \tag{b}$$

$$\int_{\overline{\Omega}} \frac{\partial g}{\partial y}(\mathbf{x}, \bar{y}(\mathbf{x})) y(\mathbf{x}) d\mu(\mathbf{x}) = 0$$
 (c)

and $y = S'(\bar{u})h$. Note that $y \in C(\overline{\Omega})$ in the two-dimensional case. Moreover, the derivative of the Hamiltonian is given by $\frac{\partial H}{\partial u}(x, \bar{y}, \bar{u}, p) = \frac{\partial l}{\partial u}(x, \bar{y}, \bar{u}) + p$.

Now, we are in the position to state the second-order sufficient conditions for (P). With the above results, in particular Lemma 6.8, the corresponding proof is completely analogous to the one presented in [12].

Theorem 6.23 Let d = 2, let Assumptions 6.1–6.3 be satisfied and suppose that $\bar{u} \in L^{\infty}(\Gamma)$ with associated state $\bar{y} \in W_{\Gamma}^{1,2}(\Omega) \cap C^{\alpha}(\Omega)$ is feasible for (P). Moreover, let

 $p \in W^{1,q'}_{\Gamma}(\Omega)$ and $\mu \in \mathcal{M}(\Omega)$ exist such that (6.20), (6.25), and (6.26) are satisfied. In addition, it is assumed that there are two constants $\omega, \tau > 0$ with

$$\frac{\partial^2 l}{\partial u^2}(\mathbf{x}, \bar{\mathbf{y}}(\mathbf{x}), \bar{u}(\mathbf{x})) \ge \omega \quad \text{if} \quad \left| \frac{\partial H}{\partial u}(\mathbf{x}, \bar{\mathbf{y}}(\mathbf{x}), \bar{u}(\mathbf{x}), p(\mathbf{x})) \right| \le \tau, \quad a.e. \text{ on } \Gamma \quad (6.27)$$

$$\frac{\partial^2 \mathcal{L}}{\partial u^2}(\bar{u},\mu)h^2 > 0 \quad \text{for all } h \in C(\bar{u}) \setminus \{0\}.$$
(6.28)

Then there exist ε , $\delta > 0$, *such that*

$$j(u) \ge j(\bar{u}) + \frac{\delta}{2} \|u - \bar{u}\|_{L^2(\Gamma)}^2$$

for all feasible $u \in L^{\infty}(\Gamma)$ with $||u - \overline{u}||_{L^{\infty}(\Gamma)} < \varepsilon$.

Note that, according to Lemma 6.19, \mathcal{L} is only continuously differentiable from $L^{\infty}(\Gamma)$ to \mathbb{R} . However, it is straightforward to see that Assumptions 6.1–6.3 ensure that $\frac{\partial \mathcal{L}}{\partial u}$ and $\frac{\partial^2 \mathcal{L}}{\partial u^2}$ can be extended from $L^{\infty}(\Gamma)$ to $L^2(\Gamma)$ using (6.19) and (6.21). This extension is also used in (6.28). Note further that the sufficient conditions (6.27) and (6.28) are natural in the sense that they are comparatively close to the necessary optimality conditions (see [12, Remark 4.2] for details).

Remark 6.24 We point out that the second-order analysis can be extended to the three dimensional case if distributed controls are applied instead of boundary control, since $L^2(\Omega) \hookrightarrow W_{\Gamma}^{-1,q}(\Omega)$ and thus continuous states are obtained with controls in $L^2(\Omega)$ (see [12, Theorem 4.1]). Nevertheless, up to the authors' best knowledge, there is no proof of second-order conditions accounting for strongly active sets in case of pointwise state constraints and boundary controls in three dimensions.

References

- 1. Alibert, J.-J., Raymond, J.-P.: Boundary control of semilinear elliptic equations with discontinuous leading coefficients and unbounded controls. Numer. Funct. Anal. Optim. **18**, 235–250 (1997)
- Amann, H.: Dynamic theory of quasilinear parabolic equations: Abstract evolution equations. Nonlinear Anal. Theory Methods Appl. 12, 895–919 (1988)
- Amann, H.: Nonhomogeneous linear and quasilinear elliptic and parabolic boundary value problems. In: Schmeisser, H.-J. (ed.) Function Spaces, Differential Operators and Nonlinear Analysis. Teubner-Texte zur Mathematik, vol. 133, pp. 9–126. Teubner, Stuttgart (1993)
- Auscher, P., Tchamitchian, P.: Square root problem for divergence operators and related topics, Astérisque, 249 (1998)
- Bandelow, U., Kaiser, H.-C., Koprucki, T., Rehberg, J.: Modeling and simulation of strained quantum wells in semiconductor lasers. In: Jäger, W., Krebs, H.-J. (eds.) Mathematics—Key Technology for the Future. Joint Projects Between Universities and Industry, pp. 377–390. Springer, Berlin/Heidelberg (2003)
- Berestycki, H., Hamel, F., Roques, L.: Analysis of the periodically fragmented environment model: I-species persistence. J. Math. Biol. 51, 75–113 (2005)
- Bonnans, J., Casas, E.: An extension of Pontryagin's principle for state-constrained optimal control of semilinear elliptic equations and variational inequalities. SIAM J. Control Optim. 33, 274–298 (1995)
- Casas, E.: Boundary control of semilinear elliptic equations with pointwise state constraints. SIAM J. Control Optim. 31, 993–1006 (1993)

- Casas, E., Mateos, M.: Uniform convergence of the FEM. Applications to state constrained control problems. Comput. Appl. Math. 21, 67–100 (2002)
- Casas, E., Mateos, M.: Second order optimality conditions for semilinear elliptic control problems with finitely many state constraints. SIAM J. Control Optim. 40, 1431–1454 (2002)
- Casas, E., Tröltzsch, F., Unger, A.: Second order sufficient optimality conditions for some stateconstrained control problems of semilinear elliptic equations. SIAM J. Control Optim. 38, 1369–1391 (2000)
- Casas, E., Tröltzsch, F., de los Reyes, J.C.: Sufficient second-order optimality conditions for semilinear control problems with pointwise state constraints. SIAM J. Optim. 19, 616–643 (2008)
- 13. Chen, Y.Z., Wu, L.C.: Second Order Elliptic Equations and Elliptic Systems. Translations of Mathematical Monographs, vol. 174. Am. Math. Soc., Providence (1998)
- Ciarlet, P.G.: The Finite Element Method for Elliptic Problems. Studies in Mathematics and Its Applications. North-Holland, Amsterdam (1979)
- De Giorgi, E.: Sulla differenziabilita e l'analiticita delle estremali degli integrali multipli regolari. Mem. Accad. Sci. Torino, P.I. III. 3, 25–43 (1957)
- Dieudonné, J.: Grundzüge der Modernen Analysis, vol. 1. VEB Deutscher Verlag der Wissenschaften, Berlin (1971)
- Duderstadt, F., Hömberg, D., Khludnev, A.M.: A mathematical model for impulse resistance welding. Math. Methods Appl. Sci. 26, 717–737 (2003)
- Elschner, J., Rehberg, J., Schmidt, G.: Optimal regularity for elliptic transmission problems including C¹ interfaces. Interfaces Free Bound. 9, 233–252 (2007)
- Evans, L.C., Gariepy, R.F.: Measure Theory and Fine Properties of Functions. Studies in Advanced Mathematics. CRC, Boca Raton (1992)
- Franzone, P.C., Guerri, L., Rovida, S.: Wavefront propagation in an activation model of the anisotropic cardiac tissue: asymptotic analysis and numerical simulation. J. Math. Biol. 28, 121–176 (1990)
- Gajewski, H.: Analysis und Numerik von Ladungstransport in Halbleitern (Analysis and numerics of carrier transport in semiconductors). Mitt. Ges. Angew. Math. Mech. 16, 35–57 (1993)
- Gajewski, H., Gröger, K., Zacharias, K.: Nichtlineare Operatorgleichungen und Operatordifferentialgleichungen. Akademie-Verlag, Berlin (1974)
- 23. Giusti, E.: Metodi Diretti nel Calcolo delle Variazioni. Unione Matematica Italiana, Bologna (1994)
- Glitzky, A., Hünlich, R.: Global estimates and asymptotics for electro-reaction-diffusion systems in heterostructures. Appl. Anal. 66, 205–226 (1997)
- Griepentrog, J.A.: Linear elliptic boundary value problems with non-smooth data: Campanato spaces of functionals. Math. Nachr. 243, 19–42 (2002)
- Griepentrog, J.A., Recke, L.: Linear elliptic boundary value problems with non-smooth data: Normal solvability on Sobolev-Campanato spaces. Math. Nachr. 225, 39–74 (2001)
- Griepentrog, J.A., Kaiser, H.-C., Rehberg, J.: Heat kernel and resolvent properties for second order elliptic differential operators with general boundary conditions on L^p. Adv. Math. Sci. Appl. 11, 87–112 (2001)
- Griepentrog, J.A., Gröger, K., Kaiser, H.C., Rehberg, J.: Interpolation for function spaces related to mixed boundary value problems. Math. Nachr. 241, 110–120 (2002)
- 29. Grisvard, P.: Elliptic Problems in Nonsmooth Domains. Pitman, Boston (1985)
- Gröger, K.: A W^{1,p}-estimate for solutions to mixed boundary value problems for second order elliptic differential equations. Math. Ann. 283, 679–687 (1989)
- 31. Gröger, K., Rehberg, J.: Resolvent estimates in $W^{-1,p}$ for second order elliptic differential operators in case of mixed boundary conditions. Math. Ann. **285**, 105–113 (1989)
- 32. Haller-Dintelmann, R., Rehberg, J.: Maximal parabolic regularity for divergence operators including mixed boundary conditions, WIAS-preprint 1288 (2008)
- Haller-Dintelmann, R., Kaiser, H.-C., Rehberg, J.: Elliptic model problems including mixed boundary conditions and material heterogeneities. J. Math. Pures Appl. (9) 89(1), 25–48 (2008)
- 34. Kato, T.: Perturbation Theory for Linear Operators, Classics in Mathematics. Springer, Berlin (1980) (Reprint of the corr. print. of the 2nd edn.)
- Kinderlehrer, D., Stampacchia, G.: An Introduction to Variational Inequalities and Their Applications. Pure and Applied Mathematics, vol. 88. Academic Press, New York (1980)
- Koprucki, T., Kaiser, H.-C., Fuhrmann, J.: Electronic states in semiconductor nanostructures and upscaling to semi-classical models. In: Mielke, A. (ed.) Analysis, Modeling and Simulation of Multiscale Problems, pp. 367–396. Springer, Berlin (2006)
- Ladyzhenskaya, O.A., Ural'tseva, N.N.: Linear and Quasilinear Elliptic Equations. Mathematics in Science and Engineering. Academic Press, New York (1968)

- Leguillon, D., Sanchez-Palenzia, E.: Computation of Singular Solutions in Elliptic Problems and Elasticity. Wiley, Chichester (1987)
- Li, Y., Liu, J., Voskoboynikov, O., Lee, C., Sze, S.: Electron energy level calculations for cylindrical narrow gap semiconductor quantum dot. Comput. Phys. Commun. 140, 399–404 (2001)
- 40. Liebermann, G.M.: Mixed boundary value problems for elliptic and parabolic differential equations of second order. J. Math. Anal. Appl. **113**, 422–440 (1986)
- Liebermann, G.M.: Optimal Hölder regularity for mixed boundary value problems. J. Math. Anal. Appl. 143, 572–586 (1989)
- 42. Maz'ya, V.: Sobolev Spaces. Springer, Berlin (1985)
- Mitrea, I., Mitrea, M.: The Poisson problem with mixed boundary conditions in Sobolev and Besov spaces in non-smooth domains. Trans. Am. Math. Soc. 359, 4143–4182 (2007)
- 44. Prignet, A.: Remarks on existence and uniqueness of solutions of elliptic problems with right-hand side measures. Rendiconti di Mat. **15**, 321–337 (1995)
- Prüss, J.: Maximal regularity for evolution equations in L^p-spaces. Conf. Semin. Mat. Univ. Bari 285, 1–39 (2002)
- 46. Selberherr, S.: Analysis and Simulation of Semiconductors. Springer, Wien (1984)
- Serrin, J.: Pathological solutions of elliptic differential equations. Ann. Scuola Norm. Sup. Pisa Cl. Sci. 18, 385–387 (1964)
- Sobolevskij, P.E.: Equations of parabolic type in a Banach space. Am. Math. Soc. Transl. II 49, 1–62 (1966)
- Sommerfeld, A.: Electrodynamics. Lectures on Theoretical Physics, vol. III. Academic Press, New York (1952)
- Sommerfeld, A.: Thermodynamics and Statistical Mechanics. Lectures on Theoretical Physics, vol. V. Academic Press, New York (1956)
- Stampacchia, G.: Problemi al contorno ellittici, con dati discontinui, dotati di soluzioni hölderiane. Ann. Mat. Pura Appl. IV 51, 1–37 (1960)
- Stampacchia, G.: Le problème de Dirichlet pour les équations elliptiques du second ordre à coefficients discontinus. Ann. Inst. Fourier 15, 189–258 (1965)
- Triebel, H.: Interpolation Theory, Function Spaces, Differential Operators. North-Holland, Amsterdam (1978)
- 54. Triebel, H.: On spaces of $B^s_{\infty,q}$ and C^s type. Math. Nachr. 85, 75–90 (1978)
- 55. Tukia, P.: The planar Schönflies theorem for Lipschitz maps. Ann. Acad. Sci. Fenn. Ser. A I 5, 49–72 (1980)
- Wang, W., Hwang, T., Lin, W., Liu, J.: Numerical methods for semiconductor heterostructures with band nonparabolicity. J. Comput. Phys. 190, 141–158 (2003)
- 57. Weisbuch, C., Vinter, B.: Quantum Semiconductor Structures: Fundamentals and Applications. Academic Press, Boston (1991)
- 58. Ziemer, W.P.: Weakly Differentiable Functions. Springer, Berlin (1989)