# $W^{2, p}$-A PRIORI ESTIMATES FOR THE NEUTRAL POINCARÉ PROBLEM* 

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To the memory of Filippo Chiarenza


#### Abstract

A degenerate oblique derivative problem is studied for uniformly elliptic operators with low regular coefficients in the framework of Sobolev's classes $W^{2, p}(\Omega)$ for arbitrary $p>1$. The boundary operator is prescribed in terms of a directional derivative with respect to the vector field $\boldsymbol{\ell}$ that becomes tangential to $\partial \Omega$ at the points of some non-empty subset $\mathcal{E} \subset \partial \Omega$ and is directed outwards $\Omega$ on $\partial \Omega \backslash \mathcal{E}$. Under quite general assumptions of the behaviour of $\ell$, we derive a priori estimates for the $W^{2, p}(\Omega)$-strong solutions for any $p \in(1, \infty)$.


## Introduction

The lecture deals with regularity in Sobolev's spaces $W^{2, p}(\Omega), \forall p \in(1, \infty)$, of the strong solutions to the oblique derivative problem

$$
\begin{cases}\mathcal{L} u:=a^{i j}(x) D_{i j} u=f(x) & \text { a.e. } \Omega,  \tag{1}\\ \mathcal{B} u:=\partial u / \partial \ell=\varphi(x) & \text { on } \partial \Omega\end{cases}
$$

where $\mathcal{L}$ is a uniformly elliptic operator with low regular coefficients and $\mathcal{B}$ is prescribed in terms of a directional derivative with respect to the unit vector field $\ell(x)=\left(\ell^{1}(x), \ldots, \ell^{n}(x)\right)$ defined on $\partial \Omega, n \geq 3$. Precisely, we are interested in the Poincaré problem (1) (cf. [17, 20, 16]), that is, a situation when $\ell(x)$ becomes tangential to $\partial \Omega$ at the points of a non-empty subset $\mathcal{E}$ of $\partial \Omega$.

From a mathematical point of view, (1) is not an elliptic boundary value problem. In fact, it follows from the general PDEs theory that (1) is a regular (elliptic) problem if and only if the Shapiro-Lopatinskij complementary condition is satisfied which means $\ell$ must be transversal to $\partial \Omega$ when $n \geq 3$ and $|\ell| \neq 0$ as $n=2$. If $\boldsymbol{\ell}$ is tangent to $\partial \Omega$ then (1) is a degenerate problem and new effects occur in contrast to the regular case. It turns out that the qualitative properties of (1) depend on the behaviour of $\boldsymbol{\ell}$ near the set of tangency $\mathcal{E}$ and especially on the way the normal component $\gamma \boldsymbol{\nu}$ of $\boldsymbol{\ell}$ (with respect to the outward normal $\boldsymbol{\nu}$ to $\partial \Omega$ ) changes or no its sign on the trajectories of $\boldsymbol{\ell}$ when these cross $\mathcal{E}$. The main results were obtained by Hörmander [6], Egorov and Kondrat'ev [2], Maz'ya [8], Maz'ya and Paneah [9], Melin and Sjöstrand [10], Paneah [15] and good surveys and details can be found in

[^0]Popivanov and Palagachev [20] and Paneah [16]. The problem (1) has been studied in the framework of Sobolev spaces $H^{s}\left(\equiv H^{s, 2}\right)$ assuming $C^{\infty}$-smooth data and this naturally involved techniques from the pseudo-differential calculus.

The simplest case arises when $\gamma:=\boldsymbol{\ell} \cdot \boldsymbol{\nu}$, even if zero on $\mathcal{E}$, conserves the sign on $\partial \Omega$. Then $\mathcal{E}$ and $\ell$ are of neutral type (a terminology coming from the physical interpretation of (1) in the theory of Brownian motion, see [20]) and (1) is a problem of Fredholm type (cf. [2]). Assume now that $\gamma$ changes the sign from "-" to "+" in positive direction along the $\ell$-integral curves passing through the points of $\mathcal{E}$. Then $\boldsymbol{\ell}$ is of emergent type and $\mathcal{E}$ is called attracting manifold. The new effect appearing now is that the kernel of (1) is infinite-dimensional ([6]) and to get a well-posed problem one has to modify (1) by prescribing the values of $u$ on $\mathcal{E}$ (cf. [2]). Finally, suppose the sign of $\gamma$ changes from " + " to "-" along the $\boldsymbol{\ell}$-trajectories. Now $\boldsymbol{\ell}$ is of submergent type and $\mathcal{E}$ corresponds to a repellent manifold. The problem (1) has infinite-dimensional cokernel ([6]) and Maz'ya and Paneah [9] were the first to propose a relevant modification of (1) by violating the boundary condition at the points of $\mathcal{E}$. As consequence, a Fredholm problem arises, but the restriction $\left.u\right|_{\partial \Omega}$ has a finite jump at $\mathcal{E}$. What is the common feature of the degenerate problems, independently of the type of $\ell$, is that the solution "loses regularity" near the set of tangency from the data of (1) in contrast to the non-degenerate case when any solution gains two derivatives from $f$ and one derivative from $\varphi$. Roughly speaking, that loss of smoothness depends on the order of contact between $\ell$ and $\partial \Omega$ and is given by the subelliptic estimates obtained for the solutions of degenerate problems (cf. $[4,5,6,9]$ ). Precisely, if $\ell$ has a contact of order $k$ with $\partial \Omega$ then the solution of (1) gains $2-k /(k+1)$ derivatives from $f$ and $1-k /(k+1)$ derivatives from $\varphi$.


(c) submergent vector field $\ell$

For what concerns the geometric structure of $\mathcal{E}$, it was supposed initially to be a submanifold of $\partial \Omega$ of codimension one. Melin and Sjöstrand [10] and Paneah [15] were the first to study the Poincaré problem (1) in a more general situation when $\mathcal{E}$ is a massive subset of $\partial \Omega$ with positive surface measure, allowing $\mathcal{E}$ to contain arcs of $\boldsymbol{\ell}$-trajectories of finite length. Their results were extended by Winzell ([21, 22]) to the framework of Hölder's spaces who studied (1) assuming $C^{1, \alpha}$-smoothness of the coefficients of $\mathcal{L}$. It is worth noting that $\ell$ has automatically an infinite order of contact with $\partial \Omega$ when $\mathcal{E}$ is a massive subset of the boundary.

To deal with non-linear Poincaré problems, however, we have to dispose of precise information on the linear problem (1) with coefficients less regular than $C^{\infty}$ (see [11, 18, 19, 20]). Indeed, a priori estimates in $W^{2, p}$ for solutions to (1) would imply easily pointwise estimates for $u$ and $D u$ for suitable values of $p>1$ through the Sobolev imbeddings. This way, we are naturally led to consider the problem (1) in a strong sense, that is, to searching for solutions lying in $W^{2, p}$ which satisfy $\mathcal{L} u=f$ almost everywhere (a.e.) in $\Omega$ and $\mathcal{B} u=\varphi$ holds in the sense of trace on $\partial \Omega$.

In the papers $[4,5]$ by Guan and Sawyer solvability and precise subelliptic estimates have been obtained for (1) in $H^{s, p}$-spaces ( $\equiv W^{s, p}$ for integer $s!$ ). However, [4] treats operators with $C^{\infty}$-coefficients and this determines the technique involved and the results obtained, while in [5] the coefficients are $C^{0, \alpha}$-smooth, but the field $\ell$ is of finite type, that is, it has a finite order of contact with $\partial \Omega$.

The main goal of this lecture is to derive a priori estimates in Sobolev's classes $W^{2, p}(\Omega)$ with any $p \in(1, \infty)$ for the solutions of the Poincaré problem (1), weakening both Winzell's assumptions on $C^{1, \alpha}$-regularity of the coefficients of $\mathcal{L}$ and these of Guan and Sawyer on the finite type of $\boldsymbol{\ell}$. We are dealing with the simpler case when $\gamma$ preserves its sign on $\partial \Omega$ which means the field $\ell$ is of neutral type. Of course, the loss of smoothness mentioned, imposes some more regularity of the data near the set $\mathcal{E}$. We assume the coefficients of $\mathcal{L}$ to be Lipschitz continuous near $\mathcal{E}$ while only continuity (and even discontinuity controlled in $V M O$ ) is allowed away from $\mathcal{E}$. Similarly, $\boldsymbol{\ell}$ is a Lipschitz vector field on $\partial \Omega$ with Lipschitz continuous first derivatives near $\mathcal{E}$, and no restrictions on the order of contact with $\partial \Omega$ are required. Regarding the tangency set $\mathcal{E}$, it may have positive surface measure and is restricted only to a sort of non-trapping condition that all trajectories of $\boldsymbol{\ell}$ through the points of $\mathcal{E}$ are non-closed and leave $\mathcal{E}$ in a finite time.

The technique adopted is based on a dynamical system approach employing the fact that $\partial u / \partial \ell$ is a local strong solution, near $\mathcal{E}$, to a Dirichlet-type problem with right-hand side depending on the solution $u$ itself. Application of the $L^{p}$-estimates for such problems leads to the functional inequality (26) for suitable $W^{2, p}$-norms of $u$ on a family of subdomains which, starting away from $\mathcal{E}$, evolve along the $\boldsymbol{\ell}$ trajectories and exhaust a sort of their tubular neighbourhoods. Fortunately, that is an inequality with advanced argument and the desired $W^{2, p}$-estimate follows by iteration with respect to the curvilinear parameter on the trajectories of $\boldsymbol{\ell}$. Another advantage of this approach is the improving-of-integrability property obtained for the solutions of (1). Roughly speaking, it asserts that the problem (1), even if a degenerate one, behaves as an elliptic problem for what concerns the degree $p$ of integrability. In other words, the second derivatives of any solution to (1) will have the same rate of integrability as $f$ and $\varphi$. We refer the reader to the paper [14] for outgrowths of the $W^{2, p}$-a priori estimates, such as uniqueness in $W^{2, p}(\Omega), \forall p>1$, of the strong solutions to (1) as well as its Fredholmness.

Concluding this introduction, we should mention the article [13] where similar results have been obtained by different technique in the particular case when the tangency set $\mathcal{E}$ contains trajectories of $\ell$ with positive, but small enough lengths.

## 1. Hypotheses and the Main Result

Hereafter $\Omega \subset \mathbb{R}^{n}, n \geq 3$, will be a bounded domain with reasonably smooth boundary and $\boldsymbol{\nu}(x)=\left(\overline{\nu^{1}}(x), \ldots, \nu^{n}(x)\right)$ stands for the unit outward normal to $\partial \Omega$ at $x \in \partial \Omega$. Consider a unit vector field $\ell(x)=\left(\ell^{1}(x), \ldots, \ell^{n}(x)\right)$ on $\partial \Omega$ and let $\boldsymbol{\ell}(x)=\boldsymbol{\tau}(x)+\gamma(x) \boldsymbol{\nu}(x)$, where $\boldsymbol{\tau}: \partial \Omega \rightarrow \mathbb{R}^{n}$ is the projection of $\boldsymbol{\ell}(x)$ on the hyperplane tangent to $\partial \Omega$ at $x \in \partial \Omega$ and $\gamma: \partial \Omega \rightarrow \mathbb{R}$ is the inner product $\gamma(x):=$ $\boldsymbol{\ell}(x) \cdot \boldsymbol{\nu}(x)$. The set of zeroes of $\gamma$,

$$
\mathcal{E}:=\{x \in \partial \Omega: \quad \gamma(x)=0\}
$$

is indeed the subset of $\partial \Omega$ where the field $\ell(x)$ becomes tangent to it.
Fix $\mathcal{N} \subset \bar{\Omega}$ to be a closed neighbourhood of $\mathcal{E}$ in $\bar{\Omega}$. We suppose $\mathcal{L}$ is a uniformly elliptic operator with measurable coefficients, satisfying

$$
\begin{equation*}
\lambda^{-1}|\xi|^{2} \leq a^{i j}(x) \xi_{i} \xi_{j} \leq \lambda|\xi|^{2} \quad \text { a.a. } x \in \Omega, \forall \xi \in \mathbb{R}^{n}, \quad a^{i j}(x)=a^{j i}(x) \tag{2}
\end{equation*}
$$

for some positive constant $\lambda$. Regarding the regularity of the data, we assume

$$
\left\{\begin{array}{l}
a^{i j} \in V M O(\Omega) \cap C^{0,1}(\mathcal{N}),  \tag{3}\\
\partial \Omega \in C^{1,1}, \quad \partial \Omega \cap \mathcal{N} \in C^{2,1}, \quad \ell^{i} \in C^{0,1}(\partial \Omega) \cap C^{1,1}(\partial \Omega \cap \mathcal{N})
\end{array}\right.
$$

with $\operatorname{VMO}(\Omega)$ being the Sarason class of functions of vanishing mean oscillation and $C^{k, 1}$ denotes the space of functions with Lipschitz continuous $k$-th order derivatives. Let us point out that (2), (3) and the Rademacher theorem give $a^{i j} \in L^{\infty}(\Omega) \cap$ $W^{1, \infty}(\mathcal{N})$. For what concerns the boundary operator $\mathcal{B}$, we assume
(4) $\left\{\begin{array}{l}\gamma(x)=\boldsymbol{\ell}(x) \cdot \boldsymbol{\nu}(x) \geq 0 \quad \forall x \in \partial \Omega, \quad \text { and } \\ \text { the arcs of the } \boldsymbol{\ell} \text {-trajectories lying in } \mathcal{E} \text { (which coincide with these of } \boldsymbol{\tau} \text { ) } \\ \text { are all non-closed and of finite lengths. }\end{array}\right.$

The first assumption simply means that $\ell(x)$ is either tangential to $\partial \Omega$ or is directed outwards $\Omega$, that is, the field $\ell$ is of neutral type on $\partial \Omega$, while the second one is a sort of non-trapping condition on the tangency set $\mathcal{E}$. It implies that the $\boldsymbol{\ell}$-integral curves leave $\mathcal{E}$ in a finite time in both directions.


Figure 1. The set of tangency $\mathcal{E}$ is the union $\mathcal{E}_{1} \cup \mathcal{E}_{2} \cup \mathcal{E}_{3}$ where $\operatorname{codim}_{\partial \Omega} \mathcal{E}_{1}=\operatorname{codim}{ }_{\partial \Omega} \mathcal{E}_{2}=1$ while meas ${ }_{\partial \Omega} \mathcal{E}_{3}>0$. The vector field $\boldsymbol{\ell}$ is transversal to $\mathcal{E}_{1}$ and tangent to $\mathcal{E}_{2}$. Actually, $\mathcal{E}_{2}$ consists of an arc of $\boldsymbol{\tau}$-trajectory, whereas $\mathcal{E}_{3}$ is union of such arcs.

Throughout the text $W^{k, p}$ stands for the Sobolev class of functions with $L^{p_{-}}$ summable weak derivatives up to order $k \in \mathbb{N}$ while $W^{s, p}(\partial \Omega)$ with $s>0$ noninteger and $p \in(1,+\infty)$, is the Sobolev space of fractional order on $\partial \Omega$. Further, we use the standard parameterization $t \mapsto \boldsymbol{\psi}_{\boldsymbol{L}}(t ; x)$ for the trajectory (equivalently, phase curve, maximal integral curve) of a given vector field $\boldsymbol{L}$ passing through a point $x$, that is, $\partial_{t} \boldsymbol{\psi}_{\boldsymbol{L}}(t ; x)=\boldsymbol{L}\left(\boldsymbol{\psi}_{\boldsymbol{L}}(t ; x)\right)$ and $\boldsymbol{\psi}_{\boldsymbol{L}}(0 ; x)=x$.

We will employ below an extension of the field $\ell$ near $\partial \Omega$ which preserves therein its regularity and geometric properties. All the results and proofs in the sequel work for such an arbitrary $\ell$-extension but, in order to make more evident some geometric constructions, we prefer to introduce a special extension as follows. For


Figure 2
each $x \in \mathbb{R}^{n}$ near $\partial \Omega$ set $d(x)=\operatorname{dist}(x, \partial \Omega)$ and define $\Gamma:=\left\{x \in \mathbb{R}^{n}: d(x) \leq d_{0}\right\}$ with small $d_{0}>0$. Letting $\Omega_{0}:=\Omega \backslash \Gamma$ and $y(x) \in \partial \Omega$ for the unique point closest to $x \in \Gamma$, we have (see [3, Chapter 14]) $y(x) \in C^{0,1}(\Gamma)$ while $y(x) \in C^{1,1}$ near $\mathcal{E}$. Regarding the distance function $d(x)=|x-y(x)|$, it is Lipschitz continuous in $\Gamma$ and inherits the regularity of $\partial \Omega$ at $y(x)$ when considered on the parts of $\Gamma$ lying in/out $\Omega$, but its normal derivative has a finite jump on $\partial \Omega$. Anyway, it is a routine to check $(d(x))^{2} \in C^{1,1}(\Gamma)$. Setting $\boldsymbol{L}(x)$ for the normalized representative of $\boldsymbol{\ell}(y(x))+(d(x))^{2} \boldsymbol{\nu}(y(x)) \forall x \in \Gamma$, it results $|\boldsymbol{L}(x)|=1,\left.\boldsymbol{L}\right|_{\partial \Omega}=\boldsymbol{\ell},\left.\boldsymbol{L}\right|_{\mathcal{E}}=\boldsymbol{\tau}$ and $\boldsymbol{L} \in C^{0,1}(\Gamma) \cap C^{1,1}(\Gamma \cap \mathcal{N})$. Moreover, the field $\boldsymbol{L}$ is strictly transversal to $\partial \Omega_{0}$.

As consequence of the non-trapping condition (4), the compactness of $\mathcal{E}$ and the semi-continuity properties of the lengths of the $\boldsymbol{\tau}$-maximal integral curves, it is not hard to get that (see [22, Proposition 3.1] and [20, Proposition 3.2.4]) under the hypotheses (3) and (4), there is a finite upper bound $\kappa_{0}$ for the arclengths of the $\boldsymbol{\tau}$-trajectories lying in $\mathcal{E}$. Moreover, each point of $\Gamma$ can be reached from $\partial \Omega_{0}$ by an L-integral curve of length at most $\kappa=$ const $>0$.

In what follows, the letter $C$ will denote a generic constant depending on known quantities defined by the data of (1), that is, on $n, p, \lambda$, the respective norms of the coefficients of $\mathcal{L}$ and $\mathcal{B}$ in $\Omega$ and $\mathcal{N}$, the regularity of $\partial \Omega$ and the constants $\kappa_{0}$ and $\kappa$.

In order to control precisely the regularity of $u$ near the tangency set $\mathcal{E}$, we have to introduce the appropriate functional spaces. For, take an arbitrary $p \in(1, \infty)$ and define the Banach spaces

$$
\mathcal{F}^{p}(\Omega, \mathcal{N}):=\left\{f \in L^{p}(\Omega): \partial f / \partial \boldsymbol{L} \in L^{p}(\mathcal{N})\right\}
$$

equipped with norm $\|f\|_{\mathcal{F}^{p}(\Omega, \mathcal{N})}:=\|f\|_{L^{p}(\Omega)}+\|\partial f / \partial \boldsymbol{L}\|_{L^{p}(\mathcal{N})}$, and

$$
\Phi^{p}(\partial \Omega, \mathcal{N}):=\left\{\varphi \in W^{1-1 / p, p}(\partial \Omega): \varphi \in W^{2-1 / p, p}(\partial \Omega \cap \mathcal{N})\right\}
$$

normed by $\|\varphi\|_{\Phi^{p}(\partial \Omega, \mathcal{N})}:=\|\varphi\|_{W^{1-1 / p, p}(\partial \Omega)}+\|\varphi\|_{W^{2-1 / p, p}(\partial \Omega \cap \mathcal{N})}$.
Our main result asserts that the couple $(\mathcal{L}, \mathcal{B})$ improves the integrability of solutions to (1) for any $p$ in the range $(1, \infty)$ and, moreover, provides for an a priori estimate in the $L^{p}$-Sobolev scales for any such solution.

Theorem 1. Under the hypotheses (2)-(4) let $u \in W^{2, p}(\Omega)$ be a strong solution of the problem (1) with $f \in \mathcal{F}^{q}(\Omega, \mathcal{N})$ and $\varphi \in \Phi^{q}(\partial \Omega, \mathcal{N})$ where $1<p \leq q<\infty$.

Then $u \in W^{2, q}(\Omega)$ and there is an absolute constant $C$ such that

$$
\begin{equation*}
\|u\|_{W^{2, q}(\Omega)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}\right) . \tag{5}
\end{equation*}
$$

Let us point out reader's attention that the directional derivative $\partial u / \partial \boldsymbol{L}$ of each $W^{2, p}$-solution to (1) belongs to $W^{2, p}(\mathcal{N})$. For, $\partial u / \partial \boldsymbol{L} \in W^{1, p}(\mathcal{N})$ and taking the difference quotients in (1) in the direction of $\boldsymbol{L}$ (cf. [3, Chapter 8 and Lemma 7.24]) gives that $\partial u / \partial \boldsymbol{L} \in W^{2, p}(\mathcal{N})$ is a strong local solution to the Dirichlet problem
(6) $\left\{\begin{array}{l}\mathcal{L}\left(\frac{\partial u}{\partial \boldsymbol{L}}\right)=\frac{\partial f}{\partial \boldsymbol{L}}+2 a^{i j} D_{j} L^{k} D_{k i} u+a^{i j} D_{i j} L^{k} D_{k} u-\frac{\partial a^{i j}}{\partial \boldsymbol{L}} D_{i j} u \quad \text { a.e. } \mathcal{N}, \\ \frac{\partial u}{\partial \boldsymbol{L}}=\varphi \quad \text { on } \partial \Omega \cap \mathcal{N}\end{array}\right.$
where $\boldsymbol{L}(x)=\left(L^{1}(x), \ldots, L^{n}(x)\right) \in C^{1,1}(\mathcal{N})$. Therefore, once having proved $u \in$ $W^{2, q}(\Omega)$ and the estimate (5), we have

$$
\|\partial u / \partial \boldsymbol{L}\|_{W^{2, q}(\tilde{\mathcal{N}})} \leq C^{\prime}\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}\right)
$$

for any closed neighbourhood $\widetilde{\mathcal{N}}$ of $\mathcal{E}$ in $\bar{\Omega}, \widetilde{\mathcal{N}} \subset \mathcal{N}$, by means of the $L^{p}$-theory of uniformly elliptic equations (see [1] or [3, Chapter 9]). In other words, if a strong solution $u$ to (1) belongs to $W^{2, q}(\Omega)$ then $\partial u / \partial \boldsymbol{L} \in W^{2, q}(\mathcal{N})$ automatically, provided $f \in \mathcal{F}^{q}(\Omega, \mathcal{N})$ and $\varphi \in \Phi^{q}(\partial \Omega, \mathcal{N})$.

## 2. Proof of Theorem 1

Fix hereafter $\mathcal{N}^{\prime} \subset \mathcal{N}^{\prime \prime} \subset \mathcal{N}$ to be closed neighbourhoods of $\mathcal{E}$ in $\bar{\Omega}$ with $\mathcal{N}^{\prime \prime}$ so "narrow" that $\mathcal{N}^{\prime \prime} \subset \Omega \backslash \Omega_{0}$ (see Figure 3). The next result is an immediate consequence of $\gamma(x)>0 \forall x \in \partial \Omega \backslash \mathcal{N}^{\prime}$ and the $L^{p}$-theory of regular oblique derivative problems for uniformly elliptic operators with $V M O$ principal coefficients (cf. [7, Theorem 2.3.1]).

Proposition 2. Assume (2), (3) and $\gamma(x)>0 \forall x \in \Omega \backslash \mathcal{E}$, and let $u \in W^{2, p}(\Omega)$ be a solution to (1) with $f \in L^{q}(\Omega)$ and $\varphi \in W^{1-1 / q, q}(\partial \Omega)$, where $1<p \leq q<\infty$.

Then $u \in W^{2, q}\left(\Omega \backslash \mathcal{N}^{\prime}\right)$ and there is a constant such that

$$
\begin{equation*}
\|u\|_{W^{2, q}\left(\Omega \backslash \mathcal{N}^{\prime}\right)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{L^{q}(\Omega)}+\|\varphi\|_{W^{1-1 / q, q}(\partial \Omega)}\right) . \tag{7}
\end{equation*}
$$

To derive the improving-of-integrability near the tangency set $\mathcal{E}$, we consider any solution of the problem (1) for which $a^{i j}, \partial a^{i j} / \partial \boldsymbol{L} \in L^{\infty}(\mathcal{N})$ in view of $(3)^{1}$ and $f$, $\partial f / \partial \boldsymbol{L} \in L^{q}(\mathcal{N})$ and $\varphi \in W^{2-1 / q, q}(\partial \Omega \cap \mathcal{N})$ by hypotheses.

Lemma 3. Under the assumptions of Theorem 1, the solution $u$ of (1) belongs to $u \in W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)$ and there is a constant such that

$$
\begin{equation*}
\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}\right) \tag{8}
\end{equation*}
$$

Proof. Take an arbitrary point $x_{0} \in \mathcal{E}$. According to (4), the $\boldsymbol{L}$-trajectory through $x_{0}$ leaves $\mathcal{E}$ in both directions for a finite time, that is, $\boldsymbol{\psi}_{\boldsymbol{L}}\left(t^{-} ; x_{0}\right) \in \mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}$, $\psi_{\boldsymbol{L}}\left(t^{+} ; x_{0}\right) \in \mathbb{R}^{n} \backslash \bar{\Omega}$ (see Figure 3) for suitable $t^{-}<0<t^{+}$.


Figure 3. $\mathcal{T}_{r}$ is the dotted set, while the double-dotted one is $\mathcal{P}_{r, T}$.
Set $\mathcal{H}$ for the $(n-1)$-dimensional hyperplane through $x_{0}$ and orthogonal to $\boldsymbol{L}\left(x_{0}\right)$, and define

$$
B_{r}\left(x_{0}\right):=\left\{x \in \mathcal{H}: \quad\left|x-x_{0}\right|<r\right\}
$$

with $r>0$ to be chosen later. It follows from the Picard inequality ${ }^{2}$ that if $r$ is small enough, then the flow of $B_{r}\left(x_{0}\right)$ along the $\boldsymbol{L}$-trajectories at time $t^{-}$,

$$
B_{r}^{\prime}\left(x_{0}\right):=\boldsymbol{\psi}_{\boldsymbol{L}}\left(t^{-} ; B_{r}\left(x_{0}\right)\right):=\left\{\boldsymbol{\psi}_{\boldsymbol{L}}\left(t^{-} ; y\right): \quad y \in B_{r}\left(x_{0}\right)\right\}
$$

is entirely contained in $\mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}$ whence $B_{r}^{\prime}\left(x_{0}\right) \cap \mathcal{E}=\emptyset$. The set

$$
\Theta_{r}:=\left\{\boldsymbol{\psi}_{\boldsymbol{L}}\left(t ; x^{\prime}\right): \quad x^{\prime} \in B_{r}^{\prime}\left(x_{0}\right), \quad t \in\left(0, t^{+}-t^{-}\right)\right\}
$$

is an $n$-dimensional neighbourhood of the $\boldsymbol{L}$-trajectory through $x_{0}$ and defining

$$
\mathcal{T}_{r}:=\Theta_{r} \cap \Omega
$$

[^1]the boundary $\partial \mathcal{T}_{r}$ is composed of the "base" $B_{r}^{\prime}\left(x_{0}\right)$ and the "lateral" components $\partial_{1} \mathcal{T}_{r}:=\partial \mathcal{T}_{r} \cap \partial \Omega$ and $\partial_{2} \mathcal{T}_{r}:=\left(\partial \mathcal{T}_{r} \cap \Omega\right) \backslash B_{r}^{\prime}\left(x_{0}\right)$. Indeed, $\mathcal{T}_{r} \subset \mathcal{N}^{\prime \prime}$ if $r>0$ is small enough.

We will derive (8) in $\mathcal{T}_{r}$ after that the desired estimate will follow by covering the compact $\mathcal{E} \subset \partial \Omega$ by a finite number of sets like $\overline{\mathcal{T}_{r}}$. Our strategy is based on a representation of $u(x)$ in $\mathcal{T}_{r}$ by means of $u\left(x^{\prime}\right)$ with $x^{\prime}=\boldsymbol{\psi}_{\boldsymbol{L}}(-\xi(x) ; x) \in B_{r}^{\prime}\left(x_{0}\right)$ for some $\xi(x)>0$, and the integral of $\partial u / \partial \boldsymbol{L}$ along the $\boldsymbol{L}$-trajectory joining $x^{\prime}$ with $x$. Thus the Sobolev norm of $u$ will be expressed by the respective norm of $\partial u / \partial \boldsymbol{L}$ and that of $u$ itself near $B_{r}^{\prime}\left(x_{0}\right)$ where we dispose of (7). Concerning $\partial u / \partial \boldsymbol{L}$, it is a local solution of Dirichlet problem near $\mathcal{E}$ with right-hand side depending on $u$.

Let $\mu: \mathcal{H} \rightarrow \mathbb{R}^{+}$be a $C^{\infty}$ cut-off function such that

$$
\mu(y)= \begin{cases}1 & y \in B_{r / 2}\left(x_{0}\right),  \tag{9}\\ 0 & y \in \mathcal{H} \backslash B_{3 r / 4}\left(x_{0}\right)\end{cases}
$$

and extend it to $\mathbb{R}^{n}$ as constant on the $\boldsymbol{L}$-trajectory through $y \in \mathcal{H}$. The function $U(x):=\mu(x) u(x)$ is a $W^{2, p}(\mathcal{N})$-solution of

$$
\left\{\begin{array}{l}
\mathcal{L} U=F(x):=\mu f+2 a^{i j} D_{j} \mu D_{i} u+u a^{i j} D_{i j} \mu \quad \text { a.e. } \mathcal{T}_{r},  \tag{10}\\
\partial U / \partial \boldsymbol{L}=\Phi:= \begin{cases}\mu \varphi & \text { on } \partial_{1} \mathcal{T}_{r}, \\
0 & \text { near } \partial_{2} \mathcal{T}_{r}, \\
\mu \partial u / \partial \boldsymbol{L} & \text { on } B_{r}^{\prime}\left(x_{0}\right) \subset \mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime} .\end{cases}
\end{array}\right.
$$

Indeed, $u \in W^{2, p}(\mathcal{N})$ implies $D u \in L^{n p /(n-p)}$ if $p<n$ and $D u \in L^{s} \forall s>1$ when $p \geq n$, whence $F \in L^{q^{\prime}}(\mathcal{N})$ with

$$
q^{\prime}:= \begin{cases}\min \left\{q, \frac{n p}{n-p}\right\} & \text { if } p<n  \tag{11}\\ q & \text { if } p \geq n\end{cases}
$$

Further, $\partial F / \partial \boldsymbol{L} \in L^{q^{\prime}}\left(\mathcal{N}^{\prime \prime}\right)$ as consequence of (6), $\partial u / \partial \boldsymbol{L} \in W^{2, q}\left(\mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}\right)$ by Proposition 2 whence $\Phi \in W^{2-1 / q, q}\left(\partial \mathcal{T}_{r}\right)$. Thus (2), (3), $\mathcal{T}_{r} \subset \mathcal{N}^{\prime \prime}$ and (6) give that

$$
V(x):=\partial U / \partial \boldsymbol{L}
$$

is a $W^{2, p}\left(\mathcal{T}_{r}\right)$-solution of the Dirichlet problem

$$
\left\{\begin{array}{l}
\mathcal{L} V=\partial F / \partial \boldsymbol{L}+2 a^{i j} D_{j} L^{k} D_{i k} U+a^{i j} D_{i j} L^{k} D_{k} U-\frac{\partial a^{i j}}{\partial \boldsymbol{L}} D_{i j} U \quad \text { a.e. } \mathcal{I}_{r}  \tag{12}\\
V=\Phi \text { on } \partial \mathcal{T}_{r} .
\end{array}\right.
$$

Now we pass from $x \in \Theta_{r}$ into the new variables $\left(x^{\prime}, \xi\right)$ with $x^{\prime}=\boldsymbol{\psi}_{\boldsymbol{L}}(-\xi(x) ; x) \in$ $B_{r}^{\prime}\left(x_{0}\right)$ and $\xi: \Theta_{r} \rightarrow\left(0, t^{+}-t^{-}\right), \xi(x) \in C^{1,1}\left(\Theta_{r}\right)$. The transform $x \mapsto\left(x^{\prime}, \xi\right)$ defines a $C^{1,1}$-diffeomorphism because the field $\boldsymbol{L}$ is transversal to $B_{r}^{\prime}\left(x_{0}\right)$. Moreover, $\partial / \partial \boldsymbol{L} \equiv \partial / \partial \xi, \boldsymbol{\psi}_{\boldsymbol{L}}\left(t ; x^{\prime}\right)=\left(x^{\prime}, t\right)$ and $V\left(x^{\prime}, \xi\right)=\partial U\left(x^{\prime}, \xi\right) / \partial \xi$ as $\left(x^{\prime}, \xi\right) \in \mathcal{T}_{r}$. Since $V\left(x^{\prime}, \xi\right)$ is an absolutely continuous function in $\xi$ for a.a. $x^{\prime} \in B_{r}^{\prime}\left(x_{0}\right)$ ) (after redefining it, if necessary, on a set of zero measure) we get

$$
\begin{equation*}
U\left(x^{\prime}, \xi\right)=U\left(x^{\prime}, 0\right)+\int_{0}^{\xi} V\left(x^{\prime}, t\right) d t \quad \text { for a.a. }\left(x^{\prime}, \xi\right) \in \mathcal{T}_{r} \tag{13}
\end{equation*}
$$

where the point $\left(x^{\prime}, 0\right) \in B_{r}^{\prime}\left(x_{0}\right)$ lies in $\mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}$ and $U\left(x^{\prime}, 0\right) \in W^{2, q}$ there by Proposition 2, the Fubini theorem and [12, Remark 2.1]. Passing to the new variables $\left(x^{\prime}, \xi\right)$ in (12), taking the derivatives of (13) up to second order and substituting them into the right-hand side of (12), this last reads

$$
\begin{cases}\mathcal{L}^{\prime} V=F_{1}\left(x^{\prime}, \xi\right)+\int_{0}^{\xi} \mathcal{D}_{2}(\xi) V\left(x^{\prime}, t\right) d t & \text { a.e. } \mathcal{T}_{r}  \tag{14}\\ V=\Phi & \text { on } \partial \mathcal{T}_{r}\end{cases}
$$

where $\mathcal{L}^{\prime}$ is the operator $\mathcal{L}$ in terms of $\left(x^{\prime}, \xi\right)=\left(x_{1}^{\prime}, \ldots, x_{n-1}^{\prime}, \xi\right)$,

$$
\begin{align*}
F_{1}\left(x^{\prime}, \xi\right) & :=\partial F / \partial \boldsymbol{L}+\mathcal{D}_{1} V\left(x^{\prime}, \xi\right)+\mathcal{D}_{1}^{\prime} U\left(x^{\prime}, \xi\right)+\mathcal{D}_{2}^{\prime} U\left(x^{\prime}, 0\right),  \tag{15}\\
\mathcal{D}_{2}(\xi) V\left(x^{\prime}, t\right) & :=\sum_{i, j=1}^{n-1} A^{i j}\left(x^{\prime}, \xi\right) D_{x_{i}^{\prime} x_{j}^{\prime}} V\left(x^{\prime}, t\right), \quad A^{i j} \in L^{\infty}
\end{align*}
$$

$\mathcal{D}_{1}, \mathcal{D}_{1}^{\prime}, \mathcal{D}_{2}^{\prime}$ are linear differential operators with $L^{\infty}$-coefficients, ord $\mathcal{D}_{1}=$ ord $\mathcal{D}_{1}^{\prime}=$ 1 , ord $\mathcal{D}_{2}^{\prime}=2$. The Sobolev imbedding theorem implies $F_{1} \in L^{q^{\prime}}\left(\mathcal{T}_{r}\right)$ with $q^{\prime}$ given by $(11)$ as consequence of $\partial F / \partial \boldsymbol{L} \in L^{q^{\prime}}\left(\mathcal{N}^{\prime \prime}\right), U\left(x^{\prime}, 0\right) \in W^{2, q}\left(B_{r}^{\prime}\left(x_{0}\right)\right)$ and $U, V \in W^{2, p}\left(\mathcal{N}^{\prime \prime}\right)$. Nevertheless the second-order operator $\mathcal{D}_{2}(\xi)$ has a quite rough characteristic form which is neither symmetric nor sign-definite, the improving-ofintegrability holds for (14) thanks to the particular structure of $\mathcal{T}_{r}$ as union of $\boldsymbol{L}$-trajectories through $B_{r}^{\prime}\left(x_{0}\right)$. Actually, we will show that if $V \in W^{2, q^{\prime}}$ on a subset of $\mathcal{T}_{r}$ with $\xi<T$, then $V$ remains a $W^{2, q^{\prime}}$-function on a larger subset with $\xi<T+r$ for small enough $r$, after that the higher integrability of $U$ will follow from Proposition 2 and (13). For, take an arbitrary $T \in\left(0, t^{+}-t^{-}\right)$and define

$$
\mathcal{P}_{r, T}:=\left\{\left(x^{\prime}, \xi\right) \in \mathcal{T}_{r}: \quad \xi<T\right\}
$$

For a fixed $r>0,\left\{\mathcal{P}_{r, T}\right\}_{T>0}$ is a non-decreasing family of domains exhausting $\mathcal{T}_{r}$ and $\mathcal{P}_{r, T} \equiv \mathcal{T}_{r}$ for values of $T$ greater than the maximal exit-time

$$
T_{\max }:=\sup _{x^{\prime} \in B_{r}^{\prime}\left(x_{0}\right)} \sup \left\{t>0: \quad \boldsymbol{\psi}_{\boldsymbol{L}}\left(t ; x^{\prime}\right) \in \Omega, x^{\prime} \in B_{r}^{\prime}\left(x_{0}\right)\right\}
$$

Proposition 4. Let $T \in\left(0, t^{+}-t^{-}\right)$and consider the solution $V \in W^{2, p}\left(\mathcal{T}_{r}\right)$ of the problem (14). Suppose $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T}\right)$ where $q^{\prime}$ is given by (11).

There exists an $r_{0}>0$ such that $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T+r}\right)$ for all $r<r_{0}$.
Proof. There are three possible cases to be distinguished.
Case A: $T+3 r<T_{\max }$. We have $\mathcal{P}_{r, T} \subset \mathcal{P}_{r, T+3 r} \subset \mathcal{T}_{r} \equiv \mathcal{P}_{r, T_{\max }}$ and consider the $C^{\infty}$-function $\eta: \mathbb{R} \rightarrow[0,1]$ such that

$$
\eta(\xi)= \begin{cases}1 & \text { as } \xi \in(-\infty, T+r]  \tag{16}\\ \text { strictly decreases } & \text { as } \xi \in(T+r, T+2 r) \\ 0 & \text { as } \xi \geq T+2 r\end{cases}
$$

Setting $\widetilde{V}\left(x^{\prime}, \xi\right):=\eta(\xi) V\left(x^{\prime}, \xi\right)$, it follows $\mathcal{L}^{\prime} \widetilde{V}=\eta\left(\mathcal{L}^{\prime} V\right)+\mathcal{L}_{1} V$ where $\mathcal{L}_{1}$ is a firstorder differential operator with $L^{\infty}$-coefficients depending on these of $\mathcal{L}^{\prime}$ and on the
derivatives of $\eta$. Therefore,

$$
\begin{align*}
\mathcal{L}^{\prime} \widetilde{V} & =\eta F_{1}+\mathcal{L}_{1} V+\eta(\xi) \int_{0}^{\xi} \mathcal{D}_{2}(\xi) V\left(x^{\prime}, t\right) d t  \tag{17}\\
& =\eta F_{1}+\mathcal{L}_{1} V+\int_{0}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) \widetilde{V}\left(x^{\prime}, t\right) d t
\end{align*}
$$

because $\mathcal{D}_{2}(\xi)$ is a second-order operator acting in the $x^{\prime}$-variables only.
We set $\Omega_{r} \subset \mathcal{P}_{r, T+3 r} \backslash \mathcal{P}_{r, T-3 r}$ for a $C^{1,1}$-smooth domain containing $\mathcal{P}_{3 r / 4, T+2 r} \backslash$ $\mathcal{P}_{3 r / 4, T-2 r}$ and such that

$$
r^{-1} \Omega_{r}:=\left\{\left(\widetilde{y}^{\prime}, \widetilde{\xi}\right): \quad \widetilde{y}^{\prime}=x^{\prime} / r, \widetilde{\xi}=(\xi-T) / r,\left(x^{\prime}, \xi\right) \in \Omega_{r}\right\} \in C^{1,1}
$$

uniformly in $r$. The boundary $\partial \Omega_{r}$ consists of the "lateral" parts $\partial_{1} \Omega_{r}:=\partial \Omega_{r} \cap \partial \Omega$ and $\partial_{2} \Omega_{r}:=\partial \Omega_{r} \cap \Omega \cap\{\xi \in(T-2 r, T+2 r)\} \subset\left(\mathcal{P}_{r, T+2 r} \backslash \mathcal{P}_{r, T-2 r}\right) \backslash\left(\mathcal{P}_{3 r / 4, T+2 r} \backslash\right.$ $\left.\mathcal{P}_{3 r / 4, T-2 r}\right)$, and of two $C^{1,1}$-smooth components $\partial \Omega_{r}^{ \pm}$lying in $\mathcal{P}_{r, T+3 r} \backslash \mathcal{P}_{r, T+2 r}$ and $\mathcal{P}_{r, T-2 r} \backslash \mathcal{P}_{r, T-3 r}$, respectively. The properties of $\mu(\mathrm{cf}$. (9)) ensure $U \equiv 0, V \equiv 0$, $\widetilde{V} \equiv 0$ on $\mathcal{T}_{r} \backslash \mathcal{T}_{3 r / 4}$ whence $\widetilde{V} \equiv 0$ near $\partial_{2} \Omega_{r}$.

For an arbitrary $\left(x^{\prime}, \xi\right) \in \Omega_{r}$, the factor $\eta(\xi) / \eta(t)$ in (17) vanishes when $\xi \geq T+2 r$ while $\eta(\xi) / \eta(t) \leq 1$ because $\eta$ decreases in $(T+r, T+2 r)$. Moreover, $|\xi-T|<3 r$ for $\left(x^{\prime}, \xi\right) \in \Omega_{r}$ and

$$
\begin{aligned}
\int_{0}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) \widetilde{V}\left(x^{\prime}, t\right) d t & =\int_{0}^{T} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) \widetilde{V}\left(x^{\prime}, t\right) d t+\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) \widetilde{V}\left(x^{\prime}, t\right) d t \\
& =\eta(\xi) \int_{0}^{T} \mathcal{D}_{2}(\xi) V\left(x^{\prime}, t\right) d t+\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) \widetilde{V}\left(x^{\prime}, t\right) d t
\end{aligned}
$$

by means of (15) and since $\eta(t)=\underset{\sim}{\eta}(T)=1$ as $t \leq T$.
We get from (14) and (17) that $\widetilde{V} \in W^{2, p}\left(\Omega_{r}\right)$ solves the Dirichlet problem

$$
\left\{\begin{array}{l}
\mathcal{L}^{\prime} \widetilde{V}=F_{2}\left(x^{\prime}, \xi\right)+\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) \widetilde{V}\left(x^{\prime}, t\right) d t \quad \text { a.a. }\left(x^{\prime}, \xi\right) \in \Omega_{r}  \tag{18}\\
\widetilde{V}=\widetilde{\Phi}:=\eta \Phi=\left\{\begin{array}{lll}
\eta \mu \varphi \in W^{2-1 / q, q} & \text { on } \partial_{1} \Omega_{r} & (\text { by }(10)) \\
0 & \text { on } \partial_{2} \Omega_{r} & (\text { by }(10)), \\
0 & \text { on } \partial \Omega_{r}^{+} & (\text {by }(16)) \\
V \in W^{2-1 / q^{\prime}, q^{\prime}} & \text { on } \partial \Omega_{r}^{-} & (\text {since } \xi<T-2 r \text { and } \\
& \left.V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T}\right)\right)
\end{array}\right.
\end{array}\right.
$$

where, recalling $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T}\right)$, we have

$$
\begin{equation*}
F_{2}\left(x^{\prime}, \xi\right):=\eta F_{1}+\mathcal{L}_{1} V+\eta(\xi) \int_{0}^{T} \mathcal{D}_{2}(\xi) V\left(x^{\prime}, t\right) d t \in L^{q^{\prime}}\left(\Omega_{r}\right) \tag{19}
\end{equation*}
$$

We are going to prove now that $\widetilde{V} \in W^{2, q^{\prime}}\left(\Omega_{r}\right)$ for small enough $r>0$, whence it will follow $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T+r}\right)$ in view of $(16)$ and $V \equiv 0$ near $\partial_{2} \Omega_{r}$. The claim is obvious if $q^{\prime}=p$ because $V \in W^{2, p}\left(\mathcal{T}_{r}\right)$. Otherwise, take an arbitrary $s \in\left[p, q^{\prime}\right]$ and denote by $W_{*}^{2, s}\left(\Omega_{r}\right)$ the Sobolev space $W^{2, s}\left(\Omega_{r}\right)$ normed with

$$
\|u\|_{W_{*}^{2, s}\left(\Omega_{r}\right)}:=\|u\|_{L^{s}\left(\Omega_{r}\right)}+r\|D u\|_{L^{s}\left(\Omega_{r}\right)}+r^{2}\left\|D^{2} u\right\|_{L^{s}\left(\Omega_{r}\right)} .
$$

Define now the operator $\mathfrak{F}: W_{*}^{2, s}\left(\Omega_{r}\right) \rightarrow W_{*}^{2, s}\left(\Omega_{r}\right)$ as follows: for any $w \in W_{*}^{2, s}\left(\Omega_{r}\right)$ the image $\mathfrak{F} w \in W_{*}^{2, s}\left(\Omega_{r}\right)$ is the unique solution of the Dirichlet problem

$$
\begin{cases}\mathcal{L}^{\prime}(\mathfrak{F} w)=F_{2}+\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi) w\left(x^{\prime}, t\right) d t \in L^{s}\left(\Omega_{r}\right) & \text { a.a. }\left(x^{\prime}, \xi\right) \in \Omega_{r}  \tag{20}\\ \mathfrak{F} w=\widetilde{\Phi} \in W^{2-1 / s, s}\left(\partial \Omega_{r}\right) & \text { on } \partial \Omega_{r}\end{cases}
$$

We will prove that $\mathfrak{F}$ is a contraction for small values of $r$. For this goal, take arbitrary $w_{1}, w_{2} \in W_{*}^{2, s}\left(\Omega_{r}\right)$. The difference $\mathfrak{F} w_{1}-\mathfrak{F} w_{2}$ solves

$$
\begin{cases}\mathcal{L}^{\prime}\left(\mathfrak{F} w_{1}-\mathfrak{F} w_{2}\right)=\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi)\left(w_{1}-w_{2}\right)\left(x^{\prime}, t\right) d t & \text { a.a. }\left(x^{\prime}, \xi\right) \in \Omega_{r}  \tag{21}\\ \mathfrak{F} w_{1}-\mathfrak{F} w_{2}=0 & \text { on } \partial \Omega_{r}\end{cases}
$$

In order to apply the $L^{s}$-a priori estimates from [1] or [3] for the solutions of (21), we have to control the dependence on $r$ therein. For, we recall that $r^{-1} \Omega_{r} \in C^{1,1}$ uniformly in $r$ and apply a standard approach consisting of dilation of $\Omega_{r}$ onto $r^{-1} \Omega_{r}$, reduction of the problem (21) to a new one in variables $\left(\widetilde{y}^{\prime}, \widetilde{\xi}\right) \in r^{-1} \Omega_{r}$, application of the $L^{s}$-estimates from [3, Theorem 9.17] and finally turning back to (21) (see the Proof of Lemma 2.2, Eq. (2.12) in [12]). This way, one gets

$$
\begin{equation*}
\left\|\mathfrak{F} w_{1}-\mathfrak{F} w_{2}\right\|_{W_{*}^{2, s}\left(\Omega_{r}\right)} \leq C r^{2}\left\|\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi)\left(w_{1}-w_{2}\right)\left(x^{\prime}, t\right) d t\right\|_{L^{s}\left(\Omega_{r}\right)} \tag{22}
\end{equation*}
$$

where the constant $C$ is independent of $r$. Jensen's integral inequality yields

$$
r^{2}\left\|\int_{T}^{\xi} \frac{\eta(\xi)}{\eta(t)} \mathcal{D}_{2}(\xi)\left(w_{1}-w_{2}\right)\left(x^{\prime}, t\right) d t\right\|_{L^{s}\left(\Omega_{r}\right)} \leq C \max _{\left(x^{\prime}, \xi\right) \in \Omega_{r}}|\xi-T|\left\|w_{1}-w_{2}\right\|_{W_{*}^{2, s}\left(\Omega_{r}\right)}
$$

and thus (22) rewrites into

$$
\left\|\mathfrak{F} w_{1}-\mathfrak{F} w_{2}\right\|_{W_{*}^{2, s}\left(\Omega_{r}\right)} \leq C \max _{\left(x^{\prime}, \xi\right) \in \Omega_{r}}|\xi-T|\left\|w_{1}-w_{2}\right\|_{W_{*}^{2, s}\left(\Omega_{r}\right)}
$$

We have $\max _{\left(x^{\prime}, \xi\right) \in \Omega_{r}}|\xi-T|<3 r, C$ is independent of $r$ and therefore $\mathfrak{F}$ will be really a contraction from $W_{*}^{2, s}\left(\Omega_{r}\right)$ into itself for any $s \in\left[p, q^{\prime}\right]$ if $r \leq r_{0}$ with $r_{0}$ under control and small enough. Fixing $r=r_{0} / 2$, there is a unique fixed point of $\mathfrak{F}$ in $W_{*}^{2, s}\left(\Omega_{r}\right)$ for all $s \in\left[p, q^{\prime}\right]$. However, $\widetilde{V} \in W^{2, p}\left(\Omega_{r}\right)$ is already a fixed point of $\mathfrak{F}$ since it solves (18) and therefore $\widetilde{V} \in W^{2, q^{\prime}}\left(\Omega_{r}\right)$. It follows $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T+r}\right)$ by means of $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T}\right), \widetilde{V} \equiv 0$ on $\mathcal{T}_{r} \backslash \mathcal{T}_{3 r / 4}$ and the properties of $\eta(\xi)$.
Case B: $T<T_{\max } \leq T+3 r$. We have $\mathcal{T}_{r} \backslash \mathcal{P}_{r, T} \neq \emptyset, \mathcal{P}_{r, T+3 r} \equiv \mathcal{T}_{r}$ now and we do not need anymore the cut-off function $\eta$ because $V=\partial U / \partial \boldsymbol{L} \equiv 0$ near the points of $\partial_{2} \mathcal{T}_{r}$ where $\xi>T$ (cf. (9)). Thus, it suffices to repeat the above arguments with $\eta(\xi) \equiv 1 \forall \xi \in \mathbb{R}$ and $\Omega_{r} \in C^{1,1}$ defined as before when $\xi \leq T$ while $\mathcal{T}_{3 r / 4} \backslash \mathcal{P}_{3 r / 4, T} \subset\left(\Omega_{r} \cap\{\xi>T\}\right) \subset \mathcal{T}_{r} \backslash \mathcal{P}_{r, T}$ (cf. (9)). We have anyway a problem like (18) for $V \equiv \widetilde{V}$ with boundary condition

$$
V=\partial U / \partial \boldsymbol{L}= \begin{cases}\mu \varphi \in W^{2-1 / q, q} & \text { on } \partial_{1} \Omega_{r}=\partial \Omega_{r} \cap \partial \Omega \\ 0 & \text { on } \partial_{2} \Omega_{r}=\partial \Omega_{r} \cap \Omega \cap\{\xi>T-3 r\} \\ V \in W^{2-1 / q^{\prime}, q^{\prime}} & \text { on } \partial \Omega_{r}^{-} \quad \text { (by hypothesis) }\end{cases}
$$

Therefore, the procedure from Case $A$ gives $V \in W^{2, q^{\prime}}\left(\mathcal{P}_{r, T+3 r}\right)$.
Case C: $T_{\max } \leq T$. We have $\mathcal{P}_{r, T+r} \equiv \mathcal{P}_{r, T} \equiv \mathcal{T}_{r}$ now and thus the claim.
Proposition 5. Suppose $r<r_{0}$ with $r_{0}$ given in Proposition 4. Then the solution $V$ of the problem (14) lies in $W^{2, q}\left(\mathcal{T}_{r}\right)$ and satisfies the estimate

$$
\begin{align*}
\|V\|_{W^{2, q}\left(\mathcal{T}_{r}\right)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}\right. & +\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}  \tag{23}\\
& \left.+\|u\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}+\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}\right)
\end{align*}
$$

Proof. We note that $V \in W^{2, q} \subseteq W^{2, q^{\prime}}$ near $B_{r}^{\prime}\left(x_{0}\right)$ in view of $B_{r}^{\prime}\left(x_{0}\right) \subset \mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}$, Proposition 2 and (6). Therefore, successive applications of Proposition 4 with increasing values of $T$ will give $V \in W^{2, q^{\prime}}\left(\mathcal{T}_{r}\right), q^{\prime}>p$. After that, in order to get $V \in W^{2, q}\left(\mathcal{T}_{r}\right)$, it suffices to put $q^{\prime}$ in the place of $p$ in (11) and to repeat finitely many times the above arguments until $q^{\prime}=q$.

To obtain (23), we take $T \in\left(0, t^{+}-t^{-}\right)$to be arbitrary, fix $r=r_{0} / 2$, and consider the domains $\Omega_{r}$ defined in the proof of Proposition 4. Let $\widetilde{V}=\eta V \in W^{2, q}\left(\mathcal{T}_{r}\right)$ solve (18) with $\eta$ given by (16) in Case $A$ and $\eta \equiv 1$ in Case $B$. Since $\widetilde{V}$ is a fixed point of the mapping $\mathfrak{F}$ : $W^{2, q}\left(\Omega_{r}\right) \rightarrow W^{2, q}\left(\Omega_{r}\right), \mathfrak{F} \widetilde{V}=\widetilde{V}$, we get

$$
\left\|D^{2} \widetilde{V}\right\|_{L^{q}\left(\Omega_{r}\right)}=\left\|D^{2}(\mathfrak{F} \widetilde{V})\right\|_{L^{q}\left(\Omega_{r}\right)} \leq\left\|D^{2}(\mathfrak{F} \widetilde{V}-\mathfrak{F} 0)\right\|_{L^{q}\left(\Omega_{r}\right)}+\left\|D^{2}(\mathfrak{F} 0)\right\|_{L^{q}\left(\Omega_{r}\right)},
$$

while

$$
\left\|D^{2}\left(\mathfrak{F} w_{1}-\mathfrak{F} w_{2}\right)\right\|_{L^{q}\left(\Omega_{r}\right)} \leq \theta\left\|D^{2}\left(w_{1}-w_{2}\right)\right\|_{L^{q}\left(\Omega_{r}\right)} \quad \forall w_{1}, w_{2} \in W^{2, q}\left(\Omega_{r}\right), \quad \theta<1
$$

because $\mathfrak{F}$ is a contraction, (22) and the fact that $\mathcal{D}_{2}(\xi)$ is a homogeneous secondorder operator (cf. (15)). This way, $\left\|D^{2}(\mathfrak{F} \widetilde{V}-\mathfrak{F} 0)\right\|_{L^{q}\left(\Omega_{r}\right)} \leq \theta\left\|D^{2}(\widetilde{V}-0)\right\|_{L^{q}\left(\Omega_{r}\right)}=$ $\theta\left\|D^{2} \widetilde{V}\right\|_{L^{q}\left(\Omega_{r}\right)}$ and therefore

$$
\begin{equation*}
\left\|D^{2} \widetilde{V}\right\|_{L^{q}\left(\Omega_{r}\right)} \leq C\left\|D^{2}(\mathfrak{F} 0)\right\|_{L^{q}\left(\Omega_{r}\right)} \tag{24}
\end{equation*}
$$

with $\mathfrak{F} 0 \in W^{2, q}\left(\Omega_{r}\right)$ being the unique solution of the Dirichlet problem

$$
\left\{\mathcal{L}^{\prime}(\mathfrak{F} 0)=F_{2} \quad \text { a.e. } \Omega_{r}, \quad \mathfrak{F} 0=\widetilde{\Phi} \quad \text { on } \partial \Omega_{r}\right.
$$

(see (20)), for which the $L^{p}$-theory (cf. [3, Chapter 9]) gives

$$
\begin{equation*}
\left\|D^{2}(\mathfrak{F} 0)\right\|_{L^{q}\left(\Omega_{r}\right)} \leq\|\mathfrak{F} 0\|_{W^{2, q}\left(\Omega_{r}\right)} \leq C\left(\left\|F_{2}\right\|_{L^{q}\left(\Omega_{r}\right)}+\|\widetilde{\Phi}\|_{W^{2-1 / q, q}\left(\partial \Omega_{r}\right)}\right) . \tag{25}
\end{equation*}
$$

Direct applications, based on (19) and (15), yield

$$
\begin{aligned}
\left\|F_{2}\right\|_{L^{q}\left(\Omega_{r}\right)}= & \left\|\eta F_{1}+\mathcal{L}_{1} V+\eta(\xi) \int_{0}^{T} \mathcal{D}_{2}(\xi) V\left(x^{\prime}, t\right) d t\right\|_{L^{q}\left(\Omega_{r}\right)} \\
\leq & C\left(\|\partial F / \partial \boldsymbol{L}\|_{L^{q}\left(\Omega_{r}\right)}+\|U\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}\right)}+\|U\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}+\|V\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}\right. \\
& \left.+\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T}\right)}\right) \\
\leq & C\left(\|\partial f / \partial \boldsymbol{L}\|_{L^{q}(\mathcal{N})}+\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime} \backslash \mathcal{N}^{\prime}\right)}+\|u\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}+\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}\right. \\
& \left.+\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T}\right)}\right)
\end{aligned}
$$

in view of (7), (10), $U=\mu u, V=\partial U / \partial \boldsymbol{L}$ and (9). Moreover,

$$
\begin{aligned}
\|\widetilde{\Phi}\|_{W^{2-1 / q, q}\left(\partial \Omega_{r}\right)} & \leq C\left(\|\varphi\|_{W^{2-1 / q, q}(\partial \Omega \cap \mathcal{N})}+\|V\|_{W^{2, q}\left(\mathcal{P}_{r, T}\right)}\right) \\
& \leq C\left(\|\varphi\|_{W^{2-1 / q, q}(\partial \Omega \cap \mathcal{N})}+\|V\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}+\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T}\right)}\right) \\
& \leq C\left(\|\varphi\|_{W^{2-1 / q, q}(\partial \Omega \cap \mathcal{N})}+\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}+\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T}\right)}\right)
\end{aligned}
$$

by (18) and $\partial \Omega_{r}^{-} \subset \mathcal{P}_{r, T}$. Further on, $\widetilde{V}=V$ on $\mathcal{P}_{r, T+r}$, whence

$$
\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T+r}\right)} \leq\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T}\right)}+\left\|D^{2} \widetilde{V}\right\|_{L^{q}\left(\Omega_{r}\right)}
$$

Therefore, setting $\zeta(T):=\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T}\right)}$ and $K:=\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+$ $\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}+\|u\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}+\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}\left(\mathcal{T}_{r}\right)}$, it follows from (24), (25) and Proposition 2 that

$$
\begin{equation*}
\zeta(T+r) \leq C(K+\zeta(T)) \quad \forall T \in\left(0, t^{+}-t^{-}\right) \tag{26}
\end{equation*}
$$

To get (23), we let $m$ to be the least integer such that $T_{\max } \leq m r$ and iterate (26) in order to obtain

$$
\begin{aligned}
\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{T}_{r}\right)} & =\left\|D^{2} V\right\|_{L^{q}\left(\mathcal{P}_{r, T_{\max }}\right)}=\zeta\left(T_{\max }\right)=\zeta(m r)=\zeta((m-1) r+r) \\
& \leq C(K+\zeta((m-1) r))=C(K+\zeta((m-2) r+r)) \\
& \leq K\left(C+C^{2}\right)+C^{2} \zeta((m-2) r) \\
& \vdots \\
& \leq K \sum_{j=1}^{m} C^{j}+C^{m} \zeta(0)=K \sum_{j=1}^{m} C^{j}
\end{aligned}
$$

This proves (23).
Remark 6. It is important to note that the constant $C$ in Proposition 5 depends on $m$ through $T_{\max }$, and therefore on the point $x_{0} \in \mathcal{E}$. Actually, that constant will have the very same value for each other point of $\mathcal{E}$ lying on the same $\boldsymbol{L}$-trajectory as $x_{0}$.

Moreover, if the improving-of-integrability property asserted in Propositions 4 and 5 holds on a set $S \subset \bar{\Omega}$ then it is guaranteed, on the base of (13), on any other set which can be reached from $S$ along $\boldsymbol{L}$-trajectories.

To complete the proof of Lemma 3, we select a finite set $\left\{\mathcal{T}_{r}^{j}\right\}_{j=1}^{N}$ of neighbourhoods covering the compact $\mathcal{E}$, each of the type $\mathcal{T}_{r}$ above with $r=r_{0} / 2$, and such that $\mathcal{T}:=$ closure $\left(\bigcup_{j=1}^{N} \mathcal{T}_{r / 2}^{j}\right) \subset \mathcal{N}^{\prime \prime}$ is a closed neighbourhood of $\mathcal{E}$ in $\bar{\Omega}$. It is clear that Proposition 2 remains true with $\mathcal{T}$ instead of $\mathcal{N}^{\prime}$ and then (7) rewrites into

$$
\begin{equation*}
\|u\|_{W^{2, q}(\Omega \backslash \mathcal{T})} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{L^{q}(\Omega)}+\|\varphi\|_{W^{1-1 / q, q}(\partial \Omega)}\right) \tag{27}
\end{equation*}
$$

The improving-of-integrability claimed in Lemma 3 then follows from (13), Proposition 5 and (27) (recall $U=u$ on $\mathcal{T}_{r / 2}^{j}$ ). Similarly, (13), (27) and (23) yield

$$
\begin{equation*}
\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)} \leq\|u\|_{W^{2, q}(\mathcal{T})}+\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime} \backslash \mathcal{T}\right)} \tag{28}
\end{equation*}
$$

$$
\begin{aligned}
\leq & C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}\right. \\
& \left.+\|u\|_{W^{1, q}(\mathcal{N})}+\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}(\mathcal{N})}\right)
\end{aligned}
$$

Later on, $\mathcal{N} \backslash \mathcal{N}^{\prime \prime} \subset \Omega \backslash \mathcal{N}^{\prime}$ and

$$
\begin{aligned}
\|u\|_{W^{1, q}(\mathcal{N})} & \leq\|u\|_{W^{1, q}\left(\mathcal{N}^{\prime \prime}\right)}+\|u\|_{W^{1, q}\left(\mathcal{N} \backslash \mathcal{N}^{\prime \prime}\right)} \\
& \leq \varepsilon\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)}+C(\varepsilon)\left(\|u\|_{L^{q}(\Omega)}+\|u\|_{W^{2, q}\left(\Omega \backslash \mathcal{N}^{\prime}\right)}\right)
\end{aligned}
$$

in view of the interpolation inequality for the $W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)$-norms with $\varepsilon>0$ under control ${ }^{3}$. In the same manner,

$$
\begin{aligned}
\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}(\mathcal{N})} & \leq\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}\left(\mathcal{N}^{\prime}\right)}+\|\partial u / \partial \boldsymbol{L}\|_{W^{1, q}\left(\mathcal{N} \backslash \mathcal{N}^{\prime}\right)} \\
& \leq \varepsilon\|\partial u / \partial \boldsymbol{L}\|_{W^{2, q}\left(\mathcal{N}^{\prime}\right)}+C(\varepsilon)\left(\|\partial u / \partial \boldsymbol{L}\|_{L^{q}\left(\mathcal{N}^{\prime}\right)}+\|u\|_{W^{2, q}\left(\Omega \backslash \mathcal{N}^{\prime}\right)}\right)
\end{aligned}
$$

while

$$
\|\partial u / \partial \boldsymbol{L}\|_{W^{2, q}\left(\mathcal{N}^{\prime}\right)} \leq C\left(\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)}+\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}\right)
$$

by means of the local a priori estimates ([3, Theorem 9.11]) for the problem (6).
A substitution of the above expressions into (28) and (7) give

$$
\begin{aligned}
\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}\right. & +\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})} \\
& \left.+\varepsilon\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)}+C(\varepsilon)\|\partial u / \partial \boldsymbol{L}\|_{L^{q}\left(\mathcal{N}^{\prime}\right)}\right)
\end{aligned}
$$

whence, choosing $\varepsilon>0$ small enough, we get

$$
\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}+\|u\|_{W^{1, q}\left(\mathcal{N}^{\prime}\right)}\right) .
$$

Similarly, another application of the interpolation inequality yields

$$
\|u\|_{W^{1, q}\left(\mathcal{N}^{\prime}\right)} \leq\|u\|_{W^{1, q}\left(\mathcal{N}^{\prime \prime}\right)} \leq \delta\|u\|_{W^{2, q}\left(\mathcal{N}^{\prime \prime}\right)}+C(\delta)\|u\|_{L^{q}\left(\mathcal{N}^{\prime \prime}\right)}
$$

and thus

$$
\|u\|_{W^{2}, q\left(\mathcal{N}^{\prime \prime}\right)} \leq C\left(\|u\|_{L^{q}(\Omega)}+\|f\|_{\mathcal{F}^{q}(\Omega, \mathcal{N})}+\|\varphi\|_{\Phi^{q}(\partial \Omega, \mathcal{N})}\right) .
$$

for small $\delta>0$. The proof of Lemma 3 is completed.
The statement of Theorem 1 follows from Proposition 2 and Lemma 3.

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[^1]:    ${ }^{1}$ It will be clear from the considerations given below that instead of Lipschitz continuity of the coefficients of $\mathcal{L}$ in $\mathcal{N}$ as (3) asks, it suffices to have essentially bounded their directional derivatives with respect to the field $\boldsymbol{L}$.
    ${ }^{2}\left|\boldsymbol{\psi}_{\boldsymbol{L}}\left(t ; x^{\prime}\right)-\boldsymbol{\psi}_{\boldsymbol{L}}\left(t ; x^{\prime \prime}\right)\right| \leq e^{t\|\boldsymbol{L}\|_{C^{1}}(\mathcal{N})}\left|x^{\prime}-x^{\prime \prime}\right|$ for all $x^{\prime}, x^{\prime \prime} \in \mathcal{N}$.

[^2]:    ${ }^{3}$ This requires some minimal smoothness of $\partial \mathcal{N}^{\prime \prime}$ and it is not restrictive to take it Lipschitz continuous at the very beginning.

