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# On the detection of a moving obstacle in an ideal fluid by a boundary measurement 

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#### Abstract

In this paper, we investigate the problem of the detection of a moving obstacle in a perfect fluid occupying a bounded domain in $\mathbb{R}^{2}$ from the measurement of the velocity of the fluid on one part of the boundary. We show that when the obstacle is a ball, we may identify the position and the velocity of its centre of mass from a single boundary measurement. Linear stability estimates are also established by using shape differentiation techniques.


## 1. Introduction

Inverse problems in fluid mechanics constitute a challenging topic with numerous potential applications, ranging from engineering, medicine, and military surveillance to fishing. In [4], the authors established that a fixed smooth convex obstacle surrounded by a real fluid modellized by Navier-Stokes equations could be identified via a localized boundary measurement of the velocity of the fluid and the Cauchy forces. Directional stability estimates were also derived in the same paper. The results in [4] strongly rested on the unique continuation property for the Stokes system due to Fabre-Lebeau [8]. In [6] the obstacle was identified by a measurement of both the gradient of the pressure and the velocity of the fluid on a part of the boundary, and the stability was established by shape differentiation. The distance from a chosen point to the obstacle was estimated in [10] from boundary measurements for a fluid governed by the stationary Stokes equation. As water is often considered as a perfect fluid on a small time-scale, it is natural to wonder whether the above results are still valid when the viscosity coefficient tends to zero, i.e., for an ideal fluid. The answer to that question is of great importance for applications.

In this paper, we shall address the issue of whether a moving obstacle surrounded by a perfect fluid may be detected by the measurement of the tangential velocity of the fluid on one part of the boundary. Assume a fixed domain $\Omega \subset \mathbb{R}^{2}$, and a rigid body $S$ occupying the set $S(t) \subset \Omega$ at time $t$. Let us denote by $h(t)$ the centre of mass of $S(t), m$ the mass of the rigid body and $J$ its moment of inertia. Then the equations modelling the dynamics of the system solid + fluid read [19] as follows:

$$
\begin{align*}
& \frac{\partial u}{\partial t}+(u \cdot \nabla) u+\nabla p=0 \quad(x, t) \in(\Omega \backslash S(t)) \times \mathbb{R}  \tag{1.1}\\
& \operatorname{div} u=0 \quad(x, t) \in(\Omega \backslash S(t)) \times \mathbb{R},  \tag{1.2}\\
& u \cdot n=\left(h^{\prime}+r(x-h)^{\perp}\right) \cdot n \quad(x, t) \in \partial S(t) \times \mathbb{R}  \tag{1.3}\\
& u \cdot n=g \quad(x, t) \in \partial \Omega \times \mathbb{R},  \tag{1.4}\\
& m h^{\prime \prime}(t)=\int_{\partial S(t)} p n \mathrm{~d} \sigma+f(t) \quad t \in \mathbb{R}  \tag{1.5}\\
& J r^{\prime}(t)=\int_{\partial S(t)}(x-h(t))^{\perp} \cdot p n \mathrm{~d} \sigma+T(t) \quad t \in \mathbb{R} \tag{1.6}
\end{align*}
$$

In these equations, $u=u(x, t)$ (resp. $p=p(x, t))$ is the velocity (resp. the pressure) of the fluid, $g$ is the flow through the boundary $\Omega$ (just assumed to be given here), $r$ is the angular velocity of the solid, $x^{\perp}=\left(-x_{2}, x_{1}\right)$ if $x=\left(x_{1}, x_{2}\right), n$ is the outward unit normal vector and $f(t)$ (resp. $T(t))$ stands for the external force (resp. the external torque) applied to the solid in addition to the contribution of the fluid pressure represented by the integral term. For a rigid body without a self-propelling mechanism (i.e. $f=T=0$ ) moving in the whole space ( $\Omega=\mathbb{R}^{2}$ ), it has been proved that system (1.1)-(1.6) admits a unique classical solution defined for all times in $[18,19]$. When $\Omega$ is a half-plane, the existence of chocks in finite time between the rigid body and the boundary of the domain has been established in [12] when $S$ is a ball and $u$ is a potential velocity.

In this paper, we focus on the determination of the position and the velocity of the obstacle from a boundary measurement of the velocity of the fluid at a given time $t$. This means that we will ignore the Newton laws (1.5) and (1.6) in our analysis. This setting is convenient in situations where the self-propelling data, namely $f$ and $T$, are not known. This is the case, for example, when we aim to localize a submarine from a pressure measurement.

In contrast to what happens for Navier-Stokes equations, the Euler equations do not exhibit any unique continuation property because of the existence of the famous ghost solutions with compact support [16]. A simple example of a ghost solution is provided by the stationary solution $v(x)=\left(\partial \psi / \partial x_{2},-\partial \psi / \partial x_{1}\right)$, where the stream function $\psi$ is given by

$$
\psi(x)=-\int_{1}^{|x|} \frac{1}{r}\left(\int_{1}^{r} s \omega(s) \mathrm{d} s\right) \mathrm{d} r
$$

and the vorticity $\omega \in C^{\infty}\left(\mathbb{R}^{+}\right)$is chosen so that $\omega(s)=0$ for $s \geqslant 1$ and $\int_{r}^{1} s \omega(s) \mathrm{d} s=0$ for $r \in\left(0, r_{0}\right)$, where $r_{0} \in(0,1)$ is a given number. As $v$ is supported in the set $\left\{r_{0} \leqslant|x| \leqslant 1\right\}$, we deduce that no obstacle contained in the ball $B_{r_{0}}(0)$ can be detected from measurements performed at a distance from the origin larger than 1 ; that is, the identifiability property fails for Eulerian flows.

However, the above obstruction to the detection disappears if we restrict ourselves to potential flows, that is, flows for which the velocity assumes the form $v=\nabla \varphi$ for a scalar function $\varphi=\varphi(x, t)$. It is well known (see, e.g. [17]) that a bidimensional Eulerian flow in


Figure 1. Moving obstacle in a pipeline.
a domain with one hole $S$ is potential if the vorticity vanishes everywhere and the circulation along $S$ is null. As it has been noticed in [13], an Eulerian flow remains potential as long as the incoming flow, located at the part of the boundary where $g<0$, has a null vorticity. We shall assume that the incoming flow fulfils that condition.

Plugging $u=\nabla \varphi$ in (1.1)-(1.4) results in the system

$$
\begin{align*}
& \nabla\left(\frac{\partial \varphi}{\partial t}+\frac{1}{2}|\nabla \varphi|^{2}+p\right)=0 \quad(x, t) \in(\Omega \backslash S(t)) \times \mathbb{R},  \tag{1.7}\\
& \Delta \varphi=0 \quad(x, t) \in(\Omega \backslash S(t)) \times \mathbb{R},  \tag{1.8}\\
& \frac{\partial \varphi}{\partial n}=\left(h^{\prime}+r(x-h)^{\perp}\right) \cdot n \quad(x, t) \in \partial S(t) \times \mathbb{R},  \tag{1.9}\\
& \frac{\partial \varphi}{\partial n}=g \quad(x, t) \in \partial \Omega \times \mathbb{R} . \tag{1.10}
\end{align*}
$$

Clearly, measuring the tangential component of the velocity on one part of the boundary amounts to measuring the function $\varphi$ itself. When the obstacle is fixed $\left(h^{\prime}=r=0\right)$, condition (1.9) simplifies to $\partial \varphi / \partial n=0$, so that the detection of the obstacle reduces to a very classical problem (see, e.g. [1-3, 5, 9, 14, 22]). Such a problem arises in different contexts including the corrosion detection by electrostatic measurements and the crack detection in nonferrous metals from electromagnetic measurements.

As far as we know, the situation where the obstacle is moving (i.e. $(h, r) \neq(0,0))$ has not yet been investigated. It turns out that this problem is more difficult to study than the stationary one for two reasons: (i) the velocity of the rigid body being unknown, the classical argument based upon the unique continuation property for the Laplace equation is not sufficient to derive the identifiability property; (ii) unlike [2, 9], we cannot use several Neumann data and apply topological arguments to identify the obstacle. Indeed, the obstacle may occupy different positions and undergo different velocities for different Neumann data.

The goal of this paper is to address the identifiability issue when the obstacle has a known form. For the sake of simplicity, we shall assume here that the obstacle is the ball $B_{1}(h(t))$ of radius one centred at the point $h(t)$. (See figure 1.) Note that, for any $x \in \partial B_{1}(h(t)), x-h(t)=-n$, hence $(x-h)^{\perp} \cdot n=0$. Setting $l=h^{\prime}$, the system reads

$$
\begin{array}{ll}
\Delta \varphi=0 & \text { in } \quad \Omega \backslash \overline{B_{1}(h(t))}, \\
\frac{\partial \varphi}{\partial n}=l \cdot n & \text { on } \quad \partial B_{1}(h(t)), \tag{1.12}
\end{array}
$$

$$
\begin{equation*}
\frac{\partial \varphi}{\partial n}=g \quad \text { on } \quad \partial \Omega \tag{1.13}
\end{equation*}
$$

and we assume that $\varphi$ is measured on a part $\Gamma_{m}$ of the boundary $\partial \Omega$. The identifiability issue is to understand whether only one pair $(h, l)$ may be associated with a given measurement.

Clearly, the data $g(x)=l \cdot x$, with $l \in \mathbb{R}^{2}$ a given fixed vector, has to be excluded, for it may lead to the situation where the ball, which is surrounded by a fluid flowing at the same velocity $(\varphi(x, t)=l \cdot x)$, is not identifiable. We shall prove in this paper that for any data $g$ which is not of this form, the identifiability problem has a positive answer, whatever be the distance between the ball and $\partial \Omega$.

The method of proof relies on a careful investigation of the singularities of the solution $\psi$ to the Dirichlet problem

$$
\begin{align*}
& \Delta \psi=0 \quad \text { in } \quad B_{1}(h) \backslash \overline{B_{1}(-h)},  \tag{1.14}\\
& \psi=y \quad \text { on } \quad \partial B_{1}(h) \cap\{y>0\},  \tag{1.15}\\
& \psi=-y \quad \text { on } \quad \partial B_{1}(-h) \cap\{y>0\}, \tag{1.16}
\end{align*}
$$

where $h=(0, \delta)$ and $0<\delta<1$. More precisely, we will show that the solution $\psi$ of (1.14)-(1.16) is not of the class $C^{2}$ at the point $M_{+}=\left(\sqrt{1-\delta^{2}}, 0\right)$ when $\delta \geqslant 1 / \sqrt{2}$, and that $\psi$ may be extended analytically on the set $B_{1}(-h) \cap\{y>0\}$ when $\delta>0$, by using a Möbius transformation and a version of Schwarz reflection principle for harmonic functions. It is likely that the function $\psi$ fails to be analytic in a neighbourhood of $M_{+}$for any $\delta>0$, but despite our efforts, we were not able to prove it.

The second main objective of the paper is to investigate the stability properties of the $\operatorname{map} \varphi_{\left.\right|_{r_{m}}} \rightarrow(h, l)$. Under the same assumption on the data $g$ as above, we shall derive a linear stability estimate. The method of proof rests on the concept of shape differentiation introduced by Simon in [21].

To summarize, we present in this paper sharp results for the identification of a moving obstacle surrounded by a potential flow via a single boundary measurement, when the obstacle is a ball in $\mathbb{R}^{2}$. It would be interesting to see whether these results can be extended to a smooth obstacle of arbitrary (known) form in dimension two. On the other hand, it is clear that more information can be collected by repeating measurements on a time interval. It would be interesting to see whether the shape of the obstacle could be identified with a measurement over a time interval. Finally, preliminary computations indicate that a single measurement of the fluid velocity on the boundary is probably not sufficient to extend the results of the paper to the dimension three. This suggests that repeating measurements over a time interval could be essential in dimension three. These issues, which are below the scope of this paper, will be investigated elsewhere.

The paper is outlined as follows. The identifiability result is stated and proved in section 2 . Section 3 is devoted to the derivation of the stability estimate. Finally, the annexe contains the proof of the fact that $\psi$ is not of class $C^{2}$ at $M_{+}$when $\delta \geqslant 1 / \sqrt{2}$.

## 2. Identifiability

Let $\Omega$ be a bounded (connected) open set in $\mathbb{R}^{2}$, with a smooth boundary $\partial \Omega$. Assume given an open set $\Gamma_{m}$ in $\partial \Omega$, and a function $g \in H^{s}(\partial \Omega)$, with $s \geqslant 0$, such that $\int_{\partial \Omega} g \mathrm{~d} \sigma=0$. We denote by $\Omega_{a}$ the set of admissible positions for the centres of the balls of radius one included in $\Omega$; i.e.,

$$
\Omega_{a}:=\{h \in \Omega, \operatorname{dist}(h, \partial \Omega)>1\} .
$$

Let $h_{1}, h_{2} \in \Omega_{a}$ and $l_{1}, l_{2} \in \mathbb{R}^{2}$. For $i=1,2$, we denote by $B_{i}$ the ball $B_{1}\left(h_{i}\right)$, and by $\varphi_{i}$ the solution (defined up to an additive constant) of the following Neumann problem:

$$
\begin{array}{lc}
\Delta \varphi_{i}=0 & \text { in } \quad \Omega \backslash \overline{B_{i}}, \\
\frac{\partial \varphi_{i}}{\partial n}=g \quad & \text { on } \quad \partial \Omega, \\
\frac{\partial \varphi_{i}}{\partial n}=l_{i} \cdot n & \text { on } \quad \partial B_{i}, \tag{2.3}
\end{array}
$$

where $n$ stands for the outward unit normal vector. We shall say that problem (2.1)-(2.3) is identifiable if, for a convenient choice of the input $g$, the following implication holds:

$$
\begin{equation*}
\varphi_{1}=\varphi_{2} \quad \text { on } \quad \Gamma_{m} \Rightarrow h_{1}=h_{2} \quad \text { and } \quad l_{1}=l_{2} \tag{2.4}
\end{equation*}
$$

Let us introduce the two-dimensional space

$$
V=\operatorname{Span}\left\{e_{1} \cdot n, e_{2} \cdot n\right\} \subset L^{\infty}(\partial \Omega)
$$

where $\left\{e_{1}, e_{2}\right\}$ denotes the canonical basis of $\mathbb{R}^{2}$.
The following result is the first main result of this paper.
Theorem 2.1. Assume that $g \in H^{s}(\partial \Omega) \backslash V$ with $s>1 / 2$. For $i=1,2$, pick any $\left(h_{i}, l_{i}\right) \in \Omega_{a} \times \mathbb{R}^{2}$ and let $\varphi_{i}$ denote the solution (defined up to a constant) of (2.1)-(2.3). Then (2.4) holds.

Proof. By standard regularity results for elliptic problems [15], we know that $\varphi_{i} \in$ $H^{s+3 / 2}\left(\Omega \backslash \overline{B_{i}}\right) \subset C^{1}\left(\bar{\Omega} \backslash B_{i}\right)$. Assume that $\varphi_{1}=\varphi_{2}$ on $\Gamma_{m}$. Since also $\frac{\partial \varphi_{1}}{\partial n}=g=\frac{\partial \varphi_{2}}{\partial n}$ on $\Gamma_{m}$, we infer from the unique continuation property that

$$
\varphi_{1}=\varphi_{2} \quad \text { on } \quad \Omega \backslash \overline{B_{1} \cup B_{2}} .
$$

Define a function $\varphi: \Omega \backslash \overline{B_{1} \cap B_{2}} \rightarrow \mathbb{R}$ by

$$
\varphi(x):=\left\{\begin{array}{lll}
\varphi_{1}(x) & \text { if } & x \in \Omega \backslash \overline{B_{1}}  \tag{2.5}\\
\varphi_{2}(x) & \text { if } & x \in \Omega \backslash \overline{B_{2}} .
\end{array}\right.
$$

Then $\varphi$ fulfils

$$
\begin{array}{ll}
\Delta \varphi=0 & \text { in } \quad \Omega \backslash \overline{B_{1} \cap B_{2}}, \\
\frac{\partial \varphi}{\partial n}=g & \text { on } \quad \partial \Omega, \\
\frac{\partial \varphi}{\partial n}=l_{1} \cdot n & \text { on } \quad \partial B_{1}, \\
\frac{\partial \varphi}{\partial n}=l_{2} \cdot n & \text { on } \quad \partial B_{2} . \tag{2.9}
\end{array}
$$

If $B_{1} \cap B_{2}=\emptyset$, then $\varphi$ is, as $\varphi_{2}$, defined and harmonic in $B_{1}$, and we infer from (2.8) that $\varphi(x)=l_{1} \cdot x+$ const on $B_{1}$. The same property holds on $\Omega$ by unique continuation. This gives $g=l_{1} \cdot n$ on $\partial \Omega$, and hence $g \in V$, which is a contradiction.

The nontrivial case is the one for which $B_{1} \cap B_{2} \neq \emptyset$, i.e. $\left\|h_{2}-h_{1}\right\|<2$. Obviously, if $h_{1}=h_{2}$, then (2.8) and (2.9) yield $l_{1}=l_{2}$. Assume from now on that $h_{1} \neq h_{2}$. If $l_{1}=l_{2}$, then introducing the domain $D_{1}=B_{1} \backslash \overline{B_{2}}$, we see that $\varphi$ solves

$$
\begin{equation*}
\Delta \varphi=0 \quad \text { in } \quad D_{1} \tag{2.10}
\end{equation*}
$$



Figure 2. The Möbius transformation from $D_{1}$ to $C_{1}$.

$$
\begin{equation*}
\frac{\partial \varphi}{\partial n}=l_{1} \cdot n \quad \text { on } \quad \partial D_{1} \tag{2.11}
\end{equation*}
$$

which gives again $\varphi(x)=l_{1} \cdot x+$ const in $D_{1}$ and $g \in V$, which is a contradiction.
We shall therefore assume that $h_{1} \neq h_{2}$ and $l_{1} \neq l_{2}$. Using Green's formula, we infer from (2.10) and (2.8)-(2.9) that

$$
\begin{equation*}
\left(l_{2}-l_{1}\right) \cdot\left(h_{2}-h_{1}\right)=0 . \tag{2.12}
\end{equation*}
$$

Translating and rotating $\Omega$ if needed, we may assume that $h_{1}=(0, \delta), h_{2}=(0,-\delta)$ with $0<\delta<1$, and $l_{2}-l_{1}=\lambda e_{1}$ for some $\lambda \neq 0$. Replacing $\varphi$ and $g$ by $(-2 / \lambda)\left(\varphi-l_{1} \cdot x\right)+e_{1} \cdot x$ and $(-2 / \lambda)\left(g-l_{1} \cdot n\right)+e_{1} \cdot n$, respectively, we may assume that $l_{1}=e_{1}$ and that $l_{2}=-e_{1}$. We are thus led to investigate the properties of a function $\varphi: \Omega \backslash \overline{B_{1} \cap B_{2}} \rightarrow \mathbb{R}$ satisfying the system

$$
\begin{array}{ll}
\Delta \varphi=0 & \text { in } \\
\frac{\partial \varphi}{\partial n}=g \quad \Omega \overline{B_{1} \cap B_{2}}, \\
\frac{\partial \varphi}{\partial n}=e_{1} \cdot n & \text { on }
\end{array} \quad \partial \Omega,
$$

We introduce the points $M_{ \pm}=\left( \pm \sqrt{1-\delta^{2}}, 0\right)$ located at the intersection of the circles $\partial B_{1}$ and $\partial B_{2}$ (see figure 2).

We shall use thereafter some complex analysis, denoting the coordinates by $(x, y)$ instead of ( $x_{1}, x_{2}$ ), and identifying a couple $(x, y)$ of real numbers with the complex number $z=x+\mathrm{i} y$.

Pick a number $\eta>0$ such that $B_{1+\eta}\left(h_{1}\right) \subset \Omega$. The function $\tilde{\varphi}(x, y)=\varphi(x, y)-x$ fulfils the system

$$
\begin{array}{lll}
\Delta \tilde{\varphi}=0 & \text { on } & B_{1+\eta}\left(h_{1}\right) \backslash \overline{B_{1}}, \\
\frac{\partial \tilde{\varphi}}{\partial n}=0 & \text { on } & \partial B_{1}
\end{array}
$$

and is of class $C^{1}$ on $\bar{\Omega} \backslash B_{1}$. By the reflection principle (see [11]), we may extend $\tilde{\varphi}$ to the annulus $\mathcal{A}_{1}=B_{1+\eta}\left(h_{1}\right) \backslash \overline{B_{(1+\eta)^{-1}}\left(h_{1}\right)}$ as a harmonic function in setting

$$
\tilde{\varphi}(z)=\tilde{\varphi}\left(\overline{(z-\mathrm{i} \delta)^{-1}}+\mathrm{i} \delta\right) \quad \text { for } \quad(1+\eta)^{-1}<|z-\mathrm{i} \delta|<1
$$

Therefore, $\varphi$ may as well be extended to $\mathcal{A}_{1}$ as a harmonic function. Analogously, $\varphi$ may be extended as a harmonic function on the annulus $\mathcal{A}_{2}=B_{1+\eta}\left(h_{2}\right) \backslash \overline{B_{(1+\eta)^{-1}}\left(h_{2}\right)}$. To obtain the contradiction, we shall prove that $\varphi$ is also analytic in $B_{1} \cap B_{2}$, so that by (2.15), $\varphi(x)=x \cdot e_{1}+$ const, and again $g \in V$, which contradicts the assumptions.

Since $\int_{\partial B_{1}}(\partial \varphi / \partial n) \mathrm{d} \sigma=0$, the function $\varphi$ possesses a harmonic conjugate function $\psi$ defined on $\mathcal{A}_{1} \cup \mathcal{A}_{2} \cup \Omega \backslash\left(\overline{B_{1} \cap B_{2}}\right)$, fulfilling $\nabla \psi=(\nabla \varphi)^{\perp}$. Let $\theta$ (resp. $\theta^{\prime}$ ) denote the angle $\left(\vec{e}_{1}, \overrightarrow{h_{1} M}\right)$ (resp. $\left(\vec{e}_{1}, \overrightarrow{h_{2} M}\right)$ ).

Then $\partial \psi / \partial \theta=\partial \varphi / \partial r=\cos \theta$ on $\partial B_{1}$, which gives upon integration $\psi=\sin \theta+C$. Similarly, $\psi=-\sin \theta^{\prime}+C^{\prime}$ on $\partial B_{2}$. Picking the constants $C$ and $C^{\prime}$ so that $\psi\left(M_{ \pm}\right)=0$, we see that $\psi$ solves

$$
\begin{array}{lc}
\Delta \psi=0 & \text { in } \quad D_{1}=B_{1} \backslash \overline{B_{2}}, \\
\psi=y & \text { on } \quad \Gamma_{1}:=\left(\partial B_{1}\right) \backslash B_{2}, \\
\psi=-y \quad & \text { on } \quad \gamma_{2}:=\left(\partial B_{2}\right) \cap B_{1} . \tag{2.19}
\end{array}
$$

A similar Dirichlet problem is satisfied by $\psi$ on $-D_{1}=B_{2} \backslash \overline{B_{1}}$, and from the uniqueness of the solution we infer that

$$
\psi(x,-y)=\psi(x, y)=\psi(-x, y)
$$

To prove that $\varphi$ has no singularity in $B_{1} \cap B_{2}$, it is therefore sufficient to check that $\psi$ does not have any singularity in the set $B_{2} \cap\{z=x+\mathrm{i} y ; y>0\}$. We first transform problem (2.17)-(2.19) into a Dirichlet problem in a corner.

Let $T_{1}: z \mapsto z_{1}=x_{1}+\mathrm{i} y_{1}=\left(z+\sqrt{1-\delta^{2}}\right)^{-1}$ denote the inversion of pole $M_{-}=-\sqrt{1-\delta^{2}}$. As $T_{1}$ is a Möbius transformation, it carries circles into circles or lines (see [11]). Since $T_{1}\left(M_{-}\right)=\infty$, we see that $l_{1}=T_{1}\left(\partial B_{1}\right)$ (resp. $\left.l_{2}=T_{1}\left(\partial B_{2}\right)\right)$ is the line passing through $T_{1}\left(M_{+}\right)=\left(2 \sqrt{1-\delta^{2}}\right)^{-1}$ and $T_{1}(\mathrm{i}(1+\delta))=(-\mathrm{i}+\sqrt{(1-\delta) /(1+\delta)}) / 2$ (resp. through $T_{1}\left(M_{+}\right)$and $T_{1}(-\mathrm{i}(1+\delta))=(\mathrm{i}+\sqrt{(1-\delta) /(1+\delta)}) / 2$ ) (see figure 2). Clearly, $T_{1}\left(\Gamma_{1}\right)$ is the half-line $l_{1}^{+} \subset l_{1}$ issuing from $T_{1}\left(M_{+}\right)$and containing $T_{1}(\mathrm{i}(1+\delta))$, while $T_{1}\left(\gamma_{2}\right)$ is the half-line $l_{2}^{+} \subset l_{2}$ issued from $T_{1}\left(M_{+}\right)$and which does not contain $T_{1}(-1(1+\delta))$. Therefore, $T_{1}\left(D_{1}\right)$ is the convex corner $C_{1}=l_{1}^{+} l_{2}^{+}$.

Note that $z=T_{1}^{-1}\left(z_{1}\right)=z_{1}^{-1}-\sqrt{1-\delta^{2}}$. Let $\psi_{1}\left(z_{1}\right):=\psi(z)$. Then $\psi_{1}$ solves the system

$$
\begin{align*}
& \Delta \psi_{1}=0 \quad \text { in } \quad C_{1}  \tag{2.20}\\
& \psi_{1}\left(z_{1}\right)=-\frac{y_{1}}{x_{1}^{2}+y_{1}^{2}} \quad \text { on } \quad l_{1}^{+},  \tag{2.21}\\
& \psi_{1}\left(z_{1}\right)=\frac{y_{1}}{x_{1}^{2}+y_{1}^{2}} \quad \text { on } \quad l_{2}^{+},  \tag{2.22}\\
& \psi_{1}\left(z_{1}\right) \rightarrow 0 \quad \text { as } \quad z_{1} \rightarrow \infty, z_{1} \in C_{1} . \tag{2.23}
\end{align*}
$$

For notational convenience, we translate and rotate the corner $C_{1}$. We let $C_{2}=T_{2}\left(C_{1}\right)$, where $T_{2}\left(z_{1}\right):=z_{2}=-\left(z_{1}-\left(2 \sqrt{1-\delta^{2}}\right)^{-1}\right)$. Then

$$
C_{2}=\left\{z_{2} \in \mathbb{C}^{*} ; \frac{\pi-\theta}{2}<\arg z_{2}<\frac{\pi+\theta}{2}\right\}
$$

where $\theta \in(0, \pi)$ stands for the angle of $C_{1}$ at $T_{1}\left(M_{+}\right)$, or of $\partial D_{1}$ at $M_{+}$by conformal invariance.

Let $\psi_{2}\left(z_{2}\right):=\psi_{1}\left(z_{1}\right)$. Then $\psi_{2}$ solves the system

$$
\begin{align*}
& \Delta \psi_{2}=0 \quad \text { in } C_{2},  \tag{2.24}\\
& \psi_{2}\left(z_{2}\right)=\frac{y_{2}}{\left(x_{2}+c\right)^{2}+y_{2}^{2}} \quad \text { on } \quad d_{-1},  \tag{2.25}\\
& \psi_{2}\left(z_{2}\right)=-\frac{y_{2}}{\left(x_{2}+c\right)^{2}+y_{2}^{2}} \quad \text { on } \quad d_{0},  \tag{2.26}\\
& \psi_{2}\left(z_{2}\right) \rightarrow 0 \quad \text { as } \quad z_{2} \rightarrow \infty, z_{2} \in C_{2}, \tag{2.27}
\end{align*}
$$

where $c:=-\left(2 \sqrt{1-\delta^{2}}\right)^{-1}$, and for any $k \in \mathbb{Z}, d_{k}$ denotes the half-line

$$
d_{k}=\left\{z_{2} \in \mathbb{C}^{*} ; \arg z_{2}=\theta_{k}:=\frac{\pi+(2 k+1) \theta}{2}\right\}
$$

Note that $\left(T_{2} \circ T_{1}\right)\left(B_{2} \cap\{y>0\}\right)$ is the corner $C=\left\{z_{2} \in \mathbb{C}^{*} ; \theta_{0}<\arg z_{2}<\pi\right\}$. To prove that $\psi$ does not have singularities in $B_{2} \cap\{y>0\}$, it is then sufficient to check that $\psi_{2}$ can be extended as a harmonic function on $C$. This is done in applying several times the following reflection principle for harmonic functions.

Lemma 2.2. Let $\theta_{0} \in \mathbb{R}$ and $\theta \in(0, \pi / 2)$. Let $l_{ \pm}=\left\{z \in \mathbb{C}^{*} ; \arg z=\theta_{0} \pm \theta\right\}$, and let $l_{0}=\left\{z \in \mathbb{C}^{*} ; \arg z=\theta_{0}\right\}$. Let $C_{+}=\left\{z \in \mathbb{C}^{*} ; \theta_{0}<\arg z<\theta_{0}+\theta\right\}$ (resp. $C_{-}=$ $\left\{z \in \mathbb{C}^{*} ; \theta_{0}-\theta<\arg z<\theta_{0}\right\}$ ) be the sectors bounded by the half-lines $l_{0}$ and $l_{+}$(resp. by the half-lines $l_{-}$and $l_{0}$ ). Let $\psi$ be a harmonic function on $C_{-}$such that

$$
\begin{array}{ll}
\lim _{z \rightarrow Z, z \in C_{-}} \psi(z)=\operatorname{Im} f_{-}(Z) & \forall Z \in l_{-}, \\
\lim _{z \rightarrow Z, z \in C_{-}} \psi(z)=\operatorname{Im} f_{0}(Z) & \forall Z \in l_{0}, \tag{2.29}
\end{array}
$$

where $f_{-}\left(\right.$resp. $\left.f_{0}\right)$ is a holomorphic function in a neighbourhood of $l_{-}\left(\right.$resp. on $\left.C_{-} \cup l_{0} \cup C_{+}\right)$. Then $\psi$ can be extended as a harmonic function on the set $C_{-} \cup l_{0} \cup C_{+}$, and

$$
\begin{equation*}
\lim _{z \rightarrow Z, z \in C_{+}} \psi(z)=-I m f_{-}\left(\mathrm{e}^{-2 \mathrm{i} \theta} Z\right)+\lim _{z \rightarrow Z, z \in C_{+}} \operatorname{Im}\left(f_{0}(z)+f_{0}\left(\mathrm{e}^{2 \mathrm{i} \theta_{0}} \bar{z}\right)\right) \tag{2.30}
\end{equation*}
$$

for each $Z \in l_{+}$for which the limit in the right-hand side of (2.30) exists.
Proof. Using the transformation $z \mapsto \mathrm{e}^{-\mathrm{i} \theta_{0}} z$, we may without loss of generality assume that $\theta_{0}=0$, hence $l_{+}=\overline{l_{-}}$and $C_{+}=\overline{C_{-}}$, ( $\cdot$ means conjugation). Pick a holomorphic function $f$ on $C_{-}$such that $\psi(z)=\operatorname{Im} f(z)$ on $C_{-}$, and let $F(z):=f(z)-f_{0}(z)$. Then $F$ is holomorphic on $C_{-}$, and for any $Z \in l_{0}$

$$
\lim _{z \rightarrow Z, z \in C_{-}} \operatorname{Im} F(z)=\lim _{z \rightarrow Z, z \in C_{-}} \psi(z)-\operatorname{Im} f_{0}(Z)=0
$$

by (2.29). Using the Schwarz reflection principle for holomorphic functions stated in [20], we infer that $F$ may be extended as an holomorphic function on $C_{-} \cup l_{0} \cup C_{+}$in setting

$$
F(z)=\overline{F(\bar{z})} \quad \forall z \in C_{+} .
$$

Letting

$$
\psi(z)=\operatorname{Im}\left(F(z)+f_{0}(z)\right) \quad \forall z \in C_{+} \cup l_{0}
$$



Figure 3. Analytic extension of $\psi_{2}$ when $\theta=\pi / 5$.
we obtain a harmonic extension of $\psi$ on $C_{-} \cup l_{0} \cup C_{+}$. For $Z \in l_{+}$, we have

$$
\begin{align*}
\lim _{z \rightarrow Z, z \in C_{+}} \psi(z) & =\lim _{z \rightarrow Z, z \in C_{+}} \operatorname{Im}\left(\overline{F(\bar{z})}+f_{0}(z)\right) \\
& =\lim _{z \rightarrow Z, z \in C_{+}} \operatorname{Im}\left(\overline{f(\bar{z})}-\overline{f_{0}(\bar{z})}+f_{0}(z)\right) \\
& =-\lim _{z \rightarrow Z, z \in C_{+}} \psi(\bar{z})+\lim _{z \rightarrow Z, z \in C_{+}} \operatorname{Im}\left(f_{0}(\bar{z})+f_{0}(z)\right) \\
& =-\operatorname{Im} f_{-}(\bar{Z})+\lim _{z \rightarrow Z, z \in C_{+}} \operatorname{Im}\left(f_{0}(z)+f_{0}(\bar{z})\right) \tag{2.31}
\end{align*}
$$

whenever the limit in (2.31) does exist.
To see that $\psi_{2}$ can be extended in an analytic way on the sector $C$, we apply inductively lemma 2.2. Starting with $\left(l_{+}, l_{0}, l_{-}\right)=\left(d_{1}, d_{0}, d_{-1}\right)$ and $f_{-}\left(z_{2}\right)=-\left(z_{2}+c\right)^{-1}, f_{0}\left(z_{2}\right)=$ $\left(z_{2}+c\right)^{-1}$, we obtain for $\arg z_{2}=\theta_{1}$

$$
\begin{equation*}
\psi_{2}\left(z_{2}\right)=\operatorname{Im}\left(\frac{2}{\mathrm{e}^{-2 i \theta} z_{2}+c}+\frac{1}{z_{2}+c}\right) . \tag{2.32}
\end{equation*}
$$

Note that $\arg \left(\mathrm{e}^{-2 \mathrm{i} \theta} z_{2}\right)=-2 \theta+\theta_{1}=\theta_{-1}>0$, and hence the right-hand side of (2.32) is well defined on $d_{1}$. (Recall that $c<0$.)

Applying again lemma 2.2 with $\left(l_{+}, l_{0}, l_{-}\right)=\left(d_{2}, d_{1}, d_{0}\right)$ and

$$
f_{-}(z)=\frac{1}{z_{2}+c}, \quad f_{0}\left(z_{2}\right)=\frac{2}{\mathrm{e}^{-2 \mathrm{i} \theta} z_{2}+c}+\frac{1}{z_{2}+c},
$$

it follows that for $\arg z_{2}=\theta_{2}$,

$$
\psi_{2}\left(z_{2}\right)=\operatorname{Im}\left(\frac{2}{\mathrm{e}^{-4 i \theta} z_{2}+c}+\frac{2}{\mathrm{e}^{-2 \mathrm{i} \theta} z_{2}+c}+\frac{1}{z_{2}+c}\right) .
$$

Assume first that $\theta=\frac{\pi}{2 N+1}$ for some $N \in \mathbb{N}^{*}$, so that $\theta_{N}=\pi$. Then, we can prove by induction on $k$ that for each $k \in\{1, \ldots, N\}$ the function $\psi_{2}$ can be extended in an analytic way on the sector $d_{-1} d_{k}$, with

$$
\begin{equation*}
\psi_{2}\left(z_{2}\right)=\operatorname{Im}\left(\frac{1}{z_{2}+c}+\sum_{l=1}^{k} \frac{2}{\mathrm{e}^{-2 l i \theta} z_{2}+c}\right) \quad \forall z_{2} \in d_{k} \tag{2.33}
\end{equation*}
$$

Note that for any $k \leqslant N$, the right-hand side of (2.33) does not present any singularity in the sector $d_{k-1} d_{k+1}$ (hence the extension at the step $k$ can be performed), as

$$
\arg \left(\mathrm{e}^{-2 k \mathrm{i} \theta} z_{2}\right)>\theta_{k-1}-2 k \theta=\frac{\pi}{2}\left(1-\frac{2 k+1}{2 N+1}\right) \geqslant 0
$$

A final application of lemma 2.2 gives that $\psi_{2}$ may be extended analytically on the sector $d_{-1} d_{N+1}$, with

$$
\begin{equation*}
\psi_{2}\left(z_{2}\right)=\operatorname{Im}\left(\frac{1}{z_{2}+c}+\sum_{l=1}^{N+1} \frac{2}{\mathrm{e}^{-2 l i \theta} z_{2}+c}\right) \tag{2.34}
\end{equation*}
$$

for any $z \in d_{N+1}$ for which the right-hand side of (2.34) is meaningful. This occurs for any point of $d_{N+1}$, except for $z_{2}=|c| \mathrm{e}^{\mathrm{i} \theta_{N+1}}$. This point is the first singularity encountered during the extension procedure of $\psi_{2}$. As it is outside $C$, since $\theta_{N+1} \in(\pi, 2 \pi)$, we are done.

Assume now that $\frac{\pi}{2 N+3}<\theta<\frac{\pi}{2 N+1}$ for some $N \in \mathbb{N}$. Then $\theta_{N}<\pi<\theta_{N+1}$. The analytic extension may be done in the sector $d_{-1} d_{N+1}$, as for $k \leqslant N$ and $\theta_{k-1}<\arg z_{2}<\theta_{k+1}$ we have

$$
\arg \left(\mathrm{e}^{-2 k i \theta} z_{2}\right)>\theta_{k-1}-2 k \theta=\frac{\pi}{2}-\frac{2 k+1}{2} \theta \geqslant \frac{\pi}{2}-\frac{2 N+1}{2} \theta>0 .
$$

Once again, the analytic extension of $\psi_{2}$ does not present any singularity in $C$. The proof of theorem 2.1 is complete.

Remark 2.3. The above proof of theorem 2.1 is still valid when $g=0$ and $l_{1} \neq 0$. Indeed, the relation $l_{1} \cdot n=0$ cannot hold everywhere on $\partial \Omega$. This means that, in the absence of flow through the boundary, the obstacle can be identified when it is moving, and only in that case.

## 3. Stability estimates

In this section, we investigate the stability properties of the map $\varphi_{\mid \Gamma_{m}} \rightarrow(h, l)$. Linear stability estimates will be established by using shape differentiation.

Fix $h_{0} \in \Omega_{a}, l_{0} \in \mathbb{R}^{2}$ and a function $g$ fulfilling

$$
\begin{equation*}
g \in H^{s}(\partial \Omega) \quad \text { for some } \quad s \geqslant 1 \quad \text { and } \quad \int_{\partial \Omega} g \mathrm{~d} \sigma=0 \tag{3.1}
\end{equation*}
$$

Write $B_{0}=B_{1}\left(h_{0}\right)$. Let $\varphi_{0}$ denote the solution of the reference Neumann problem

$$
\begin{array}{lc}
\Delta \varphi_{0}=0 & \text { in } \quad \Omega \backslash \overline{B_{0}} \\
\frac{\partial \varphi_{0}}{\partial n}=g \quad \text { on } \quad \partial \Omega \\
\frac{\partial \varphi_{0}}{\partial n}=l_{0} \cdot n \quad \quad \text { on } \quad \partial B_{0} . \tag{3.4}
\end{array}
$$

Pick any $(h, l) \in \Omega_{a} \times \mathbb{R}^{2}$, and let $\varphi$ denote the solution of the perturbed Neumann problem

$$
\begin{equation*}
\Delta \varphi=0 \quad \text { in } \quad \Omega \backslash \bar{B}, \tag{3.5}
\end{equation*}
$$

$$
\begin{align*}
& \frac{\partial \varphi}{\partial n}=g \quad \text { on } \quad \partial \Omega  \tag{3.6}\\
& \frac{\partial \varphi}{\partial n}=l \cdot n \quad \text { on } \quad \partial B \tag{3.7}
\end{align*}
$$

where $B=B_{1}(h)$. Once again, the functions $\varphi$ and $\varphi_{0}$ are defined up to an additive constant. By standard regularity results for elliptic problems, we know that $\varphi \in H^{s+3 / 2}(\Omega \backslash \bar{B})$, hence $\varphi_{\mid \partial \Omega} \in H^{s+1}(\partial \Omega)$. We may therefore define a map $\Lambda: \Omega_{a} \times \mathbb{R}^{2} \times\left\{g \in H^{s}(\partial \Omega) ; \int_{\partial \Omega} g \mathrm{~d} \sigma=\right.$ $0\} \rightarrow H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}$ by

$$
\begin{equation*}
\Lambda(h, l, g)=\varphi_{\mid \Gamma_{m}} . \tag{3.8}
\end{equation*}
$$

Recall that the quotient space $H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}$ is a Banach space for the norm

$$
\|g\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}}:=\inf _{t \in \mathbb{R}}\|g+t\|_{H^{s+1}\left(\Gamma_{m}\right)}
$$

Proceeding as in [4], one may prove that this map is of class $C^{1}$. We are now in a position to state the second main result of the paper.

Theorem 3.1. Let $g$ fulfilling (3.1) and let $\left(h_{0}, l_{0}\right) \in \Omega_{a} \times \mathbb{R}^{2}$. If $g \notin V$, then there exist two constants $\rho>0$ and $C>0$ depending only on $\left(h_{0}, l_{0}, g\right)$ such that for any $(h, l) \in B_{\rho}\left(h_{0}, l_{0}\right)$ we have

$$
\begin{equation*}
\left\|\Lambda(h, l, g)-\Lambda\left(h_{0}, l_{0}, g\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \geqslant C\left\|\left(h-h_{0}, l-l_{0}\right)\right\|_{\mathbb{R}^{4}} \tag{3.9}
\end{equation*}
$$

Proof. Let $\mathrm{d} \Lambda\left(h_{0}, l_{0}, g\right)$ denote the differential of $\Lambda$ at the point $\left(h_{0}, l_{0}, g\right)$, and let $L=\mathrm{d} \Lambda\left(h_{0}, l_{0}, g\right)_{\mid \mathbb{R}^{2} \times \mathbb{R}^{2} \times\{0\}}$. We need the following result, whose proof will be postponed.

Proposition 3.2. Let $g$ be as in theorem 3.1. Then the map $L: \mathbb{R}^{4} \rightarrow H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}$ is one-to-one.

By the compactness of the unit sphere in $\mathbb{R}^{4}$, we infer from proposition 3.2 the existence of two positive constants $C_{1}, C_{2}$ such that

$$
\begin{equation*}
C_{1}\|(\hat{h}, \hat{l})\| \leqslant\|L(\hat{h}, \hat{l})\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \leqslant C_{2}\|(\hat{h}, \hat{l})\| \quad \forall(\hat{h}, \hat{l}) \in \mathbb{R}^{4} \tag{3.10}
\end{equation*}
$$

On the other hand, we can write

$$
\begin{equation*}
\Lambda\left(h_{0}+\hat{h}, l_{0}+\hat{l}, g\right)=\Lambda\left(h_{0}, l_{0}, g\right)+L(\hat{h}, \hat{l})+\|(\hat{h}, \hat{l})\| \varepsilon(\hat{h}, \hat{l}) \tag{3.11}
\end{equation*}
$$

where $\varepsilon(h, l)$ is a function such that $\varepsilon(h, l) \rightarrow 0$ as $(h, l) \rightarrow 0$. Pick $\rho>0$ so that $\|\varepsilon(\hat{h}, \hat{l})\|<C_{1} / 2$ whenever $\|(\hat{h}, \hat{l})\|<\rho$. Then we infer from (3.10) and (3.11) that

$$
\left\|\Lambda\left(h_{0}+\hat{h}, l_{0}+\hat{l}, g\right)-\Lambda\left(h_{0}, l_{0}, g\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \geqslant\left(C_{1} / 2\right)\|(\hat{h}, \hat{l})\|
$$

for $\|(\hat{h}, \hat{l})\|<\rho$. The proof of theorem 3.1 is achieved.
It remains to prove proposition 3.2.
Proof of proposition 3.2. Let $h_{0}, l_{0}$ and $g$ be as in the statement of theorem 3.1. Without loss of generality, we may assume that $h_{0}=(0,0)$.

If $(\hat{h}, \hat{l}) \in \mathbb{R}^{4}$ is given, then by a classical result due to Simon (see [21]) we have that $L(\hat{h}, \hat{l})=\psi_{\mid \Gamma_{m}}$, where $\psi$ denotes the solution (defined up to a constant) of

$$
\begin{array}{ll}
\Delta \psi=0 & \text { in } \quad \Omega \backslash \overline{B_{0}} \\
\frac{\partial \psi}{\partial n}=0 & \text { on } \quad \partial \Omega \tag{3.13}
\end{array}
$$

$$
\begin{equation*}
\frac{\partial \psi}{\partial n}=-\hat{h} \cdot n \frac{\partial^{2} \varphi_{0}}{\partial n^{2}}+\left(\nabla \varphi_{0}-l_{0}\right) \cdot \operatorname{grad}_{\partial \Omega}(\hat{h} \cdot n)+\hat{l} \cdot n \quad \text { on } \quad \partial B_{0}, \tag{3.14}
\end{equation*}
$$

$\varphi_{0}$ denoting the solution of (3.2)-(3.4), and $\operatorname{grad}_{\partial \Omega}$ standing for the tangential gradient, defined as

$$
\operatorname{grad}_{\partial \Omega} f:=\nabla f-(\nabla f \cdot n) n .
$$

To prove the proposition, we argue by contradiction. If the map $L$ is not one-to-one, then we can pick a pair $(\hat{h}, \hat{l}) \neq(0,0)$ such that $L(\hat{h}, \hat{l})=0$, i.e. $\psi_{\mid \Gamma_{m}}=$ const. Since $\left.\frac{\partial \psi}{\partial n}\right|_{\Gamma_{m}}=0$ and $\Delta \psi=0$ in $\Omega \backslash \overline{B_{0}}$, we infer that $\psi \equiv$ const in $\Omega \backslash \overline{B_{0}}$ by unique continuation. Therefore, (3.14) gives

$$
\begin{equation*}
0=-\hat{h} \cdot n \frac{\partial^{2} \varphi_{0}}{\partial n^{2}}+\left(\nabla \varphi_{0}-l_{0}\right) \cdot \operatorname{grad}_{\partial \Omega}(\hat{h} \cdot n)+\hat{l} \cdot n \quad \text { on } \quad \partial B_{0} \tag{3.15}
\end{equation*}
$$

Note that $\hat{h} \neq 0$, otherwise $\hat{l}=0$ by (3.15). Let $(r, \theta)$ denote the polar coordinates with respect to the origin, and let $e_{r}:=(\cos \theta, \sin \theta)$ and $e_{\theta}:=e_{r}^{\perp}=(-\sin \theta, \cos \theta)$. Then $e_{r}=-n$ on $\partial B_{0}$, so $\partial^{2} \varphi_{0} / \partial n^{2}=\partial^{2} \varphi_{0} / \partial r^{2}$. Since $g \in H^{1}(\partial \Omega)$, we see that $\varphi_{0} \in H^{\frac{5}{2}}(\Omega)$, hence all the second derivatives of $\varphi_{0}$ possess traces in $L^{2}\left(\partial B_{0}\right)$. In particular,

$$
\begin{equation*}
0=\Delta \varphi_{0}=\frac{\partial^{2} \varphi_{0}}{\partial r^{2}}+\frac{1}{r} \frac{\partial \varphi_{0}}{\partial r}+\frac{1}{r^{2}} \frac{\partial^{2} \varphi_{0}}{\partial \theta^{2}} \quad \text { on } \quad \partial B_{0} . \tag{3.16}
\end{equation*}
$$

On the other hand,
$\nabla \varphi_{0}=\frac{\partial \varphi_{0}}{\partial r} e_{r}+\frac{1}{r} \frac{\partial \varphi_{0}}{\partial \theta} e_{\theta} \quad$ and $\quad \operatorname{grad}_{\partial \Omega}(\hat{h} \cdot n)=\frac{\partial \hat{h} \cdot n}{\partial \theta} e_{\theta} \quad$ on $\quad \partial B_{0}$.
Using (3.15)-(3.17), we obtain

$$
\begin{aligned}
0 & =-\hat{h} \cdot n\left(-\frac{\partial \varphi_{0}}{\partial r}-\frac{\partial^{2} \varphi_{0}}{\partial \theta^{2}}\right)+\left(\frac{\partial \varphi_{0}}{\partial r} e_{r}+\frac{\partial \varphi_{0}}{\partial \theta} e_{\theta}-l_{0}\right) \cdot \frac{\partial \hat{h} \cdot n}{\partial \theta} e_{\theta}+\hat{l} \cdot n \\
& =-\hat{h} \cdot e_{r}\left(l_{0} \cdot e_{r}+\frac{\partial^{2} \varphi_{0}}{\partial \theta^{2}}\right)-\left(\frac{\partial \varphi_{0}}{\partial \theta}-l_{0} \cdot e_{\theta}\right) \frac{\partial \hat{h} \cdot e_{r}}{\partial \theta}-\hat{l} \cdot e_{r}
\end{aligned}
$$

and hence, gathering together the second derivatives with respect to $\theta$,

$$
\begin{equation*}
\frac{\partial}{\partial \theta}\left(\left(\hat{h} \cdot e_{r}\right) \frac{\partial \varphi_{0}}{\partial \theta}\right)=\left(\hat{h} \cdot e_{\theta}\right)\left(l_{0} \cdot e_{\theta}\right)-\left(\hat{h} \cdot e_{r}\right)\left(l_{0} \cdot e_{r}\right)-\hat{l} \cdot e_{r} . \tag{3.18}
\end{equation*}
$$

Let $M_{0}$ and $M_{0}^{\prime}$ be the two points $M \in \partial B_{0}$ at which $\hat{h} \cdot e_{r}(M)=0$, and let $\theta(M)$ denote the angle $\left(\overrightarrow{O M_{0}}, \overrightarrow{O M}\right)$. Then, by (3.18),

$$
\begin{aligned}
\hat{h} \cdot e_{r}(M) \frac{\partial \varphi_{0}}{\partial \theta}(M) & =\int_{0}^{\theta(M)} \frac{\partial}{\partial \theta}\left(\hat{h} \cdot e_{r} \frac{\partial \varphi_{0}}{\partial \theta}\right) \mathrm{d} \theta \\
& =\int_{0}^{\theta(M)}\left[\left(\hat{h} \cdot e_{\theta}\right)\left(l_{0} \cdot e_{\theta}\right)-\left(\hat{h} \cdot e_{r}\right)\left(l_{0} \cdot e_{r}\right)\right] \mathrm{d} \theta-\int_{0}^{\theta(M)} \hat{l} \cdot e_{r} \mathrm{~d} \theta .
\end{aligned}
$$

As $\mathrm{d} e_{r} / \mathrm{d} \theta=e_{\theta}$ and $\mathrm{d} e_{\theta} / \mathrm{d} \theta=-e_{r}$, we obtain at once that

$$
\int_{0}^{\theta(M)}\left[\left(\hat{h} \cdot e_{\theta}\right)\left(l_{0} \cdot e_{\theta}\right)-\left(\hat{h} \cdot e_{r}\right)\left(l_{0} \cdot e_{r}\right)\right] \mathrm{d} \theta=\left[\left(\hat{h} \cdot e_{r}\right)\left(l_{0} \cdot e_{\theta}\right)\right]_{0}^{\theta(M)}
$$

and

$$
\int_{0}^{\theta(M)} \hat{l} \cdot e_{r} \mathrm{~d} \theta=-\hat{l} \cdot\left[e_{\theta}\right]_{0}^{\theta(M)}
$$

We conclude that

$$
\hat{h} \cdot e_{r} \frac{\partial \varphi_{0}}{\partial \theta}=\left[\left(\hat{h} \cdot e_{r}\right)\left(l_{0} \cdot e_{\theta}\right)+\hat{l} \cdot e_{\theta}\right]_{0}^{\theta(M)}
$$

Pick now $M=M_{0}^{\prime}$, so that $\hat{h} \cdot e_{r}\left(M_{0}^{\prime}\right)=0$. We obtain

$$
0=\hat{l} \cdot\left(e_{\theta}\left(M_{0}^{\prime}\right)-e_{\theta}\left(M_{0}\right)\right)=-2 \hat{l} \cdot e_{\theta}\left(M_{0}\right)
$$

It follows that $\hat{l} \in \operatorname{Span}\left\{e_{r}\left(M_{0}\right)\right\}$, so

$$
\hat{l} \cdot \hat{h}=0
$$

Writing $\hat{l}=\lambda \hat{h}^{\perp}$ for some constant $\lambda$, we have that $\hat{l} \cdot e_{\theta}=\lambda \hat{h} \cdot e_{r}$, and thus

$$
\hat{h} \cdot e_{r} \frac{\partial \varphi_{0}}{\partial \theta}=\left(\hat{h} \cdot e_{r}\right)\left(l_{0} \cdot e_{\theta}\right)+\lambda \hat{h} \cdot e_{r}
$$

Dividing by $\hat{h} \cdot e_{r}$ (which is non-null for $M \neq M_{0}, M_{0}^{\prime}$ ) and integrating over $\theta$, we obtain $\varphi_{0}=l_{0} \cdot e_{r}+\lambda \theta+\mu$, where $\mu$ denotes another constant. The function $\varphi_{0}$ is therefore a solution to the system

$$
\begin{align*}
& \Delta \xi=0 \quad \text { on } \quad \Omega \backslash \overline{B_{0}},  \tag{3.19}\\
& \xi=l_{0} \cdot x+\lambda \theta+\mu \quad \text { on } \quad \partial B_{0},  \tag{3.20}\\
& \frac{\partial \xi}{\partial n}=l_{0} \cdot n \quad \text { on } \quad \partial B_{0} . \tag{3.21}
\end{align*}
$$

An obvious solution of (3.19)-(3.21) on $\mathbb{C} \backslash\left(\mathbb{R}^{-} \cup \overline{B_{0}}\right)$ is given by $\xi=l_{0} \cdot x+\lambda \theta+\mu$. By the unique continuation property, we conclude that $\varphi_{0}=\xi$. As $\varphi_{0}$ is continuous on $\bar{\Omega} \backslash B_{0}$, we infer that $\lambda=0$, and that $g=\partial \varphi_{0} / \partial n=l_{0} \cdot n$ on $\partial \Omega$, which is a contradiction. The proof of proposition 3.2 is complete.

Combining theorems 2.1 and 3.1, we can state a semi-global stability result.
Corollary 3.3. Let $g$ be as in theorem 3.1, and let $\Lambda$ be the map defined in (3.8). Let $K \subset \Omega_{a}$ be a compact set, and let $R>0$ be a given number. Then there exists a constant $C=C(K, R, g)>0$ such that for all $\left(h_{1}, l_{1}\right),\left(h_{2}, l_{2}\right) \in K \times \overline{B_{R}(0)}$ it holds

$$
\begin{equation*}
\left\|\Lambda\left(h_{1}, l_{1}, g\right)-\Lambda\left(h_{2}, l_{2}, g\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \geqslant C\left\|\left(h_{1}-h_{2}, l_{1}-l_{2}\right)\right\|_{\mathbb{R}^{4}} \tag{3.22}
\end{equation*}
$$

Proof. For any $\left(h_{0}, l_{0}\right) \in K \times \overline{B_{R}(0)}$, let $L_{h_{0}, l_{0}}: \mathbb{R}^{4} \rightarrow H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}$ denote the linear map $L_{h_{0}, l_{0}}=\left.\mathrm{d} \Lambda\left(h_{0}, l_{0}, g\right)\right|_{\mathbb{R}^{2} \times \mathbb{R}^{2} \times\{0\}}$. Using the continuity of the map $\left(h_{0}, l_{0}, \hat{h}, \hat{l}\right) \mapsto L_{h_{0}, l_{0}}(\hat{h}, \hat{l})$, the compactness of $K \times \overline{B_{R}(0)} \times S^{3}$ and proposition 3.2, we infer the existence of two constants $C_{1}, C_{2}$ such that for any $\left(h_{0}, l_{0}\right) \in K \times \overline{B_{R}(0)}$,

$$
\begin{equation*}
C_{1}\|(\hat{h}, \hat{l})\| \leqslant\left\|L_{h_{0}, k_{0}}(\hat{h}, \hat{l})\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \leqslant C_{2}\|(\hat{h}, \hat{l})\| \quad \forall(\hat{h}, \hat{l}) \in \mathbb{R}^{4} \tag{3.23}
\end{equation*}
$$

On the other hand, the map $\left(h_{0}, l_{0}\right) \mapsto L_{h_{0}, l_{0}}$ being uniformly continuous on the compact set $K \times \overline{B_{R}(0)}$, we can find a small number $\delta>0$ such that if $\left(h_{0}, l_{0}\right),\left(h_{0}^{\prime}, l_{0}^{\prime}\right) \in K \times \overline{B_{R}(0)}$ satisfy $\left\|\left(h_{0}, l_{0}\right)-\left(h_{0}^{\prime}, l_{0}^{\prime}\right)\right\|<\delta$, then

$$
\begin{equation*}
\left\|\left(L_{h_{0}, l_{0}}-L_{h_{0}^{\prime}, l_{0}}\right)(\hat{h}, \hat{l})\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \leqslant \frac{C_{1}}{2}\|(\hat{h}, \hat{l})\| \quad \forall(\hat{h}, \hat{l}) \in \mathbb{R}^{4} \tag{3.24}
\end{equation*}
$$

Pick two pairs $\left(h_{1}, l_{1}\right),\left(h_{2}, l_{2}\right)$ in $K \times \overline{B_{R}(0)}$. Assume first that $\left\|\left(h_{1}-h_{2}, l_{1}-l_{2}\right)\right\|<\delta$, and set $\hat{h}=h_{2}-h_{1}, \hat{l}=l_{2}-l_{1}$. Then we have

$$
\begin{aligned}
\Lambda\left(h_{2}, l_{2}, g\right)-\Lambda\left(h_{1}, l_{1}, g\right) & =\int_{0}^{1} \frac{\mathrm{~d}}{\mathrm{~d} s} \Lambda\left(h_{1}+s \hat{h}, l_{1}+s \hat{l}, g\right) \mathrm{d} s \\
& =\int_{0}^{1} L_{h_{1}+s \hat{h}, l_{1}+s \hat{l}}(\hat{h}, \hat{l}) \mathrm{d} s \\
& =L_{h_{1}, l_{1}}(\hat{h}, \hat{l})+\int_{0}^{1}\left[L_{h_{1}+s \hat{h}, l_{1}+s \hat{l}}-L_{h_{1}, l_{1}}\right](\hat{h}, \hat{l}) \mathrm{d} s
\end{aligned}
$$

Using (3.23) and (3.24), we infer that

$$
\left\|\Lambda\left(h_{2}, l_{2}, g\right)-\Lambda\left(h_{1}, l_{1}, g\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \geqslant \frac{C_{1}}{2}\|(\hat{h}, \hat{l})\| .
$$

If $\left\|\left(h_{1}-h_{2}, l_{1}-l_{2}\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \geqslant \delta$, then

$$
\left\|\Lambda\left(h_{2}, l_{2}, g\right)-\Lambda\left(h_{1}, l_{1}, g\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}} \geqslant \frac{m}{d}\|(\hat{h}, \hat{l})\|,
$$

where $m$ denotes the minimum of the continuous map

$$
\left(h_{1}, l_{1}, h_{2}, l_{2}\right) \mapsto\left\|\Lambda\left(h_{2}, l_{2}, g\right)-\Lambda\left(h_{1}, l_{1}, g\right)\right\|_{H^{s+1}\left(\Gamma_{m}\right) / \mathbb{R}}
$$

taken over the compact set $\left\{\left(h_{1}, l_{1}, h_{2}, l_{2}\right) \in\left(K \times \overline{B_{R}(0)}\right)^{2},\left\|\left(h_{2}-h_{1}, l_{2}-l_{1}\right)\right\| \geqslant \delta\right\}$ and $d$ denotes the diameter of that compact set. Note that $m>0$, according to theorem 2.1. The proof is completed by taking $C=\min \left(m / d, C_{1} / 2\right)$.

Remark 3.4. The constant $C$ in theorem 3.1 is not given explicitly, as the method of proof relies upon compactness arguments. One way to give explicit values for this constant would be to estimate precisely the first eigenvalue of the linear map $L$ in proposition 3.2.

## Annexe

We prove in this annexe that the solution $\psi$ of (2.17)-(2.19) cannot be of class $C^{2}$ at the point $M_{+}$when $\delta \geqslant 1 / \sqrt{2}$. Using the transformation $T_{2} \circ T_{1}$ which is analytic in a neighbourhood of $M_{+}$, its inverse being also analytic near the origin, this is equivalent to show that the solution $\psi_{2}$ of (2.24)-(2.27) is not of class $C^{2}$ at the origin. This is a direct consequence of the following result.

Proposition 3.5. Assume that $\delta \geqslant 1 / \sqrt{2}$. Then the second derivative $\partial^{2} \psi_{2} / \partial y_{2}^{2}(0,0)$ fails to exist.

Proof. Recall that $d_{-1}=\left\{z_{2} \in \mathbb{C}^{*}\right.$; $\left.\arg z_{2}=\theta^{*}\right\}$, where we denote $\theta^{*}=\theta_{-1}=\frac{\pi-\theta}{2}$. Obviously, $\delta \geqslant 1 / \sqrt{2}$ if and only if $\theta^{*} \leqslant \pi / 4$. Proving that $\psi_{2}$ is not of class $C^{2}$ in a neighbourhood of the origin is quite easy when $\theta^{*}=\pi / 4$. Indeed, in that case we may use the system of orthogonal coordinates $(u, v)=(\sqrt{2})^{-1}\left(x_{2}+y_{2},-x_{2}+y_{2}\right)$, with $d_{-1}$ and $d_{0}$ as the coordinate axes, to compute the Laplacian of $\psi_{2}$ at the origin. If $\psi_{2}$ were of class $C^{2}$ near the origin, we should have

$$
0=\Delta \psi_{2}(0,0)=\frac{\partial^{2} \psi_{2}}{\partial u^{2}}(0,0)+\frac{\partial^{2} \psi_{2}}{\partial v^{2}}(0,0)
$$

But straightforward computations give

$$
\frac{\partial^{2} \psi_{2}}{\partial u^{2}}(0,0)=\frac{\partial^{2} \psi_{2}}{\partial v^{2}}(0,0)=-2 c^{-3} \neq 0
$$

When $\theta^{*} \neq \pi / 4$, the value of the Laplacian of $\psi_{2}$ at the origin cannot be deduced from the knowledge of the second derivatives of $\psi_{2}$ along the axes $d_{-1}, d_{0}$. An exact computation of $\psi_{2}$ is therefore required. This is done in the following step.

Step 1. Reduction to a Dirichlet problem in the half-plane $\mathbb{C}^{+}$.
Let us introduce the number $\alpha>1$ defined by

$$
\alpha^{-1}:=1-\frac{\theta^{*}}{\frac{\pi}{2}}
$$

Introduce the transformation $T_{3}\left(z_{2}\right):=z_{3}=x_{3}+\mathrm{i} y_{3}=\left(\mathrm{e}^{-\mathrm{i} \theta^{*}} z_{2}\right)^{\alpha}$. Then
$T_{3}\left(d_{-1}\right)=\mathbb{R}^{+}, \quad T_{3}\left(d_{0}\right)=\mathbb{R}^{-} \quad$ and $\quad T_{3}\left(C_{2}\right)=\mathbb{C}^{+}:=\{z=x+\mathrm{i} y ; y>0\}$.
Define $\psi_{3}$ on $\mathbb{C}^{+}$by $\psi_{3}\left(z_{3}\right)=\psi_{2}\left(z_{2}\right)$. To determine the values of $\psi_{3}$ on $\partial \mathbb{C}^{+}=\mathbb{R}$, we need to express $z_{2}$ as a function of $z_{3}$ when $z_{3}=x_{3} \in \mathbb{R}$.

If $z_{3}=x_{3} \in \mathbb{R}^{+}$, then $z_{2}=\mathrm{e}^{\mathrm{i} \theta^{*}} x_{3}^{1 / \alpha}$, and hence

$$
\psi_{3}\left(z_{3}\right)=\frac{y_{2}}{\left(x_{2}+c\right)^{2}+y_{2}^{2}}=\frac{x_{3}^{1 / \alpha} \sin \theta^{*}}{\left(c+x_{3}^{1 / \alpha} \cos \theta^{*}\right)^{2}+\left(x_{3}^{1 / \alpha} \sin \theta^{*}\right)^{2}}
$$

If $z_{3}=x_{3} \in \mathbb{R}^{-}$, then $z_{2}=\mathrm{e}^{\mathrm{i} \theta^{*}}\left|x_{3}\right|^{1 / \alpha} \mathrm{e}^{\mathrm{i} \pi / \alpha}=\left|x_{3}\right|^{1 / \alpha} \mathrm{e}^{\mathrm{i}\left(\pi-\theta^{*}\right)}$, and hence

$$
\psi_{3}\left(z_{3}\right)=-\frac{y_{2}}{\left(x_{2}+c\right)^{2}+y_{2}^{2}}=-\frac{\left|x_{3}\right|^{1 / \alpha} \sin \theta^{*}}{\left(c-\left|x_{3}\right|^{1 / \alpha} \cos \theta^{*}\right)^{2}+\left(\left|x_{3}\right|^{1 / \alpha} \sin \theta^{*}\right)^{2}} .
$$

Let us set $x^{\{\alpha\}}:=\operatorname{sgn}(x)|x|^{\alpha}$ for any $x \in \mathbb{R}$ and any $\alpha \in \mathbb{R}^{+}$. Then $\psi_{3}$ solves the system

$$
\begin{align*}
& \Delta \psi_{3}=0 \quad \text { in } \quad \mathbb{C}^{+},  \tag{3.25}\\
& \psi_{3}\left(x_{3}\right)=\frac{x_{3}^{\{1 / \alpha\}} \sin \theta^{*}}{\left(c+x_{3}^{\{1 / \alpha\}} \cos \theta^{*}\right)^{2}+\left(x_{3}^{\{1 / \alpha\}} \sin \theta^{*}\right)^{2}} \quad \text { on } \quad \mathbb{R},  \tag{3.26}\\
& \psi_{3}\left(z_{3}\right) \rightarrow 0 \quad \text { as } \quad z_{3} \rightarrow \infty, z_{3} \in \mathbb{C}^{+} . \tag{3.27}
\end{align*}
$$

Using the Poisson formula (see, e.g. [7]), we obtain that

$$
\psi_{3}\left(z_{3}\right)=\frac{y_{3}}{\pi} \int_{-\infty}^{\infty} \frac{\psi_{3}(t)}{\left|z_{3}-t\right|^{2}} \mathrm{~d} t
$$

Going back to the variable $z_{2}=r_{2} \mathrm{e}^{\mathrm{i} \varphi_{2}}$, and using the fact that $z_{3}=\left(\mathrm{e}^{-\mathrm{i} \theta^{*}} z_{2}\right)^{\alpha}=r_{2}^{\alpha} \mathrm{e}^{\mathrm{i} \alpha\left(\varphi_{2}-\theta^{*}\right)}$, we conclude that

$$
\begin{gathered}
\psi_{2}\left(z_{2}\right)=\frac{r_{2}^{\alpha} \sin \alpha\left(\varphi_{2}-\theta^{*}\right)}{\pi} \int_{-\infty}^{\infty} \frac{t^{\{1 / \alpha\}} \sin \theta^{*}}{\left(c+t^{\{1 / \alpha\}} \cos \theta^{*}\right)^{2}+\left(t^{\{1 / \alpha\}} \sin \theta^{*}\right)^{2}} \\
\times \frac{\mathrm{d} t}{\left(t-r_{2}^{\alpha} \cos \alpha\left(\varphi_{2}-\theta^{*}\right)\right)^{2}+\left(r_{2}^{\alpha} \sin \alpha\left(\varphi_{2}-\theta^{*}\right)\right)^{2}}
\end{gathered}
$$

For the value $\varphi_{2}=\pi / 2$, we obtain
$\psi_{2}\left(0, y_{2}\right)=\psi_{2}\left(r_{2} \mathrm{e}^{\mathrm{i} \pi / 2}\right)=\frac{r_{2}^{\alpha}}{\pi} \int_{-\infty}^{\infty} \frac{t^{\{1 / \alpha\}} \sin \theta^{*}}{\left(c+t^{\{1 / \alpha\}} \cos \theta^{*}\right)^{2}+\left(t^{\{1 / \alpha\}} \sin \theta^{*}\right)^{2}} \frac{\mathrm{~d} t}{t^{2}+\left|r_{2}\right|^{2 \alpha}}$.
In a second step, we show that $\lim _{y_{2} \rightarrow 0^{+}} \frac{\partial^{2} \psi_{2}}{\partial y_{2}^{2}}\left(0, y_{2}\right)$ exists and is finite if and only if $\alpha>2$ (i.e. $\delta<1 / \sqrt{2}$ ).

Step 2. Estimation of the second derivative of $\psi_{2}$ with respect to $y_{2}$ at the origin.

Let us introduce the functions

$$
f(t)=\frac{t^{\{1 / \alpha\}}}{\left(c+t^{\{1 / \alpha\}} \cos \theta^{*}\right)^{2}+\left(t^{\{1 / \alpha\}} \sin \theta^{*}\right)^{2}}, \quad g(r)=r \int_{-\infty}^{\infty} f(t) \frac{\mathrm{d} t}{t^{2}+r^{2}}
$$

defined for $t \in \mathbb{R}$ and $r \in(0,+\infty)$, respectively. Note that $f \in C(\mathbb{R}) \cap L^{\infty}(\mathbb{R})$ with $|f(t)| \leqslant \operatorname{const}\left(1+|t|^{1 / \alpha}\right)^{-1}$. We aim to prove that for $\alpha \leqslant 2,\left|\mathrm{~d}^{2}\left[g\left(r^{\alpha}\right)\right] / \mathrm{d} r^{2}\right| \rightarrow \infty$ as $r \rightarrow 0^{+}$. Using the Lebesgue theorem and a change of variables, we have

$$
\begin{aligned}
g^{\prime}(r) & =\int_{-\infty}^{\infty} \frac{f(t)}{t^{2}+r^{2}} \mathrm{~d} t-2 r^{2} \int_{-\infty}^{\infty} \frac{f(t)}{\left(t^{2}+r^{2}\right)^{2}} \mathrm{~d} t \\
& =r^{-1} \int_{-\infty}^{\infty} \frac{f(r s)}{s^{2}+1} \mathrm{~d} s-2 r^{-1} \int_{-\infty}^{\infty} \frac{f(r s)}{\left(s^{2}+1\right)^{2}} \mathrm{~d} s \\
g^{\prime \prime}(r) & =-6 r \int_{-\infty}^{\infty} \frac{f(t)}{\left(t^{2}+r^{2}\right)^{2}} \mathrm{~d} t+8 r^{3} \int_{-\infty}^{\infty} \frac{f(t)}{\left(t^{2}+r^{2}\right)^{3}} \mathrm{~d} t \\
& =-6 r^{-2} \int_{-\infty}^{\infty} \frac{f(r s)}{\left(s^{2}+1\right)^{2}} \mathrm{~d} s+8 r^{-2} \int_{-\infty}^{\infty} \frac{f(r s)}{\left(s^{2}+1\right)^{3}} \mathrm{~d} s
\end{aligned}
$$

It may be seen that $\mathrm{d}\left[g\left(r^{\alpha}\right)\right] / \mathrm{d} r \rightarrow 0$ as $r \rightarrow 0$; that is, $\partial \psi_{2} / \partial y_{2}(0,0)=$ $\lim _{y_{2} \rightarrow 0^{+}} \partial \psi_{2} / \partial y_{2}\left(0, y_{2}\right)=0$. On the other hand,

$$
\begin{aligned}
\frac{\mathrm{d}^{2}\left[g\left(r^{\alpha}\right)\right]}{\mathrm{d} r^{2}} & =\left(\alpha r^{\alpha-1}\right)^{2} g^{\prime \prime}\left(r^{\alpha}\right)+\alpha(\alpha-1) r^{\alpha-2} g^{\prime}\left(r^{\alpha}\right) \\
& =\alpha r^{-1}\left((\alpha-1) I_{1}(r)+(-8 \alpha+2) I_{2}(r)+8 \alpha I_{3}(r)\right)
\end{aligned}
$$

where
$I_{j}(r):=\int_{-\infty}^{\infty} \frac{s^{\{1 / \alpha\}}}{c^{2}+2 c\left(\cos \theta^{*}\right) r s^{\{1 / \alpha\}}+\left(r s^{\{1 / \alpha\}}\right)^{2}} \frac{\mathrm{~d} s}{\left(s^{2}+1\right)^{j}} \quad$ for $\quad j \in\{1,2,3\}$.
As $s \mapsto s^{\{1 / \alpha\}}$ is odd, we see that $I_{j}(0)=0$ for each $j$. Since

$$
\left|\frac{2 c\left(\cos \theta^{*}\right) s^{\{1 / \alpha\}}+2 r|s|^{2 / \alpha}}{\left(c^{2}+2 c\left(\cos \theta^{*}\right) r s^{\{1 / \alpha\}}+r^{2}|s|^{2 / \alpha}\right)^{2}} \frac{s^{\{1 / \alpha\}}}{\left(s^{2}+1\right)^{j}}\right| \leqslant \text { const } \frac{|s|^{2 / \alpha}}{\left(s^{2}+1\right)^{j}},
$$

we infer from the Lebesgue theorem that

$$
I_{j}^{\prime}(r)=-\int_{-\infty}^{\infty} \frac{2 c\left(\cos \theta^{*}\right) s^{\{1 / \alpha\}}+2 r|s|^{2 / \alpha}}{\left(c^{2}+2 c\left(\cos \theta^{*}\right) r s^{\{1 / \alpha\}}+r^{2}|s|^{2 / \alpha}\right)^{2}} \frac{s^{\{1 / \alpha\}}}{\left(s^{2}+1\right)^{j}} \mathrm{~d} s \quad \forall r>0
$$

for any $\alpha>1$ if $j \geqslant 2$, and for any $\alpha>2$ if $j=1$.
If $\alpha>2$, then $\lim _{r \rightarrow 0^{+}} I_{j}(r) / r=I_{j}^{\prime}(0)$ for $j=1,2,3$, thus $\frac{\partial^{2} \psi_{2}}{\partial y_{2}^{2}}(0,0)=$ $\lim _{y_{2} \rightarrow 0^{+}} \frac{\partial^{2} \psi_{2}}{\partial y_{2}^{2}}\left(0, y_{2}\right)$ exists.

Assume now that $1<\alpha \leqslant 2$. Once again, $\lim _{r \rightarrow 0^{+}} I_{j}(r) / r=I_{j}^{\prime}(0) \in \mathbb{R}$ for $j=2,3$. For $j=1$, we claim that $I_{1}(r) / r \rightarrow-\infty$ as $r \rightarrow 0^{+}$.

Claim. $\lim _{r \rightarrow 0^{+}} I_{1}(r) / r=-\infty$.
Indeed, letting $\sigma=r^{\alpha} s$, we have that

$$
I_{1}(r)=r^{-1-\alpha} \int_{-\infty}^{\infty} \frac{\sigma^{\{1 / \alpha\}}}{c^{2}+2 c\left(\cos \theta^{*}\right) \sigma^{\{1 / \alpha\}}+|\sigma|^{2 / \alpha}} \frac{\mathrm{d} \sigma}{\left(\sigma / r^{\alpha}\right)^{2}+1},
$$

hence

$$
\begin{aligned}
\frac{I_{1}(r)}{r^{\alpha-1}} & =\int_{-\infty}^{\infty} \frac{\sigma^{\{1 / \alpha\}}}{c^{2}+2 c\left(\cos \theta^{*}\right) \sigma^{\{1 / \alpha\}}+|\sigma|^{2 / \alpha}} \frac{\mathrm{d} \sigma}{\sigma^{2}+r^{2 \alpha}} \\
& =-4 c\left(\cos \theta^{*}\right) J(r)
\end{aligned}
$$

where
$J(r):=\int_{0}^{\infty} \frac{\sigma^{2 / \alpha}}{\left(c^{2}+2 c\left(\cos \theta^{*}\right) \sigma^{1 / \alpha}+\sigma^{2 / \alpha}\right)\left(c^{2}-2 c\left(\cos \theta^{*}\right) \sigma^{1 / \alpha}+\sigma^{2 / \alpha}\right)} \frac{\mathrm{d} \sigma}{\sigma^{2}+r^{2 \alpha}}$.
An application of the monotone convergence theorem yields $J(r) \rightarrow J(0)$ as $r \rightarrow 0^{+}$, where $J(0) \in(0,+\infty)$ for $\alpha \in(1,2)$ and $J(0)=+\infty$ for $\alpha=2$. The claim follows at once. The proof of proposition 3.5 is complete.

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