# INVERSE PROBLEMS FOR ELLIPTIC EQUATIONS WITH FRACTIONAL POWER TYPE NONLINEARITIES 

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#### Abstract

We study inverse problems for semilinear elliptic equations with fractional power type nonlinearities. Our arguments are based on the higher order linearization method, which helps us to solve inverse problems for certain nonlinear equations in cases where the solution for a corresponding linear equation is not known. By using a fractional order adaptation of this method, we show that the results of [LLLS20a, LLLS20b] remain valid for general power type nonlinearities.


Keywords. Inverse boundary value problem, Calderón problem, partial data, semilinear elliptic equations, higher order linearization, transversally anisotropic manifold.

## Contents

1. Introduction ..... 1
2. Preliminaries ..... 7
3. Global uniqueness in Euclidean space ..... 14
4. Global uniqueness in Riemannian manifolds ..... 21
References ..... 23

## 1. Introduction

In this work we study inverse problems for semilinear elliptic equations with fractional power type nonlinearities, extending the earlier results in [LLLS20a, LLLS20b] from integer powers to fractional powers. Here, when we say $r$ is fractional we mean $r \in \mathbb{R} \backslash \mathbb{Z}$. Let $r>1$ be fractional and let $\Omega \subset \mathbb{R}^{n}$ be a bounded domain with $C^{\infty}$-smooth boundary $\partial \Omega$, for $n \geq 2$. Consider the semilinear elliptic equation

$$
\begin{cases}\Delta u+q(x)|u|^{r-1} u=0 & \text { in } \Omega  \tag{1.1}\\ u=f & \text { on } \partial \Omega\end{cases}
$$

where $q \in C^{\alpha}(\bar{\Omega})$ is a potential function and $C^{\alpha}$ is the space of $\alpha$-Hölder continuous functions. By assuming a suitable smallness condition on the boundary data $f$, one can obtain the well-posedness of the Dirichlet problem (1.1) for small solutions (see Section 2). One can then define the corresponding Dirichlet-to-Neumann (DN) $\operatorname{map} \Lambda_{q}$ of (1.1) by

$$
\Lambda_{q}: C^{2, \alpha}(\partial \Omega) \rightarrow C^{1, \alpha}(\partial \Omega),\left.\quad f \mapsto \partial_{\nu} u_{f}\right|_{\partial \Omega}
$$

for some $0<\alpha<1$, where $u_{f} \in C^{2, \alpha}(\bar{\Omega})$ is the unique small solution of (1.1), and $\nu$ is the unit outer normal on $\partial \Omega$. We will consider the following problem:

- Inverse Problem 1: Determine the potential $q$ from the knowledge of $\Lambda_{q}$.

A typical method in the study of inverse boundary value problems for nonlinear elliptic equations was initiated by Isakov [Isa93], where he introduced the first linearization of the given (nonlinear) DN map. More precisely, the first linearization allows one to reduce the nonlinear equations to the linear equations, and one can adapt some known results for the linear equations to solve certain inverse problems for the nonlinear equations. Meanwhile, the second order linearization has been successfully applied in solving inverse problems, see [AZ17, CNV19, KN02, Sun96, SU97].

Throughout this paper the number $r>1$ is fractional, and the solution $u$ is real valued but may change sign, so it is natural to consider $q(x)|u|^{r-1} u$ instead of $q(x) u^{r}$ to have well-defined nonlinear term. Note also that at least when $n=1$ the case $0<r<1$ would roughly correspond to the second order differential equation $u^{\prime \prime}=F(u)$, where $F$ is not Lipschitz. In this case, it is well-known that uniqueness of solutions can fail, so the assumption $r>1$ is reasonable. Let us write $r=k+\alpha>1$ for some $k \in \mathbb{N}$ and $\alpha \in(0,1)$ in the rest of this work.

In case of $r=m \in \mathbb{N}$ and nonlinear term $q(x) u^{m}$, corresponding inverse problems were first investigated in [FO20, LLLS20a], and related problems have been further studied in many works. For example, the articles [LLLS20b, KU20c, KU20b] studied related inverse problems for semilinear elliptic equations with partial data. In [LL20, Lin20, LO20], the authors studied inverse problems for fractional semilinear elliptic equations. In [LZ20, KU20a, CF20, KKU20], the authors studied partial data inverse problems for the nonlinear magnetic Schrödinger and conductivity equations. The nonlinearities in these articles are typically integer power type, or holomorphic in $u$ and $\nabla u$ (i.e. sums of integer powers).

The main tool in solving these inverse problems is based on the higher order linearization technique, where one introduces extra small parameters for the Dirichlet data to reduce inverse problems for nonlinear elliptic equations into statements involving solutions of simpler linear elliptic equations. In the case of nonlinearity $q(x) u^{m}$ where $m \in \mathbb{N}$, this just means that we are looking at the $m$ th order Fréchet derivative of the nonlinear measurement operator. For a nonlinearity of fractional order $r=k+\alpha$, we will in some sense need to use the $\alpha$ th fractional derivative of the $k$ th Fréchet derivative instead. A somewhat related method was used in [CK20] for a $p$-Laplace type equation. Thanks to the higher order linearization method, one may solve related inverse problems for certain semilinear elliptic equations in cases where the analogous problems for the corresponding linear equations still remain open.

Let us state our first main result to answer Inverse Problem 1:
Theorem 1.1 (The Calderón problem with full data). Let $\Omega \subset \mathbb{R}^{n}$ be a connected bounded domain with $C^{\infty}$-smooth boundary $\partial \Omega$, for $n \geq 2$. Let $r>1$ be a fractional number, $q_{j} \in C^{\alpha}(\bar{\Omega})$ for some $0<\alpha<1$, and $\Lambda_{q_{j}}$ be the DN map of

$$
\begin{cases}\Delta u_{j}+q_{j}\left|u_{j}\right|^{r-1} u_{j}=0 & \text { in } \Omega \\ u_{j}=f & \text { on } \partial \Omega\end{cases}
$$

for $j=1,2$. Assume that $\Lambda_{q_{1}}(f)=\Lambda_{q_{2}}(f)$, for all $f \in C^{2, \alpha}(\partial \Omega)$ with $\|f\|_{C^{2, \alpha}(\partial \Omega)}<$ $\delta$, where $\delta>0$ is a sufficiently small number. Then

$$
q_{1}=q_{2} \text { in } \Omega
$$

Moreover, in dimensions $n \geq 3$ the statement holds true if we only assume that $\Lambda_{q_{1}}(f)=\Lambda_{q_{2}}(f)$ whenever $\|f\|_{C^{2, \alpha}(\partial \Omega)}<\delta$ and $f \geq 0$.

We remark that in certain applications it is natural to consider nonnegative Dirichlet data (see e.g. [RZ18]). Theorem 1.1 applies in this case when $n \geq 3$.

However, the methods for proving the other main theorems in this paper require sign-changing solutions, and we do not know if those results are valid if one only has access to measurements for nonnegative Dirichlet data.

We briefly explain the higher order linearization in the fractional power case. Let $(M, g)$ be a compact $C^{\infty}$ Riemannian manifold with a $C^{\infty}$ smooth boundary $\partial M$. Recall that $\Delta_{g}$ is the Laplace-Beltrami operator, given in local coordinates by

$$
\Delta_{g} u=\frac{1}{\operatorname{det}(g)^{1 / 2}} \sum_{a, b=1}^{n} \frac{\partial}{\partial x_{a}}\left(\operatorname{det}(g)^{1 / 2} g^{a b} \frac{\partial u}{\partial x_{b}}\right)
$$

where $g=\left(g_{a b}(x)\right)$ and $g^{-1}=\left(g^{a b}(x)\right)$. Throughout this work, we assume that $g=\left(g_{a b}\right)$ is uniformly elliptic. Let $q \in C^{\alpha}(M)$. In Proposition 2.3 we will see that by setting the Dirichlet data as

$$
f=\epsilon_{0} f_{0}+\ldots+\epsilon_{k} f_{k}
$$

and differentiating the equation (1.1) with respect to $\epsilon^{\prime}=\left(\epsilon_{1}, \ldots, \epsilon_{k}\right)$ we obtain a new equation

$$
\begin{equation*}
\Delta_{g} w^{\epsilon_{0}}(x)=-\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}}\left(q(x)\left|u_{f}\right|^{r-1} u_{f}\right)\right|_{\epsilon^{\prime}=0} \text { in } M \tag{1.2}
\end{equation*}
$$

where $w^{\epsilon_{0}}:=\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}} u_{f}\right|_{\epsilon^{\prime}=0}$ and $\left.w^{\epsilon_{0}}\right|_{\partial M}=\left.\epsilon_{0} f_{0}\right|_{\partial M}$.
Furthermore, eliminating $\epsilon_{0}^{\alpha}$ on the both sides of (1.2), by taking the limit $\epsilon_{0} \rightarrow 0$, we get

$$
\epsilon_{0}^{-\alpha} w^{\epsilon_{0}} \rightarrow w \text { in } C^{2, \alpha}(M), \quad \text { as } \epsilon_{0} \rightarrow 0
$$

where $w$ solves

$$
\Delta_{g} w=c_{r} q(x) \operatorname{sgn}\left(v_{0}\right)^{k-1}\left|v_{0}\right|^{\alpha} v_{1} \cdots v_{k} \text { in } M
$$

Here $c_{r}$ is the constant given by $c_{r}=-r(r-1) \cdots(r-(k-1)), \operatorname{sgn}\left(v_{0}(x)\right)$ is the sign of $v_{0}(x)$, and the functions $v_{\ell}$ are harmonic in $M$ with the corresponding boundary values $f_{\ell}$, for $\ell=0,1, \ldots, k$. Moreover, we will multiply this equation by an extra auxiliary harmonic function $v_{k+1}$ in $M$ with its boundary data $\left.v_{k+1}\right|_{\partial M}=f_{k+1}$. Now integrating over $M$ and using integration by parts, we see that from the knowledge of the DN map for the equation $\Delta_{g} u+q(x)|u|^{r-1} u=0$ in $M$ it is possible to determine the integrals

$$
c_{r} \int_{M} q(x) \operatorname{sgn}\left(v_{0}\right)^{k-1}\left|v_{0}\right|^{\alpha} v_{1} \cdots v_{k+1} d V
$$

It thus suffices to choose the boundary data $f_{\ell}$ for $\ell=0,1, \ldots, k$, so that $v_{0} \neq 0$ in $M$ and the scalar products $v_{1} \cdots v_{k+1}$ become dense in a suitable function space. This recovers the function $q$ (see Sections 3 and 4).

Next we study the Calderón problem with partial data for elliptic equations with fractional power type nonlinearities. Let $\Omega \subset \mathbb{R}^{n}$ be a connected bounded domain, and $\Gamma \subset \partial \Omega$ be a nonempty relatively open subset. By using the well-posedness of (1.1) (Proposition 2.1), one can define the corresponding partial DN map $\Lambda_{q}^{\Gamma}$ of (1.1) by

$$
\Lambda_{q}^{\Gamma}: C_{0}^{2, \alpha}(\Gamma) \rightarrow C^{1, \alpha}(\Gamma),\left.\quad f \mapsto \partial_{\nu} u_{f}\right|_{\Gamma}
$$

for some $0<\alpha<1$, where $u_{f} \in C^{2, \alpha}(\bar{\Omega})$ is the unique (small) solution of (1.1) (see Section 2) with $f \in C_{0}^{2, \alpha}(\Gamma)$. Then our second question is:

- Inverse Problem 2: Determine the potential $q$ from the knowledge of $\Lambda_{q}^{\Gamma}$.

Our second main result is to solve Inverse Problem 2:

Theorem 1.2 (Partial data). Let $\Omega \subset \mathbb{R}^{n}$ be a connected bounded domain with $C^{\infty}$-smooth boundary $\partial \Omega$, for $n \geq 2$, and $\Gamma \subset \partial \Omega$ be a nonempty relatively open subset. Let $r>1$ be a fractional number, $q_{j} \in C^{\alpha}(\bar{\Omega})$ for some $0<\alpha<1$, and $\Lambda_{q_{j}}^{\Gamma}$ be the $D N$ map of

$$
\begin{cases}\Delta u_{j}+q_{j}\left|u_{j}\right|^{r-1} u_{j}=0 & \text { in } \Omega \\ u_{j}=f & \text { on } \partial \Omega\end{cases}
$$

for $j=1,2$. If $\Lambda_{q_{1}}^{\Gamma}(f)=\Lambda_{q_{2}}^{\Gamma}(f)$, for all $f \in C_{0}^{2, \alpha}(\Gamma)$ with $\|f\|_{C_{0}^{2, \alpha}(\Gamma)}<\delta$, where $\delta>0$ is a sufficiently small number, then

$$
q_{1}=q_{2} \text { in } \Omega
$$

Moreover, one can consider more general nonlinear terms that are (asymptotic) sums of homogeneous functions. Let $\Omega \subset \mathbb{R}^{n}$ be a bounded domain with $C^{\infty}$ _ smooth boundary $\partial \Omega$.

Definition 1.1. Let $r_{l}, l \geq 1$, be real numbers with $1<r_{1}<r_{2}<\ldots$, and let $0<\alpha<1$. A function $a=a(x, y): \bar{\Omega} \times \mathbb{R} \rightarrow \mathbb{R}$ is polyhomogeneous, written

$$
a(x, y) \sim \sum_{l=1}^{\infty} b_{l}(x, y)
$$

if each $b_{l}(\cdot, y) \in C^{\alpha}(\bar{\Omega})$ is positively homogeneous of degree $r_{l}$ with respect to the $y$ variable, and if for any $N \geq 1$ there is $C_{N}>0$ so that the function $\beta_{N}:=a-\sum_{l=1}^{N-1} b_{l}$ (with $\beta_{1}=a$ ) is in $C_{\operatorname{loc}}^{1, \alpha}\left(\mathbb{R}, C^{\alpha}(\bar{\Omega})\right)$ and satisfies

$$
\begin{equation*}
\left\|\beta_{N}(\cdot, y)\right\|_{C^{\alpha}(\bar{\Omega})}+|y|\left\|\partial_{y} \beta_{N}(\cdot, y)\right\|_{C^{\alpha}(\bar{\Omega})} \leq C_{N}|y|^{r_{N}}, \quad|y| \leq 1 \tag{1.3}
\end{equation*}
$$

We will assume that $1+\alpha \leq r_{1}$ (this can be arranged by decreasing $\alpha$ ).
Note that the above definition (using $N=1$ ) implies that

$$
\begin{equation*}
a(x, 0)=\partial_{y} a(x, 0)=0 . \tag{1.4}
\end{equation*}
$$

A typical example of polyhomogeneous function $a(x, y)$ is a finite sum

$$
a(x, y)=\sum_{l=1}^{m} q_{l}(x) f_{l}(y)
$$

where $q_{l}(x) \in C^{\alpha}(\bar{\Omega})$ and $f_{l}(y)$ is positively homogeneous of degree $r_{l}$, i.e. $f_{l}(\lambda y)=$ $\lambda^{r_{l}} f_{l}(y)$ for $y \in \mathbb{R}$ and $\lambda>0$. One could also consider infinite sums of this type. In fact, functions $a(x, y)$ that are $C^{\alpha}$ in $x$, holomorphic or antiholomorphic in $y$, and satisfy (1.4) are polyhomogeneous with $r_{l}=l+1$ just by using Taylor expansions. It is worth emphasizing that since we are always considering small solutions, only the behaviour for small $|y|$ plays a role.

We also mention that the function $f(y)=|y|^{r-1} y$, at least roughly speaking, encompasses all positively homogeneous functions. Indeed, if $f$ is positively homogeneous of degree $r>0$, then $f$ is of the form

$$
f(y)= \begin{cases}y^{r} f(1), & \text { if } y \geq 0 \\ f(-|y|)=|y|^{r} f(-1), & \text { if } y<0\end{cases}
$$

The case $f(y)=|y|^{r-1} y$ is obtained by taking $f(1)=1$ and $f(-1)=-1$. This computation also shows that if $r=k+\alpha$ where $k \geq 1$ and $\alpha \in(0,1)$, then $f(y)$ is $C^{k}$ and $f^{(k)}(y)$ is $C^{\alpha}$.

Let us consider the following Dirichlet problem in a bounded smooth domain $\Omega \subset \mathbb{R}^{n}$

$$
\begin{cases}\Delta u+a(x, u)=0 & \text { in } \Omega  \tag{1.5}\\ u=f & \text { on } \partial \Omega\end{cases}
$$

where $a=a(x, y)$ is a polyhomogeneous function given by Definition 1.1. By Proposition 2.1, for any sufficiently small Dirichlet data $f \in C_{0}^{2, \alpha}(\Gamma)$ with $\Gamma \subset \partial \Omega$, one can define the corresponding (partial) DN map via

$$
\Lambda_{a}^{\Gamma}: C_{0}^{2, \alpha}(\Gamma) \rightarrow C^{1, \alpha}(\Gamma),\left.\quad f \mapsto \partial_{\nu} u_{f}\right|_{\Gamma}
$$

for some $0<\alpha<1$, where $u_{f} \in C^{2, \alpha}(\bar{\Omega})$ is the unique small solution of (1.5). The inverse problem is to determine the unknown function $a(x, y)$.

Theorem 1.3 (Partial data for general coefficients). Let $\Omega \subset \mathbb{R}^{n}$ be a connected bounded domain with $C^{\infty}$-smooth boundary $\partial \Omega$, for $n \geq 2$, and $\Gamma \subset \partial \Omega$ be a nonempty relatively open subset. Let us consider the equations

$$
\begin{equation*}
\Delta u+a_{j}(x, u)=0 \text { in } \Omega \tag{1.6}
\end{equation*}
$$

for $j=1,2$, where $a_{j}(x, y) \sim \sum_{l=1}^{\infty} b_{j, l}(x, y)$ is polyhomogeneous in the sense of Definition 1.1 where the orders $1<r_{1}<r_{2}<\ldots$ are the same for $j=1,2$. Let $\Lambda_{a_{j}}^{\Gamma}: C_{0}^{2, \alpha}(\Gamma) \rightarrow C^{1, \alpha}(\Gamma)$ be the (partial) DN maps of (1.6), for $j=1,2$. Assume that

$$
\Lambda_{a_{1}}^{\Gamma}(f)=\Lambda_{a_{2}}^{\Gamma}(f)
$$

for all $f \in C_{0}^{2, \alpha}(\Gamma)$ with $\|f\|_{C_{0}^{2, \alpha}(\Gamma)}<\delta$, where $\delta>0$ is a sufficiently small number. Then we have

$$
b_{1, l}(x, y)=b_{2, l}(x, y), \quad \text { for } x \in \Omega, y \in \mathbb{R} \text { and } l \in \mathbb{N} .
$$

In particular, if $b_{j, l}$ is of the form $b_{j, l}(x, y)=q_{j, l}|y|^{r_{l}-1} y$, where $q_{j, l}(x) \in C^{\alpha}(\bar{\Omega})$, then

$$
q_{1, l}(x)=q_{2, l}(x) \text { in } \Omega, \quad \text { for } l \in \mathbb{N} .
$$

Theorem 1.3 corresponds to the recovery of the coefficients of the asymptotic series expansion of $a(x, y)$ in the $y$-variable. Note that numbers $r_{1}, r_{2}, \ldots$ could also be integers $\geq 2$. Therefore, we can regard Theorem 1.3 as a generalization of the corresponding Euclidean results in [LLLS20a, LLLS20b].

Inspired by the partial data results of inverse problems for semilinear elliptic equations [LLLS20b, KU20b], one can also consider the inverse boundary value problem of recovering an obstacle and coefficients simultaneously. Let $\Omega \subset \mathbb{R}^{n}$ be a bounded domain with a connected $C^{\infty}$-smooth boundary $\partial \Omega$. Let $D \Subset \Omega$ be an open set with $C^{\infty}$-smooth boundary $\partial D$ such that $\Omega \backslash \bar{D}$ is connected. Consider the boundary value problem

$$
\begin{cases}\Delta u+a(x, u)=0 & \text { in } \Omega \backslash \bar{D}  \tag{1.7}\\ u=0 & \text { on } \partial D \\ u=f & \text { on } \partial \Omega\end{cases}
$$

where $a=a(x, y)$ is a polyhomogeneous function defined via Definition 1.1, for $x \in \Omega \backslash \bar{D}$.

As shown in Proposition 2.1, given any Dirichlet data $f \in C^{2, \alpha}(\partial \Omega)$ with $\|f\|_{C^{2, \alpha}(\partial \Omega)}<\delta$, for some sufficiently small number $\delta>0$, the equation (1.7) is well-posed and admits a unique (small) solution $u \in C^{2, \alpha}(\bar{\Omega} \backslash D)$. Let $\Gamma \subset \partial \Omega$ be
an arbitrarily nonempty relatively open subset, then we can define the corresponding partial DN map $\Lambda_{a, D}^{\Gamma}$ by

$$
\Lambda_{a, D}^{\Gamma}: C^{2, \alpha}(\Gamma) \rightarrow C^{1, \alpha}(\Gamma),\left.\quad f \mapsto \partial_{\nu} u_{f}\right|_{\Gamma},
$$

for any $f \in C_{0}^{2, \alpha}(\Gamma)$ with sufficiently small $\|f\|_{C_{0}^{2, \alpha}(\Gamma)}$, where $u_{f} \in C^{2, \alpha}(\bar{\Omega} \backslash D)$ is the unique solution of (1.7). The following result is analogous to [LLLS20b, Theorem 1.2] and [KU20b, Theorem 1.6].

Theorem 1.4 (Simultaneous recovery: Unknown obstacle and coefficient). Let $\Omega \subset \mathbb{R}^{n}, n \geq 2$ be a bounded connected domain with connected $C^{\infty}$ boundary $\partial \Omega$. Let $D_{1}, D_{2} \Subset \Omega$ be nonempty open subsets with $C^{\infty}$ boundaries such that $\Omega \backslash \overline{D_{j}}$ are connected. For $j=1,2$, let $a_{j}=a_{j}(x, y)$ be polyhomogeneous functions in $y \in \mathbb{R}$, for $x \in \bar{\Omega} \backslash D_{j}$. Denote by $\Lambda_{a_{j}, D_{j}}^{\Gamma}$ the partial DN maps of the following Dirichlet problems

$$
\begin{cases}\Delta u_{j}+a_{j}\left(x, u_{j}\right)=0 & \text { in } \Omega \backslash \overline{D_{j}} \\ u_{j}=0 & \text { on } \partial D_{j} \\ u_{j}=f & \text { on } \partial \Omega\end{cases}
$$

defined for any $f \in C_{0}^{2, \alpha}(\Gamma)$ with $\|f\|_{C_{0}^{2, \alpha}(\Gamma)}<\delta$, where $\delta>0$ is a sufficiently small number. Assume that

$$
\Lambda_{a_{1}, D_{1}}^{\Gamma}(f)=\Lambda_{a_{2}, D_{2}}^{\Gamma}(f), \text { for any }\|f\|_{C_{0}^{2, \alpha}(\Gamma)}<\delta
$$

Then

$$
D:=D_{1}=D_{2},
$$

and

$$
b_{1, l}(x, y)=b_{2, l}(x, y), \quad \text { for } x \in \Omega \backslash \bar{D}, y \in \mathbb{R} \text { and } l \in \mathbb{N}
$$

Remark 1.2. It is worth emphasizing that the simultaneous recovery of an embedded obstacle and the surrounding potentials in the linear setting, for example, the linear Schrödinger equation (i.e., for the case $r=1$ in Theorem 1.4) is an open problem. We refer readers to [Isa90, LLLS20b] for further discussions and [CLL19] for arguments in a linear nonlocal setting.

The proof of Theorem 1.4 is similar to the proof of Theorem 1.3, and the only difference is that we need to recover the unknown obstacle first. The method to recover the unknown obstacle has been investigated in [LLLS20b, Theorem 1.2]. We will give the proof in Section 4.

We are also able to extend the geometric results in [LLLS20a] to fractional power type nonlinearities. We refer to [LLLS20a] for the introduction of these problems.

Theorem 1.5 (Simultaneous recovery of metric and potential in the plane). Let $\left(M_{1}, g_{1}\right)$ and $\left(M_{2}, g_{2}\right)$ be two compact connected $C^{\infty}$ Riemannian manifolds with mutual $C^{\infty}$ boundary $\partial M$ and $\operatorname{dim}\left(M_{1}\right)=\operatorname{dim}\left(M_{2}\right)=2$. For $j=1,2$, let $\Lambda_{M_{j}, g_{j}, q_{j}}$ be the DN maps of

$$
\begin{equation*}
\Delta_{g_{j}} u+q_{j}|u|^{r-1} u=0 \text { in } M_{j} \tag{1.8}
\end{equation*}
$$

where $r>1$ is a fractional number. Let $0<\alpha<1$ and assume that

$$
\Lambda_{M_{1}, g_{1}, q_{1}}(f)=\Lambda_{M_{2}, g_{2}, q_{2}}(f) \text { on } \partial M,
$$

for any $f \in C^{2, \alpha}(\partial M)$ with $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta$, where $\delta>0$ is a sufficiently small number. Then:
(1) There exists a conformal diffeomorphism $J: M_{1} \rightarrow M_{2}$ and a positive smooth function $\sigma \in C^{\infty}\left(M_{1}\right)$ such that

$$
\sigma J^{*} g_{2}=g_{1} \text { in } M_{1}
$$

with $\left.J\right|_{\partial M}=\mathrm{Id}$ and $\left.\sigma\right|_{\partial M}=1$.
(2) Moreover, one can also recover the potential up to a natural gauge invariance in the sense that

$$
\sigma q_{1}=q_{2} \circ J \text { in } M_{1} .
$$

Furthermore, as shown in [LLLS20a] for integer power type nonlinearities, one can also consider the corresponding Calderón type inverse problem on a transversally anisotropic manifold. Let us consider inverse problems for the semilinear Schrödinger equation on transversally anisotropic manifold with fractional power type nonlinearities. The definition of a transversally anisotropic manifold is given as follows.

Definition 1.3. Let $(M, g)$ be a compact oriented manifold with a $C^{\infty}$ boundary and with $\operatorname{dim} M \geq 3$. $(M, g)$ is called transversally anisotropic if $(M, g) \Subset(T, g)$, where $T=\mathbb{R} \times M_{0}$ and $g(x)=g\left(x_{1}, x^{\prime}\right)=e\left(x_{1}\right) \oplus g_{0}\left(x^{\prime}\right)$ for $x_{1} \in \mathbb{R}$ and $x^{\prime} \in$ $M_{0}$. Here $(\mathbb{R}, e)$ denotes the Euclidean line and $\left(M_{0}, g_{0}\right)$ stands for an $(n-1)$ dimensional compact manifold with a smooth boundary.

Theorem 1.6. Let $(M, g)$ be a transversally anisotropic manifold, let $q_{j} \in C^{\infty}(M)$, and let $\Lambda_{q_{j}}$ be the DN maps for the equations

$$
\Delta_{g} u+q_{j}|u|^{r-1} u=0 \text { in } M
$$

for $j=1,2$, where we further assume the fractional number satisfies

$$
r>3
$$

Suppose that the DN maps satisfy

$$
\Lambda_{q_{1}}(f)=\Lambda_{q_{2}}(f) \text { on } \partial M,
$$

for all $f$ with $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta$, for a sufficiently small number $\delta>0$ and for some $0<\alpha<1$. Then $q_{1}=q_{2}$ in $M$.

Theorems 1.5 and 1.6 follow from the corresponding arguments in [LLLS20a] if we use the integral identity (2.11) with the choice $v_{0}=1$ in $M$ (by taking $f_{0}=1$ on $\partial M$ ).

The structure of this article is given as follows. In Section 2, we give wellposedness results for the relevant semilinear elliptic equations and derive the integral identity which plays a crucial role in the study of our inverse problems. In Section 3, we prove global uniqueness and simultaneous recovery in the Euclidean case, i.e., Theorems 1.1-1.4. Finally, we prove Theorems 1.5-1.6 in Section 4.

## 2. Preliminaries

First, let us recall the definition of Hölder spaces. Let $U \subset \mathbb{R}^{n}$ be an open set, let $k \in \mathbb{N} \cup\{0\}$, and let $0<\alpha<1$. The function space $C^{k, \alpha}(\bar{U})$ consists of those real valued functions $u \in C^{k}(\bar{U})$ for which the norm

$$
\|f\|_{C^{k, \alpha}(\bar{U})}:=\sum_{|\gamma| \leq k}\left\|\partial^{\gamma} f\right\|_{L^{\infty}(U)}+\sup _{x \neq y, x, y \in \bar{U}} \sum_{|\gamma|=k} \frac{\left|\partial^{\gamma} f(x)-\partial^{\gamma} f(y)\right|}{|x-y|^{\alpha}},
$$

is finite. Here $\gamma=\left(\gamma_{1}, \cdots, \gamma_{n}\right)$ is a multi-index with $\gamma_{i} \in \mathbb{N} \cup\{0\}$ and $|\gamma|=$ $\gamma_{1}+\cdots+\gamma_{n}$. Furthermore, we also denote the space

$$
C_{0}^{k, \alpha}(\bar{U}):=\text { closure of } C_{c}^{\infty}(U) \text { in } C^{k, \alpha}(\bar{U}) .
$$

In short, we only use $C^{\alpha}(\bar{U})$ to denote $C^{0, \alpha}(\bar{U})$ when $k=0$. In addition, one can define Hölder spaces on any Riemannian manifold ( $M, g$ ) using the Riemannian distance or via local coordinates, see e.g. [Tay11, Section 13.8 in vol. III].
2.1. Well-posedness. Let $(M, g)$ be a $C^{\infty}$ compact Riemannian manifold with $C^{\infty}$-smooth boundary $\partial M$. We study the well-posedness of the following boundary value problem

$$
\begin{cases}\Delta_{g} u+a(x, u)=0 & \text { in } M  \tag{2.1}\\ u=f & \text { on } \partial M\end{cases}
$$

for any sufficiently small Dirichlet data $f \in C^{2, \alpha}(\partial M)$, for some $0<\alpha<1$. Let us assume that the nonlinear coefficient $a=a(x, y) \in C_{\mathrm{loc}}^{k, \alpha}\left(\mathbb{R}, C^{\alpha}(M)\right)$ for some $k \geq 1$, meaning that $y \mapsto \partial_{y}^{j} a(\cdot, y)$ is a continuous map $\mathbb{R} \rightarrow C^{\alpha}(M)$ for $0 \leq j \leq k$ and for any $R>0,\left\|\partial_{y}^{k} a(\cdot, y)-\partial_{y}^{k} a(\cdot, z)\right\|_{C^{\alpha}} \leq C_{R}|y-z|^{\alpha}$ whenever $|y|,|z| \leq R$. Also assume that the following two conditions hold:

$$
\begin{equation*}
a(x, 0)=0, \quad \text { for } x \in M \tag{2.2}
\end{equation*}
$$

The map $v \mapsto \Delta_{g} v+\partial_{y} a(\cdot, 0) v$ is injective on $H_{0}^{1}(M)$.
We prove the well-posedness of (2.1) for small Dirichlet data $f \in C^{2, \alpha}(\partial M)$.
Proposition 2.1 (Well-posedness). Let $(M, g)$ be a compact Riemannian manifold with $C^{\infty}$ boundary $\partial M$ and let $Q$ be the semilinear elliptic operator

$$
Q(u):=\Delta_{g} u+a(x, u),
$$

where $a \in C_{\mathrm{loc}}^{k, \alpha}\left(\mathbb{R}, C^{\alpha}(M)\right)$ for some $k \geq 1$, $\alpha \in(0,1)$, and (2.2) and (2.3) are satisfied. There exist $\delta, C>0$ such that for any $f$ in the set

$$
U_{\delta}:=\left\{h \in C^{2, \alpha}(\partial M) ;\|h\|_{C^{2, \alpha}(\partial M)} \leq \delta\right\}
$$

there is a solution $u=u_{f}$ of

$$
\begin{cases}\Delta_{g} u+a(x, u)=0 & \text { in } M  \tag{2.4}\\ u=f & \text { on } \partial M\end{cases}
$$

which satisfies

$$
\begin{equation*}
\|u\|_{C^{2, \alpha}(M)} \leq C\|f\|_{C^{2, \alpha}(\partial M)} \tag{2.5}
\end{equation*}
$$

The solution $u_{f}$ is unique within the class $\left\{w \in C^{2, \alpha}(M) ;\|w\|_{C^{2, \alpha}(M)} \leq C \delta\right\}$. In addition, there are $C^{k}$ Frechét differentiable maps

$$
\begin{aligned}
& S: U_{\delta} \rightarrow C^{2, \alpha}(M), \quad f \mapsto u_{f} \\
& \Lambda: U_{\delta} \rightarrow C^{1, \alpha}(\partial M),\left.\quad f \mapsto \partial_{\nu} u_{f}\right|_{\partial M}
\end{aligned}
$$

In particular, if $a(x, u)=q(x)|u|^{r-1} u$ for a fractional number $r>1$ and $q \in$ $C^{\alpha}(M)$, then the function $q(x)|u|^{r-1} u$ satisfies the condition $a(x, 0)=\partial_{y} a(x, 0)=$ 0 , which implies that the conditions (2.2) and (2.3) hold automatically (due to the well-posedness of the Laplace equation). Hence, Proposition 2.1 implies the well-posedness of the Dirichlet problem (1.1) immediately.

For the proof of Proposition 2.1, we will need a lemma that will also be useful later.

Lemma 2.2. Let $(M, g)$ be a compact Riemannian manifold with $C^{\infty}$ boundary $\partial M$, let $0<\alpha<1$, and let $b(x, y) \in C_{\mathrm{loc}}^{\alpha}\left(\mathbb{R}, C^{\alpha}(M)\right)$. For any $u \in C^{1}(M)$ one has $b(x, u(x)) \in C^{\alpha}(M)$, and

$$
\begin{equation*}
\|b(x, u+v)-b(x, u)\|_{C^{\alpha}(M)}=o(1), \text { as }\|v\|_{C^{1}(M)} \rightarrow 0 \tag{2.6}
\end{equation*}
$$

Proof. The assumption that $t \mapsto b(\cdot, t)$ is a $C_{\mathrm{loc}}^{\alpha}$ function $\mathbb{R} \rightarrow C^{\alpha}(M)$ means that for any $R>0$ there is $C_{R}>0$ such that

$$
\begin{aligned}
|b(x, t)| & \leq C_{R}, \\
|b(x, t)-b(y, t)| & \leq C_{R} d_{g}(x, y)^{\alpha}, \\
|b(x, t)-b(x, s)| & \leq C_{R}|t-s|^{\alpha}, \\
|b(x, t)-b(x, s)-(b(y, t)-b(y, s))| & \leq\|b(\cdot, t)-b(\cdot, s)\|_{C^{\alpha}(M)} d_{g}(x, y)^{\alpha} \\
& \leq C_{R} d_{g}(x, y)^{\alpha}|t-s|^{\alpha},
\end{aligned}
$$

whenever $x, y \in M$ and $|t|,|s| \leq R$.
Now if $u \in C^{1}(M)$ with $\|u\|_{L^{\infty}(M)} \leq R$, one has $|b(x, u(x))| \leq C_{R}$ and

$$
\begin{aligned}
|b(x, u(x))-b(y, u(y))| & \leq|b(x, u(x))-b(y, u(x))|+|b(y, u(x))-b(y, u(y))| \\
& \leq C_{R}\left[1+\|u\|_{C^{1}(M)}^{\alpha}\right] d_{g}(x, y)^{\alpha} .
\end{aligned}
$$

This shows that $b(x, u(x)) \in C^{\alpha}(M)$.
Let now $u, v \in C^{1}(M)$ with $\|u\|_{L^{\infty}} \leq R$ and $\|u+v\|_{L^{\infty}} \leq R$. Then

$$
\|b(x, u+v)-b(x, u)\|_{L^{\infty}(M)} \leq C_{R}\|v\|_{L^{\infty}(M)}^{\alpha}
$$

Let us next estimate the $C^{\alpha}$ norm of $b(x, u+v)-b(x, u)$. Writing $h(x, u):=b(x, u)$ and $w_{t}(x):=u(x)+t v(x)$, we have

$$
\begin{align*}
& \left|h\left(x, w_{1}(x)\right)-h\left(x, w_{0}(x)\right)-\left[h\left(y, w_{1}(y)\right)-h\left(y, w_{0}(y)\right)\right]\right| \\
& \leq\left|h\left(x, w_{1}(x)\right)-h\left(x, w_{0}(x)\right)-\left[h\left(y, w_{1}(x)\right)-h\left(y, w_{0}(x)\right)\right]\right|  \tag{2.7}\\
& \quad+\left|h\left(y, w_{1}(x)\right)-h\left(y, w_{0}(x)\right)-\left[h\left(y, w_{1}(y)\right)-h\left(y, w_{0}(y)\right)\right]\right| .
\end{align*}
$$

The first absolute value on the right of $(2.7)$ is $\leq C_{R} d_{g}(x, y)^{\alpha}|v(x)|^{\alpha}$. The second absolute value on the right of (2.7) can be estimated by grouping the terms in two different ways and using the triangle inequality: it is either $\leq C_{R}\|v\|_{L^{\infty}(M)}^{\alpha}$ or $\leq C_{R}\left(\|u\|_{C^{1}(M)}+\|v\|_{C^{1}(M)}\right)^{\alpha} d_{g}(x, y)^{\alpha}$.

By interpolation, this shows that for any $\beta<\alpha$ one has

$$
\|b(x, u+v)-b(x, u)\|_{C^{\beta}(M)}=o(1), \text { as }\|v\|_{C^{1}(M)} \rightarrow 0
$$

This estimate is also true for $\beta=\alpha$. This can be seen by writing

$$
b=b_{\epsilon}+r_{\epsilon},
$$

where

$$
b_{\epsilon}(x, t)=\int_{\mathbb{R}} \varphi_{\epsilon}(t-s) b(x, s) d s
$$

Here $\varphi_{\epsilon}(t)=\epsilon^{-n} \varphi(t / \epsilon)$ is a standard mollifier with $\varphi \in C_{c}^{\infty}((-1,1)), 0 \leq \varphi \leq 1$, and $\int_{\mathbb{R}} \varphi(t) d t=1$. Repeating the argument above for $b_{\epsilon}$ using a higher Hölder exponent in $t$, and using the estimate $\left\|r_{\epsilon}(\cdot, t)\right\|_{C^{\alpha}(M)} \leq C_{R} \epsilon^{\alpha}$ for $|t| \leq R$ which follows from the regularity of $b$, finally yields the estimate

$$
\|b(x, u+v)-b(x, u)\|_{C^{\alpha}(M)}=o(1), \text { as }\|v\|_{C^{1}(M)} \rightarrow 0
$$

Proof of Proposition 2.1. We prove the existence of solutions by using the implicit function theorem in Banach spaces [Zei86, Theorem 4.B]. Let

$$
X=C^{2, \alpha}(\partial M), \quad Y=C^{2, \alpha}(M), \quad Z=C^{\alpha}(M) \times C^{2, \alpha}(\partial M)
$$

Consider the map

$$
F: X \times Y \rightarrow Z, \quad F(f, u)=\left(Q(u),\left.u\right|_{\partial M}-f\right)
$$

Now $F$ indeed maps to $Z$, since by Lemma 2.2 the map $u \mapsto a(x, u)$ takes $C^{2, \alpha}(M)$ to $C^{\alpha}(M)$. Thus $F$ is well defined.

We next show that $F$ is a $C^{k}$ map. Let $0<m \leq k$ be an integer. If $u, v \in$ $C^{2, \alpha}(M)$ we use the Taylor formula

$$
\begin{align*}
& a(x, u+v)  \tag{2.8}\\
= & \sum_{j=0}^{m-1} \frac{\partial_{u}^{j} a(x, u)}{j!} v^{j}+\int_{0}^{1} \frac{\partial_{u}^{m} a(x, u+t v)}{(m-1)!} v^{m}(1-t)^{m-1} d t \\
= & \sum_{j=0}^{m} \frac{\partial_{u}^{j} a(x, u)}{j!} v^{j}-\frac{v^{m}}{m!} \partial_{u}^{m} a(x, u)+\int_{0}^{1} \frac{\partial_{u}^{m} a(x, u+t v)}{(m-1)!} v^{m}(1-t)^{m-1} d t \\
= & \sum_{j=0}^{m} \frac{\partial_{u}^{j} a(x, u)}{j!} v^{j}+\frac{v^{m}}{(m-1)!} \int_{0}^{1}\left[\partial_{u}^{m} a(x, u+t v)-\partial_{u}^{m} a(x, u)\right](1-t)^{m-1} d t .
\end{align*}
$$

We study the remainder term. From (2.6) with $b=\partial_{u}^{m} a$ we obtain the estimate

$$
\left\|\partial_{u}^{m} a(x, u+t v)-\partial_{u}^{m} a(x, u)\right\|_{C^{\alpha}(M)}=o(1), \text { if } t \in[0,1] \text { and }\|v\|_{C^{2, \alpha}(M)} \rightarrow 0
$$

Inserting this in the Taylor formula computation (2.8) yields

$$
\left\|a(x, u+v)-\sum_{j=0}^{m} \frac{\partial_{u}^{j} a(x, u)}{j!} v^{j}\right\|_{C^{\alpha}(M)}=o\left(\|v\|_{C^{2, \alpha}(M)}^{m}\right), \text { as }\|v\|_{C^{2, \alpha}(M)} \rightarrow 0 .
$$

This shows that $u \mapsto a(x, u)$ is a $C^{k} \operatorname{map} C^{2, \alpha}(M) \rightarrow C^{\alpha}(M)$. Since the other parts of $F$ are linear, $F$ is a $C^{k}$ map.

Note that $F(0,0)=0$ by (2.2). The linearization of $F$ at $(0,0)$ in the $u$-variable is

$$
\left.D_{u} F\right|_{(0,0)}(v)=\left(\Delta_{g} v+\partial_{u} a(x, 0) v,\left.v\right|_{\partial M}\right)
$$

This is a homeomorphism $Y \rightarrow Z$ by (2.3). To see this, let $(w, \phi) \in Z=C^{\alpha}(M) \times$ $C^{2, \alpha}(\partial M)$, and consider the Dirichlet problem

$$
\begin{cases}\left(\Delta_{g}+\partial_{u} a(x, 0)\right) v=w & \text { in } M  \tag{2.9}\\ v=\phi & \text { on } \partial M\end{cases}
$$

The solution of (2.9), if it exists, is unique by (2.3), and by using the Fredholm alternative and Schauder estimates the solution $v \in Y=C^{2, \alpha}(M)$ exists (see e.g. [Tay11, Exercise 1 in Section 13.8]) and depends continuously on the data ( $w, \phi$ ). Thus the implicit function theorem in Banach spaces [Zei86, Theorem 4.B] yields that there is $\delta>0$, a closed ball $U_{\delta}=\overline{B_{X}(0, \delta)} \subset X$, and a $C^{k} \operatorname{map} S: U \rightarrow Y$ such that whenever $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta$ we have

$$
F(f, S(f))=(0,0)
$$

Since $S$ is Lipschitz continuous and $S(0)=0, u=S(f)$ satisfies

$$
\|u\|_{C^{2, \alpha}(M)} \leq C\|f\|_{C^{2, \alpha}(\partial M)}
$$

Moreover, by redefining $\delta$ if necessary $u=S(f)$ is the only solution to $F(f, u)=$ $(0,0)$ whenever $\|u\|_{C^{2, \alpha}(M)} \leq C \delta$. We have proven the existence of unique small solutions of the Dirichlet problem (2.4) and the fact that the solution operator $S: U_{\delta} \rightarrow C^{2, \alpha}(M)$ is a $C^{k}$ map. Since the normal derivative is a linear map $C^{2, \alpha}(M) \rightarrow C^{1, \alpha}(\partial M)$, it follows that also $\Lambda$ is a well defined $C^{k}$ map $U_{\delta} \rightarrow$ $C^{1, \alpha}(\partial M)$.

In the next proposition we present an integral identity involving the $k$ th linearization the DN map $\Lambda_{q}$. Below, we write

$$
\left(D^{k} f\right)_{x}\left(y_{1}, \ldots, y_{k}\right)
$$

to denote the $k$ th derivative at $x$ of a $C^{k}$ map $f$ between Banach spaces, considered as a symmetric $k$-linear form acting on $\left(y_{1}, \ldots, y_{k}\right)$. We refer to [Hor85, Section 1.1], where the notation $f^{(k)}\left(x ; y_{1}, \ldots, y_{k}\right)$ is used instead of $\left(D^{k} f\right)_{x}\left(y_{1}, \ldots, y_{k}\right)$.

Proposition 2.3 (Integral identity). Let $(M, g)$ be a compact $C^{\infty}$ Riemannian manifold with a $C^{\infty}$ smooth boundary $\partial M$. Let $q \in C^{\alpha}(M)$, and let $\Lambda_{q}$ be the $D N$ map for the semilinear elliptic equation

$$
\begin{equation*}
\Delta_{g} u+q|u|^{r-1} u=0 \text { in } M, \tag{2.10}
\end{equation*}
$$

where

$$
r=k+\alpha, \quad k \geq 1 \text { and } \alpha \in(0,1)
$$

Let $f_{0} \in C^{2, \alpha}(\partial M)$. Then the $k$ th linearization $\left(D^{k} \Lambda_{q}\right)_{\epsilon_{0} f_{0}}$ of $\Lambda_{q}$ at $\epsilon_{0} f_{0}$ satisfies the following identity: For any $f_{1}, \ldots, f_{k+1} \in C^{2, \alpha}(\partial M)$ one has

$$
\begin{align*}
& \lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha} \int_{\partial M}\left(D^{k} \Lambda_{q}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right) f_{k+1} d S \\
& \quad=c_{r} \int_{M} q\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k+1} d V \tag{2.11}
\end{align*}
$$

where $c_{r}$ is the constant given by

$$
c_{r}=-r(r-1) \cdots(r-(k-1)) .
$$

Here each $v_{\ell}, \ell=0, \ldots, k+1$, is a harmonic function satisfying

$$
\begin{cases}\Delta_{g} v_{\ell}=0 & \text { in } M  \tag{2.12}\\ v_{\ell}=f_{\ell} & \text { on } \partial M\end{cases}
$$

Proof. Let $f_{0} \in C^{2, \alpha}(\partial M)$ and denote $h_{0}=\epsilon_{0} f_{0}$, where $\epsilon_{0}$ is small. The nonlinearity $a(x, u)=q(x)|u|^{r-1} u$ satisfies the conditions in Proposition 2.1, and thus the DN $\operatorname{map} \Lambda_{q}=\left.\partial_{\nu} S\right|_{\partial M}$ is well defined for boundary data $f$ with $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta$. Here $S: f \mapsto u_{f}$ is the solution operator for the Dirichlet problem of the equation (2.10).

We first compute the derivatives of $\Lambda_{q}$ at $h_{0}$. For this it is enough to consider the derivatives of $S$. Let us write

$$
\widetilde{f}=\widetilde{f}\left(x ; \epsilon_{1}, \ldots, \epsilon_{k}\right):=\epsilon_{1} f_{1}(x)+\ldots+\epsilon_{k} f_{k}(x)
$$

Let $f=h_{0}+\widetilde{f}$, then the solution

$$
u_{f}:=S(f)=S\left(h_{0}+\epsilon_{1} f_{1}+\cdots+\epsilon_{k} f_{k}\right) \in C^{2, \alpha}(M)
$$

is $k$ times continuously differentiable with respect to the parameters $\epsilon_{1}, \ldots, \epsilon_{k}$ by Proposition 2.1. Let us denote

$$
\epsilon:=\left(\epsilon_{0}, \epsilon^{\prime}\right), \quad \epsilon^{\prime}:=\left(\epsilon_{1}, \ldots, \epsilon_{k}\right)
$$

Applying $\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}}\right|_{\epsilon^{\prime}=0}$ to the Taylor formula for $C^{k}$ maps (see e.g. [Hor85, equation (1.1.8)])

$$
u_{f}=S\left(h_{0}+\widetilde{f}\right)=\sum_{m=0}^{k} \frac{\left(D^{m} S\right)_{h_{0}}(\tilde{f}, \ldots, \tilde{f})}{m!}+o\left(\|\widetilde{f}\|_{C^{2, \alpha}(\partial M)}^{k}\right)
$$

implies that $\left(D^{m} S\right)_{h_{0}}$ for $0 \leq m \leq k$ may be computed using the formula

$$
\begin{equation*}
\left(D^{m} S\right)_{h_{0}}\left(f_{1}, \ldots, f_{m}\right)=\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{m}} u_{f}\right|_{\epsilon^{\prime}=0} \tag{2.13}
\end{equation*}
$$

Moreover, since $S$ is $C^{k} \operatorname{map} C^{2, \alpha}(\partial M) \rightarrow C^{2, \alpha}(M)$, since $u \mapsto q(x)|u|^{r-1} u$ is a $C^{k} \operatorname{map} C^{2, \alpha}(M) \rightarrow C^{\alpha}(M)$ by the argument in Proposition 2.1, and since $\Delta_{g}$ is linear, we may differentiate the equation

$$
\begin{cases}\Delta_{g} u_{f}+q(x)\left|u_{f}\right|^{r-1} u_{f}=0 & \text { in } M  \tag{2.14}\\ u_{f}=f=h_{0}+\widetilde{f} & \text { on } \partial M\end{cases}
$$

up to $k$ times in the $\epsilon_{\ell}$ variables at $\epsilon^{\prime}=0$ (recalling that $\widetilde{f}=f\left(x ; \epsilon^{\prime}\right)=\widetilde{f}\left(x ; \epsilon_{1}, \ldots, \epsilon_{k}\right)$ ).
Let $\ell \in\{1, \ldots, k\}$. Then for any $\beta>0$ we have the identity

$$
\partial_{\epsilon_{\ell}}\left(\left|u_{f}\right|^{\beta} u_{f}\right)=\left(\beta\left|u_{f}\right|^{\beta-2} u_{f}^{2}+\left|u_{f}\right|^{\beta}\right) \partial_{\epsilon_{\ell}} u_{f}=(\beta+1)\left|u_{f}\right|^{\beta} \partial_{\epsilon_{\ell}} u_{f}
$$

so that

$$
\begin{cases}\Delta_{g}\left(\left.\partial_{\epsilon_{\ell}} u_{f}\right|_{\epsilon^{\prime}=0}\right)+\left.q(x) r\left|u_{f}\right|^{r-1} \partial_{\epsilon_{\ell}} u_{f}\right|_{\epsilon^{\prime}=0}=0 & \text { in } M,  \tag{2.15}\\ \left.\partial_{\epsilon_{\ell}} u_{f}\right|_{\epsilon^{\prime}=0}=f_{\ell} & \text { on } \partial M .\end{cases}
$$

Thus the first linearization of the map $S$ at $h_{0}$ is

$$
\begin{equation*}
v_{\ell}^{\epsilon_{0}}:=(D S)_{h_{0}}\left(f_{\ell}\right)=\left.\partial_{\epsilon_{\ell}} u_{f}\right|_{\epsilon^{\prime}=0} \tag{2.16}
\end{equation*}
$$

where $v_{\ell}^{\epsilon_{0}}$ satisfies (2.15). For $\ell=1,2, \ldots, k$, we also claim that

$$
\begin{equation*}
\lim _{\epsilon_{0} \rightarrow 0} v_{\ell}^{\epsilon_{0}}=v_{\ell} \text { in } C^{2, \alpha}(M) \tag{2.17}
\end{equation*}
$$

where $v_{\ell}$ is the harmonic function satisfying (2.12) with Dirichlet data $f_{\ell}$. To prove (2.17), note by the Schauder estimates we have

$$
\begin{aligned}
\left\|v_{\ell}^{\epsilon_{0}}-v_{\ell}\right\|_{C^{2, \alpha}(M)} & \leq C\left(\left\|\Delta_{g}\left(v_{\ell}^{\epsilon_{0}}-v_{\ell}\right)\right\|_{C^{\alpha}(M)}+\left\|\epsilon_{0} f_{0}+f_{\ell}-f_{\ell}\right\|_{C^{2, \alpha}(\partial M)}\right) \\
& =C\left(\left\|\left.q\left[r\left|u_{f}\right|^{r-1} \partial_{\epsilon_{\ell}} u_{f}\right]\right|_{\epsilon^{\prime}=0}\right\|_{C^{\alpha}(M)}+\left\|\epsilon_{0} f_{0}\right\|_{C^{2, \alpha}(\partial M)}\right) \\
& \leq C\left(\left\|\left|u_{\epsilon_{0} f_{0}}\right|^{r-1}\right\|_{C^{\alpha}(M)}+\epsilon_{0}\right)
\end{aligned}
$$

Now $\left\|u_{\epsilon_{0} f_{0}}\right\|_{C^{2, \alpha}(M)} \leq C \epsilon_{0}\left\|f_{0}\right\|_{C^{2, \alpha}(\partial M)}$ by (2.5). Then (2.6) with $b(x, t)$ replaced by $|t|^{r-1}$ implies that $\left\|\left|u_{\epsilon_{0} f_{0}}\right|^{r-1}\right\|_{C^{\alpha}(M)} \rightarrow 0$ as $\epsilon_{0} \rightarrow 0$, proving (2.17).

Let now $2 \leq j \leq k$. Applying $\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}}\right|_{\epsilon^{\prime}=0}$ to (2.14) gives that

$$
\begin{cases}\Delta_{g}\left(\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}} u_{f}\right|_{\epsilon^{\prime}=0}\right)=-\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}}\left(q(x)|u|^{r-1} u\right)\right|_{\epsilon^{\prime}=0} & \text { in } M \\ \left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}} u_{f}\right|_{\epsilon^{\prime}=0}=0 & \text { on } \partial M\end{cases}
$$

Since $r>k$, the fact that $u_{f}$ is $k$ times continuously Frechét differentiable in $\epsilon^{\prime}$ gives that

$$
\left.\lim _{\epsilon_{0} \rightarrow 0} \partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}}\left(q(x)|u|^{r-1} u\right)\right|_{\epsilon^{\prime}=0}=0
$$

By an argument similar to the one above using Schauder estimates we obtain

$$
\left.\lim _{\epsilon_{0} \rightarrow 0} \partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{j}} u_{f}\right|_{\epsilon^{\prime}=0}=0
$$

Let us consider the $k$ th mixed derivative $w^{\epsilon_{0}}:=\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}} u_{f}\right|_{\epsilon^{\prime}=0}$ further. It satisfies the equation

$$
\begin{cases}\Delta_{g} w^{\epsilon_{0}}=-\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}}\left(q(x)|u|^{r-1} u\right)\right|_{\epsilon^{\prime}=0} & \text { in } M  \tag{2.18}\\ w^{\epsilon_{0}}=0 & \text { on } \partial M\end{cases}
$$

We wish to multiply (2.18) by $\epsilon_{0}^{-\alpha}$ and take the limit as $\epsilon_{0} \rightarrow 0$. Since $f(t)=|t|^{r-1} t$ for $r=k+\alpha$ satisfies the homogeneity relation $f(\lambda t)=\lambda^{r} f(t)$ for $\lambda>0$, we have that

$$
\frac{d^{k}}{d y^{k}}\left(|y|^{r-1} y\right)=r(r-1) \cdots(r-(k-1))|y|^{r-1} y^{1-k}=-c_{r}|y|^{r-1} y^{1-k}
$$

Using Faà di Bruno's formula, see [Har06], we find that

$$
\begin{align*}
\left.\partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}}\left(\left|u_{f}\right|^{r-1} u_{f}\right)\right|_{\epsilon^{\prime}=0}= & \left.\sum_{\sigma \in P} c_{\sigma}\left|u_{f}\right|^{r-1} u_{f}^{1-|\sigma|} \prod_{\delta \in \sigma} \partial_{\epsilon^{\prime}}^{\delta} u_{f}\right|_{\epsilon^{\prime}=0} \\
= & \left.c_{r}\left|u_{f}\right|^{r-1} u_{f}^{1-k}\left(\partial_{\epsilon_{1}} u_{f}\right) \cdots\left(\partial_{\epsilon_{k}} u_{f}\right)\right|_{\epsilon^{\prime}=0}  \tag{2.19}\\
& +\left.\sum_{\substack{\sigma \in P,|\sigma|<k}} c_{\sigma}\left|u_{f}\right|^{r-1} u_{f}^{1-|\sigma|} \prod_{\delta \in \sigma} \partial_{\epsilon^{\prime}}^{\delta} u_{f}\right|_{\epsilon^{\prime}=0},
\end{align*}
$$

where $P$ contains all partitions of $\{1, \ldots, k\}$ and the product over $\delta \in \sigma$ runs over all sets in the partition $\sigma$. The number $|\sigma|$ denotes the cardinality of the set $\sigma$ and $\partial_{\epsilon^{\prime}}^{\delta}$ is the usual multi-index notation for partial derivatives in $\epsilon^{\prime}$.

Observe that $\left.u_{f}\right|_{\epsilon^{\prime}=0}$ solves the nonlinear equation (2.10) with boundary value $h_{0}=\epsilon_{0} f_{0}$. By continuity and uniqueness of solutions, we have that

$$
\begin{equation*}
\left.\epsilon_{0}^{-1} u_{f}\right|_{\epsilon^{\prime}=0} \rightarrow v_{0} \text { in } C^{2, \alpha}(M), \quad \text { as } \epsilon_{0} \rightarrow 0 . \tag{2.20}
\end{equation*}
$$

Then note that $|\sigma|<k$ implies that the products

$$
\left.\prod_{\delta \in \sigma} \partial_{\epsilon^{\prime}}^{\delta} u_{f}\right|_{\epsilon^{\prime}=0}
$$

are bounded in $C^{\alpha}(M)$ as $\epsilon_{0} \rightarrow 0$, because the solution operator $S$ is continuously $k$-Fréchet differentiable and the Hölder space $C^{\alpha}(M)$ is an algebra. Next, since the function $g(y)=|y|^{r-1} y^{1-|\sigma|}$ is homogeneous of degree $k-|\sigma|+\alpha \geq 1+\alpha$, Euler's homogeneous function theorem shows that it belongs to $C^{1}(\mathbb{R})$. Since the composition of $C^{1}(\mathbb{R})$ function with a $C^{2, \alpha}(M)$ function is at least $C^{\alpha}(M)$, we have that

$$
\begin{equation*}
\left.\epsilon^{-\alpha}\left|u_{f}\right|^{r-1} u_{f}^{1-|\sigma|}\right|_{\epsilon^{\prime}=0}=\left.\epsilon_{0}^{k-|\sigma|}\left|\frac{u_{f}}{\epsilon_{0}}\right|^{r-1}\left(\frac{u_{f}}{\epsilon_{0}}\right)^{1-|\sigma|}\right|_{\epsilon^{\prime}=0} \rightarrow 0 \quad \text { in } C^{\alpha}(M) \tag{2.21}
\end{equation*}
$$

as $\epsilon_{0} \rightarrow 0$. By using (2.17), (2.20) and (2.21), we see that after multiplying (2.19) by $\epsilon_{0}^{-\alpha}$ and taking the limit $\epsilon_{0} \rightarrow 0$, only the first term on the right hand side of (2.19) survives. To analyze this first term in the right-hand side of (2.19), observe that $g(y)=|y|^{r-1} y^{1-k}$ belongs to $C^{\alpha}(\mathbb{R})$ and $u_{f}$ is in $C^{2, \alpha}(M)$, so the composition $\left|u_{f}\right|^{r-1} u_{f}^{1-k}$ is in $C^{\alpha}(M)$. Recall again from (2.16) that $\left.\partial_{\epsilon_{\ell}} u_{f}\right|_{\epsilon^{\prime}=0} \rightarrow v_{\ell}$ in $C^{2, \alpha}(M)$ as $\epsilon_{0} \rightarrow 0$ for all $\ell=1,2, \ldots, k$. Due to the continuity of the solution map $S$, we finally have in $C^{\alpha}$ the limit

$$
\begin{equation*}
\left.\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha} \partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}}\left(q\left|u_{f}\right|^{r-1} u_{f}\right)\right|_{\epsilon^{\prime}=0}=-c_{r} q\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k} . \tag{2.22}
\end{equation*}
$$

Integrating the equation (2.18) against the harmonic function $v_{k+1}$, we have

$$
\int_{\partial M}\left(\partial_{\nu} w^{\epsilon_{0}}\right) f_{k+1} d S=-\left.\int_{M} \partial_{\epsilon_{1}} \cdots \partial_{\epsilon_{k}}\left(q(x)\left|u_{f}\right|^{r-1} u_{f}\right)\right|_{\epsilon^{\prime}=0} v_{k+1} d V
$$

Since $\Lambda_{q}=\partial_{\nu} S$ where $\partial_{\nu}$ is linear, the formula (2.13) gives that $\left.\partial_{\nu} w^{\epsilon_{0}}\right|_{\partial M}=$ $\left(D^{k} \Lambda_{q}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right)$. Now (2.22) yields

$$
\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha} \int_{\partial M}\left(D^{k} \Lambda_{q}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right) f_{k+1} d S=c_{r} \int_{M} q\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k} d V
$$

as required.
It is easy to see that the integral identity also holds for any $f \in C_{0}^{2, \alpha}(\Gamma)$, for any open subset $\Gamma \subset \partial M$. The following result is an easy consequence of the preceding proposition. For simplicity we only state the result in Euclidean domains.

Corollary 2.4 (Integral identity with partial data). Let $\Omega \subset \mathbb{R}^{n}$ be a bounded domain with $C^{\infty}$-smooth boundary $\partial \Omega$, for $n \geq 2$, and let $\Gamma \subset \partial \Omega$ be a nonempty relatively open subset. Let $q \in C^{\alpha}(\bar{\Omega})$ for some $0<\alpha<1$, and let $\Lambda_{q}^{\Gamma}$ be the partial data $D N$ map for the semilinear elliptic equation

$$
\begin{cases}\Delta u+q|u|^{r-1} u=0 & \text { in } \Omega \\ u=f & \text { on } \partial \Omega\end{cases}
$$

where $r=k+\alpha$ with $k \geq 1$ and $\alpha \in(0,1)$. The $k$ th linearization $D^{k} \Lambda_{q}^{\Gamma}$ of $\Lambda_{q}^{\Gamma}$ satisfies the following identity: For any $f_{0}, f_{1}, \ldots, f_{k+1} \in C_{0}^{2, \alpha}(\Gamma)$, one has

$$
\begin{align*}
& \lim _{\epsilon_{0} \rightarrow 0} \int_{\partial \Omega} \epsilon_{0}^{-\alpha}\left(D^{k} \Lambda_{q}^{\Gamma}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right) f_{k+1} d S  \tag{2.23}\\
& \quad=c_{r} \int_{\Omega} q\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k+1} d x,
\end{align*}
$$

where $c_{r}=-r(r-1) \cdots(r-(k-1))$. Here each $v_{\ell}, \ell=0, \ldots, k+1$, is a harmonic function satisfying

$$
\Delta v_{\ell}=0 \text { in } \Omega \quad \text { and } \quad v_{\ell}=f_{\ell} \text { on } \partial \Omega .
$$

The result follows immediately from Proposition 2.3, even if the Dirichlet data is supported in a relatively open subset $\Gamma \subset \partial \Omega$.

It is worth mentioning that even in the case $1<r<2$ we can use two boundary functions $f_{0}$ and $f_{1}$. A suitable choice of the Dirichlet data $f_{0}$ allows us to get rid of the nonlinear term $\left|v_{0}\right|^{\alpha}$, if necessary, while still retaining the ability to choose $f_{1}$ and the auxiliary function $f_{2}$ in an appropriate way.

Remark 2.5. We mention that for nonlinearities $a(x, u)=q(x)|u|^{\alpha} u$ where $q \in$ $C^{\alpha}(M)$ and $\alpha \in(0,1)$, one can prove that the solution of

$$
\begin{cases}\Delta u_{\epsilon}+q\left|u_{\epsilon}\right|^{\alpha} u_{\epsilon}=0 & \text { in } M \\ u_{\epsilon}=\epsilon f & \text { on } \partial M\end{cases}
$$

where $f \in C^{2, \alpha}(\partial M)$ and $\epsilon>0$ is small, has the asymptotic expansion

$$
u_{\epsilon}=\epsilon v+\epsilon^{1+\alpha} w+O\left(\epsilon^{1+2 \alpha}\right)
$$

where $v$ is the harmonic function satisfying

$$
\begin{cases}\Delta v=0 & \text { in } M \\ v=f & \text { on } \partial M\end{cases}
$$

and $w$ is the solution of

$$
\begin{cases}\Delta w=-q|v|^{\alpha} v & \text { in } M \\ w=0 & \text { on } \partial M\end{cases}
$$

One could use such one-parameter asymptotic expansions to give alternative proofs of some of our full data inverse problems. However, we will instead use Proposition 2.3 and Corollary 2.4, which are based on multiparameter expansions and will lead to more general results. For our proof of Theorem 1.6 it is crucial to use Proposition 2.3 with $k \geq 3$.

## 3. Global uniqueness in Euclidean space

In this section, let us prove our main Euclidean results. Recall that we are considering real-valued solutions. In order to apply the density results [FKSU09, LLLS20a] involving products of complex-valued harmonic functions, let us start with the following simple lemma also used in [LLLS20b]:

Lemma 3.1. Let $\Omega \subset \mathbb{R}^{n}$ be a bounded domain with $C^{\infty}$-smooth boundary $\partial \Omega$, for $n \geq 2$. Let $f \in L^{\infty}(\Omega), v_{1}, v_{2} \in L^{2}(\Omega)$, and $v_{3}, \ldots, v_{k} \in L^{\infty}(\Omega)$ be complex valued functions where $k \geq 2$. Then

$$
\int_{\Omega} f v_{1} \cdots v_{k} d x=\sum_{j=1}^{2^{k}} \int_{\Omega} c_{j} f w_{1}^{(j)} \cdots w_{k}^{(j)} d x
$$

where $c_{j} \in\{ \pm 1, \pm i\}$ and $w_{1}^{(j)} \in\left\{\operatorname{Re}\left(v_{1}\right), \operatorname{Im}\left(v_{1}\right)\right\}, \cdots, w_{k}^{(j)} \in\left\{\operatorname{Re}\left(v_{k}\right), \operatorname{Im}\left(v_{k}\right)\right\}$ for $1 \leq j \leq 2^{k}$.
Proof. The result follows by writing

$$
\int_{M} f v_{1} \cdots v_{k} d x=\int_{M} f\left(\operatorname{Re}\left(v_{1}\right)+i \operatorname{Im}\left(v_{1}\right)\right) \cdots\left(\operatorname{Re}\left(v_{k}\right)+i \operatorname{Im}\left(v_{k}\right)\right) d x
$$

and by multiplying out the right hand side.
Lemma 3.1 also holds on Riemannian manifolds $(M, g)$, which will be applied in Section 4.

Proof of Theorem 1.1. Since $\Lambda_{q_{1}}(f)=\Lambda_{q_{2}}(f)$ for all small $f$ and since $\Lambda_{q_{j}}$ is a $C^{k}$ map by Proposition 2.1, one has

$$
\left(D^{k} \Lambda_{q_{1}}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right)=\left(D^{k} \Lambda_{q_{2}}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right)
$$

for all $f_{0}, \ldots, f_{k+1} \in C^{2, \alpha}(\partial \Omega)$ and for $\epsilon_{0}$ small. The integral identity (2.23) applied with $q_{1}$ and $q_{2}$ implies that

$$
\int_{\Omega}\left(q_{1}-q_{2}\right)\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k+1} d x=0
$$

for any real-valued harmonic functions $v_{0}, \ldots, v_{k+1} \in C^{2, \alpha}(\bar{\Omega})$. Let $v_{0}=v_{3}=\ldots=$ $v_{k+1}=1$ be constant functions in $\Omega$. Then

$$
\begin{equation*}
\int_{\Omega}\left(q_{1}-q_{2}\right) v_{1} v_{2} d x=0 \tag{3.1}
\end{equation*}
$$

whenever $v_{j} \in C^{2, \alpha}(\bar{\Omega})$ are real-valued and harmonic. Since the real and imaginary parts of a complex valued harmonic function are harmonic, it follows from Lemma 3.1 that (3.1) remains true for complex valued harmonic functions.

Now let $v_{1}(x)=e^{(-\zeta+i \xi) \cdot x}$ and $v_{2}(x)=e^{(\zeta+i \xi) \cdot x}$ be Calderón's exponential solutions (see [Cal80]), which are harmonic, and where $\zeta, \xi \in \mathbb{R}^{n}$ with $|\zeta|=|\xi|$ and $\zeta \cdot \xi=0$. Then we have

$$
\begin{align*}
0 & =\int_{\Omega}\left(q_{1}-q_{2}\right) v_{1} v_{2} d x \\
& =\int_{\Omega}\left(q_{1}-q_{2}\right) e^{(-\zeta+i \xi) \cdot x} e^{(\zeta+i \xi) \cdot x} d x  \tag{3.2}\\
& =\int_{\Omega}\left(q_{1}-q_{2}\right) e^{2 i \xi \cdot x} d x
\end{align*}
$$

Thus, via (3.2), we obtain that the Fourier transform of the difference $q_{1}-q_{2}$ at $-2 \xi$ is zero. Since $\xi \in \mathbb{R}^{n}$ can be chosen arbitrarily, we must have $q_{1}=q_{2}$ as desired.

Let us give another proof of this result when $n \geq 3$ and when we only assume that $\Lambda_{q_{1}}(f)=\Lambda_{q_{2}}(f)$ for all small $f$ with $f \geq 0$. As before, let $f_{0}=f_{3}=\ldots=f_{k+1}=1$ so that $v_{0}=v_{3}=\ldots=v_{k+1}=1$ in $\Omega$. Then (3.1) holds whenever $f_{1}, f_{2} \geq 0$. Let $x \notin \bar{\Omega}$ and choose the boundary values $f_{1}, f_{2}$ so that $v_{1}(y)=v_{2}(y)=|x-y|^{2-n}$. Then $v_{1}, v_{2}>0$ are harmonic in $\Omega$. Inserting these solutions to (3.1) and writing $q=q_{1}-q_{2}$, we see that

$$
\int_{\Omega}|x-y|^{4-2 n} q(y) d y=0
$$

for $x \notin \bar{\Omega}$. By [Isa90, page 79], the knowledge of the Riesz potential

$$
I_{\beta} \mu(x)=\int_{\Omega}|x-y|^{\beta} d \mu(y)
$$

for $x \notin \bar{\Omega}$ uniquely determines the measure $\mu(y)$ in $\Omega$, when $\beta \neq 2 k$ and $\beta+n \neq$ $2 k+2$ for all $k=0,1, \ldots$. Since these conditions are satisfied for $\beta=4-2 n$, we see that $q=0$ by setting $d \mu(y)=q(y) d y$ above. Isakov [Isa90] credits M. Riesz [Rie38] and M. M. Lavrentiev [Lav67] for the first results about determination of a measure from the Riesz potential.

Proof of Theorem 1.2. Since the DN maps satisfy $\Lambda_{q_{1}}^{\Gamma}(f)=\Lambda_{q_{2}}^{\Gamma}(f)$ for any sufficiently small Dirichlet data $f \in C_{0}^{2, \alpha}(\Gamma)$, we have for any $f_{0}, \ldots, f_{k+1} \in C_{0}^{2, \alpha}(\Gamma)$

$$
\begin{equation*}
\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha} \int_{\partial \Omega}\left(D^{k} \Lambda_{q_{1}}^{\Gamma}-D^{k} \Lambda_{q_{2}}^{\Gamma}\right)_{\epsilon_{0} f_{0}}\left(f_{1}, \ldots, f_{k}\right) f_{k+1} d S=0 \tag{3.3}
\end{equation*}
$$

Therefore, by subtracting the integral identity (2.23) for $q=q_{1}, q_{2}$ and inserting (3.3), one has

$$
\int_{\Omega}\left(q_{1}-q_{2}\right)\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \ldots v_{k+1} d x=0
$$

where $v_{\ell}$ are the solutions of (2.12) in $\Omega$ for $\ell=0,1, \ldots, k+1$ with $\left.v_{\ell}\right|_{\partial \Omega}=f_{\ell}$. Write $F:=\left(q_{1}-q_{2}\right)\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{3} \ldots v_{k+1}$, so that we have

$$
\int_{\Omega} F v_{1} v_{2} d x=0
$$

By applying Lemma 3.1, we see that the last identity is valid for complex-valued harmonic functions $v_{1}, v_{2} \in C^{2, \alpha}(\bar{\Omega})$ with $\operatorname{supp}\left(\left.v_{\ell}\right|_{\partial \Omega}\right) \subset \Gamma$. On the other hand, via the density result of [FKSU09], one can choose $\left\{v_{1} v_{2}\right\}$ to form a dense subset in $L^{1}(\Omega)$ with $\operatorname{supp}\left(\left.v_{1}\right|_{\partial \Omega}\right), \operatorname{supp}\left(\left.v_{2}\right|_{\partial \Omega}\right) \subset \Gamma$. This implies that $F=0$ in $\Omega$. Finally, by choosing $f_{0}, f_{3}, \ldots, f_{k+1} \not \equiv 0$ to be nonnegative Dirichlet data supported in $\Gamma$, we see that $v_{0}, v_{3}, \ldots, v_{k+1}$ are positive in $\Omega$ by the maximum principle. Thus one can conclude that $q_{1}=q_{2}$ in $\Omega$.

Next we prove Theorem 1.3.
Proof of Theorem 1.3. Via Proposition 2.1, let $u_{j} \in C^{2, \alpha}(\bar{\Omega})$, for $j=1,2$, be the unique (small) solutions to

$$
\begin{cases}\Delta u_{j}+a_{j}\left(x, u_{j}\right)=0 & \text { in } \Omega  \tag{3.4}\\ u_{j}=\epsilon_{0} f_{0}+\epsilon_{1} f_{1} & \text { on } \partial \Omega\end{cases}
$$

where $\epsilon_{\ell} \geq 0$ are small parameters and $f_{\ell} \in C_{0}^{2, \alpha}(\Gamma)$, for $\ell=0,1$. Then, as in equation (2.16) in the proof of Proposition 2.3, we have that the first linearization of the solution map $S_{j}$ to (3.4), $j=1,2$, at $h_{0}:=\epsilon_{0} f_{0}$ satisfies

$$
v_{j, 1}^{\epsilon_{0}}:=\left(D S_{j}\right)_{h_{0}}\left(f_{1}\right)=\left.\partial_{\epsilon_{1}} u_{j}\right|_{\epsilon_{1}=0}
$$

where $v_{j, 1}^{\epsilon_{0}}$ satisfies

$$
\begin{cases}\Delta v_{j, 1}^{\epsilon_{0}}=-\partial_{y} a_{j}\left(x,\left.u_{j}\right|_{\epsilon_{1}=0}\right) v_{j, 1}^{\epsilon_{0}} & \text { in } \Omega  \tag{3.5}\\ v_{j, 1}^{\epsilon_{0}}=f_{1} & \text { on } \partial \Omega\end{cases}
$$

for $j=1,2$. Analogously to (2.17) in the proof of Proposition 2.3, one has

$$
v_{j, 1}^{\epsilon_{0}} \rightarrow v_{1} \text { in } C^{2, \alpha}(\bar{\Omega}), \quad \text { as } \epsilon_{0} \rightarrow 0
$$

where $v_{1}$ solves $\Delta v_{1}=0$ in $\Omega$ and $\left.v_{1}\right|_{\partial \Omega}=f_{1}$.

Fix $f_{2} \in C_{0}^{2, \alpha}(\Gamma)$ and let $v_{2}$ solve $\Delta v_{2}=0$ in $\Omega$ with $\left.v_{2}\right|_{\partial \Omega}=f_{2}$. Since $\Lambda_{a_{1}}^{\Gamma}(f)=$ $\Lambda_{a_{2}}^{\Gamma}(f)$ for any sufficiently small $f \in C_{0}^{2, \alpha}(\Gamma)$, integration by parts and (3.5) yield that

$$
\begin{align*}
0= & \left.\partial_{\epsilon_{1}}\right|_{\epsilon_{1}=0}\left(\int_{\partial \Omega} f_{2}\left(\Lambda_{a_{1}}^{\Gamma}-\Lambda_{a_{2}}^{\Gamma}\right)\left(\epsilon_{0} f_{0}+\epsilon_{1} f_{1}\right) d S\right) \\
= & \left.\partial_{\epsilon_{1}}\right|_{\epsilon_{1}=0}\left(\int_{\Omega} v_{2}\left(\Delta u_{1}-\Delta u_{2}\right) d x\right)+\left.\partial_{\epsilon_{1}}\right|_{\epsilon_{1}=0}\left(\int_{\Omega} \nabla v_{2} \cdot \nabla\left(u_{1}-u_{2}\right) d x\right) \\
= & -\left.\int_{\Omega} v_{2} \partial_{\epsilon_{1}}\right|_{\epsilon_{1}=0}\left(a_{1}\left(x, u_{1}\right)-a_{2}\left(x, u_{2}\right)\right) d x \\
& +\left.\partial_{\epsilon_{1}}\right|_{\epsilon_{1}=0}\left(\int_{\partial \Omega} \partial_{\nu} v_{2}\left(\left.u_{1}\right|_{\partial \Omega}-\left.u_{2}\right|_{\partial \Omega}\right) d S\right)  \tag{3.6}\\
= & -\int_{\Omega} v_{2}\left(\partial_{y} a_{1}\left(x,\left.u_{1}\right|_{\epsilon_{1}=0}\right) v_{1,1}^{\epsilon_{0}}-\partial_{y} a_{2}\left(x,\left.u_{2}\right|_{\epsilon_{1}=0}\right) v_{2,1}^{\epsilon_{0}}\right) d x \\
& +\int_{\partial \Omega} \partial_{\nu} v_{2}\left(f_{1}-f_{1}\right) d S \\
= & -\int_{\Omega} v_{2}\left(\partial_{y} a_{1}\left(x,\left.u_{1}\right|_{\epsilon_{1}=0}\right) v_{1,1}^{\epsilon_{0}}-\partial_{y} a_{2}\left(x,\left.u_{2}\right|_{\epsilon_{1}=0}\right) v_{2,1}^{\epsilon_{0}}\right) d x .
\end{align*}
$$

For $j=1,2$, the function

$$
w_{j}:=\left.u_{j}\right|_{\epsilon_{1}=0}
$$

now solves

$$
\begin{cases}\Delta w_{j}+a_{j}\left(x, w_{j}\right)=0 & \text { in } \Omega, \\ w_{j}=\epsilon_{0} f_{0} & \text { on } \partial \Omega .\end{cases}
$$

By (2.5) we have

$$
\left\|w_{j}\right\|_{C^{2, \alpha}(\bar{\Omega})} \leq C \epsilon_{0}\left\|f_{0}\right\|_{C^{2, \alpha}(\partial \Omega)}
$$

Since $\Delta\left(w_{j}-\epsilon_{0} v_{0}\right)=-a_{j}\left(x, w_{j}\right)$ in $\Omega$ with $w_{j}-\left.\epsilon_{0} v_{0}\right|_{\partial \Omega}=0$, Schauder estimates imply that

$$
\left\|w_{j}-\epsilon_{0} v_{0}\right\|_{C^{2, \alpha}(\bar{\Omega})} \leq C\left\|a_{j}\left(x, w_{j}\right)\right\|_{C^{\alpha}(\bar{\Omega})} .
$$

Using the Taylor formula as in (2.8) together with the conditions

$$
a_{j}(x, 0)=\partial_{y} a_{j}(x, 0)=0
$$

gives that

$$
a_{j}\left(x, w_{j}(x)\right)=w_{j}(x) \int_{0}^{1}\left(\partial_{y} a_{j}\left(x, t w_{j}(x)\right)-\partial_{y} a_{j}(x, 0)\right) d t
$$

We may now apply (2.6) with $b$ replaced by $a_{j}$ to obtain that

$$
\begin{align*}
\left\|w_{j}-\epsilon_{0} v_{0}\right\|_{C^{2, \alpha}(\bar{\Omega})} & \leq C\left\|w_{j}\right\|_{C^{\alpha}(\bar{\Omega})} \int_{0}^{1}\left\|\partial_{y} a_{j}\left(x, t w_{j}\right)-\partial_{y} a_{j}(x, 0)\right\|_{C^{\alpha}(\bar{\Omega})} d t  \tag{3.7}\\
& =o\left(\epsilon_{0}\right)
\end{align*}
$$

as $\epsilon_{0} \rightarrow 0$.
We have by assumption $a_{j}(x, y) \sim \sum_{l=1}^{\infty} b_{j, l}(x, y)$, where each $b_{j, l}(\cdot, y) \in C^{\alpha}(\bar{\Omega})$ is homogeneous of order $r_{l}>1$ with respect to the variable $y \in \mathbb{R}$, for $l \geq 1$. Let us also write $\beta_{j, N}:=a_{j}-\sum_{l=1}^{N-1} b_{j, l}$ for $j=1,2$ and $N \geq 1$, with $\beta_{j, 1}=a_{j}$. Then $\beta_{j, N}$
is in $C_{\text {loc }}^{1, \alpha}\left(\mathbb{R}, C^{\alpha}(\bar{\Omega})\right)$ as in Definition 1.1. It follows from (1.3) that, in particular,

$$
\left\|\partial_{y} a_{j}(\cdot, y)-\sum_{l=1}^{N-1} \partial_{y} b_{j, l}(\cdot, y)\right\|_{L^{\infty}(\Omega)} \leq C_{N}|y|^{r_{N}-1}, \quad|y| \leq 1
$$

for $j=1,2$.
We apply the above with $N=2$ and $y=w_{j}(x)=\left.u_{j}(x)\right|_{\epsilon_{1}=0}$ to have for $x \in \bar{\Omega}$, for $j=1,2$ that

$$
\left|\partial_{y} a_{j}\left(x, w_{j}\right)-\partial_{y} b_{j, 1}\left(x, w_{j}\right)\right| \leq C_{2}\left|w_{j}\right|^{r_{2}-1} \leq C \epsilon_{0}^{r_{2}-1}
$$

Multiplying this by $\epsilon_{0}^{-r_{1}+1}$ and using the facts that $r_{2}>r_{1}$ and $\partial_{y} b_{j, 1}(x, y)$ is homogeneous of order $r_{1}-1$ in $y$, we obtain in $L^{\infty}(\Omega)$ that

$$
\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-r_{1}+1} \partial_{y} a_{j}\left(x, w_{j}\right)=\lim _{\epsilon_{0} \rightarrow 0} \partial_{y} b_{j, 1}\left(x, \epsilon_{0}^{-1} w_{j}\right)=\partial_{y} b_{j, 1}\left(x, v_{0}\right)
$$

Here in the last equality we additionally used (3.7). Recall that we also have that the limit $\lim _{\epsilon_{0} \rightarrow 0} v_{j, 1}^{\epsilon_{0}}=v_{1}$ in $C^{2, \alpha}(\bar{\Omega})$, for both $j=1,2$. Hence, we obtain

$$
\begin{aligned}
0 & =\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-r_{1}+1} \int_{\Omega} v_{2}\left[\partial_{y} a_{1}\left(x,\left.u_{1}\right|_{\epsilon_{1}=0}\right) v_{1,1}^{\epsilon_{0}}-\partial_{y} a_{2}\left(x,\left.u_{2}\right|_{\epsilon_{1}=0}\right) v_{2,1}^{\epsilon_{0}}\right] d x \\
& =\int_{\Omega}\left[\partial_{y} b_{1,1}\left(x, v_{0}\right)-\partial_{y} b_{2,1}\left(x, v_{0}\right)\right] v_{1} v_{2} d x
\end{aligned}
$$

Via the density result of [FKSU09], products $v_{1} v_{2}$ of pairs of harmonic functions with boundary values supported in $\Gamma \subset \partial \Omega$ are dense in $L^{1}(\Omega)$. Therefore, we must have

$$
\partial_{y} b_{1,1}\left(x, v_{0}\right)=\partial_{y} b_{2,1}\left(x, v_{0}\right), \text { for } x \in \Omega
$$

In addition, notice that the boundary value $f_{0} \in C_{0}^{2, \alpha}(\Gamma)$ has been arbitrary so far. Let $x_{0} \in \Omega$, let $y_{0} \in \mathbb{R}$ and let us choose by Runge approximation (see e.g. [LLS19, Proposition A.2]) a boundary value $f_{0}=f_{0, x_{0}} \in C_{0}^{\infty}(\Gamma)$ so that

$$
\begin{equation*}
v_{0}\left(x_{0}\right)=y_{0} . \tag{3.8}
\end{equation*}
$$

We deduce that

$$
\partial_{y} b_{1,1}\left(x_{0}, y_{0}\right)=\partial_{y} b_{2,1}\left(x_{0}, y_{0}\right)
$$

for any $x_{0} \in \Omega$ and any $y_{0}$. Thus we have $\partial_{y} b_{1,1}=\partial_{y} b_{2,1}$. By Euler's homogeneous function theorem, we have

$$
b_{1,1}(x, y)=\frac{y}{r_{1}} \partial_{y} b_{1,1}(x, y)=\frac{y}{r_{1}} \partial_{y} b_{2,1}(x, y)=b_{2,1}(x, y)
$$

where $r_{1}>1$ is the degree of homogeneity for $b_{j, 1}(x, y)$ with respect to the $y$ variable, for $j=1,2$. Thus $b_{1,1}=b_{2,1}$.

We proceed by induction on the index $l \in \mathbb{N}$ of $b_{j, l}, j=1,2$, to show that $b_{1, l}=b_{2, l}$ for any $l \in \mathbb{N}$. We have already shown the case $l=1$. Let us then make the induction assumption that $b_{1, l}=b_{2, l}$ for $l=1, \ldots, L$, for some $L \in \mathbb{N}$. Then, we have that

$$
\begin{aligned}
& \left|\left(\partial_{y} a_{1}(x, y)-\partial_{y} a_{2}(x, y)\right)-\left(\partial_{y} b_{1, L+1}(x, y)-\partial_{y} b_{2, L+1}(x, y)\right)\right| \\
= & \mid\left(\partial_{y} a_{1}(x, y)-\partial_{y} a_{2}(x, y)\right)-\sum_{l=1}^{L} \partial_{y} b_{1, l}(x, y)+\sum_{l=1}^{L} \partial_{y} b_{2, l}(x, y) \\
& \quad-\left(\partial_{y} b_{1, L+1}(x, y)-\partial_{y} b_{2, L+1}(x, y)\right) \mid \\
= & \left|\left(\partial_{y} a_{1}(x, y)-\sum_{l=1}^{L+1} \partial_{y} b_{1, l}(x, y)\right)-\left(\partial_{y} a_{2}(x, y)-\sum_{l=1}^{L+1} \partial_{y} b_{2, l}(x, y)\right)\right| \\
= & \left|\partial_{y} \beta_{1, L+2}(x, y)-\partial_{y} \beta_{2, L+2}(x, y)\right| \leq 2 C_{L+2}|y|^{r_{L+2}-1} .
\end{aligned}
$$

Here we used the induction assumption in the first equality. Applying this for $y=w_{j}(x)=\left.u_{j}(x)\right|_{\epsilon_{1}=0}$ we have for $x \in \bar{\Omega}$, and for $j=1,2$, that

$$
\left|\left(\partial_{y} a_{1}\left(x, w_{j}\right)-\partial_{y} a_{2}\left(x, w_{j}\right)\right)-\left(\partial_{y} b_{1, L+1}\left(x, w_{j}\right)-\partial_{y} b_{2, L+1}\left(x, w_{j}\right)\right)\right| \leq C \epsilon_{0}^{r_{L+2}-1},
$$

for some constant $C>0$. Here we used again $\left\|w_{j}\right\|_{C^{2, \alpha}(\bar{\Omega})} \leq C \epsilon_{0}\left\|f_{0}\right\|_{C^{2, \alpha}(\partial \Omega)}$
Therefore, by using (3.7), homogeneity and $r_{L+2}>r_{L+1}$, we obtain in $L^{\infty}(\Omega)$ that

$$
\begin{aligned}
& \lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-r_{L+1}+1}\left(\partial_{y} a_{1}\left(x,\left.u_{1}\right|_{\epsilon_{1}=0}\right)-\partial_{y} a_{2}\left(x,\left.u_{2}\right|_{\epsilon_{1}=0}\right)\right) \\
= & \lim _{\epsilon_{0} \rightarrow 0}\left(\partial_{y} b_{1, L+1}\left(x, \epsilon_{0}^{-1} w_{1}\right)-\partial_{y} b_{2, L+1}\left(x, \epsilon_{0}^{-1} w_{2}\right)\right) \\
= & \partial_{y} b_{1, L+1}\left(x, v_{0}\right)-\partial_{y} b_{2, L+1}\left(x, v_{0}\right) .
\end{aligned}
$$

By repeating the arguments we used to prove the special case $N=2$, which especially use the integral identity (3.6) and [FKSU09], we obtain

$$
\partial_{y} b_{1, L+1}=\partial_{y} b_{2, L+1}
$$

By Euler's homogeneous function theorem again, we then have $b_{1, L+1}=b_{2, L+1}$ in $\Omega$ as desired, which concludes the induction step and the proof of the theorem.

Remark 3.2. In the previous proof we recovered the expansion coefficients $b_{l}(x, y)$ of the potential $a \sim \sum_{l=1}^{\infty} b_{l}$ at arbitrary point $\left(x_{0}, y_{0}\right) \in \Omega \times \mathbb{R}$. This was done by using Runge approximation (see (3.8)) to select a boundary value $f_{0}$ so that the corresponding solution $v_{0}$ satisfies $v_{0}\left(x_{0}\right)=y_{0}$. This is slightly different from earlier results in [LLLS20a, LLLS20b, KU20b], where one recovers the Taylor coefficients $\tilde{b}_{l}(x, y):=\partial_{y}^{l} \tilde{a}(x, y)$ of an unknown smooth potential $\tilde{a}(x, y)$ only at $y=0, x \in \Omega$.

In the end of this section, let us prove the simultaneous recovery of an obstacle and a potential.

Proof of Theorem 1.4. For $\ell=0,1$, let $\epsilon_{\ell} \geq 0$ be sufficiently small parameters, and $f_{\ell} \in C_{0}^{2, \alpha}(\Gamma)$. Consider the Dirichlet data $f=\epsilon_{0} f_{0}+\epsilon_{1} f_{1}$ and let $u_{j}=u_{j}(x)$ be the solution of

$$
\begin{cases}\Delta u_{j}+a_{j}\left(x, u_{j}\right)=0 & \text { in } \Omega  \tag{3.9}\\ u_{j}=0 & \text { on } \partial D_{j} \\ u_{j}=f & \text { on } \partial \Omega\end{cases}
$$

for $j=1,2$, where $a_{j}=a_{j}(x, z)$ are polyhomogeneous in the sense of Definition 1.1 with $x \in \Omega \backslash \overline{D_{j}}$. We first show that $D_{1}=D_{2}$ and then recover the coefficients similarly as in the proof of Theorem 1.3.

## Step 1. Recovering the obstacle.

As in the proof of Proposition 2.3, see (2.16), we have that the first linearization of the solution map $S_{j}$ to (3.9), $j=1,2$, at $h_{0}:=\epsilon_{0} f_{0}$ satisfies

$$
v_{j, \ell}^{\epsilon_{0}}:=\left(D S_{j}\right)_{h_{0}}\left(f_{\ell}\right)=\left.\partial_{\epsilon_{\ell}} u_{j}\right|_{\epsilon^{\prime}=0}
$$

where $v_{j, \ell}^{\epsilon_{0}}$ is the solution of

$$
\begin{cases}\Delta_{g} v_{j, \ell}^{\epsilon_{0}}=-\partial_{y} a_{j}\left(x,\left.u_{j}\right|_{\epsilon^{\prime}=0}\right) v_{j, \ell}^{\epsilon_{0}} & \text { in } \Omega \\ v_{j, \ell}^{\epsilon_{0}}=0 & \text { on } \partial D_{j} \\ v_{j, \ell}^{\epsilon_{0}}=f_{\ell} & \text { on } \partial \Omega\end{cases}
$$

Analogously to (2.17) in the proof of Proposition 2.3, one has

$$
v_{j, \ell}^{\epsilon_{0}} \rightarrow v_{j}^{(\ell)} \text { in } C^{2, \alpha}\left(\bar{\Omega} \backslash D_{j}\right), \quad \text { as } \epsilon_{0} \rightarrow 0
$$

where

$$
\begin{cases}\Delta v_{j}^{(\ell)}=0 & \text { in } \Omega \backslash \overline{D_{j}}, \\ v_{j}^{(\ell)}=0 & \text { on } \partial D_{j}, \\ v_{j}^{(\ell)}=f_{\ell} & \text { on } \partial \Omega\end{cases}
$$

for $j=1,2$ and $\ell=0,1$. The rest of the proof is the analogous to the proof of [LLLS20b, Theorem 1.2]. (See also [KU20b, Theorem 1.6].) For the sake of completeness, we offer details of the proof below.

Let $G$ be the connected connected component of $\Omega \backslash\left(\overline{D_{1} \cup D_{2}}\right)$, whose boundary contains $\partial \Omega$. Consider the function $\widetilde{v}^{(\ell)}:=v_{1}^{(\ell)}-v_{2}^{(\ell)}$, which solves

$$
\begin{cases}\Delta \widetilde{v}^{(\ell)}=0 & \text { in } G, \\ \widetilde{v}^{(\ell)}=\partial_{\nu} \widetilde{v}^{(\ell)}=0 & \text { on } \Gamma,\end{cases}
$$

where we have used that $\Lambda_{a_{1}, D_{1}}^{\Gamma}(f)=\Lambda_{a_{2}, D_{2}}^{\Gamma}(f)$, which holds for all sufficiently small Dirichlet data $f \in C_{0}^{2, \alpha}(\Gamma)$. By the unique continuation of harmonic functions this yields that $\widetilde{v}^{(\ell)}=0$ in $G$. That is, for $\ell=0$, 1 , we have

$$
\begin{equation*}
v_{1}^{(\ell)}=v_{2}^{(\ell)} \text { in } G \tag{3.10}
\end{equation*}
$$

We use a contradiction argument to prove $D_{1}=D_{2}$. For this, let us assume that $D_{1} \neq D_{2}$. Note that the connected component $G \neq \emptyset$. By using [LLLS20b, Lemma A.3], there exists

$$
x_{1} \in \partial G \cap\left(\Omega \backslash \overline{D_{1}}\right) \cap \partial D_{2} .
$$

Since $x_{1} \in \partial D_{2}$, we have $v_{2}^{(\ell)}\left(x_{1}\right)=0$. By (3.10) and continuity, we also have that $v_{1}^{(\ell)}\left(x_{1}\right)=0$. Note that $x_{1}$ is an interior point of the open set $\Omega \backslash \overline{D_{1}}$.

We next fix one of the boundary values $f_{\ell}$ to be non-negative and not identically 0 . Since $v_{1}^{(\ell)}\left(x_{1}\right)=0$, the maximum principle implies that $v_{1}^{(\ell)} \equiv 0$ in $\Omega \backslash \overline{D_{1}}$, which contradicts to the assumption that $v_{1}^{(\ell)}=f_{\ell}$ on $\partial \Omega$ is not identically zero (because the harmonic function $v_{1}^{(\ell)}$ is continuous up to boundary). This shows that

$$
D:=D_{1}=D_{2} .
$$

Step 2. Recovering the coefficient.
Since we have proved that $D_{1}=D_{2}=D$, it follows that the partial data Dirichlet-to-Neumann maps for the equations $\Delta u+a_{j}(x, u)=0$ in $\Omega \backslash \bar{D}$ agree on $\Gamma$. Applying Theorem 1.3 in the connected set $\Omega \backslash \bar{D}$ then implies that $b_{1, l}=b_{2, l}$ for all $l \in \mathbb{N}$. This concludes the proof.

## 4. Global uniqueness in Riemannian manifolds

In this last section of this paper, we prove Theorem 1.5 and Theorem 1.6. In our earlier work [LLLS20a], we proved similar theorems for power type nonlinearities, with integer exponents. We begin with the proof of Theorem 1.5.

Proof of Theorem 1.5. The proof is similar to the proof of [LLLS20a, Theorem 1.2]. We first recover the manifold and the its conformal class by the first linearization. After that we use the integral identity (2.11) to recover the potential.

Step 1. Recovering the conformal manifold.
By using Proposition 2.1, the equality $\Lambda_{M_{1}, g_{1}, q_{1}}(f)=\Lambda_{M_{2}, g_{2}, q_{2}}(f)$, for all $f \in$ $C^{2, \alpha}(\partial M)$ with $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta$, where $\delta>0$ is a sufficiently small number, implies

$$
\left(D \Lambda_{M_{1}, g_{1}, q_{1}}\right)_{0}=\left(D \Lambda_{M_{2}, g_{2}, q_{2}}\right)_{0}
$$

Here, for $j=1,2$, the maps $\left(D \Lambda_{M_{j}, g_{j}, q_{j}}\right)_{0}$ are the DN maps of the linearizations of the equations $\Delta_{g_{j}} u_{j}+q_{j}\left|u_{j}\right|^{r-1} u_{j}=0$ in $M_{j}$ at a boundary value $f=0$. This implies that the DN maps on $\partial M$ of the first linearized equation

$$
\begin{cases}\Delta_{g_{j}} v_{j}=0 & \text { in } M_{j} \\ v_{j}=f & \text { on } \partial M\end{cases}
$$

agree on $\partial M$. That is, we know the DN maps on $\partial M$ of the anisotropic Calderón problem on two-dimensional Riemannian manifolds. Thus, as noted in the proof of [LLLS20a, Theorem 1.2], we may use [LLS19, Theorem 5.1] to determine the manifold and the Riemannian metric up to a conformal transformation: There exists a $C^{\infty}$ smooth diffeomorphism $J: M_{1} \rightarrow M_{2}$ such that

$$
\sigma J^{*} g_{2}=g_{1} \text { in } M_{1}
$$

with $\left.J\right|_{\partial M}=$ Id. Here the function $\sigma \in C^{\infty}\left(M_{1}\right)$ is positive with $\left.\sigma\right|_{\partial M}=1$.
Step 2. Recovering the potential.
Let us transform the equation $\Delta_{g_{2}} u_{2}+q_{2}\left|u_{2}\right|^{r-1} u_{2}=0$ from the manifold $\left(M_{2}, g_{2}\right)$ into the manifold ( $M_{1}, g_{1}$ ) as follows. We denote in $M_{1}$

$$
\widetilde{q}_{2}=\sigma^{-1}\left(q_{2} \circ J\right) \equiv \sigma^{-1} J^{*} q_{2}
$$

Let $u_{2}$ be the solution to

$$
\begin{cases}\Delta_{g_{2}} u_{2}+q_{2}\left|u_{2}\right|^{r-1} u_{2}=0 & \text { in } M_{2}  \tag{4.1}\\ u_{2}=f & \text { on } \partial M\end{cases}
$$

where $f \in C^{2, \alpha}(\partial M)$ with $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta, \delta>0$ sufficiently small. Let us define

$$
\widetilde{u}_{2}:=J^{*} u_{2} \equiv u_{2} \circ J,
$$

in $M_{1}$. Then $\widetilde{u}_{2}$ satisfies in $M_{1}$

$$
\begin{aligned}
& \Delta_{g_{1}} \widetilde{u}_{2}+\widetilde{q}_{2}\left|\widetilde{u}_{2}\right|^{r-1} \widetilde{u}_{2} \\
= & \Delta_{\sigma J^{*} g_{2}} \widetilde{u}_{2}+\widetilde{q}_{2}\left|\widetilde{u}_{2}\right|^{r-1} \widetilde{u}_{2} \\
= & \sigma^{-1} \Delta_{J^{*} g_{2}} \widetilde{u}_{2}+\sigma^{-1}\left(J^{*} q_{2}\right)\left|\widetilde{u}_{2}\right|^{r-1} \widetilde{u}_{2} \\
= & \sigma^{-1} J^{*}\left(\Delta_{g_{2}} u_{2}\right)+\sigma^{-1}\left(J^{*} q_{2}\right)\left|J^{*} u_{2}\right|^{r-1} J^{*} u_{2} \\
= & \sigma^{-1} J^{*}\left(\Delta_{g_{2}} u_{2}+q_{2}\left|u_{2}\right|^{r-1} u_{2}\right) .
\end{aligned}
$$

Here we used the conformal invariance of the Laplace-Beltrami operator in two dimensions and the coordinate invariance of Laplace-Beltrami operator in the second and third equality respectively. Therefore, one has

$$
\begin{cases}\Delta_{g_{1}} \widetilde{u}_{2}+\widetilde{q}_{2}\left|\widetilde{u}_{2}\right|^{r-1} \widetilde{u}_{2}=0 & \text { in } M_{1}  \tag{4.2}\\ \widetilde{u}_{2}=f & \text { on } \partial M\end{cases}
$$

where we have used that $u_{2}$ is the solution of (4.1), $f \in C^{2, \alpha}(\partial M)$ and $\left.J\right|_{\partial M}=\mathrm{Id}$.
Let $u_{1}$ be the solution to the nonlinear equation $\Delta_{g_{1}} u_{1}+q_{1}\left|u_{1}\right|^{r-1} u_{1}=0$ in $M_{1}$ with potential $q_{1}$ and boundary data $f$. We show next that

$$
\begin{equation*}
\partial_{\nu_{1}} u_{1}=\partial_{\nu_{1}} \widetilde{u}_{2} \text { on } \partial M \tag{4.3}
\end{equation*}
$$

Via the assumption that $\Lambda_{M_{1}, g_{1}, q_{1}}(f)=\Lambda_{M_{2}, g_{2}, q_{2}}(f)$, it follows that if $u_{1}=u_{2}=$ $f \in C^{2, \alpha}(\partial M)$ on $\partial M$, then

$$
\begin{equation*}
\partial_{\nu_{1}} u_{1}=\partial_{\nu_{2}} u_{2} \text { on } \partial M \tag{4.4}
\end{equation*}
$$

We compute that

$$
\begin{equation*}
\partial_{\nu_{2}} u_{2}=\nu_{2} \cdot d u_{2}=\nu_{2} \cdot d\left(u_{2} \circ J \circ J^{-1}\right)=\left(J_{*}^{-1} \nu_{2}\right) \cdot d \widetilde{u}_{2}=\nu_{1} \cdot d \widetilde{u}_{2}=\partial_{\nu_{1}} \widetilde{u}_{2}, \tag{4.5}
\end{equation*}
$$

where - denotes the canonical pairing between vectors and covectors, and $d$ is the exterior derivative of a function. For example $\nu_{2} \cdot d u_{2}=g\left(\nu_{2}, \nabla u_{2}\right)=\sum_{k=1}^{2} \nu_{2}^{k} \partial_{k} u_{2}$. We used that $J: M_{1} \rightarrow M_{2}$ is conformal diffeomorphism, $\sigma J^{*} g_{2}=g_{1}$, with $\left.J\right|_{\partial M}=$ Id and $\left.\sigma\right|_{\partial M}=1$ in (4.5). Combining (4.4) and (4.5), we have (4.3) as claimed.

We have by (4.3) that

$$
\begin{equation*}
\Lambda_{M_{1}, g_{1}, q_{1}}(f)=\partial_{\nu_{1}} u_{1}=\partial_{\nu_{1}} \widetilde{u}_{2}=\widetilde{\Lambda}_{M_{1}, g_{1}, \widetilde{q}_{2}}(f) \tag{4.6}
\end{equation*}
$$

for all $f \in C^{2, \alpha}(\partial M)$ with $\|f\|_{C^{2, \alpha}(\partial M)} \leq \delta$, where $\widetilde{\Lambda}_{M_{1}, g_{1}, \widetilde{q}_{2}}$ denotes the DN map of the Dirichlet problem (4.2) on $\partial M$.

We apply Proposition 2.3 on $\left(M_{1}, g_{1}\right)$, the DN maps $\Lambda_{M_{1}, g_{1}, q_{1}}$ and $\widetilde{\Lambda}_{M_{1}, g_{1}, \widetilde{q}_{2}}$, which agree by (4.6). By Proposition 2.1 we have

$$
\left.\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha}\left(D^{k} \Lambda_{M_{1}, g_{1}, q_{1}}\right)\right|_{\epsilon_{0} f_{0}}=\left.\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha}\left(D^{k} \widetilde{\Lambda}_{M_{1}, g_{1}, \widetilde{q}_{2}}\right)\right|_{\epsilon_{0} f_{0}} \text { on } \partial M
$$

and by Proposition 2.3

$$
\int_{M_{1}}\left(q_{1}-\widetilde{q}_{2}\right)\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k+1} d V=0
$$

where $v_{0}, v_{1}, \cdots, v_{k} \in C^{2, \alpha}\left(M_{1}\right)$ are harmonic functions in $\left(M_{1}, g_{1}\right)$ with $r=k+$ $\alpha>1$. We can choose $v_{0}=v_{1}=\cdots=v_{k-2}=1$ in $M_{1}$, hence

$$
\int_{M_{1}}\left(q_{1}-\widetilde{q}_{2}\right) v_{k-1} v_{k} d V=0
$$

for any harmonic functions $v_{k-1}$ and $v_{k}$ in $M_{1}$.
By choosing $v_{k-1}$ and $v_{k}$ to be complex geometrical optics solutions constructed in [GT11] (see the proof of Proposition 5.1 in [GT11]), we conclude that

$$
q_{1}=\widetilde{q}_{2} \text { in } M_{1} .
$$

We point out that the construction in [GT11] can be simplified in our case where $v_{k-1}$ and $v_{k}$ are harmonic. In such case, Carleman estimates are not needed and the construction in [GST19] would suffice. We have proven the claim.

Proof of Theorem 1.6. Let us write $r=k+\alpha, k \in \mathbb{N}, k \geq 3$ and $\alpha \in(0,1)$. For $j=$ 1,2 , consider $\Lambda_{q_{j}}$ to be the DN map for the equation $\Delta_{g} u_{j}+q_{j}\left|u_{j}\right|^{r-1} u_{j}=0$ in $M$. If $\Lambda_{q_{1}}(f)=\Lambda_{q_{2}}(f)$ for any sufficiently small $f \in C^{2, \alpha}(\partial M)$, then by Proposition 2.1

$$
\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha}\left(D^{k} \Lambda_{q_{1}}\right)_{\epsilon_{0} f_{0}}=\lim _{\epsilon_{0} \rightarrow 0} \epsilon_{0}^{-\alpha}\left(D^{k} \Lambda_{q_{2}}\right)_{\epsilon_{0} f_{0}}
$$

Hence, by Proposition 2.3, we have

$$
\int_{M}\left(q_{1}-q_{2}\right)\left|v_{0}\right|^{r-1} v_{0}^{1-k} v_{1} \cdots v_{k+1} d V=0
$$

where $v_{j} \in C^{2, \alpha}(M)$ are harmonic functions in $M$. Therefore, by choosing $v_{0} \equiv$ 1 and by using [LLLS20a, Proposition 5.1], one obtains that $q_{1}=q_{2}$ in $M$, as desired.

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