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# THE INVERSE PROBLEM FOR THE DIRICHLET-TO-NEUMANN MAP <br> ON LORENTZIAN MANIFOLDS 

# THE INVERSE PROBLEM FOR THE DIRICHLET-TO-NEUMANN MAP ON LORENTZIAN MANIFOLDS 

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We consider the Dirichlet-to-Neumann map $\Lambda$ on a cylinder-like Lorentzian manifold related to the wave equation related to the metric $g$, the magnetic field $A$ and the potential $q$. We show that we can recover the jet of $g, A, q$ on the boundary from $\Lambda$ up to a gauge transformation in a stable way. We also show that $\Lambda$ recovers the following three invariants in a stable way: the lens relation of $g$, and the light ray transforms of $A$ and $q$. Moreover, $\Lambda$ is an FIO away from the diagonal with a canonical relation given by the lens relation. We present applications for recovery of $A$ and $q$ in a logarithmically stable way in the Minkowski case, and uniqueness with partial data.

## 1. Introduction and main results

Let $(M, g)$ be a Lorentzian manifold of dimension $1+n, n \geq 2$; i.e., $g$ is a metric with signature $(-1,1, \ldots, 1)$. Suppose a part of $\partial M$ is timelike. An example of $M$ is a cylinder-like domain representing a moving and shape-changing compact manifold in the $x$-space (if we have fixed time and space variables) with the requirement that the normal speed of the boundary is less than 1 ; see Section 5 .

Denote the wave operator by $\square_{g}$; in local coordinates $x=\left(x^{0}, \ldots, x^{n}\right)$ it takes the form

$$
\square_{g}:=\frac{1}{\sqrt{|\operatorname{det} g|}} \partial_{j}\left(\sqrt{|\operatorname{det} g|} g^{j k} \partial_{k}\right)
$$

Consider the following operator $P=P_{g, A, q}$, which is a first-order perturbation of $\square_{g}$ :

$$
\begin{equation*}
P=P_{g, A, q}:=\frac{1}{\sqrt{|\operatorname{det} g|}}\left(\partial_{j}-i A_{j}\right) \sqrt{|\operatorname{det} g|} g^{j k}\left(\partial_{k}-i A_{k}\right)+q . \tag{1}
\end{equation*}
$$

Here $i=\sqrt{-1}, A$ is a smooth 1-form on $M$, and $q$ is a smooth function on $M$.
The goal of this work is to study the inverse problem of recovery of $g, A$ and $q$, up to a data-preserving gauge transformation, from the outgoing Dirichlet-to-Neumann (DN) $\Lambda$ map on a timelike boundary associated with the wave equation

$$
\begin{equation*}
P u=0 \quad \text { in } M . \tag{2}
\end{equation*}
$$

We are motivated by applications in relativity but also in applications to classical wave-propagation problems with media moving and/or changing at a speed not negligible compared to the wave speed. We

[^0]are interested in possible stability results even though some steps in the recovery are inherently unstable. This problem remains widely open. The results we prove are the following. First, we show that one can recover the jet of $g, A, q$ at the boundary (up to a gauge transform) in a Hölder-stable way. Next, we show that one can extract the natural geometric invariants of $g, A, q$ from $\Lambda$ in a Hölder-stable way. More precisely, $\Lambda$ recovers the lens relation $\mathcal{L}$ related to $g$, in stable way. If we know $g$, the light ray transform $L_{1} A$ of $A$ is recovered stably. If $g$ and $A$ are known, the light ray transform $L_{0} q$ of $q$ is recovered stably. The lens relation $\mathcal{L}$ is the canonical relation of the Fourier integral operator (FIO) $\Lambda$ away from the diagonal, and the light ray transforms $L_{1} A$ and $L_{0} q$ are in fact encoded in the principal and the subprincipal symbols of it. In fact, $\mathcal{L}$ is directly measurable from $\Lambda$.

Since the results we prove are local or semilocal (near a fixed lightlike geodesic), and the proofs are microlocal, we do not formulate a global mixed problem for the wave equation at the beginning but we do consider one in Section 5. In fact, existence of solutions of such problems depends on global properties of $(M, g)$, one of which is global hyperbolicity, which are not needed for our weaker formulation and for the proofs. Instead, we define the DN map up to smoothing operators only. In the case when one can prove the existence of a global solution, the true DN map would coincide with ours up to a smoothing error, see Section 5, and our results are not affected by adding smoothing operators.

This problem has a long history in the stationary Riemannian setting, i.e., when $M=[0, T] \times M_{0}$, where $\left(M_{0}, g\right)$ is a compact Riemannian manifold with boundary, and the metric is $-\mathrm{d} t^{2}+g_{i j}(x) \mathrm{d} x^{i} \mathrm{~d} x^{j}$. The boundary control method [Belishev 1987] and Tataru's uniqueness continuation theorem [1995; 1999] give uniqueness provided that $T$ is greater than a certain sharp critical value $T$, as shown by Belishev and Kurylev [1992]; see also the survey [Belishev 2007]. Stability however does not follow from such arguments. Stability results for recovering of the metric and lower-order terms appeared in [Stefanov and Uhlmann 1998; 2005b; Montalto 2014; Bellassoued and Dos Santos Ferreira 2011; Bao and Zhang 2014], with [Montalto 2014] covering the general case. A main assumption in those works is that the metric is simple, i.e., that there are no conjugate points and the boundary is strictly convex (a not so essential assumption) and the main technical tool for recovery of the metric is to reduce it to stability for the boundary/lens rigidity problem; see, e.g., [Stefanov and Uhlmann 2005a]. For related results, we refer to [Isakov and Sun 1992; Sun 1990]. Recently, the progress in treating the local rigidity problem allowed results under the more general foliation condition [Stefanov et al. 2016], which allows conjugate points. In any case, some condition is believed to be necessary for stability. It is worth noticing that all inverse (hyperbolic) scattering problems for compactly supported perturbations are equivalent to inverse DN map problems.

Recently, there has been increased interest in this problem or in related inverse scattering problems in time-space. Recovery of lower-order time-dependent terms for the Minkowski metric has been studied in [Stefanov 1989; Ramm and Sjöstrand 1991; Ramm and Rakesh 1991; Waters 2014; Salazar 2013; Ben Aïcha 2015; Bellassoued and Ben Aïcha 2017], and for $-\mathrm{d} t^{2}+g_{i j}(x) \mathrm{d} x^{i} \mathrm{~d} x^{j}$ in [Kian et al. 2018]. Eskin [2017] proved that one can recover $g, A, q$ up to a gauge transformation, assuming existence of a global time variable $t$ and analyticity of all coefficients with respect to it. The proof is based on an adaptation of the boundary control methods and the analyticity is needed so that one can still use the
unique continuation results in [Tataru 1999]. Stability does not follow from such arguments. Other inverse problems on Lorentzian manifolds are studied in [Kurylev et al. 2014a; 2014b; Lassas et al. 2016]. The inverse scattering problem of recovering a moving boundary is studied in [Cooper and Strauss 1984; Stefanov 1991; Eskin and Ralston 2010]. The first author showed in [Stefanov 1989] that in the case where $g$ is Minkowski and $A=0$, the problem of recovery of $q$ reduces to the inversion of the X-ray transform in time-space over light rays, which was shown there to be injective for functions tempered in time and uniformly compactly supported in space. In [Lassas et al. 2017], it was shown that the linearized metric problem leads to the inversion of a light ray transform of tensor fields. Such light ray transforms are inherently unstable however because they are smoothing on the timelike cone. They require specialized tools for analyzing the singularities near the lightlike cone, not fully developed in the geodesic case; see [Greenleaf and Uhlmann 1989; 1990a; 1990b]. The light ray transform has been also studied in [Boman and Quinto 1993; Begmatov 2001; Stefanov 2017; Kian 2016].

We describe the main results below. Let $x_{0} \in \partial M$ and assume that $\partial M$ is timelike near $x_{0}$. Then $\partial M$ with the induced metric is a Lorentzian manifold as well and we choose (locally) one of the two time orientations, which we call future pointing.

Let $f \in \mathcal{E}^{\prime}(\partial M)$ be supported near $x_{0}$ with $\mathrm{WF}(f)$ close to a fixed timelike $\left(x_{0}, \xi^{0 \prime}\right) \in T^{*} \partial M \backslash 0$. We define the local outgoing solution operator $f \mapsto u$, defined up to a smoothing operator, as the operator mapping $f$ to the outgoing solution $u$ of

$$
\begin{equation*}
P u \in C^{\infty} \quad \text { in } M \text { near } x_{0},\left.\quad u\right|_{\partial M}=f \bmod C^{\infty} . \tag{3}
\end{equation*}
$$

The term "outgoing" here refers to the following. We chose that microlocal solution (parametrix) for which the singularities of the solution are required to propagate along future-pointing bicharacteristics. We refer to the subsection on page 1385 for more details. On the other hand, it is "local" because it solves (3) near $x_{0}$ only and this keeps the singularities close enough to $\partial M$ without allowing them to hit $\partial M$ again.

Define the associated local outgoing Dirichlet-to-Neumann map as

$$
\begin{equation*}
\Lambda_{g, A, q}^{\mathrm{loc}} f=\left.\left(\partial_{\nu} u-i\langle A, v\rangle u\right)\right|_{\partial M}, \tag{4}
\end{equation*}
$$

where $v$ denotes the unit outer normal vector field to $\partial M$, and the equality is modulo smoothing operators applied to $f$. By definition, the $\Lambda_{g, A, q}^{\mathrm{loc}}$ is defined near $x_{0}$ only, and in fact, in some conic neighborhood of the timelike $\left(x_{0}, \xi^{0 \prime}\right)$. Since the latter is arbitrary, $\Lambda_{g, A, q}^{\text {loc }}$ extends naturally to the whole timelike cone on $\partial M$ but we keep it microlocalized near $\left(x_{0}, \xi^{0 \prime}\right)$ to emphasize what we can recover given microlocal data only.

As we show in Theorem 3.1, $\Lambda_{g, A, q}^{\text {loc }}$ is actually a $\Psi D O$ (pseudodifferential operator) on the timelike cone bundle near $x_{0}$. The main result about $\Lambda_{g, A, q}^{\mathrm{loc}}$ is Theorem 3.2: a stability estimate about the recovery of the boundary jets of the coefficients.

Let $f \in \mathcal{E}^{\prime}(\partial M)$ have $\mathrm{WF}(f)$ as above. Let $u$, as in (3), be the parametrix in a neighborhood of the future-pointing null bicharacteristic issued from the unique future-pointing lightlike covector $\left(x_{0}, \xi^{0}\right) \in T^{*} M \backslash 0$ with orthogonal projection $\left(x_{0}, \xi^{0 \prime}\right)$. Note that the direction of ( $x_{0}, \xi^{0}$ ) and that of the bicharacteristic might be the same or opposite. Assume that this bicharacteristic hits $\partial M$ again,
transversely, at point $y_{0}$ in the codirection $\eta^{0}$ and let $\eta^{0 \prime}$ be the corresponding orthogonal tangential projection on $T_{y_{0}}^{*} \partial M$. Then $\left(y_{0}, \eta^{0 \prime}\right)$ is timelike, as well. Let $\mathcal{U}$ and $\mathcal{V}$ be two small conic timelike neighborhoods in $T^{*} \partial M \backslash 0$ of $\left(x_{0}, \xi^{0 \prime}\right)$ and ( $y_{0}, \eta^{0 \prime}$ ), respectively. If $\mathcal{U}$ is small enough, for every timelike $\left(x, \xi^{\prime}\right) \in \mathcal{U}$ close to $\left(x_{0}, \xi^{0 \prime}\right)$, we can define $\left(y, \eta^{\prime}\right)$ in the same way. This defines the lens relation

$$
\begin{equation*}
\mathcal{L}: \mathcal{U} \longrightarrow \mathcal{V}, \quad \mathcal{L}\left(x, \xi^{\prime}\right)=\left(y, \eta^{\prime}\right) \tag{5}
\end{equation*}
$$

see Figure 1. By definition, $\mathcal{L}$ is an even map in the second variable; i.e., $\mathcal{L}\left(x,-\xi^{\prime}\right)=\left(y,-\eta^{\prime}\right)$. If $\left(x, \xi^{\prime}\right)$ is future-pointing (i.e., if the associated vector by the metric is such), then ( $x,-\xi^{\prime}$ ) is past-pointing but we can interpret $\left(y,-\eta^{\prime}\right)$ as the end point of the null geodesic with initial point projecting to $\left(y,-\eta^{\prime}\right)$ but moving "backward" with respect to the parameter over it. This property correlates well with Theorem 4.2 since the wave equation has two wave "speeds" of opposite signs.

The map $\mathcal{L}$ is positively homogeneous of order 1 in its second variable. Now, for $f$ as above, let $u$ be an outgoing microlocal solution to (3) near (the projection on the base of) the bicharacteristic $\gamma_{0}$ issued from $\left(x_{0}, \xi^{0}\right)$ all the way to its second contact with $\partial M$ at $y_{0}$. In other words, $u$ is a distribution defined near $\gamma_{0}$ and solving (3) there rather than just near $x_{0}$, having future-propagating singularities only. It is unique up to a function smooth near $\gamma_{0}$; see Proposition 4.1. At this point, we assume that $\left(x_{0}, \xi^{0 \prime}\right)$ is not a fixed point for $\mathcal{L}$, which means that the reflected bicharacteristic does not become a periodic one after the first reflection. This solution is constructed as a solution with the boundary removed, in some neighborhood of $\gamma_{0}$. Since $f$ is smooth near $\left(y_{0}, \eta^{0 \prime}\right)$ that means there is no singularity of the solution $u$ at $\left(y_{0}, \eta^{0 \prime}\right)$; therefore, the singularity reflects at $y_{0}$. We extend the solution microlocally over a small segment of the reflected ray before reaching $\partial M$ again; see Proposition 4.1 for details. Then we define the global outgoing $D N$ map $\Lambda_{g, A, q}^{\mathrm{gl}}$ by (4) again but with the right-hand side localized to $V$, the projection of $\mathcal{V}$ to the base. In fact, by propagation of singularities, $\Lambda_{g, A, q}^{\mathrm{gl}} f$ has a wave-front set in $\mathcal{V}$ only and we can cut smoothly outside some neighborhood of $y_{0}$. The map $\Lambda_{g, A, q}^{\mathrm{gl}}$ is actually just semiglobal because it is the DN map restricted to a solution near one geodesic segment connecting boundary points. Also, it is only defined up to an operator smoothing near $\left(x_{0}, \xi^{0}\right)$. If a global initial boundary value problem is well defined, $\Lambda_{g, A, q}^{\mathrm{gl}}$ coincides with the associated DN operator up to a smoothing operator; see Section 5. In Theorem 4.2, we prove that $\Lambda_{g, A, q}^{\mathrm{gl}}$ is an FIO associated with the graph of $\mathcal{L}$. In Theorem 4.3, we show that $\Lambda_{g, A, q}^{\mathrm{gl}}$ recovers $\mathcal{L}$ in a stable way, which is also a general property of FIOs associated to a local canonical diffeomorphism.

Another fundamental object is the light ray transform $L$ which integrates functions, or more generally tensor fields, along lightlike geodesics. We define $L$ on functions by

$$
\begin{equation*}
L_{0} f(\gamma)=\int f(\gamma(s)) \mathrm{d} s \tag{6}
\end{equation*}
$$

and on covector fields of order 1 by

$$
\begin{equation*}
L_{1} f(\gamma)=\int\langle f(\gamma(s)), \dot{\gamma}(s)\rangle \mathrm{d} s \tag{7}
\end{equation*}
$$

where $\langle f(\gamma(s)), \dot{\gamma}(s)\rangle=f_{j}(\gamma(s)) \dot{\gamma}^{j}(s)$ in local coordinates and $\gamma$ runs over a given set of lightlike geodesics, and we always assume that supp $f$ is such that the integral is taken over a finite interval. In
our results below, $\gamma$ 's in $L_{0}$ and $L_{1}$ are the maximal geodesics through $M$ connecting boundary points. Unlike the Riemannian case, lightlike geodesics do not have a natural speed-1 parametrization and every rescaling of the parameter along them (even if that rescaling changes from geodesic to geodesic) keeps them being lightlike. The transform $L_{1}$ is invariant under reparametrization of the geodesics and can be considered as an integral of $\langle f, \mathrm{~d} \gamma\rangle$ over the geodesics. On the other hand, $L_{0}$ is not. Despite that freedom, the property $L_{0} f=0$ does not change. One way to parametrize it is to define it locally near a lightlike geodesic hitting a timelike surface at $s=0$, in our case, $\partial M$. Then the orthogonal projection $\dot{\gamma}^{\prime}(0)$ of each such $\gamma$ on $T \partial M$ (the prime stands for projection) determines $\dot{\gamma}(0)$, and therefore $\gamma$, uniquely. To normalize the projections on $T \partial M$, we can choose a timelike covector field $Z$ on $T \partial M$ locally and require $g(\dot{\gamma}, Z)=\mp 1$ for future-/past-pointing directions.

In Theorem 4.4, we show that given $g$, one can recover $L_{1} A$ in a Hölder-stable way; and if we are given $g$, $A$, one can recover $L_{0} q$ in a Hölder-stable way. Notice that we do not require absence of conjugate points and we do not use Gaussian beams. Instead, we use standard microlocal tools including Egorov's theorem. In Section 5, we consider some cases where $L_{1}$ and $L_{0}$ can be inverted to derive uniqueness results. As we mentioned above, those transforms are unstable. The reason is that they are microlocally smoothing in the spacelike cone; see, e.g., [Greenleaf and Uhlmann 1990b; Stefanov 2017; Lassas et al. 2017]. Therefore, stable recovery of $L_{1} A$ and $L_{0} q$ does not imply Hölder-stable recovery of $A_{1}$ (up to a gauge transform) and $q$ but allows for weaker logarithmic estimates using the estimate for recovery of $q$ from $L_{0} q$ in the Minkowski case proven in [Begmatov 2001], for example. We discuss some of those possible corollaries in Section 5. Recovery of $g$ from $\mathcal{L}$ is an open problem, with some results about the linearized problems obtained recently in [Lassas et al. 2017].

## 2. Preliminaries

Notation and terminology. In what follows, we denote by $U$ and $V$ the projections of $\mathcal{U}$ and $\mathcal{V}$ onto the base $\partial M$. We freely assume that $\mathcal{U}$ and $\mathcal{V}$, and therefore, $U$ and $V$ are small enough to satisfy the needed requirements below.

If $\xi$ is a covector based at a point $x$ on $\partial M$, we denote by $\xi^{\prime}$ its orthogonal projection to $T_{x}^{*} \partial M$. We routinely denote covectors on $T_{x}^{*} \partial M$ by placing primes, like $\xi^{\prime}$, etc., even if a priori such a covector is not a projection of a given one.

Timelike/spacelike/lightlike vectors $v$ are the ones satisfying $g(v, v)<0$, or $g(v, v)>0$, or $g(v, v)=0$, respectively. We identify vectors and covectors by the metric. We choose an orientation in $U$ that we call future pointing (FP). More precisely, we choose some smooth timelike vector $Z$ in $U$ (identified with an open set in the tangent bundle) and we call future pointing those timelike vectors $v$ for which $g(v, Z)>0$. If we have a time variable $t$, for example, such a choice could be $Z=\partial / \partial t$. In semigeodesic coordinates $\left(x^{0}=t, x\right)$ near a spacelike hypersurface, see (9) after Lemma 2.3, FP $v=\left(v^{0}, v^{\prime}\right)$ means $v^{0}>0$. Notice that for the associated covector $(\tau, \xi)=g v$, we have $\tau<0$.

Given a timelike $\left(x, \xi^{\prime}\right) \in \mathcal{U}$, assume first that $\xi^{\prime}$ is FP. Let $\xi$ be the lightlike covector pointing into $M$ with orthogonal projection $\xi^{\prime}$, identified with the vector $v=g^{-1} \xi$. The geodesic $\gamma_{x, \xi^{\prime}}(s)$ issued from $(x, v)$, for $s \geq 0$ will be called the FP geodesic issued from $\left(x, \xi^{\prime}\right)$. In Figure 1 , left, $v=v_{\text {int }}$ and


Figure 1. A tangent timelike future-pointing (FP) vector $v^{\prime}$ on the left, and a past-pointing on the right; and the two lightlike vectors $v_{\text {int }}$ and $v_{\text {ext }}$ with the same projection, pointing to $M$ and outside $M$, respectively. The FP geodesic $\gamma=\gamma_{x, \xi^{\prime}}(s)$ in both cases propagates to the future but on the right, it is determined by negative values of the parameter over it. The corresponding covectors $\xi^{\prime}, \xi_{\text {int }}$ and $\xi_{\text {ext }}$ are plotted, as well. The lens relation is $\mathcal{L}\left(x, \xi^{\prime}\right)=\left(y, \eta^{\prime}\right)$.
$\gamma_{x, \xi^{\prime}}(s)=\gamma$. If $\left(x, \xi^{\prime}\right)$ is past pointing, then we choose $v$ to be the lightlike vector projecting to $v^{\prime}$ pointing to the exterior ( $v_{\mathrm{ext}}$ in Figure 1, right) and take $\gamma_{x, \xi^{\prime}}(s)$ for $s \leq 0$. By propagation of singularities, a boundary singularity ( $x, \xi^{\prime}$ ) as above would propagate either along the FP geodesics chosen above, or along the past-pointing ones (or both) that we did not choose. The choice we made reflects the requirement that singularities should propagate to the future only. We call such microlocal solutions outgoing. We borrow that term from scattering theory. In the case of the classical formulation of the Riemannian version of this problem, this is guaranteed by the condition $u=0$ for $t<0$.

Gauge invariance. There exist some gauge transformations which leave the local and the global versions of the Dirichlet-to-Neumann map $\Lambda_{g, A, q}$ invariant; thus one can only expect to recover the corresponding gauge-equivalence class. To simplify the formulations, we assume that the DN map $\Lambda_{g, A, q}$ is well defined globally on $M$. In our main theorems, we will apply this to the $\Psi D O$ part of $\Lambda_{g, A, q}$ first, and then $\Phi$ below needs to be the identity near a fixed point only. For the semiglobal one, we need $\Phi$ to be identity near both ends of the fixed lightlike geodesic only. Since the computations below are purely algebraic, the lemmas remain true for the localized maps with obvious modifications.

We will consider two types of gauge transformations in this part. The first one is a diffeomorphism in $M$ which fixes $\partial M$.

Lemma 2.1. Let $(M, g)$ be a Lorentzian manifold with boundary as above, let A be a smooth 1-form and $q$ be a smooth function on $M$. If $\Phi: M \rightarrow M$ is a diffeomorphism with $\left.\Phi\right|_{\partial M}=\mathrm{Id}$, then

$$
\Lambda_{g, A, q}=\Lambda_{\Phi^{*} g, \Phi^{*} A, \Phi^{*} q}
$$

Here Id is the identity map, and $\Phi^{*} g, \Phi^{*} A, \Phi^{*} q$ are the pullbacks of $g, A, q$ under $\Phi$, respectively. Proof. For any $f \in C^{\infty}(\partial M)$, let $u$ be the solution of $\mathcal{L}_{g, A, q} u=0$ on $M$ with $\left.u\right|_{\partial M}=f$. Define $v:=\Phi^{*} u$ as the pull-back of $u$; then a simple calculation in local coordinates shows that $\mathcal{L}_{\Phi^{*} g, \Phi^{*} A, \Phi^{*} q} v=0$ and
$\left.v\right|_{\partial M}=f$. If we write $y=\Phi(x)$ as a local coordinate representation of $\Phi$, then

$$
\begin{aligned}
\Lambda_{g, A, q} f(y) & =v^{j}(y) \frac{\partial u}{\partial y^{j}}(y)-\left.i v^{j}(y) A_{j}(y) u(y)\right|_{\partial M} \\
& =v^{j}(x) \frac{\partial x^{l}}{\partial y^{j}} \frac{\partial v}{\partial y^{l}}-\left.i \frac{\partial x^{l}}{\partial y^{j}} v^{j}(x) \frac{\partial y^{k}}{\partial x^{l}} A_{k}(x) v(x)\right|_{\partial M} \\
& =\tilde{v}^{j}(x) \frac{\partial v}{\partial x^{j}}(x)-\left.i \tilde{v}^{j}(x)\left(\Phi^{*} A\right)_{j}(x) v(x)\right|_{\partial M}=\Lambda_{\Phi^{*} g, \Phi^{*} A, \Phi^{*} q} f
\end{aligned}
$$

where $v$ and $\tilde{v}$ are the unit normals in the $y$-and $x$-variables, respectively. The above calculation essentially verifies that $\Lambda_{g, A, q}$ is defined invariantly. Therefore, $\Lambda_{g, A, q}=\Lambda_{\Phi^{*} g, \Phi^{*} A, \Phi^{*} q}$.

Another type of gauge invariance occurs when one makes a conformal change of the metric $g$. This type of gauge invariance also occurs when $g$ is a Riemannian metric and $\Lambda_{g, A, q}$ is the corresponding Dirichlet-to-Neumann map for the magnetic Schrödinger equation; see [Dos Santos Ferreira et al. 2009, Proposition 8.2].

Lemma 2.2. Let $(M, g)$ be a Lorentzian manifold with boundary as above, let A be a smooth 1-form and $q$ be a smooth function on $M$. If $\varphi$ and $\psi$ are smooth functions such that

$$
\left.\varphi\right|_{\partial M}=\left.\partial_{\nu} \varphi\right|_{\partial M}=0,\left.\quad \psi\right|_{\partial M}=0
$$

then we have

$$
\Lambda_{g, A, q}=\Lambda_{e^{-2 \varphi} g, A-d \psi, e^{2 \varphi}\left(q-q_{\varphi}\right)}
$$

where $q_{\varphi}:=e^{\frac{n-2}{2} \varphi} \square_{g} e^{\frac{2-n}{2} \varphi}$.
Proof. A direct computation in local coordinates shows that

$$
\begin{aligned}
e^{\frac{n+2}{2} \varphi} P_{g, A, q}\left(e^{\frac{2-n}{2} \varphi} u\right) & =P_{e^{-2 \varphi}} g, A, e^{2 \varphi}\left(q-q_{\varphi}\right) \\
e^{-i \psi} P_{g, A, q}\left(e^{i \psi} u\right) & =P_{g, A-d \psi, q} u .
\end{aligned}
$$

For any $f \in C^{\infty}(\partial M)$, let $u$ be the solution of $P_{g, A, q} u=0$ on $M$ with $\left.u\right|_{\partial M}=f$. Setting $v:=$ $e^{\frac{n-2}{2} \varphi} e^{-i \psi} u$, we have

$$
\begin{aligned}
P_{e^{-2 \varphi} g, A-d \psi, e^{2 \varphi}\left(q-q_{\varphi}\right)} v & =P_{e^{-2 \varphi} g, A-d \psi, e^{2 \varphi}\left(q-q_{\varphi}\right)}\left(e^{\frac{n-2}{2} \varphi} e^{-i \psi} u\right) \\
& =e^{\frac{n+2}{2} \varphi} P_{g, A-d \psi, q}\left(e^{-i \psi} u\right)=e^{\frac{n+2}{2} \varphi} e^{-i \psi} P_{g, A, q} u=0
\end{aligned}
$$

Furthermore, notice that $v_{e^{-2 \varphi}}=v_{g}$ by the assumption on $\varphi$; thus

$$
\begin{aligned}
\Lambda_{e^{-2 \varphi} g, A-d \psi, e^{2 \varphi}\left(q-q_{\varphi}\right)} f & =v^{j} \frac{\partial v}{\partial x^{j}}-\left.i v^{j}\left(A_{j}-\frac{\partial \psi}{\partial x^{j}}\right) v\right|_{\partial M} \\
& =v^{j} \frac{\partial\left(e^{\frac{n-2}{2} \varphi} e^{-i \psi} u\right)}{\partial x^{j}}-\left.i v^{j}\left(A_{j}-\frac{\partial \psi}{\partial x^{j}}\right)\left(e^{\frac{n-2}{2} \varphi} e^{-i \psi} u\right)\right|_{\partial M} \\
& =v^{j}\left(-i \frac{\partial \psi}{\partial x^{j}} u+\frac{\partial u}{\partial x^{j}}\right)-\left.i v^{j} i v^{j}\left(A_{j}-\frac{\partial \psi}{\partial x^{j}}\right) u\right|_{\partial M} \\
& =v^{j} \frac{\partial u}{\partial x^{j}}-\left.i v^{j} A_{j} u\right|_{\partial M}=\Lambda_{g, A, q} f
\end{aligned}
$$

Gauge equivalent modifications of $\boldsymbol{g}, \boldsymbol{A}, \boldsymbol{q}$. It is convenient to work in semigeodesic normal coordinates on a Lorentzian manifold. These coordinates are the Lorentzian counterparts of the well-known Riemannian semigeodesic coordinates for Riemannian manifolds with boundary. We formulate the existence of such coordinates in the following lemma.
Lemma 2.3. Let $S$ be a timelike hypersurface in $M$. For every $x_{0} \in S$, there exist $\varepsilon>0$, a neighborhood $N$ of $x_{0}$ in $M$, and a diffeomorphism $\Psi: S \cap N \times[0, T) \rightarrow N$ such that
(i) $\Psi\left(x^{\prime}, 0\right)=x^{\prime}$ for all $x^{\prime} \in S \cap N$;
(ii) $\Psi\left(x^{\prime}, x^{n}\right)=\gamma_{x^{\prime}}\left(x^{n}\right)$, where $\gamma_{x^{\prime}}\left(x^{n}\right)$ is the unit speed geodesic issued from $x^{\prime}$ normal to $S$.

Moreover, if $\left(x^{0}, \ldots, x^{n-1}\right)$ are local boundary coordinates on $S$, in the coordinate system $\left(x^{0}, \ldots, x^{n}\right)$, the metric tensor $g$ takes the form

$$
\begin{equation*}
g=g_{\alpha \beta} d x^{\alpha} \otimes d x^{\beta}+d x^{n} \otimes d x^{n}, \quad \alpha, \beta \leq n-1 \tag{8}
\end{equation*}
$$

Clearly, $g_{\alpha \beta}$ has a Lorentzian signature as well. If $M$ has a boundary, then $S$ can be $\partial M$ and $x^{n}$ is restricted to $[0, \varepsilon]$. A proof of the lemma can be found in [Petrov 1969] and is based on the fact that the lines $x^{\prime}=$ const. and $x^{n}=s$ are unit speed geodesics; therefore the Christoffel symbols $\Gamma_{n n}^{i}$ vanish for all $i$. We will call such coordinates the semigeodesic normal coordinates. The lemma remains true if $S$ is spacelike with a negative sign in front of $d x^{n} \otimes d x^{n}$ in (8) (we replace the index $n$ by 0 below), and this gives us a way to define a time function $t=x^{0}$ locally, and put the metric in the block form

$$
\begin{equation*}
g=-d t^{2}+g_{i j}(t, x) d x^{i} \otimes d x^{j}, \quad 1 \leq i, j \leq n \tag{9}
\end{equation*}
$$

with $g_{i j}$ Riemannian.
Now we use the gauge invariance of $\Lambda_{g, A, q}$ to alter $g, A, q$ without changing the DN map. Three types of modifications are made in the following, labeled as (M1)-(M3) respectively.

Firstly, given two metrics $g$ and $\tilde{g}$, one can choose diffeomorphisms as in Lemma 2.1 to obtain common semigeodesic normal coordinates. In fact, let $\Psi$ and $\widetilde{\Psi}$ be diffeomorphisms like in Lemma 2.3 with respect to $g$ and $\tilde{g}$ respectively; then $\tilde{\Psi} \circ \Psi^{-1}$ is a diffeomorphism near $\partial M$ which fixes $\partial M$. Extend $\tilde{\Psi} \circ \Psi^{-1}$ as in [Palais 1960] to be a global diffeomorphism on $M$. The properties of $\Psi$ and $\tilde{\Psi}$ ensure that the two metrics $g$ and $\left(\tilde{\Psi} \circ \Psi^{-1}\right)^{*} \tilde{g}$ have common semigeodesic normal coordinates near $\partial M$. Therefore, we may assume: (M1) If ( $x^{\prime}, x^{n}$ ) are the semigeodesic normal coordinates for $g$, they are also the semigeodesic normal coordinates for $\tilde{g}$.
Secondly, we employ the conformal gauge invariance to replace $\tilde{g}$ with a gauge-equivalent one to obtain some identities which later will help simplify the calculations.
Lemma 2.4. Let $S$ be either a timelike or a spacelike hyperplane near some point $p_{0} \in S$. Given smooth functions $r_{2}, r_{3}, \ldots$ on $S$ near $p_{0}$, there exists a smooth function $\mu$ near $p_{0}$ with $\mu=0, \partial_{\nu} \mu=0$ on $S$ so that if $\hat{\Psi}$ is the diffeomorphism in Lemma 2.3 related to the metric $\hat{g}:=e^{\mu} g$, then

$$
\partial_{n}^{j} \operatorname{det}\left(\hat{\Psi}^{*} \hat{g}\right)=r_{j}, \quad j=2,3, \ldots
$$

on $S$ near $p_{0}$. Here $\partial_{n}=\partial / \partial x^{n}$ with $\left(x^{0}, \ldots, x^{n}\right)$ the semigeodesic normal coordinates for $g$.

Before giving the proof of the lemma, we remark that $\left(x^{0}, \ldots, x^{n}\right)$ may not be the semigeodesic normal coordinates for $\hat{g}$.

Proof. The statement of the theorem is invariant under replacing $g$ by $\Psi^{*} g$ for any local diffeomorphism $\Phi$ which preserves the boundary pointwise. Therefore, we may assume that $g$ is replaced by $\Psi^{*} g$, i.e., that $x=\left(x^{\prime}, x^{n}\right)$ are semigeodesic coordinates for $g$.

Note first that the conformal factor does not change the property of a covector being normal to $S$ but rescales the normal derivative and may change the higher-order ones because $\gamma_{x^{\prime}}$ may change its curvature with respect to the old metric. More precisely, for the vector $e_{n}=(0, \ldots, 0,1)$ we have $g\left(e_{n}, e_{n}\right)=\mp 1$ but $\hat{g}\left(e_{n}, e_{n}\right)=\mp e^{\mu}$. Therefore, for the corresponding normal derivatives we have $\hat{\partial}_{v}=e^{-\frac{\mu}{2}} \partial_{\nu}=\partial_{n}$ on $x^{n}=0$. Let $\hat{\gamma}_{x^{\prime}}(s)$ be the normal geodesic at $x^{\prime} \in S$ with $\dot{\hat{\gamma}}_{x^{\prime}}$ consistent with the orientation of $S$, normalized by $\hat{g}\left(\dot{\hat{\gamma}}_{x^{\prime}}(s), \dot{\hat{\gamma}}_{x^{\prime}}(s)\right)=\mp 1$. Then for every smooth function $f$,

$$
\left.\partial_{n}^{j} \hat{\Psi}^{*} f\left(x^{\prime}\right)\right|_{x^{n}=0}=\left.\partial_{n}^{j}\right|_{x^{n}=0} f\left(\hat{\gamma}_{x^{\prime}}\left(x^{n}\right)\right)
$$

For $j=0,1$, the results are not affected by the conformal factor and we get

$$
\left.\widehat{\Psi}^{*} f\left(x^{\prime}\right)\right|_{x^{n}=0}=f\left(x^{\prime}, 0\right),\left.\quad \partial_{n} \hat{\Psi}^{*} f\left(x^{\prime}\right)\right|_{x^{n}=0}=f_{n}\left(x^{\prime}, 0\right)
$$

To compute the higher-order normal derivatives, we write

$$
\begin{equation*}
\partial_{n}^{2} \hat{\Psi}^{*} f\left(x^{\prime}\right)=f_{i j} \dot{\hat{\gamma}}_{x^{\prime}}^{i} \dot{\hat{\gamma}}_{x^{\prime}}^{j}+f_{i} \ddot{\hat{\gamma}}_{x^{\prime}}^{i} \quad \text { on } x^{n}=0 \tag{10}
\end{equation*}
$$

Under the conformal change of the metric, the Christoffel symbols are transformed by the law

$$
\hat{\Gamma}_{j k}^{k}=\Gamma_{i j}^{k}+\frac{1}{2} \delta_{i}^{k} \partial_{j} \mu+\frac{1}{2} \delta_{j}^{k} \partial_{i} \mu-g_{i j} \nabla^{k} \mu
$$

In particular,

$$
\begin{equation*}
\hat{\Gamma}_{n n}^{k}=\Gamma_{n n}^{k}+\frac{1}{2} \delta_{n}^{k} \partial_{n} \mu+\frac{1}{2} \delta_{n}^{k} \partial_{n} \mu-g_{n n} \nabla^{k} \mu=\delta_{n}^{k} \partial_{n} \mu-\frac{1}{2} g^{k l} \partial_{l} \mu . \tag{11}
\end{equation*}
$$

Therefore, $\hat{\Gamma}_{n n}^{k}=0$ on $x^{n}=0$ and (10) reduces to

$$
\begin{equation*}
\partial_{n}^{2} \hat{\Psi}^{*} f\left(x^{\prime}\right)=f_{n n} \quad \text { on } x^{n}=0 \tag{12}
\end{equation*}
$$

In a similar way, we may compute $\partial_{n}^{j} \hat{\Psi}^{*} f\left(x^{\prime}\right)$ on $x^{n}=0$. The result is $\partial_{n}^{j} f$ plus normal derivatives of $f$ of order $j-1$ and less, with coefficients depending on the normal derivatives of $\mu$ up to order $j-1$. For our purposes, the exact expression does not matter.

The metric $\hat{g}$ has the form

$$
\left(\hat{\Psi}^{*} \hat{g}\right)_{k l}=\left(\hat{g}_{i j} \circ \hat{\Psi}\right) \frac{\partial \hat{\Psi}^{i}}{\partial x^{k}} \frac{\partial \hat{\Psi}^{j}}{\partial x^{l}}=\left(\hat{g}_{\alpha \beta} \circ \hat{\Psi}\right) \frac{\partial \hat{\Psi}^{\alpha}}{\partial x^{k}} \frac{\partial \hat{\Psi}^{\beta}}{\partial x^{l}}+\frac{\partial \hat{\Psi}^{n}}{\partial x^{k}} \frac{\partial \hat{\Psi}^{n}}{\partial x^{l}}
$$

where the Greek indices range from 0 to $n-1$ (but not $n$ ). In particular,

$$
\begin{equation*}
\operatorname{det} \hat{\Psi}^{*} \hat{g}=(\operatorname{det} d \hat{\Psi})^{2} \operatorname{det}(\hat{g} \circ \hat{\Psi}) \tag{13}
\end{equation*}
$$

We need to understand the structure of $\left.\partial_{n}^{k}(\operatorname{det} \mathrm{~d} \hat{\Psi})\right|_{x^{n}=0}$ now. For $k=0$, we have $\left.\mathrm{d} \widehat{\Psi}\right|_{x^{n}=0}=\mathrm{Id}$. Notice next that

$$
\begin{equation*}
\mathrm{d} \hat{\Psi}=\left(\partial_{0} \hat{\Psi}, \ldots, \partial_{n-1} \hat{\Psi}, \partial_{n} \hat{\Psi}\right) \tag{14}
\end{equation*}
$$

where each partial derivative is a vector. Since by (11), $\partial_{n}^{2} \widehat{\Psi}^{i}=-\widehat{\Gamma}_{n n}^{i}=0$ for $x^{n}=0$,

$$
\left.\partial_{n}(\operatorname{det} d \hat{\Psi})\right|_{x^{n}=0}=0
$$

To analyze $k=2$, we notice first that

$$
\partial_{n}^{3} \hat{\Psi}^{i}=-\partial_{n} \widehat{\Gamma}_{n n}^{i}=-\partial_{n}\left(\delta_{n}^{i} \partial_{n} \mu-\frac{1}{2} g^{i l} \partial_{l} \mu\right)=-\delta_{n}^{i} \mu_{n n}+\cdots
$$

where the dots represent a term involving lower-order $\partial_{n}$-derivatives of $\mu$. Using this in (14), we get

$$
\left.\partial_{n}^{2}(\operatorname{det} \mathrm{~d} \hat{\Psi})\right|_{x^{n}=0}=-\left.\mu_{n n}\right|_{x^{n}=0}
$$

Reasoning as above, we see that

$$
\begin{equation*}
\left.\partial_{n}^{j}(\operatorname{det} \mathrm{~d} \hat{\Psi})\right|_{x^{n}=0}=-\left.\partial_{n}^{j} \mu\right|_{x^{n}=0}+\cdots \tag{15}
\end{equation*}
$$

where the dots represent terms involving normal derivatives of $\mu$ (possibly differentiated tangentially) up to order $j-1$.

We will analyze the normal derivatives of $\operatorname{det}(\hat{g} \circ \hat{\Psi})$ in (13) now. Since $\operatorname{det} \hat{g}=e^{n \mu} \operatorname{det} g$, we get

$$
\begin{align*}
\partial_{n} \operatorname{det}(\hat{g} \circ \hat{\Psi}) & =\partial_{n}\left(e^{(n+1) \mu \circ \hat{\Psi}} \operatorname{det} g \circ \hat{\Psi}\right) \\
& =(n+1) \mu_{n} \operatorname{det} g+\partial_{n} \operatorname{det} g=\partial_{n} \operatorname{det} g \quad \text { on } \partial M \tag{16}
\end{align*}
$$

We used the fact that $\mathrm{d} \Psi=\mathrm{Id}$ on $\partial M$ and that $\partial_{n} \mathrm{~d} \Psi=0$ since $\mathrm{d} \mu=0$ on $\partial M$. Therefore, $\partial_{n}^{j} \operatorname{det} \hat{g} \circ \hat{\Psi}=$ $\partial_{n}^{j} \operatorname{det} g$ on $x^{n}=0$ for $j=0,1$.

For the highest-order derivatives, notice that $\partial_{n}^{j} \hat{\Psi}$ involves $\partial_{n}^{j-1} \mu$ as its highest-order normal $\mu$ derivative, as the arguments leading to (15) show. Differentiating (16), we therefore get

$$
\begin{align*}
\partial_{n}^{j} \operatorname{det}(\hat{g} \circ \hat{\Psi}) & =\partial_{n}^{j}\left(e^{(n+1) \mu \circ \hat{\Psi}} \operatorname{det} g \circ \hat{\Psi}\right) \\
& =(n+1)\left(\partial_{n}^{j} \mu\right) \operatorname{det} g+\cdots \quad \text { on } \partial M \tag{17}
\end{align*}
$$

where the dots have the same meaning as in (15).
Using (13) in combination with (15) and (17), we get

$$
\begin{equation*}
\left.\partial_{n}^{j}\left(\operatorname{det} \hat{\Psi}^{*} \hat{g}\right)\right|_{x^{n}=0}=(n-1)\left(\partial_{n}^{j} \mu\right) \operatorname{det} g+\cdots \tag{18}
\end{equation*}
$$

To complete the proof of the lemma, we determine the normal derivatives of $\mu$ on $x^{n}=0$ for $j=2, \ldots$. We get first $\left.\partial_{n}^{2}\left(\operatorname{det} \widehat{\Psi}^{*} \hat{g}\right)\right|_{x^{n}=0}=\left.(n-1) \mu_{n n}\right|_{x^{n}=0}$, which needs to be equal to $r_{2}$, and can be solved for $\mu_{n n}$. Then we can determine the tangential derivatives of the latter. After that, we can solve (17) with $j=3$ for $\mu_{n n n}$, etc. To complete the proof, we use Borel's lemma.

Let $g$ and $\tilde{g}$ be two metrics satisfying (M1) with the two diffeomorphisms $\Psi$ and $\widetilde{\Psi}$ respectively, as in Lemma 2.3. Applying Lemma 2.4 to $S=\partial M$ and $p=x_{0}$, we can find a metric $\hat{g}:=e^{\mu} g$ with $\mu=0$,
$\partial_{\nu} \mu=0$ on $\partial M$ such that under the semigeodesic normal coordinates $\left(x^{0}, \ldots, x^{n}\right)$ for $g$ we have

$$
\partial_{n}^{j} \operatorname{det}\left(\hat{\Psi}^{*} \hat{g}\right)=\partial_{n}^{j} \operatorname{det}\left(\tilde{\Psi}^{*} \tilde{g}\right), \quad j=2,3, \ldots
$$

on $\partial M$. Notice that $\left(x^{0}, \ldots, x^{n}\right)$ are also semigeodesic normal coordinates for $\tilde{g}$ by (M1).
Now consider the metrics $\left(\hat{\Psi} \circ \tilde{\Psi}^{-1}\right)^{*} \hat{g}$ and $\tilde{g}$. These metrics have common semigeodesic normal coordinates (see the argument following Lemma 2.3), which are $\left(x^{0}, \ldots, x^{n}\right)$. In these coordinates the choice of $\hat{g}$ yields

$$
\partial_{n}^{j} \operatorname{det}\left(\tilde{\Psi}^{*} \circ\left(\hat{\Psi} \circ \tilde{\Psi}^{-1}\right)^{*} \hat{g}\right)=\partial_{n}^{j} \operatorname{det}\left(\hat{\Psi}^{*} \hat{g}\right)=\partial_{n}^{j} \operatorname{det}\left(\tilde{\Psi}^{*} \tilde{g}\right)
$$

Thus we may replace $g$ by $\left(\hat{\Psi} \circ \tilde{\Psi}^{-1}\right)^{*} \hat{g}$ and change $A, q$ accordingly, as in Lemma 2.1 and Lemma 2.2, without affecting $\Lambda_{g, A, q}$. We therefore can assume that $g$ and $\tilde{g}$ satisfy not only (M1), but also:
(M2) In the common semigeodesic normal coordinates $\left(x^{\prime}, x^{n}\right)$,

$$
\partial_{n}^{j} \operatorname{det} g\left(x^{\prime}, 0\right)=\partial_{n}^{j} \operatorname{det} \tilde{g}\left(x^{\prime}, 0\right), \quad j=2,3, \ldots
$$

Here we have identified the metrics with their coordinate representations under $\tilde{\Psi}$.
Thirdly, we make modifications to the 1-form $A$. Again the modification does not change the gaugeequivalence class of $\Lambda_{g, A, q}$ due to Lemma 2.2.
Lemma 2.5. Let $(M, g)$ be a Lorentzian manifold with boundary as above, let $A$ be a smooth 1-form and $q$ be a smooth function on $M$. There exists a smooth function $\psi$ with $\left.\psi\right|_{\partial M}=0$ such that in the semigeodesic normal coordinates $\left(x^{\prime}, x^{n}\right), B:=A-d \psi$ satisfy

$$
\begin{equation*}
\partial_{n}^{j} B_{n}\left(x^{\prime}, 0\right)=0, \quad j=0,1,2, \ldots \tag{19}
\end{equation*}
$$

Proof. We can find a smooth function $\psi$ with

$$
\psi\left(x^{\prime}, 0\right)=0, \quad \partial_{n}^{j+1} \psi\left(x^{\prime}, 0\right)=\partial_{n}^{j} A_{n}\left(x^{\prime}, 0\right), \quad j=0,1,2, \ldots
$$

Extend it in a suitable manner so that $\psi \in C^{\infty}(M)$ with $\left.\psi\right|_{\partial M}=0$. Then $B=A-d \psi$ satisfies (19).
As a result we may further assume:
(M3) In the common semigeodesic normal coordinates $\left(x^{\prime}, x^{n}\right)$ of $g$ and $\tilde{g}$,

$$
\partial_{n}^{j} A_{n}\left(x^{\prime}, 0\right)=\partial_{n}^{j} \tilde{A}_{n}\left(x^{\prime}, 0\right)=0, \quad j=0,1,2, \ldots
$$

## 3. Boundary stability

We choose the semigeodesic coordinates $\left(x^{\prime}, x^{n}\right)$ near $x_{0}$ so that $x_{0}=0, \partial M$ locally is given by $x^{n}=0$, and the interior of $M$ is given by $x^{n}>0$. Let $\xi^{0 \prime}$ be a future-pointing timelike covector in $T_{x_{0}}^{*} \partial M$ at $x_{0}$. On Figure 1, the associated vector would look like $v^{\prime}$ on the left, while the covector $\xi^{0 \prime}$ would have the opposite time direction, like the figure on the right. Let $\chi\left(x^{\prime}, \xi^{\prime}\right)$ be a smooth cutoff function with small enough support in $\mathcal{U}$ that is equal to 1 in a smaller conic timelike neighborhood of $\left(x_{0}, \xi^{0 \prime}\right)$. Assume also that $\chi$ is homogeneous in $\xi^{\prime}$ of order 0 .

For

$$
\begin{equation*}
f\left(x^{\prime}\right)=e^{i \lambda x^{\prime} \cdot \xi^{\prime}} \chi\left(x^{\prime}, \xi^{\prime}\right) \tag{20}
\end{equation*}
$$

and for every $N>0$, we would like to construct a geometric optics approximation of the outgoing solution $u$ of (3) near $x_{0}$ in $M$ of the form

$$
\begin{equation*}
u_{N}(x):=e^{i \lambda \phi\left(x, \xi^{\prime}\right)} \sum_{j=0}^{N} \frac{1}{\lambda^{j}} a_{j}\left(x, \xi^{\prime}\right) \tag{21}
\end{equation*}
$$

The eikonal and the transport equations below are based on the identity

$$
e^{-i \lambda \phi} P e^{i \lambda \phi}=-\lambda^{2} g^{j k}\left(\partial_{j} \phi\right)\left(\partial_{k} \phi\right)+i \lambda \square_{g} \phi+2 i \lambda g^{j k} \partial_{j} \phi\left(\partial_{k}-i A_{k}\right)+P .
$$

In $M$ near $x_{0}$, the phase function $\phi\left(x, \xi^{\prime}\right)$ solves the eikonal equation, which in the semigeodesic coordinates takes the form

$$
\begin{equation*}
g^{\alpha \beta} \partial_{\alpha} \phi \partial_{\beta} \phi+\left(\partial_{n} \phi\right)^{2}=0,\left.\quad \phi\right|_{x^{n}=0}=x^{\prime} \cdot \xi^{\prime} \tag{22}
\end{equation*}
$$

With the extra condition $\left.\partial_{\nu} \phi\right|_{\partial M}<0$, (22) is locally uniquely solvable. Moreover, (22) implies

$$
\begin{equation*}
\partial_{n} \phi\left(x^{\prime}, 0\right)=\xi_{n}\left(x^{\prime}, \xi^{\prime}\right)>0 \quad \text { for any }\left(x^{\prime}, \xi^{\prime}\right) \in \mathcal{U} \tag{23}
\end{equation*}
$$

where

$$
\begin{equation*}
\xi_{n}\left(x^{\prime}, \xi^{\prime}\right):=\sqrt{-g^{\alpha \beta}\left(x^{\prime}, 0\right) \xi_{\alpha} \xi_{\beta}} \tag{24}
\end{equation*}
$$

Notice that the choice of the sign of $\xi_{n}$ makes $\xi$ a lightlike future-pointing covector, pointing into $M$. In Figure 1, the associated vector $v=g^{-1} \xi$ looks like $v_{\text {int }}$ on the left.

We recall briefly the method of characteristics for solving the eikonal equation. We first determine $\partial \phi$ on $x^{n}=0$ to get (23) or the same equation with a negative square root. We choose one of them, and in this case our choice is determined by the requirement that $\partial \phi$ points into $M$; see Figure 1. Let now $\left(q_{x^{\prime}, \xi^{\prime}}(s), p_{x^{\prime}, \xi^{\prime}}(s)\right)$ be the null bicharacteristic with $q_{x^{\prime}, \xi^{\prime}}(0)=x^{\prime}, p_{x^{\prime}, \xi^{\prime}}(0)=\left(\xi^{\prime}, \xi_{n}\right)$. We think of $\left(x^{\prime}, s\right)$ as local coordinates and set $\phi\left(x^{\prime}, s\right)=x^{\prime} \cdot \xi^{\prime}$. More precisely, $\phi$ is uniquely determined locally by the requirement to be constant along the null bicharacteristics $q_{x^{\prime}, \xi^{\prime}}$. Moreover,

$$
\begin{equation*}
p(s)=\nabla_{x} \phi\left(q(s), \xi^{\prime}\right) \tag{25}
\end{equation*}
$$

Since by the Hamilton equations, $\dot{q}^{i}(s)=g^{i j} p_{j}(s)$, we get in particular that $g^{i j} \partial_{j} \phi \partial_{i}$ is just the derivative $\partial / \partial s$ along the null bicharacteristic.

In $M$ near $x_{0}$, the amplitudes $a_{0}$ and $a_{j}, j=1,2, \ldots$, solve the following transport equations:

$$
\begin{array}{rlrl}
T a_{0} & =0, & \left.a_{0}\right|_{x_{n}=0}=\chi \\
i T a_{j} & =-P a_{j-1}, & \left.a_{j}\right|_{x_{n}=0}=0, & j \geq 1 \tag{27}
\end{array}
$$

where the operator $T$ is defined as

$$
\begin{equation*}
T:=2 g^{j k} \partial_{j} \phi\left(\partial_{k}-i A_{k}\right)+\square_{g} \phi . \tag{28}
\end{equation*}
$$

We prefer to express the bicharacteristics through the geodesics

$$
\Gamma(s):=\left(q_{x^{\prime}, \xi^{\prime}}(s), p_{x^{\prime}, \xi^{\prime}}(s)\right)=\left(\gamma_{x^{\prime}, \xi^{\prime}}(s), g \dot{\gamma}_{x^{\prime}, \xi^{\prime}}(s)\right)
$$

Then along the bicharacteristics, we have

$$
\begin{equation*}
T=2 \partial_{s}-2 i\langle A, p(s)\rangle+\square_{g} \phi=2 \mu \partial_{s} \mu^{-1} \tag{29}
\end{equation*}
$$

with the integrating factor $\mu$ given by

$$
\begin{equation*}
\mu(\Gamma(s))=\exp \left\{-\frac{1}{2} \int_{0}^{s}\left(\square_{g} \phi\right)\left(\gamma_{x^{\prime}, \xi^{\prime}}(\sigma)\right) \mathrm{d} \sigma\right\} \exp \left\{i \int_{0}^{s}\left\langle A\left(\gamma_{x^{\prime}, \xi^{\prime}}(\sigma)\right), \dot{\gamma}_{x^{\prime}, \xi^{\prime}}(\sigma)\right\rangle \mathrm{d} \sigma\right\} \tag{30}
\end{equation*}
$$

The amplitudes $a_{j}, j=0,1, \ldots$, are supported in a neighborhood of the characteristics issued from $x_{0} \in \partial M$ in the codirection $\xi\left(x_{0}\right)$. As a result, on some neighborhood of $x_{0}$, we have $u_{N}$ solves $P u_{N}=O\left(\lambda^{-N}\right),\left.u\right|_{\partial M}=f$.

Theorem 3.1. $\Lambda_{g, A, q}^{\mathrm{loc}}$ is an elliptic $\Psi D O$ of order 1 in $\mathcal{U}$.
Proof. Given $f \in \mathcal{E}^{\prime}(U)$, not related to (20), with a wave-front set as in the theorem, we are looking for an outgoing solution $u$ of $P u=0$ near $x_{0}, u=f$ on $U$, of the form

$$
\begin{equation*}
u(x)=(2 \pi)^{-n} \int e^{i \phi\left(x, \xi^{\prime}\right)} a\left(x, \xi^{\prime}\right) \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime} \tag{31}
\end{equation*}
$$

The phase $\phi$ solves the eikonal equation (22) and therefore coincides with $\phi$ there. We chose the solution which guarantees a locally outgoing $u$, which corresponds to the positive square root in (24). We are looking for an amplitude $a$ of the form $a \sim \sum_{j=0}^{\infty} a_{j}\left(x, \xi^{\prime}\right)$, where $a_{j}$ is homogeneous in the $\xi^{\prime}$-variable of degree $-j$. The standard geometric optics construction leads to the transport equations (26), (27). Using the standard Borel lemma argument, we construct a convergent series for $a$. Then $u$ is the microlocal solution (up to a microlocally smoothing operator applied to $f$ ) that we used to define $\Lambda_{g, A, q}^{\mathrm{loc}}$. Then $\Lambda_{g, A, q}^{\mathrm{loc}} f=\partial u /\left.\partial \nu\right|_{U}$. Since $\phi=x^{\prime} \cdot \xi^{\prime}$ on $U$, we get that $\Lambda_{g, A, q}^{\mathrm{loc}}$ is a $\Psi \mathrm{DO}$ with symbol

$$
-i \xi_{n}\left(x^{\prime}, \xi^{\prime}\right)-\left.\partial_{n} a\right|_{x^{n}=0}
$$

In particular, for the principal symbol we get

$$
\begin{equation*}
\sigma_{p}\left(\Lambda_{g, A, q}^{\mathrm{loc}}\right)\left(x^{\prime}, \xi^{\prime}\right)=-i \xi_{n}=-i \sqrt{-g^{\alpha \beta}\left(x^{\prime}\right) \xi_{\alpha} \xi_{\beta}} \tag{32}
\end{equation*}
$$

We proceed in the same way if $\left(x^{\prime}, \xi^{\prime}\right)$ is past pointing.
It remains to show that if we use another locally outgoing solution $\tilde{u}$, the resulting $\widetilde{\Lambda}_{g, A, q}^{\text {loc }}$ would differ by a smoothing operator. This follows by considering $v:=u-\tilde{u}$ which is a locally outgoing solution with smooth boundary data, which therefore must be smooth. We omit the details.

We prove a stable determination result on the boundary next. Let $(g, A, q)$ and $(\tilde{g}, \tilde{A}, \tilde{q})$ be two triples. Define

$$
\begin{equation*}
\delta=\left\|\Lambda_{g, A, q}^{\mathrm{loc}}-\Lambda_{\tilde{g}, \tilde{A}, \tilde{q}}^{\mathrm{loc}}\right\|_{H^{1}(U) \rightarrow L^{2}(U)} \tag{33}
\end{equation*}
$$

where, as above, $\Lambda_{g, A, q}^{\text {loc }}$ and $\Lambda_{\tilde{g}, \tilde{A}, \tilde{q}}^{\text {loc }}$ are the local DN maps associated with $(g, A, q)$ and $(\tilde{g}, \tilde{A}, \tilde{q})$ respectively microlocally restricted to a fixed conic neighborhood $\mathcal{U}$ of a timelike future-pointing $\left(x_{0}, \xi^{0 \prime}\right) \in T^{*} U$ with $x_{0} \in U \subset \partial M$. As above, we assume that $\xi^{0 \prime}$ is future pointing and timelike for both $g$ and $\tilde{g}$, and that $\mathcal{U}$ is small enough so that it is included in the future timelike cone on $T^{*} U$ for both metrics. Therefore, in the theorem below, we need to know the DN map microlocally only near a fixed timelike covector on $T^{*} \partial M$.

Theorem 3.2. Let $(g, A, q)$ and $(\tilde{g}, \tilde{A}, \tilde{q})$ be replaced by their gauge equivalent triples satisfying (M1)(M3). Then for any $\mu<1$ and $m \geq 0$, and some open neighborhood $U_{0} \Subset U$ of $x_{0}$,
(1) $\sup _{x \in \bar{U}_{0},|\gamma| \leq m}\left|\partial^{\gamma}(g-\tilde{g})\right| \leq C \delta^{\frac{\mu}{2^{m}}}$;
(2) $\sup _{x \in \bar{U}_{0},|\gamma| \leq m}\left|\partial^{\gamma}(A-\tilde{A})\right| \leq C \delta^{\frac{\mu}{2^{m+1}}}$;
(3) $\sup _{x \in \bar{U}_{0},|\gamma| \leq m}\left|\partial^{\gamma}(q-\tilde{q})\right| \leq C \delta^{\frac{\mu}{2^{m+2}}}$
are valid whenever $g, \tilde{g}, A, \tilde{A}, q, \tilde{q}$ are bounded in a certain $C^{k}$ norm in the semigeodesic normal coordinates near $x_{0}$ with a constant $C>0$ depending on that bound with $k=k(m, \mu)$.

Proof. We adapt the proofs in [Montalto 2014; Stefanov and Uhlmann 2005b] in the Riemannian setting. Let $\Gamma_{0}$ be a small conic neighborhood of $\xi^{0 \prime}$. We can assume that $\chi=1$ on $U_{0} \times \Gamma_{0}$. Let $f$ be as in (20). We restrict ( $x^{\prime}, \xi^{\prime}$ ) to $U_{0} \times \Gamma_{0}$ below. In addition, we normalize $\xi^{\prime}$ to have unit Euclidean length (in that coordinate system). Since $\partial_{\nu}=-\partial_{n}$, the formal Dirichlet-to-Neumann map in the boundary normal coordinates $\left(x^{\prime}, x^{n}\right)$ is given by

$$
\begin{equation*}
\Lambda_{g, A, q}^{\mathrm{loc}} f\left(x^{\prime}\right)=-e^{i \lambda x^{\prime} \cdot \xi^{\prime}}\left(i \lambda \partial_{n} \phi\left(x^{\prime}, 0, \xi^{\prime}\right)+\sum_{j=0}^{N} \frac{1}{\lambda^{j}}\left(\partial_{n}-i A_{n}\right) a_{j}\left(x^{\prime}, 0, \xi\right)\right)+O\left(\lambda^{-N-1}\right) \tag{34}
\end{equation*}
$$

The expression for $\Lambda_{\tilde{g}, \tilde{A}, \tilde{q}} f$ is similar, with $\phi$ and $a_{j}$ replaced by $\tilde{\phi}$ and $\tilde{a}_{j}$, respectively.
The representation (34) could be derived from (21) but since $u$ there is an approximate solution only, and we defined $\Lambda_{g, A, q}^{\text {loc }}$ microlocally, we need to go back to its definition. To justify (34), notice that by [Taylor 1981, Chapter VIII.7], on the set $\chi=1$, we have $e^{-i \lambda x^{\prime} \cdot \xi^{\prime}} \Lambda_{g, A, q}^{\text {loc }} f$ is equal to the full symbol of $\Lambda_{g, A, q}^{\mathrm{loc}}$ with $\lambda=|\xi|$ and $\xi$ in (34) unit.

In the following, $C$ denotes various constants depending only on $M, \chi$ in (20), on the choice of $k \gg 1$ and on the a priori bounds of the coefficients of $P$ in $C^{k}$. Solving for $\partial_{n} \phi$ (resp. $\partial_{n} \tilde{\phi}$ ) in (34) and taking the difference we obtain

$$
\partial_{n} \phi-\partial_{n} \tilde{\phi}=\frac{1}{i \lambda}\left(\Lambda_{g, A, q}^{\mathrm{loc}} f-\Lambda_{\tilde{g}, \tilde{A}, \tilde{q}}^{\mathrm{loc}} f\right)+\frac{1}{i \lambda} \sum_{j=0}^{N} \frac{1}{\lambda^{j}}\left[\left(\partial_{n} a_{j}-\partial_{n} \tilde{a}_{j}\right)-i\left(A_{n} a_{j}-\tilde{A}_{n} \tilde{a}_{j}\right)\right]+O\left(\lambda^{-N-1}\right)
$$

in $L^{2}\left(U_{0}\right)$. Integrating in $U_{0}$ yields

$$
\begin{equation*}
\left\|\partial_{n} \phi-\partial_{n} \tilde{\phi}\right\|_{L^{2}\left(U_{0}\right)} \leq \frac{C}{\lambda} \delta\|f\|_{H^{1}\left(U_{0}\right)}+\frac{C}{\lambda} . \tag{35}
\end{equation*}
$$

The choice of $f$ in (20) indicates that $\|f\|_{H^{1}\left(U_{0}\right)} \leq C \lambda$. Thus, taking the limit $\lambda \rightarrow \infty$ yields

$$
\begin{equation*}
\left\|\xi_{n}-\tilde{\xi}_{n}\right\|_{L^{2}\left(U_{0}\right)}=\left\|\partial_{n} \phi-\partial_{n} \tilde{\phi}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta \tag{36}
\end{equation*}
$$

From relation (24) we have

$$
\begin{equation*}
\left\|\left(g^{\alpha \beta}-\tilde{g}^{\alpha \beta}\right) \xi_{\alpha} \xi_{\beta}\right\|_{L^{2}\left(U_{0}\right)}=\left\|\xi_{n}^{2}-\tilde{\xi}_{n}^{2}\right\|_{L^{2}\left(U_{0}\right)}=\left\|\left(\partial_{n} \phi\right)^{2}-\left(\partial_{n} \tilde{\phi}\right)^{2}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta \tag{37}
\end{equation*}
$$

We use the following argument here and in several places below: a quadratic form $h^{\alpha \beta} \xi_{\alpha} \xi_{\beta}$ is uniquely determined for $\xi^{\prime}$ in any fixed-in-advance open set $\Gamma$ on the unit sphere. In fact, one can choose $n(n-1) / 2$ vectors $\xi^{\prime}$ in $\Gamma$ and then the recovery is done by inverting an isomorphism on $\mathbb{R}^{\frac{n(n-1)}{2}}$, and is therefore stable; see [Dairbekov et al. 2007, Lemma 3.3]. Therefore, (37) implies $\|g-\tilde{g}\|_{L^{2}\left(U_{0}\right)} \leq C \delta$. By interpolation estimates in Sobolev space and Sobolev embedding theorems, we have for any $m \geq 0$ and $\mu<1$ that

$$
\begin{equation*}
\|g-\tilde{g}\|_{C^{m}\left(\bar{U}_{0}\right)} \leq C \delta^{\mu} \tag{38}
\end{equation*}
$$

provided $k \gg 1$ is sufficiently large.
Second, we show that the first-order normal derivatives of $g$ and the 1 -form can be stably determined on the boundary. From (34) we have

$$
\begin{aligned}
\left(\partial_{n}-i \tilde{A}_{n}\right) \tilde{a}_{0}-\left(\partial_{n}-i A_{n}\right) a_{0}= & e^{-i \lambda x^{\prime} \cdot \xi^{\prime}}\left(\Lambda_{g, A, q} f-\Lambda_{\tilde{g}, \tilde{A}, \tilde{q}} f\right) \\
& +i \lambda\left(\partial_{n} \phi-\partial_{n} \tilde{\phi}\right)+\sum_{j=1}^{N} \frac{1}{\lambda^{j}}\left(\partial_{n} a_{j}-\partial_{n} \tilde{a}_{j}\right)+O\left(\frac{1}{\lambda^{N+1}}\right) \quad \text { in } L^{2}\left(U_{0}\right) .
\end{aligned}
$$

Estimate as in (35) to obtain

$$
\left\|\left(\partial_{n}-i A_{n}\right) a_{0}-\left(\partial_{n}-i \tilde{A}_{n}\right) \tilde{a}_{0}\right\|_{L^{2}\left(U_{0}\right)} \leq C\left(\delta+\lambda \delta+\frac{1}{\lambda}\right)
$$

which holds for all $\lambda>0$. In particular, we may choose $\lambda=\delta^{-\frac{1}{2}}$ to minimize the right-hand side; then

$$
\begin{equation*}
\left\|\left(\partial_{n}-i A_{n}\right) a_{0}-\left(\partial_{n}-i \tilde{A}_{n}\right) \tilde{a}_{0}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{2}} \tag{39}
\end{equation*}
$$

In order to estimate the difference of first-order normal derivatives of the metrics, we consider the transport equation in (26). Since $\chi \equiv 1$ for $x \in U_{0}$, it follows from the boundary condition in (26) that $\partial_{\alpha} a_{0}=\partial_{\alpha} \chi=0$ for $\alpha=0, \ldots, n-1$. Moreover, $g^{n j}=\delta^{n j}$ in the semigeodesic coordinates; thus the transport equation in (26) becomes

$$
\begin{equation*}
2 \xi_{n}\left(\partial_{n}-i A_{n}\right) a_{0}-2 i A^{\alpha} \xi_{\alpha}+\frac{1}{\sqrt{-\operatorname{det} g}} \partial_{n}\left(\sqrt{-\operatorname{det} g} \partial_{n} \phi\right)+Q(g)=0 \tag{40}
\end{equation*}
$$

where, as before, Greek indices range from 0 to $n-1$ (but not $n$ ). Here $A^{\alpha}:=g^{\alpha \beta} A_{\beta}$, and $Q(g)$ is a linear combination of tangential derivatives of $g$, which is defined as follows:

$$
Q(g):=\frac{1}{\sqrt{-\operatorname{det} g}} \partial_{\alpha}\left(\sqrt{-\operatorname{det} g} g^{\alpha \beta}\right) \xi_{\beta}
$$

where we have used that $\partial_{\beta} \phi=\xi_{\beta}$ in $U_{0}, \beta=0, \ldots, n-1$. As a consequence of (38),

$$
\begin{equation*}
Q(g)-Q(\tilde{g})=O\left(\delta^{\frac{1}{2}}\right) \tag{41}
\end{equation*}
$$

Therefore, combining (39), (40) and (41) we obtain

$$
\frac{1}{\sqrt{-\operatorname{det} g}} \partial_{n}\left(\sqrt{-\operatorname{det} g} \partial_{n} \phi\right)-\frac{1}{\sqrt{-\operatorname{det} \tilde{g}}} \partial_{n}\left(\sqrt{-\operatorname{det} \tilde{g}} \partial_{n} \tilde{\phi}\right)-2 i\left(A^{\alpha}-\widetilde{A}^{\alpha}\right) \xi_{\alpha}=O\left(\delta^{\frac{1}{2}}\right)
$$

Notice that

$$
\frac{1}{\sqrt{-\operatorname{det} g}} \partial_{n}\left(\sqrt{-\operatorname{det} g} \partial_{n} \phi\right)=\frac{1}{2 \operatorname{det} g} \partial_{n} \operatorname{det} g \partial_{n} \phi+\partial_{n}^{2} \phi=\frac{\xi_{n}}{2 \operatorname{det} g} \partial_{n} \operatorname{det} g-\frac{1}{2 \xi_{n}} \partial_{n} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}
$$

is an even function of $\xi^{\prime}$. Here in the computation, $\partial_{n} \phi$ is substituted by $\xi_{n}$ due to (23) and $\partial_{n}^{2} \phi$ is calculated by differentiating the eikonal equation (22). Separating the even and odd parts in $\xi^{\prime}$ we conclude

$$
\begin{gather*}
\left(\frac{\xi_{n}}{2 \operatorname{det} g} \partial_{n} \operatorname{det} g-\frac{1}{2 \xi_{n}} \partial_{n} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}\right)-\left(\frac{\tilde{\xi}_{n}}{2 \operatorname{det} \tilde{g}} \partial_{n} \operatorname{det} \tilde{g}-\frac{1}{2 \tilde{\xi}_{n}} \partial_{n} \tilde{g}^{\alpha \beta} \xi_{\alpha} \xi_{\beta}\right)=O\left(\delta^{\frac{1}{2}}\right)  \tag{42}\\
\left(A^{\alpha}-\tilde{A}^{\alpha}\right) \xi_{\alpha}=O\left(\delta^{\frac{1}{2}}\right) \tag{43}
\end{gather*}
$$

From the odd part (43), varying $\xi^{\prime}$ locally, we get

$$
\begin{equation*}
\|A-\tilde{A}\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{2}} \tag{44}
\end{equation*}
$$

To deal with the even part, notice (42) states that

$$
\frac{\xi_{n}}{2 \operatorname{det} g} \partial_{n} \operatorname{det} g-\frac{1}{2 \xi_{n}} \partial_{n} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}
$$

is stably determined of order $O\left(\delta^{\frac{1}{2}}\right)$. As $\xi_{n}$ is stably determined on $U_{0}$, see (36), their product

$$
\begin{aligned}
\frac{\xi_{n}^{2}}{2 \operatorname{det} g} \partial_{n} \operatorname{det} g-\frac{1}{2} \partial_{n} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta} & =-\frac{1}{2 \operatorname{det} g}\left(\partial_{n} \operatorname{det} g\right) g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}-\frac{1}{2} \partial_{n} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta} \\
& =-\frac{1}{2} \frac{1}{\operatorname{det} g} \partial_{n}\left(\operatorname{det} g \cdot g^{\alpha \beta}\right) \xi_{\alpha} \xi_{\beta}
\end{aligned}
$$

is also stably determined. Since det $g$ is known to be stable and away from zero, it follows that $\partial_{n} h^{\alpha \beta}$ is stable where $h^{\alpha \beta}:=(\operatorname{det} g) g^{\alpha \beta}$. Hence, the normal derivative of $g=(\operatorname{det} h)^{\frac{1}{1-n}} h$ is also stably determined; that is,

$$
\begin{equation*}
\left\|\partial_{n} g-\partial_{n} \tilde{g}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{2}} \tag{45}
\end{equation*}
$$

Using interpolation and Sobolev embedding theorems, we obtain from (45) and (44) that for any $m \geq 0$ and $\mu<1$,

$$
\begin{equation*}
\left\|\partial_{n} g-\partial_{n} \tilde{g}\right\|_{C^{m}\left(\bar{U}_{0}\right)}+\|A-\widetilde{A}\|_{C^{m}\left(\bar{U}_{0}\right)} \leq C \delta^{\frac{\mu}{2}} \tag{46}
\end{equation*}
$$

provided $k \gg 1$ is sufficiently large.

Next we show that the second-order normal derivatives of $g$, the first-order normal derivatives of $A$, and the values of $q$ can be stably determined on the boundary. By (34) up to $\lambda^{-1}$ we obtain

$$
\left\|\left(\partial_{n}-i A_{n}\right) a_{1}-\left(\partial_{n}-i \tilde{A}_{n}\right) \tilde{a}_{1}\right\|_{L^{2}\left(U_{0}\right)} \leq C\left(\lambda^{2} \delta+\lambda \delta^{\frac{1}{2}}+\lambda^{-1}\right)
$$

Choose $\lambda=\delta^{-\frac{1}{4}}$ to minimize the right-hand side. Then

$$
\begin{equation*}
\left\|\left(\partial_{n}-i A_{n}\right) a_{1}-\left(\partial_{n}-i \tilde{A}_{n}\right) \tilde{a}_{1}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{4}} \tag{47}
\end{equation*}
$$

Consider the transport equation (27) for $a_{1}$. In the semigeodesic coordinates this equation takes the form

$$
\begin{equation*}
2 i \xi_{n}\left(\partial_{n}-i A_{n}\right) a_{1}=-\partial_{n}^{2} a_{0}+q+O\left(\delta^{\frac{1}{2}}\right) \tag{48}
\end{equation*}
$$

where $O\left(\delta^{\frac{1}{2}}\right)$ represents the stably determined terms of order $O\left(\delta^{\frac{1}{2}}\right.$. (In fact, $a_{1}=0$ in these expressions by the boundary condition in (27), but it is left here for the convenience of tracking the corresponding terms.) From the estimates (36), (45) and (48) it follows that

$$
\begin{equation*}
\left(-\partial_{n}^{2} a_{0}+\partial_{n}^{2} \tilde{a}_{0}\right)+(q-\tilde{q})=O\left(\delta^{\frac{1}{4}}\right) \tag{49}
\end{equation*}
$$

To obtain an expression of $\partial_{n}^{2} a_{0}$, we differentiate the transport equation in (26) and evaluate it on $U_{0}$ :

$$
\begin{aligned}
\partial_{n}^{2} a_{0} & =-\frac{1}{4 \operatorname{det} g} \partial_{n}^{2} \operatorname{det} g-\frac{1}{2 \xi_{n}} \partial_{n}^{3} \phi+\frac{i}{\xi_{n}} g^{\alpha \beta} \partial_{n} A_{\alpha} \xi_{\beta}+O\left(\delta^{\frac{1}{2}}\right) \\
& =-\frac{1}{4 \operatorname{det} g} \partial_{n}^{2} \operatorname{det} g+\frac{1}{4 \xi_{n}^{2}} \partial_{n}^{2} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}+\frac{i}{\xi_{n}} g^{\alpha \beta} \partial_{n} A_{\alpha} \xi_{\beta}+O\left(\delta^{\frac{1}{2}}\right)
\end{aligned}
$$

where the $O\left(\delta^{\frac{1}{2}}\right)$ terms are estimated by (38) and (46) and we have used that $\partial_{n} A_{n}\left(x^{\prime}, 0\right)=0$ in (M3). Inserting this into (49) and separating the even and odd parts in $\xi^{\prime}$ gives (notice that $\xi_{n}=\sqrt{-g^{\alpha \beta}} \xi_{\alpha} \xi_{\beta}$ is an even function of $\xi^{\prime}$ ):

$$
\begin{gather*}
\left(\frac{1}{4 \operatorname{det} g} \partial_{n}^{2} \operatorname{det} g-\frac{1}{4 \operatorname{det} \tilde{g}} \partial_{n}^{2} \operatorname{det} \tilde{g}-\frac{1}{4 \xi_{n}^{2}} \partial_{n}^{2} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}+\frac{1}{4 \tilde{\xi}_{n}^{2}} \partial_{n}^{2} \tilde{g}^{\alpha \beta} \xi_{\alpha} \xi_{\beta}\right)+(q-\tilde{q})=O\left(\delta^{\frac{1}{4}}\right)  \tag{50}\\
-\frac{i}{\xi_{n}} g^{\alpha \beta} \partial_{n} A_{\alpha} \xi_{\beta}+\frac{i}{\tilde{\xi}_{n}} \tilde{g}^{\alpha \beta} \partial_{n} \tilde{A}_{\alpha} \xi_{\beta}=O\left(\delta^{\frac{1}{4}}\right) \tag{51}
\end{gather*}
$$

To deal with (51), we multiply the two terms by $\xi_{n}$ and $\tilde{\xi}_{n}$ respectively. This is valid since $\xi_{n}$ is stably determined in (36). By the argument following (37),

$$
\left\|\partial_{n} A_{\alpha}-\partial_{n} \tilde{A}_{\alpha}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{2}}
$$

To deal with (50), recall the following matrix identity which is valid for any invertible matrix $S$ :

$$
\partial \log |\operatorname{det} S|=\operatorname{tr}\left(S^{-1} \partial S\right)
$$

Taking $S=g^{\alpha \beta}$ and applying $\partial_{n}^{j-1}$ we see that

$$
\partial_{n}^{j} \log \left(-\operatorname{det} g^{\alpha \beta}\right)=\partial_{n}^{j-1}\left(g_{\alpha \beta} \partial_{n} g^{\alpha \beta}\right), \quad j=1,2, \ldots
$$

For $j=2$, it gives

$$
g_{\alpha \beta} \partial_{n}^{2} g^{\alpha \beta}=\partial_{n}^{2} \log \left(-\operatorname{det} g^{\alpha \beta}\right)-\partial_{n} g_{\alpha \beta} \partial_{n} g^{\alpha \beta}
$$

The right-hand side is stably determined by (M2) and (45); we thus get on $U_{0}$ that

$$
\begin{equation*}
g_{\alpha \beta} \partial_{n}^{2} g^{\alpha \beta}-\tilde{g}_{\alpha \beta} \partial_{n}^{2} \tilde{g}^{\alpha \beta}=O\left(\delta^{\frac{1}{2}}\right) \tag{52}
\end{equation*}
$$

On the other hand, remember that the two metrics $g$ and $\tilde{g}$ have been modified to satisfy (M2); thus by (37)

$$
\frac{1}{4 \operatorname{det} g} \partial_{n}^{2} \operatorname{det} g-\frac{1}{4 \operatorname{det} \tilde{g}} \partial_{n}^{2} \operatorname{det} \tilde{g}=\left(\frac{1}{4 \operatorname{det} g}-\frac{1}{4 \operatorname{det} \tilde{g}}\right) \partial_{n}^{2} \operatorname{det} g=O(\delta)
$$

This together with (50) gives

$$
\begin{equation*}
\left(-\frac{1}{4 \xi_{n}^{2}} \partial_{n}^{2} g^{\alpha \beta} \xi_{\alpha} \xi_{\beta}+\frac{1}{4 \tilde{\xi}_{n}^{2}} \partial_{n}^{2} \tilde{g}^{\alpha \beta} \xi_{\alpha} \xi_{\beta}\right)+(q-\tilde{q})=O\left(\delta^{\frac{1}{4}}\right) \tag{53}
\end{equation*}
$$

Again we multiply the terms without the tilde by $\xi_{n}^{2}$ and those with it by $\tilde{\xi}_{n}^{2}$; using (24) we have

$$
\left(\partial_{n}^{2} g^{\alpha \beta}+4 q g^{\alpha \beta}-\partial_{n}^{2} \tilde{g}^{\alpha \beta}-4 \tilde{q} \tilde{g}^{\alpha \beta}\right) \xi_{\alpha} \xi_{\beta}=O\left(\delta^{\frac{1}{4}}\right)
$$

By the argument following (37),

$$
\left(\partial_{n}^{2} g^{\alpha \beta}+4 q g^{\alpha \beta}\right)-\left(\partial_{n}^{2} \tilde{g}^{\alpha \beta}+4 \tilde{q} \tilde{g}^{\alpha \beta}\right)=O\left(\delta^{\frac{1}{4}}\right)
$$

Multiplying those terms without the tilde by $g_{\alpha \beta}$, those with the tilde by $\tilde{g}_{\alpha \beta}$, and then summing up in $\alpha, \beta$ yields

$$
\left(g_{\alpha \beta} \partial_{n}^{2} g^{\alpha \beta}+4 n q\right)-\left(\tilde{g}_{\alpha \beta} \partial_{n}^{2} \tilde{g}^{\alpha \beta}+4 n \tilde{q}\right)=O\left(\delta^{\frac{1}{4}}\right)
$$

From (52) we come to the conclusion that

$$
\|q-\tilde{q}\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{4}}
$$

Inserting this into (53) and using the argument following (37),

$$
\left\|\partial_{n}^{2} g^{\alpha \beta}-\partial_{n}^{2} \tilde{g}^{\alpha \beta}\right\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{4}}
$$

Putting the estimates on $g, A, q$ together, we have established

$$
\left\|\partial_{n}^{2} g-\partial_{n}^{2} \tilde{g}\right\|_{L^{2}\left(U_{0}\right)}+\left\|\partial_{n} A_{\alpha}-\partial_{n} \tilde{A}_{\alpha}\right\|_{L^{2}\left(U_{0}\right)}+\|q-\tilde{q}\|_{L^{2}\left(U_{0}\right)} \leq C \delta^{\frac{1}{4}}
$$

As before, interpolation and the Sobolev embedding theorem lead to

$$
\left\|\partial_{n}^{2} g-\partial_{n}^{2} \tilde{g}\right\|_{C^{m}\left(\bar{U}_{0}\right)}+\left\|\partial_{n} A_{\alpha}-\partial_{n} \tilde{A}_{\alpha}\right\|_{C^{m}\left(\bar{U}_{0}\right)}+\|q-\tilde{q}\|_{C^{m}\left(\bar{U}_{0}\right)} \leq C \delta^{\frac{\mu}{4}}
$$

for $m>0$ and $\mu<1$. Repeating this type of argument will establish the stability for higher-order derivatives of $g, A, q$ on $U_{0}$.


Figure 2. The solution $u$.

## 4. Interior stability

The semiglobal microlocal solution. We construct the semiglobal microlocal solution $u$ sketched in the Introduction in the paragraph following (5) and used to define $\Lambda_{g, A, q}^{\mathrm{gl}}$. We recall the assumptions. We fix a timelike $\left(x_{0}, \xi^{0 \prime}\right) \in T^{*} \partial M \backslash 0$ and a small conic neighborhood $\mathcal{U}$ of it. We choose a local orientation so that $\left(x_{0}, \xi^{0 \prime}\right)$ is future pointing. Then there is a unique lightlike $\left(x_{0}, \xi^{0}\right) \in T^{*} M \backslash 0$ which projects orthogonally to $\left(x_{0}, \xi^{0 \prime}\right)$. Let $\gamma_{0}$ be the zero bicharacteristic issued from ( $x_{0}, \xi^{0 \prime}$ ) extended until it hits $T^{*} \partial M$ again, transversally, by assumption, at some $\left(y_{0}, \eta^{0}\right)$ with projection $\left(y_{0}, \eta^{0 \prime}\right)=\mathcal{L}\left(x_{0}, \xi^{0 \prime}\right) \in T^{*} \partial M$. Let $\mathcal{V}=\mathcal{L}(\mathcal{U})$ and denote by $U$ and $V$ the projections $\pi(\mathcal{U}), \pi(\mathcal{V})$ of $\mathcal{U}$ and $\mathcal{V}$ onto the base, i.e., their " $x$-parts". Denote by $\Gamma$ the union of all zero bicharacteristics issued "from $\mathcal{U}$ ", i.e., from all futurepointing ( $x, \xi$ ) with $x \in \partial M$ which have projections on the boundary in $\mathcal{U}$. Let $\Gamma_{0}:=\pi(\Gamma) \subset M$ be the projection of $\Gamma$ onto the base; see Figure 2. We assume below, for convenience, that $(M, g)$ is embedded in a slightly larger manifold.

Next proposition says that the microlocal solution $u$ used in the Introduction to define $\Lambda_{g, A, q}^{\mathrm{gl}}$ is well defined.

Proposition 4.1. If $\mathcal{U}$ is small enough, then for every $f \in \mathcal{E}^{\prime}(\partial M)$ with $\operatorname{WF}(f) \in \mathcal{U}$ there exists a distribution $u$ defined in a neighborhood of $\bar{\Gamma}_{0}$ so that $P u \in C^{\infty}\left(\Gamma_{0}\right),\left.u\right|_{U}-f \in C^{\infty}(U)$ and $\left.u\right|_{V} \in C^{\infty}$. Moreover, $u$ is unique up to a smooth function in $\Gamma_{0}$.

Proof. We are looking for a solution $u^{\text {inc }}$ of the form (31) with $f$ having a wave-front set in $\mathcal{U}$. Pastpointing codirections can be handled the same way. The solution is the same as in Theorem 3.1 but we are now trying to extend it as far as possible away from $\partial M$. We know that microlocally, $u^{\mathrm{inc}}$ is supported in a small neighborhood of the null bicharacteristic (projecting to a null geodesic on $M$ ) issued from $\left(x_{0}, \xi^{0}\right)$ with $\xi^{0}$ future pointing with a projection $\xi^{0^{\prime}}$ on the boundary; i.e., $\xi^{0}=\left(\xi^{0^{\prime}}, \xi_{n}\left(x_{0}, \xi^{0}\right)\right.$ ), where $\xi_{n}$ is given by (24). See Figure 1. This follows from the general propagation of singularities theory but in this particular case it can be derived from the fact that $T$ in (28) has its principal part a vector field along such null geodesics, and $\mathrm{WF}\left(u^{\mathrm{inc}}\right)$ can be analyzed directly with the aid of (31).

Such a solution is guaranteed to exist only near some neighborhood of $x_{0}$ because the eikonal equation may not be globally solvable. On the other hand, the solution is still a global FIO applied to the boundary
data $f$. Indeed, it can also be viewed as a superposition of a finite number of local FIOs, each one having a representation of the kind (31). We construct $u^{\text {inc }}$ first near $\partial M$; call it $u_{1}$. Then we restrict it to a timelike hypersurface $S_{1}$ intersecting the null geodesic $\pi\left(\gamma_{0}\right)$ transversely and we chose $S_{1}$ so that the geometric optics construction is still valid along $\pi\left(\gamma_{0}\right)$ until it hits $S_{1}$, and a bit beyond it. We take the boundary data at $S_{1}$, and solve a new similar problem, by taking the outgoing solution (the future-pointing cone on $S_{1}$ is the one determined by $\left.\mathrm{WF}\left(u_{1} \mid S_{1}\right)\right)$, etc. By compactness arguments, we can cover the whole null geodesic (the projection of $\gamma_{0}$ to the base) until it hits $\partial M$ again. This construction provides solutions (modulo smooth terms) $u_{1}, \ldots, u_{k}$, each one defined in an open set $\Gamma_{k}$, where $\bigcup_{k} \Gamma_{k}$ covers $\bar{\Gamma}_{0}$. Without loss of generality we may assume that the only intersections of the $\Gamma_{k}$ 's happen among consecutive ones. Then on $S_{k}$, near the intersection with $\pi\left(\gamma_{0}\right)$, we have two microlocal solutions: $u_{k}$ and $u_{k+1}$. They have the same traces on $S_{k}$ modulo a smooth function.

Next, in their common domain of definition, $u_{k}$ and $u_{k+1}$ coincide up to a smooth function. Indeed, the difference $v$ has smooth trace on $S_{k}$ and it is outgoing. By the last paragraph of the proof of Theorem 3.1, $v$ is smooth near $S_{k}$.

We choose a partition of unity $1=\sum_{k} \chi_{k}$ near $\bar{\Gamma}_{0}$ subordinate to that cover and set $u^{\mathrm{inc}}=\sum_{k} \chi_{k} u_{k}$. The latter is a microlocal solution (i.e., a solution up to smooth errors) in a neighborhood of $\bar{\Gamma}_{0}$. Indeed, this is not completely obvious only when supp $\chi_{k}$ and supp $\chi_{k+1}$ intersect but then $u_{k+1}=u_{k}$ modulo $C^{\infty}$ and therefore, near such a point, $u^{\text {inc }}=\chi_{k} u_{k}+\chi_{k+1} u_{k+1}=u_{k}$ modulo $C^{\infty}$, which is a microlocal solution.

We use this argument several times below. This construction is similar to that in [Duistermaat 1996], where it is shown that the Cauchy problem on a spacelike surface gives rise to a global FIO. As a result, one gets a microlocal solution $u^{\text {inc }}$ in a neighborhood of $\bar{\Gamma}_{0}$ (not satisfying the needed boundary conditions on $V$ yet) as a composition of a finite number of FIOs.

We need to reflect $u^{\text {inc }}$ at $V$ to satisfy the zero boundary condition. We write the solution $u$ as the sum of the incident wave $u^{\mathrm{inc}}$ and the reflected wave $u^{\mathrm{ref}}: u=u^{\mathrm{inc}}+u^{\text {ref }}$. The construction of $u^{\text {ref }}$ is similar we start with boundary data $-\left.u^{\mathrm{inc}}\right|_{V}$ on $V$ and singularities which propagate into $M$ into the future (the past-future orientation near $V$ is determined by declaring the singularities of $u^{\text {inc }}$ on $\mathcal{V}$ coming from the past). We refer to (57) below and the construction following it for more details. The solution $u^{\text {ref }}$ needs to be extended to a small neighborhood of the geodesics near $\gamma_{0}$ reflected at $V$ until they leave $\Gamma_{0}$. By choosing $\mathcal{U}$ small enough, we guarantee that the reflected geodesics do not hit $\partial M$ again.

Finally, we prove the uniqueness statement. If $u_{1}$ and $u_{2}$ are two such solutions, then $v:=u_{1}-u_{2}$ is smooth on both $U$ and $V$. A priori, $v$ can be only singular along bicharacteristics close to $\gamma_{0}$ or its reflection from $V$. By the argument we used above, $v$ must be smooth in $\Gamma_{0}$ with the possible exception of some neighborhood of $V$ in $M$, where $u^{\text {ref }}$ might be nontrivial. Near $V$, we know $v$ has smooth Cauchy data. An (easier) adaptation of the same argument shows that $v$ has to be smooth near $V$ as well. Indeed, otherwise, for $v$, extended as zero outside $M$, we would get that $P v$ has singularities conormal to $V$ only, and the microlocal propagation of singularities theorem then would yield that $v$ has no singularities near $\gamma_{0}$ or its reflection.

Having constructed $u$, then we define $\Lambda_{g, A, q}^{\mathrm{gl}}$ as in (4) but with the so-constructed $u$. The uniqueness part of the proposition shows that $\Lambda_{g, A, q}^{\mathrm{gl}}$ is defined up to a smoothing operator.
$\Lambda_{g, A, q}^{\mathrm{gl}}$ recovers the lens relation $\mathcal{L}$ in a stable way.
Theorem 4.2. Under the assumptions in the Introduction, $\Lambda_{g, A, q}^{\mathrm{gl}}$ is an elliptic FIO of order 1 associated with the (canonical) graph of $\mathcal{L}$.

Note that we excluded lightlike covectors in $\mathrm{WF}(f)$. This excludes bicharacteristics (geodesics) tangent to $\partial M$ carrying singularities of $u$. This is where the two Lagrangians (one of them being the diagonal) intersect. We also restricted $u$ to the first reflection shortly after that. Without that, the canonical relations would contain powers of $\mathcal{L}$. The theorem is a direct consequence of the geometric optics construction and propagation of singularities results for the wave equation and can be considered as essentially known.

As a consequence of Theorem 4.2, for every $s$, we have $\Lambda_{g, A, q}^{\text {loc }}$ maps $H^{s}(U)$ into $H^{s-1}(U)$ and $\Lambda_{g, A, q}^{\mathrm{gl}}$ maps $H^{s}(U)$ into $H^{s-1}(V)$. Fixing $s=1$, one may conclude that the natural norms for those two operators are the $H^{1} \rightarrow L^{2}$ ones. While both operators are bounded in those norms, their dependence on the metric $g$ is not necessarily continuous if we stay in those norms. For $\Lambda_{g, A, q}^{\text {loc }}$, we will see that the principal symbol (and the whole one, in fact) depends continuously on $g$; and in fact the whole operator does, as well. On the other hand, while the canonical relation of $\Lambda_{g, A, q}^{\mathrm{gl}}$ depends continuously on $g$, the operator itself does not. This observation was used in [Bao and Zhang 2014]; see also [Stefanov et al. 2016] for a discussion.

Proof of Theorem 4.2. We will analyze first the map $F:\left.f \mapsto u^{\mathrm{inc}}\right|_{S}$, where $S$ is a timelike surface as in the proof of Proposition 4.1, and (31) for $u=u_{\mathrm{inc}}$ is valid all the way to it, and a bit beyond it.

Change the coordinates $x$ so that $S=\left\{x^{n}=1\right\}$. This can be done if $S$ is close enough to $\partial M$. Then (31) with $x=\left(x^{\prime}, 1\right)$ is a local representation of the FIO $F$ and its canonical relation is given by (see, e.g., [Taylor 1981, Chapter VIII])

$$
\left(\left.\nabla_{\xi^{\prime}} \phi\right|_{x^{n}=1}, \xi^{\prime}\right) \mapsto\left(x^{\prime},\left.\nabla_{x^{\prime}} \phi\right|_{x^{n}=1}\right)
$$

By (25), with the momentum $p$ projected to $T^{*}\left\{x^{n}=1\right\}$, we get that this is the lens relation $\mathcal{L}_{1}$ from $\mathcal{U} \subset T^{*} \partial M$ to $T^{*} S$ (instead of the image being on $T^{*} \partial M$ again).

We can repeat this finitely many times by choosing $S_{1}, S_{2}$, etc., to get a composition of finitely many canonical relations, starting with $\mathcal{L}_{1}$; then $\mathcal{L}_{2}$ maps data on $T^{*} S_{1}$ to $T^{*} S_{2}$, etc. That composition of, say $m$ of them, gives the lens relation from $\partial M$ to $S_{m}$. In the final step, we need to take the normal derivative. This shows that the map $\left.f \mapsto \partial_{\nu} u^{\mathrm{inc}}\right|_{V}$ is an FIO of the claimed type.

To prove this for $\Lambda_{g, A, q}^{\mathrm{gl}}$, we need to add $\left.\partial_{\nu} u^{\text {ref }}\right|_{V}$. The latter has an oscillatory representation of the same kind with a different phase; see (57). Its normal derivative on $V$ is the same however and the principal symbol is the same as that of $\left.\partial_{\nu} u^{\text {ref }}\right|_{V}$; see (58) below.

To prove stable recovery of the lens relation $\mathcal{L}$, we recall that the $H^{1} \rightarrow L^{2}$ norm of the DN maps is not suitable for measuring how close the canonical relations $\mathcal{L}$ and $\widetilde{\mathcal{L}}$ of the FIOs $\Lambda_{g, A, q}^{\mathrm{gl}}$ and $\Lambda_{\tilde{g}, \tilde{A}, \tilde{q}}^{\mathrm{gl}}$ are. Instead, we formulate stability based on measuring propagation of singularities. Given a properly supported $\Psi D O R$ on $\partial M$ near $\left(y_{0}, \eta^{0}\right)$, with a principal symbol $r_{0}$, we consider $\Lambda^{*} R \Lambda$, where $\Lambda=\Lambda_{g, A, q}^{\mathrm{gl}}$. By the Egorov theorem, this is actually a $\Psi$ DO near $\left(x_{0}, \xi_{0}\right)$ with a principal symbol $\left(r_{0} \circ \mathcal{L}\right) \lambda_{0}$, where $\lambda_{0}$ is the principal symbol of $\Lambda \Lambda^{*}$ which depends on $g$. In this way, we do not recover $\mathcal{L}$ directly; instead
we recover functions of $\mathcal{L}$ for various choices of $r_{0}$, multiplied by $\lambda_{0}$. Choosing a finite number of $R$ 's satisfying some nondegeneracy assumption, we can apply the implicit function theorem to recover $\mathcal{L}$ locally. In fact, we choose below the differential operators

$$
\begin{equation*}
\left\{R_{j}\right\}=\left\{1, y^{0}, \ldots, y^{n-1}, \partial / \partial y^{0}, \ldots, \partial / \partial y^{n-1}\right\} \tag{54}
\end{equation*}
$$

Theorem 4.3. $\operatorname{Let}\left(y^{0}, \ldots, y^{n-1}\right)$ be local coordinates on $\partial M$ near $y_{0}$. Let

$$
\begin{gather*}
\sum_{j=0}^{n-1}\left\|\Lambda^{*} y^{j} \Lambda-\tilde{\Lambda}^{*} y^{j} \tilde{\Lambda}\right\|_{H^{2}(U) \rightarrow L^{2}(U)} \leq \delta, \quad\left\|\Lambda^{*} \Lambda-\tilde{\Lambda}^{*} \tilde{\Lambda}\right\|_{H^{2}(U) \rightarrow L^{2}(U)} \leq \delta  \tag{55}\\
\sum_{j=0}^{n-1}\left\|\Lambda^{*} \frac{\partial}{\partial y^{j}} \Lambda-\tilde{\Lambda}^{*} \frac{\partial}{\partial y^{j}} \tilde{\Lambda}\right\|_{H^{3}(U) \rightarrow L^{2}(U)} \leq \delta
\end{gather*}
$$

with $\Lambda:=\Lambda_{g, A, q}^{\mathrm{gl}}, \tilde{\Lambda}:=\tilde{\Lambda}_{\tilde{g}, \tilde{A}, \tilde{q}}^{\mathrm{gl}}$. Assume that $(g, A, q)$ and $(\tilde{g}, \tilde{A}, \tilde{q})$ are $\varepsilon$-close to a fixed triple
 there exist $k>0$ and $\mu \in(0,1)$ so that

$$
\begin{equation*}
\left|(\mathcal{L}-\widetilde{\mathcal{L}})\left(x, \xi^{\prime}\right)\right| \leq C \delta^{\mu} \sqrt{-g\left(\xi^{\prime}, \xi^{\prime}\right)} \quad \text { for all }\left(x, \xi^{\prime}\right) \in \mathcal{U} \tag{56}
\end{equation*}
$$

if $\mathcal{U}$ and $\varepsilon>0$ are small enough.
A few remarks:
(a) The square-root term is just a homogeneity factor.
(b) The cotangent bundle $T^{*} \partial M$ is not a linear space; therefore the difference $\mathcal{L}-\widetilde{\mathcal{L}}$ makes sense in fixed coordinates only.
(c) The norms in (55) are the natural one since the operators we subtract there are $\Psi$ DOs of orders 2 and 3 , respectively.
(d) The norms in (55) are equivalent to studying the quadratic forms $\left(\Lambda f, R_{j} \Lambda f\right)-\left(\tilde{\Lambda} f, R_{j} \tilde{\Lambda} f\right)$.
(e) One could reduce the number of the $R_{j}$ 's to $2 n-2$; in fact, $R_{0}=1$ in (55) is not needed, as it follows from Remark 4.8, since we can recover $\eta^{\prime} / \eta_{n}$ and use the fact $\eta=\left(\eta^{\prime}, \eta_{n}\right)$ is a null covector.

We prove Theorem 4.3 at the end of this section.
Stable recovery of the light ray transforms of $A$ and $q$. Let, as in the Introduction, $\xi^{0} \in T_{x_{0}} M \backslash 0$ be the future-pointing lightlike covector whose projection on $T^{*} \partial M \backslash 0$ is the timelike covector $\xi^{0 \prime}$ as in the definition of the semiglobal DN map. Let $\gamma_{0}:=\gamma_{x_{0}, \xi^{0,}}$ be the lightlike geodesic issued from $\left(x_{0}, \xi^{0}\right)$ which intersects $\partial M$ at another point $y_{0}$. Let $V$ be a neighborhood of $y_{0}$ containing all endpoints of future-pointing geodesics issued from $\overline{\mathcal{U}}$. Choose and fix any parametrization of the lightlike geodesics close to $\gamma_{0}$ by normalizing $\xi^{\prime}$. This defines a hypersurface $\mathcal{U}_{0}$ in $\mathcal{U}$. The theorem below holds if $\mathcal{U}$ is a small enough neighborhood of $\left(x_{0}, \xi^{0 \prime}\right)$, and therefore $\mathcal{U}_{0}$ is small enough, as well. Then $L_{1}$ and $L_{0}$ are well defined on $\mathcal{U}_{0}$.

Theorem 4.4. Fix a Lorentzian metric $g$ and $\left(x_{0}, \xi^{0}\right)$ satisfying the assumptions above. Let $(A, q)$ and $(\tilde{A}, \tilde{q})$ be two pairs of magnetic and electric potentials. Define $\delta:=\left\|\Lambda_{g, A, q}^{\mathrm{gl}}-\Lambda_{g, \tilde{A}, \tilde{q}}^{\mathrm{gl}}\right\|_{H^{1}(U) \rightarrow L^{2}(V)}$. Then:
(a) For any $\mu<1$ and $m \geq 0$, the following estimates are valid for some integer $N$ whenever $g, A, \tilde{A}, q, \tilde{q}$ are bounded in a certain $C^{k}$ norm:

$$
\left\|L_{1}(A-\tilde{A})-2 \pi N\right\|_{C^{m}\left(\overline{\mathcal{U}}_{0}\right)} \leq C \delta^{\mu}
$$

(b) Under the a priori condition $\|A-\widetilde{A}\|_{C^{1}(M)} \leq \delta_{1}$ for some $\delta_{1}>0$, for any $0<\mu<\mu^{\prime}$ and $m \geq 0$, the following estimate is valid whenever $g, A, \tilde{A}, q, \tilde{q}$ are bounded in a certain $C^{k}$ norm:

$$
\left\|L_{0}(q-\tilde{q})\right\|_{C^{m}\left(\overline{\mathcal{U}}_{0}\right)} \leq C\left(\delta^{\mu}+\delta_{1}^{\mu}\right)
$$

If there are no conjugate points along $\gamma_{0}$, the proof can be done using a geometric optics construction of the kind (21) but with a different phase in (21) all the way along that geodesic and taking the normal derivative in $V$. Since we do not want to assume there are no conjugate points along $\gamma_{0}$, we will proceed in a somewhat different way.

The fact that we cannot rule out the case $N \neq 0$ based on those arguments can be considered as a manifestation of the Aharonov-Bohm effect. If $\tilde{A}$ and $A$ are a priori close, then $N=0$.

We start with a preparation for the proof of the theorem. Consider first the geometric optics parametrix of the kind (31) of the outgoing solution $u$ like in the previous section. We assume that the boundary condition $f$ has a wave-front set in the timelike cone on the boundary, and for simplicity, assume that it is in the futurepointing one ( $\tau<0$ in local coordinates for which $\partial / \partial t$ is future pointing). Assume at this point that the construction is valid in some neighborhood of the maximal $\gamma_{0}$. We microlocalize all calculations below there. All inverses like $D^{-1}$, etc., below are microlocal parametrices and the equalities between operators are modulo smoothing operators in the corresponding conic microlocal neighborhoods depending on the context.

The construction is the same as that in the previous section, but this time the outgoing solution $u$ is constructed near the bicharacteristic issued from $\left(x_{0}, \xi^{0 \prime}\right)$ all the way to $y_{0}$. Since the solution can reach the other side of the boundary, we need to reflect it at the boundary to satisfy the zero boundary condition. We write the solution $u$ as the sum of the incident wave $u^{\text {inc }}$ and the reflected wave $u^{\text {ref. }}: u=u^{\text {inc }}+u^{\text {ref }}$ where

$$
\begin{align*}
& u^{\mathrm{inc}}(x)=(2 \pi)^{-n} \int e^{i \phi\left(x, \xi^{\prime}\right)}\left(a_{0}^{\mathrm{inc}}+a_{1}^{\mathrm{inc}}+R^{\mathrm{inc}}\right)\left(x, \xi^{\prime}\right) \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime} \\
& u^{\mathrm{ref}}(x)=(2 \pi)^{-n} \int e^{i \phi^{\mathrm{ref}}\left(x, \xi^{\prime}\right)}\left(a_{0}^{\mathrm{ref}}+a_{1}^{\mathrm{ref}}+R^{\mathrm{ref}}\right)\left(x, \xi^{\prime}\right) \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime} \tag{57}
\end{align*}
$$

Here the phase function $\phi^{\text {ref }}$ solves the same eikonal equation as $\phi$ does but satisfies the boundary condition $\left.\phi^{\text {ref }}\right|_{V}=\phi$. It differs from $\phi$ by the sign of its (exterior) normal derivative $\partial \phi / \partial v=-\partial \phi^{\text {ref }} / \partial v>0$ on $V$. The amplitudes are of orders 0 and -1 , respectively, and satisfy

$$
\begin{array}{rlrl}
T^{\mathrm{inc}} a_{0}^{\mathrm{inc}} & =0, & & \left.a_{0}^{\mathrm{inc}}\right|_{U}=\chi \\
T^{\mathrm{ref}} a_{0}^{\mathrm{ref}} & =0, & & \left.a_{0}^{\mathrm{ref}}\right|_{V}=-\left.a_{0}^{\mathrm{inc}}\right|_{V}, \\
i T^{\mathrm{inc}} a_{1}^{\mathrm{inc}} & =-P a_{0}^{\mathrm{inc}}, & & \left.a_{1}^{\mathrm{inc}}\right|_{U}=0 \\
i T^{\mathrm{ref}} a_{1}^{\mathrm{ref}} & =-P a_{0}^{\mathrm{ref}}, & \left.a_{1}^{\mathrm{ref}}\right|_{V}=-\left.a_{1}^{\mathrm{inc}}\right|_{V}
\end{array}
$$

where $T^{\mathrm{inc}}$ and $T^{\text {ref }}$ are the transport operators defined in (28), related to the corresponding phase function, and the remainder terms are of order -2 . Replace $A$ and $\widetilde{A}$ with their gauge-equivalent field satisfying (M3) on $V$. This does not change their light ray transforms. A direct computation, which can be justified as (34), yields

$$
\begin{equation*}
\Lambda_{g, A, q}^{\mathrm{gl}} f=(2 \pi)^{-n} \int e^{i \phi\left(x, \xi^{\prime}\right)}\left(2 i\left(\partial_{\nu} \phi\right) a_{0}^{\mathrm{inc}}+2 i\left(\partial_{\nu} \phi\right) a_{1}^{\mathrm{inc}}+\partial_{\nu}\left(a_{0}^{\mathrm{inc}}+a_{0}^{\mathrm{ref}}\right)+a_{-1}\right) \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime} \tag{58}
\end{equation*}
$$

where $a_{-1}$ is of order -1 and $\phi$ and the amplitudes are restricted to $x \in V$.
The expression (58) allows us to factorize $\Lambda_{g, A, q}^{\mathrm{gl}}$ as $\Lambda_{g, A, q}^{\mathrm{gl}}=2 N_{0} D$ modulo FIOs of order 0 associated with the same canonical relation, where $D f$ is the trace of $u^{\text {inc }}$ on $V$ (a "Dirichlet-to-Dirichlet map") and $N_{0}$ is the DN map $\Lambda_{g, 0,0}^{\text {loc }}$ but localized in $V$. Note that replacing $A$ and $q$ in $N_{0}$ by zeros or not contributes to lower-order error terms. Let $D_{0}$ be the operator $D$ related to $A=0, q=0$. Let $N_{0}^{-1}$ and $D_{0}^{-1}$ be microlocal parametrices of those operators which are actually parametrices of the local Neumann-to-Dirichlet map and the incoming Dirichlet-to-Dirichlet one from $V$ to $U$. Then

$$
\begin{equation*}
D_{0}^{-1} N_{0}^{-1} \Lambda_{g, A, q}^{\mathrm{gl}}=2 D_{0}^{-1} D \bmod S^{-1} \tag{59}
\end{equation*}
$$

is a $\Psi \mathrm{DO}$ of order 0 .
In the next lemma, we do not assume that the geometric optics construction is valid along the whole $\gamma_{0}$.
Lemma 4.5. The operator $D_{0}^{-1} N_{0}^{-1} \Lambda_{g, A, q}^{\mathrm{gl}}$ is a $\Psi D O$ of order 0 in $\mathcal{U}$ with principal symbol

$$
\begin{equation*}
2 \exp \left\{i L_{1} A\left(\gamma_{x^{\prime}, \xi^{\prime}}\right)\right\} \tag{60}
\end{equation*}
$$

where $\gamma_{x^{\prime}, \xi^{\prime}}$ is the future-pointing lightlike geodesic issued from $x^{\prime}$ in direction $\xi$ with projection $\xi^{\prime}$. Proof. By (59), we need to find the principal symbol of $D_{0}^{-1} D$.

The transport equation for $a_{0}^{\text {inc }}$ is

$$
\left[2 g^{j k}\left(\partial_{j} \phi\right)\left(\partial_{k}-i A_{k}\right)+\square_{g} \phi\right] a_{0}^{\mathrm{inc}}=0,\left.\quad a_{0}^{\mathrm{inc}}\right|_{U}=1
$$

As explained right after (25), $g^{j k}\left(\partial_{j} \phi\right) \partial_{k}$ is the tangent vector field along the null geodesic $\gamma_{x^{\prime}, \xi^{\prime}}$. Therefore, with $\Gamma(s):=\left(\gamma_{x^{\prime}, \xi^{\prime}}(s), g \dot{\gamma}_{x^{\prime}, \xi^{\prime}}(s)\right)$, as before, on the set $\chi=1$ we get $a_{0}^{\text {inc }}=\mu$; see (30). That is,

$$
\begin{equation*}
a_{0}^{\mathrm{inc}}(\Gamma(s))=\exp \left\{-\frac{1}{2} \int_{0}^{s}\left(\square_{g} \phi\right)(\Gamma(\sigma)) \mathrm{d} \sigma\right\} \exp \left\{i \int_{0}^{s} A_{k} \circ \gamma_{x^{\prime}, \xi^{\prime}}(\sigma) \dot{\gamma}_{x^{\prime}, \xi^{\prime}}^{k}(\sigma) \mathrm{d} \sigma\right\} . \tag{61}
\end{equation*}
$$

Take $s=s(x, \xi)$ so that $\gamma_{x^{\prime}, \xi^{\prime}}(s) \in V$ to get

$$
\left(a_{0}^{\mathrm{inc}} \circ \mathcal{L}\right)\left(x^{\prime}, \xi^{\prime}\right)=\exp \left\{-\frac{1}{2} L\left(\square_{g} \phi\right)\left(x^{\prime}, \xi^{\prime}\right)+i L_{1} A\left(x^{\prime}, \xi^{\prime}\right)\right\}
$$

where we use the coordinates $\left(x^{\prime}, \xi^{\prime}\right)$ to parametrize the lightlike geodesics locally, and the definition of $L\left(\square_{g} \phi\right)$ is clear from (61).

To construct a representation for $D_{0}^{-1}$, note first that when $A=0$, the term involving $L_{1} A$ is missing above. We look for a parametrix of the incoming solution of $\square_{g} u=0$ with boundary data $u=h$ on $V$ with $\mathrm{WF}(h) \subset \mathcal{V}$ of the form

$$
\begin{equation*}
u(x)=(2 \pi)^{-n} \int e^{i \phi\left(x, \xi^{\prime}\right)} b\left(x, \xi^{\prime}\right) \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime} \tag{62}
\end{equation*}
$$

where $\phi$ is the same phase as in the first equation in (57) and $f$, not related to (20), depends on $h$ as below. The amplitude $b$ solves the transport equation along the same bicharacteristics (with different coefficients since $A=0, q=0$ ) with the initial condition

$$
\left.b\right|_{V}=\left.a^{\mathrm{inc}}\right|_{V}
$$

where $a^{\text {inc }}$ is the full amplitude in the first equation in (57). Restricted to $V$, the map $\left.f \rightarrow u\right|_{V}$ is just $D f$. Then to satisfy $u=h$ on $V$, we need to solve $D f=h$, i.e., to take $f=D^{-1} h$ microlocally.

To illustrate the argument below better, suppose that we are solving the ODE

$$
y^{\prime}+a y=0, \quad y(0)=1
$$

from $t=0$ to $t=1$, where $a=a(t)$. Then we solve

$$
y_{1}^{\prime}+a_{1} y_{1}=0, \quad y_{1}(1)=y(1)
$$

where $a_{1}=a_{1}(t)$. A direct calculation yields

$$
y(t)=\exp \left\{-\int_{0}^{t} a(s) \mathrm{d} s\right\}, \quad y_{1}(t)=\exp \left\{-\int_{1}^{t} a_{1}(s) \mathrm{d} s\right\} y(1)
$$

In particular,

$$
y_{1}(0)=\exp \left\{-\int_{0}^{1}\left(a_{1}(s)-a(s)\right) \mathrm{d} s\right\}
$$

We apply those argument to the transport equation to get

$$
\left.b\right|_{U}=\exp \left\{i L_{1} A\left(\gamma_{x^{\prime}, \xi^{\prime}}\right)\right\}
$$

Then

$$
D_{0}^{-1} D f=(2 \pi)^{-n} \int e^{x^{\prime} \cdot \xi^{\prime}} \exp \left\{i L_{1} A\left(\gamma_{x^{\prime}, \xi^{\prime}}\right)\right\} \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime}
$$

This proves the lemma under the assumption that the geometric optics construction is valid in a neighborhood of $\gamma_{0}$.

To prove the theorem in the general case, we use the partition argument we used in Proposition 4.1. Let $S_{1}, \ldots, S_{k}$ be small timelike surfaces intersecting $\gamma_{0}$ in increasing order from $U$ to $V$ so that the geometric optics construction is valid in a neighborhood of each segment of $\gamma_{0}$ cut by two consecutive surfaces of the sequence $\left\{U, S_{1}, \ldots, S_{k}, V\right\}$. This determines Dirichlet-to-Dirichlet maps $D_{1}$ from $U$ to $S_{1}$, then $D_{2}$, from $S_{1}$ to $S_{2}$, etc., until $D_{k+1}$ from $S_{k}$ to $V$. Then $D=D_{k+1} D_{k} \cdots D_{1}$. Similarly, $D_{0}=D_{0, k+1} D_{0, k} \cdots D_{0,1}$. Then (59) is still valid and takes the form

$$
D_{0}^{-1} N_{0}^{-1} \Lambda_{g, A, q}^{\mathrm{gl}}=2 D_{0,1}^{-1} \cdots D_{0, k}^{-1} D_{0, k+1}^{-1} D_{k+1} D_{k} \cdots D_{1} \bmod S^{-1}
$$

By the first part of the proof, $D_{0, k+1}^{-1} D_{k+1}$ is a $\Psi$ DO on $V$ with principal symbol $\exp \left\{i L_{1}^{(k+1)} A\right\}$, where $L_{1}^{(k+1)}$ is the light ray transform $L_{1}$ restricted to geodesics between $S_{k}$ and $V$. Then we apply Egorov's theorem, see [Hörmander 1985b, Theorem 25.2.5], to conclude that $D_{0, k}^{-1}\left(D_{0, k+1}^{-1} D_{k+1}\right) D_{k}$ is a $\Psi$ DO with a principal symbol that of $D_{0, k+1}^{-1} D_{k+1}$, pulled back by $\mathcal{L}_{k+1}$, the canonical relation between $S_{k}$
and $V$, multiplied by the principal symbol of $D_{0, k}^{-1} D_{k}$. The result is then (60) without the factor of 2 with the integration between $S_{k}$ (through $S_{k+1}$ ) to $V$. Repeating this argument several times, we complete the proof of the lemma.

## Stability of the light ray transform of the magnetic field.

Proof of Theorem 4.4(a). We have

$$
\begin{equation*}
\left\|D_{0}^{-1} N_{0}^{-1}\left(\Lambda_{g, A, q}^{\mathrm{gl}}-\Lambda_{g, \tilde{A}, \tilde{q}}^{\mathrm{gl}}\right)\right\|_{H^{1}(U)} \leq C\left\|\Lambda_{g, A, q}^{\mathrm{gl}}-\Lambda_{g, \tilde{A}, \tilde{q}}^{\mathrm{gl}}\right\|_{H^{1}(U) \rightarrow L^{2}(V)}=C \delta . \tag{63}
\end{equation*}
$$

Set $R:=D_{0}^{-1} N_{0}^{-1}\left(\Lambda_{g, A, q}^{\mathrm{gl}}-\Lambda_{g, \tilde{A}, \tilde{q}}^{\mathrm{gl}}\right)$. By Lemma $4.5, R$ is a $\Psi$ DO in $\mathcal{U}$ of order 0 with principal symbol

$$
r_{0}\left(x^{\prime}, \xi^{\prime}\right)=2 \exp \left\{i L_{1}(\tilde{A}-A)\left(\gamma_{x^{\prime}, \xi^{\prime}}\right)\right\}
$$

and we have $\|R\|_{H^{1}(V)} \leq C \delta$, by (63). We need to derive that $r_{0}$ is "small" in $\mathcal{U}$, as well. We essentially did that in the proof of Theorem 3.2. Choose $f$ as in (20). By [Taylor 1981, Chapter VIII.7], on the set $\chi=1$, we know $e^{-i \lambda x^{\prime} \cdot \xi^{\prime}} R f$ is equal to the full symbol of $\Lambda_{g, A, q}^{\text {loc }}$ with $\lambda=|\xi|$ and $\xi$ in (34) bounded, say, unit. Therefore,

$$
\begin{equation*}
r_{0}\left(x^{\prime}, \xi^{\prime}\right)=e^{-i \lambda x^{\prime} \cdot \xi^{\prime}} R f+O(1 / \lambda) \tag{64}
\end{equation*}
$$

in $C^{k}$ for every $k$. Since $\|f\|_{L^{2}}=C$ and $\|f\|_{H^{1}} \sim \lambda$, (63) yields

$$
\left\|r_{0}\left(\cdot, \xi^{\prime}\right)\right\|_{H^{1}(U)} \leq C \lambda \delta+C / \lambda
$$

uniformly for $\xi$ in some neighborhood of $\xi^{0 \prime}$. With a little more effort one can remove $\lambda$ from $C \lambda \delta$ but this is not needed. Take $\lambda=\delta^{-\frac{1}{2}}$ to get

$$
\left\|\exp \left\{i L_{1}(\tilde{A}-A)\left(\gamma_{x^{\prime}, \xi^{\prime}}\right)\right\}\right\|_{H^{1}\left(\overline{\mathcal{U}}^{\prime}\right)} \leq C \delta^{\frac{1}{2}}
$$

Using interpolation estimates, we can replace the $H^{1}$ norm by any other one at the expense of lowering the exponent on the right from $\frac{1}{2}$ to another positive one, if $k$ in Theorem 4.4 is large enough. Since $\left|e^{i z}-1\right|<\varepsilon$ implies $|z-2 \pi N|<C \varepsilon$ for some integer $N$, this proves part (a) of the theorem.

## Stability of the light ray transform of the potential.

Proof of Theorem 4.4(b). First, we will reduce the problem to the case $\tilde{A}=A$. For $\Lambda_{g, \tilde{A}, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}$, we get a representation as in (58) with a principal symbol with seminorms $O\left(\delta_{1}^{\mu^{\prime}}\right)$, since we can use interpolation estimates to estimate the higher derivatives of $\tilde{A}-A$. Apply a parametrix $\left(\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}\right)^{-1}$ to that difference to get a $\Psi D O Q$ of order 0 microlocally supported in $\mathcal{U}$. If the geometric optics construction is valid all the way from $U$ to $V$, we get as in the proof of (a) that $Q f=O\left(\delta_{1}^{\mu^{\prime}}\right)+O(1 / \lambda)$ in $H^{1}$. This implies the same estimate for $\left\|\left(\Lambda_{g, \tilde{A}, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}\right) f\right\|_{L^{2}}$. In the general case, we can prove the same estimate as in the proof of (a). We will use this later and for now, we assume $\tilde{A}=A$.
Lemma 4.6. The operator $D^{-1} N_{0}^{-1}\left(\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, q}^{\mathrm{gl}}\right)$ is a $\Psi D O$ of order -1 on $U$ with principal symbol

$$
\begin{equation*}
2\left[L_{0}(\tilde{q}-q)\right] \circ \gamma_{x^{\prime}, \xi^{\prime}} \tag{65}
\end{equation*}
$$

where $\gamma_{x^{\prime}, \xi^{\prime}}$ is the future-pointing lightlike geodesic issued from $x^{\prime}$ in direction $\xi$ with projection $\xi^{\prime}$.

Proof. Assume first that the geometric optics construction is valid in a neighborhood of the whole $\gamma_{0}$. In the amplitude

$$
-2 i\left(\partial_{\nu} \phi\right) a_{0}^{\mathrm{inc}}-2 i\left(\partial_{\nu} \phi\right) a_{1}^{\mathrm{inc}}+\partial_{\nu}\left(a_{0}^{\mathrm{inc}}+a_{0}^{\mathrm{ref}}\right)+a_{-1}
$$

in (58), the terms $-2 i\left(\partial_{\nu} \phi\right) a_{0}^{\mathrm{inc}}$ and $\partial_{\nu}\left(a_{0}^{\mathrm{inc}}+a_{0}^{\text {ref }}\right)$ do not depend on $q$; see (61). The other two terms depend on $q$ but they are of different orders. Therefore,

$$
\begin{equation*}
\left(\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, q}^{\mathrm{gl}}\right) f=\left.(2 \pi)^{-n} \int e^{i \phi\left(x, \xi^{\prime}\right)}\left(-2 i\left(\partial_{\nu} \phi\right)\left(\tilde{a}_{1}^{\mathrm{inc}}-a_{1}^{\mathrm{inc}}\right)+\left(a_{-1}-\tilde{a}_{-1}\right)\right) \hat{f}\left(\xi^{\prime}\right) \mathrm{d} \xi^{\prime}\right|_{V} \tag{66}
\end{equation*}
$$

The order of the FIO above is zero. As in the previous proof, we can represent this as a composition of $2 N_{0}$ with the operator $\widetilde{D}-D$ (the difference of two such Dirichlet-to-Dirichlet maps):

$$
\begin{equation*}
\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, q}^{\mathrm{gl}}=2 N_{0}(\tilde{D}-D) \tag{67}
\end{equation*}
$$

modulo FIOs of order -2 . That operator $\widetilde{D}-D$ is an FIO with a symbol, compare with (66),

$$
\begin{equation*}
\sigma(\tilde{D}-D)=-2 i\left(\tilde{a}_{1}^{\mathrm{inc}}-a_{1}^{\mathrm{inc}}\right)+a_{-2} \tag{68}
\end{equation*}
$$

with $a_{-2}$ of order -2 .
To compute $a_{1}^{\text {inc }}$, recall the transport equation for $a_{1}^{\text {inc }}$

$$
\begin{equation*}
\left[2 g^{j k} \partial_{j} \phi\left(\partial_{k}-i A_{k}\right)+\square_{g} \phi\right] a_{1}^{\mathrm{inc}}=i P a_{0}^{\mathrm{inc}},\left.\quad a_{1}^{\mathrm{inc}}\right|_{U}=0 \tag{69}
\end{equation*}
$$

where

$$
i P a_{0}^{\mathrm{inc}}=i P_{g, A, 0} a_{0}^{\mathrm{inc}}+i q a_{0}^{\mathrm{inc}}
$$

The first term on the right is independent of $q$. By (29), (30), with $\Gamma(s)$ as in (61), we get

$$
\begin{align*}
a_{1}^{\mathrm{inc}}(\Gamma(s)) & =\frac{i a_{0}^{\mathrm{inc}}}{2} \int_{0}^{s} \frac{1}{a_{0}^{\mathrm{inc}}}\left[P_{g, A, 0} a_{0}^{\mathrm{inc}}+q a_{0}^{\mathrm{inc}}\right] \circ \Gamma(\sigma) \mathrm{d} \sigma \\
& =\frac{i a_{0}^{\mathrm{inc}}}{2} \int_{0}^{s}\left[\frac{1}{a_{0}^{\mathrm{inc}}} P_{g, A, 0} a_{0}^{\mathrm{inc}}+q\right] \circ \Gamma(\sigma) \mathrm{d} \sigma . \tag{70}
\end{align*}
$$

The potential $q$ depends on $x$ only, so $q \circ \Gamma(s)=q \circ \gamma(s)$. In (70), only the last term depends on $q$ and is an integral of $q$ over lightlike geodesics multiplied by an elliptic factor. Note that the integral, as well as $a_{1}^{\text {inc }}$, are homogeneous of order -1 in $\xi^{\prime}$, as they should be.

We go back to (68) now. Using (70), the terms involving $P_{g, A, 0}$ and $P_{g, \tilde{A}, 0}$ cancel below and we get

$$
\begin{equation*}
\sigma(\widetilde{D}-D) \circ \mathcal{L}=i a_{0}^{\mathrm{inc}} L_{0}(\tilde{q}-q)+a_{-2} \tag{71}
\end{equation*}
$$

where $a_{-2}$ is a symbol of order -2 , different from the one above.
Similarly to (59), we have

$$
\begin{equation*}
D^{-1} N_{0}^{-1}\left(\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, q}^{\mathrm{gl}}\right)=2 D^{-1}(\widetilde{D}-D) \bmod S^{-2} \tag{72}
\end{equation*}
$$

Therefore, we need to compute the principal symbol of $2 D^{-1}(\tilde{D}-D)$. Let $R$ be a $\Psi D O$ in $U$ with principal symbol $r_{-1}$ given by (65). Then, in $\mathcal{U}, D R$ is an FIO of the type (62) with $x \in V$ with the same phase function and a principal amplitude $b_{0}$ solving $T b_{0}=0,\left.b_{0}\right|_{U}=r_{-1}$. By (29), the solution
restricted to $x \in V$ is given by $\left.\mu r_{-1} \circ \mathcal{L}^{-1}\right|_{V}$. Recall that $\mu=a_{0}^{\text {inc }}$. By (71), this is $2 \sigma(\widetilde{D}-D)$ modulo symbols of order -2 . Therefore, $D R=2(\tilde{D}-D)$ modulo FIOs of order -2 . This proves the lemma under the assumption that the geometric optic construction is valid along the whole $\gamma_{0}$.

In the general case, we repeat the arguments of Lemma 4.5. We represent $D$ and $\widetilde{D}$ as a composition $D=D_{k+1} \cdots D_{1}$, and similarly for $\widetilde{D}$. We will do the first step. Consider $2\left(D_{2} D_{1}\right)^{-1}\left(\widetilde{D}_{2} \widetilde{D}_{1}-D_{2} D_{1}\right)$. We have

$$
\begin{aligned}
2\left(D_{2} D_{1}\right)^{-1}\left(\widetilde{D}_{2} \widetilde{D}_{1}-D_{2} D_{1}\right) & =2 D_{1}^{-1} D_{2}^{-1}\left(\left(\widetilde{D}_{2}-D_{2}\right) \widetilde{D}_{1}+D_{2}\left(\widetilde{D}_{1}-D_{1}\right)\right) \\
& =D_{1}^{-1} R_{2} \widetilde{D}_{1}+R_{1}=D_{1}^{-1} R_{2} D_{1}+R_{1}
\end{aligned}
$$

modulo FIOs of order -2 , where $R_{j}=2 D_{j}^{-1}\left(\widetilde{D}_{j}-D_{j}\right), j=1,2$. We apply Egorov's theorem to $D_{1}^{-1} R_{2} D_{1}$ to conclude that it is a $\Psi D O$ on $U$ with a principal symbol equal to the sum of two terms as in (65) with $L_{0}$ taken over the geodesic segments between $U$ and $S_{1}$ first, and $S_{1}$ and $S_{2}$ second. The sum is equal to (65) with $L_{0}$ taken over the union of those segments. Repeating this argument to include $D_{2}$, etc., completes the proof of the lemma.

We finish the proof of part (b) as we did that for part (a). Set

$$
R=D^{-1} N_{0}^{-1}\left(\Lambda_{g, A, \tilde{q}}^{\mathrm{gl}}-\Lambda_{g, A, q}^{\mathrm{gl}}\right)
$$

It is a $\Psi D O$ of order -1 rather than of order 0 as in (a). The analog of (63) is still true. If, as above, $r_{-1}$ is the principal symbol of $R$, then by Lemma 4.6,

$$
r_{-1}\left(x^{\prime}, \xi^{\prime}\right)=-2 i\left[L_{0}(\tilde{q}-q)\right] \circ \gamma_{x^{\prime}, \xi^{\prime}}=\lambda e^{-i \lambda x^{\prime} \cdot \xi^{\prime}} R f+O(1 / \lambda)
$$

with $f$ as in (20); compare with (64). Then

$$
\left\|r_{-1}\left(\cdot, \xi^{\prime}\right)\right\|_{H^{1}(U)} \leq C \lambda^{2} \delta+C / \lambda
$$

Choose $\lambda=\delta^{-\frac{1}{3}}$ to get $\left\|r_{-1}\left(\cdot, \xi^{\prime}\right)\right\|_{H^{1}(U)} \leq C \delta^{\frac{1}{3}}$. This completes the proof of the theorem.

## Proof of the stable recovery of the lens relation.

Proof of Theorem 4.3. We use the notation above. Recall the remark preceding Theorem 4.3 above. The operator $\Lambda^{*} P \Lambda$ is a $\Psi D O$ with a principal symbol $\left(p_{0} \circ \mathcal{L}\right) \lambda_{0}$. Take $p=p_{0}=1$ as in (54) to recover $\lambda_{0}$ first. Knowing the latter, we recover $p_{j} \circ \mathcal{L}$ for $j=1, \ldots, 2 n-1$; see (54). That gives us ( $y, \eta^{\prime}$ ) in (5) as functions of $\left(x, \xi^{\prime}\right)$. Therefore, we reduce the stability problem to the following: show that the principal symbol of a $\Psi$ DO $A$ of order $m$ is determined by $A: H^{m} \rightarrow L^{2}$ in a stable way which is resolved by the lemma below, see also (35), (36). Note that the lemma is a bit more general than what we need since $\left\{P_{j}\right\}$ are simple multiplication and differentiation operators.
Lemma 4.7. Let $Q$ be $\Psi D O$ in $\mathbb{R}^{n}$ with kernel supported in $K \times K$, where $K \subset \mathbb{R}^{n}$ is compact. Let $q_{m}$ be its principal symbol homogeneous of order m. Then

$$
\left\|q_{m}(\cdot, \xi)\right\|_{L^{2}} \leq C|\xi|^{m}\|Q\|_{H^{m} \rightarrow L^{2}}
$$

for all $\xi \neq 0$ with $C>0$ depending on $K$ only.

Proof. Take $f=e^{i x \cdot \xi} \chi(x)$, where $\chi \in C_{0}^{\infty}$ equals 1 in a neighborhood $K$. Then for $x$ in a neighborhood of $K$, we have $Q f(x)=e^{i x \cdot \xi}\left(q_{m}(x, \xi)+r(x, \xi)\right)$ with $r \in S^{m-1}$. We have

$$
\frac{|\xi|^{m}}{C} \leq\|f\|_{H^{m}} \leq C|\xi|^{m}
$$

for $|\xi| \geq 1$. Therefore, for such $\xi$,

$$
\frac{C_{1}\|Q f\|_{L^{2}}}{\|f\|_{H^{m}}} \geq\left\|\frac{q_{m}(\cdot, \xi)}{|\xi|^{m}}\right\|_{L^{2}}-\frac{C_{2}}{|\xi|}
$$

Take the limit $|\xi| \rightarrow \infty$ along radial rays to complete the proof.
We complete the proof of Theorem 4.3 with the aid of Lemma 4.7. We recover first the $L^{2}$-norms with respect to $x$ of $\mathcal{L}(x, \xi)-\widetilde{\mathcal{L}}(x, \xi)$ uniformly in $\xi$ (in fixed coordinates); we can choose $\mu=1$ then. Using standard interpolation estimates, we can estimate the $L^{\infty}$ norm of $\mathcal{L}(x, \xi)-\widetilde{\mathcal{L}}(x, \xi)$ with $\mu<1$ in (56), using the a priori bounds on $g$ and $\tilde{g}$ in some $C^{k}, k \gg 1$, which imply similar bounds on $\mathcal{L}$ and $\tilde{\mathcal{L}}$.

Remark 4.8. The symbol $\lambda_{0}$ can be computed. Since we do not use this formula, we will sketch the proof only. Using Green's formula, as in the proof of [Stefanov and Uhlmann 1998, Proposition 2.1], we can show that $2 N_{\mathcal{V}} \cong D^{*} \Lambda$, where $\cong$ stands for equality modulo lower-order terms, and $N_{\mathcal{V}}$ is $N$ above with the subscript $\mathcal{V}$ indicating that it acts microlocally in that set. The same proof implies that $\Lambda^{*}$ is the DN map associated with the incoming solution, i.e., the one which starts from $\mathcal{V}$ microlocally and hits $\mathcal{U}$. Therefore, $\Lambda^{*} \cong 2 N_{\mathcal{U}} D^{-1}$, where $N_{\mathcal{U}}$ now acts in $\mathcal{U}$. Those two identities and Egorov's theorem imply $\lambda_{0}=-4\left(\xi_{n} \circ \mathcal{L}\right) \xi_{n}$, where $\xi_{n}$ is the function defined in (32).

## 5. Applications and examples

We start with a partial but still general enough case. We follow [Hörmander 1985a, §24.1]. Let $M$ be a Lorentzian manifolds with a timelike boundary $\partial M$. Assume that $t$ is a real-valued smooth function on $M$ so that the level surfaces $t=$ const. are compact and spacelike. For every $a<b$, the (compact) "cylinder" $M_{a b}=\{a \leq t \leq b\}$ (assuming $[a, b]$ is in the range of $t$ ) has a boundary consisting of the spacelike surfaces $t^{-1}(a), t^{-1}(b)$ and $\partial M \cap M_{a b}$ which intersect transversely. This is a generalization of $[0, T] \times \Omega$ in the Riemannian case. By [Hörmander 1985a, Theorem 24.1.1], the following problem is well posed:

$$
P u=0 \quad \text { in } M,\left.\quad u\right|_{t<a}=0,\left.\quad u\right|_{\partial M}=f
$$

with $f \in H^{s}(\partial M), s \geq 1, f=0$ for $t<a$, with a unique solution $u \in H^{s}(M)$ vanishing for $t<a$. Moreover, the map $f \mapsto u$ is continuous. Then the Dirichlet-to-Neumann map $\Lambda_{g, A, q}$ defined as in (4), is well defined.

Let $x_{0} \in U_{0} \Subset U \subset \partial M$ be as in Theorem 3.2. Let $\chi$ be a properly supported $\Psi$ DO cutoff of order 0 localizing near some timelike covector over $x_{0} \in U_{0}$. Since there is a globally defined time function, there are no periodic lightlike geodesics. Then $\chi \Lambda_{g, A, q} \chi$ can be taken as $\Lambda_{g, A, q}^{\text {loc }}$ and Theorem 3.2 applies. If we know a priori that $\Lambda_{g, A, q}: H_{(0)}^{1}(\partial M) \rightarrow L^{2}(\partial M)$ is continuous, where the subscript (0) indicates functions vanishing for $t=0$, then we can replace $\Lambda_{g, A, q}^{\mathrm{loc}}$ by $\Lambda_{g, A, q}$ in (33) and therefore, in Theorem 3.2.


Figure 3. The DN map, with $g$ Minkowski, on the lateral boundary of the cylinder determines a potential and a magnetic field up to $\mathrm{d} \phi$ inside the cylinder but outside the two characteristic cones.

Similarly, with suitable $\Psi D O$ cutoffs $\chi_{1}$ and $\chi_{2}$, we can take $\Lambda_{g, A, q}^{\mathrm{gl}}=\chi_{1} \Lambda_{g, A, q} \chi_{2}$, under the assumptions of Theorem 4.4. And again, if we know that $\Lambda_{g, A, q}: H_{(0)}^{1}(\partial M) \rightarrow L^{2}(\partial M)$ is continuous, we can remove the cutoffs. The results with the cutoffs are actually stronger.

Some special subcases are discussed below. They recover and extend the uniqueness results in [Stefanov 1989; Ramm and Sjöstrand 1991; Ramm and Rakesh 1991; Waters 2014; Salazar 2013; Ben Aïcha 2015; Bellassoued and Ben Aïcha 2017], and some of the stability results there. Using the results in this paper with the support theorems about the light ray transform in [Stefanov 2017; RabieniaHaratbar 2017], we can get new partial data results.

Example 5.1. Let $q$ be a unknown potential but assume that the metric and the magnetic fields are known. Restrict the DN map to $M_{a b}$ for some $a<b$. Then we can recover $L_{0} q$ in a stable way as in Theorem 4.4 over all timelike geodesics intersecting the lateral boundary transversely at their endpoints. If $g$ is real-analytic, then we can apply the results in [Stefanov 2017] to recover $q$ in the set covered by those geodesics under an additional foliation condition. Note that in contrast, the results in [Eskin 2017] require $A$ and $q$ to be analytic in time.

Example 5.2. In the example above, assume that $g$ is Minkowski, and $M_{a b}=[0, T] \times \bar{\Omega}$ for some bounded smooth $\Omega \subset \mathbb{R}^{n}$. By Theorem 3.2, we can recover $L_{1} A$ and $L_{0} q$ over all lightlike geodesics (lines) $l_{z, \theta}=\{(t, x)=(s, z+s \theta): s \in \mathbb{R}\},(z, \theta) \in \mathbb{R}^{n} \times S^{n-1}$, not intersecting the top and the bottom of the cylinder. By [Stefanov 2017], we can recover $q$ in the set covered by those lines. By [RabieniaHaratbar 2017], we can recover $A$ up to $\mathrm{d} \phi, \phi=0$ on $[0, T] \times \partial \Omega$ in that set as well.

For example, if $\Omega$ is the ball $B(0,1)=\{x:|x|<1\}$, the DN map recovers uniquely $q$ and $A$, up to a gauge transform, in the cylinder $[0, T] \times \bar{B}(0,1)$ with the upward characteristic cone with base $\{0\} \times B(0,1)$ and the downward with base $\{T\} \times B(0,1)$ removed; see Figure 3. If $T \leq 2$, those two cones intersect; otherwise they do not but the result holds in both cases. This is the possibly reachable region from $[0, T] \times \partial \Omega$; thus the results are sharp since no information about the complement can be obtained by the finite speed of propagation.

This extends further the uniqueness part of the results in [Stefanov 1989; Ramm and Sjöstrand 1991; Ramm and Rakesh 1991; Waters 2014; Salazar 2013; Ben Aïcha 2015; Bellassoued and Ben Aïcha 2017]. Using the stability estimate in [Begmatov 2001] about $L_{0}$, and the logarithmic estimate for $L_{1}$ in [Salazar 2014], we can use Theorem 4.4 to recover the results in [Salazar 2014]. One important improvement however is that for uniqueness, we do not assume that $A$ and $q$ are known outside $[0, T]$, or that $T=\infty$ because the uniqueness results in [Stefanov 2017; RabieniaHaratbar 2017] do not require the function or the vector field to be compactly supported in time.
Example 5.3. A partial-data case of Example 5.2 is the following. Let $\Gamma \subset \partial \Omega$ be relatively open, and assume that $\partial \Omega$ is strictly convex. Assume that we know the DN map for $f$ supported in $[0, T] \times \Gamma$, and we measure $\Lambda f$ there, as well. Then we can recover $q$ (for all $n \geq 2$ ) and $A$ for $n \geq 3$, up to a gauge transform, in the set covered by the lightlike lines hitting $[0, T] \times \partial \Omega$ in $[0, T] \times \Gamma$ at their both endpoints. When $n=2$, the recovery of $A$ up to a potential $\mathrm{d} \phi$ requires that if we know $L_{1} A$ for some lightlike $l_{z, \theta}$, we also know it for $l_{z,-\theta}$, see [RabieniaHaratbar 2017], and this is the reason we restricted $n$ to $n \geq 3$. Those local uniqueness results for the DN maps are new.
Example 5.4. In a recent work [Bellassoued and Ben Aïcha 2017], an inverse problem for the wave operator

$$
P:=\partial_{t}^{2}+a(t, x) \partial_{t}-\Delta+b(t, x)
$$

with real-valued $a, b$ is studied. The coefficient $b$ causes absorption. We do not restrict $A$ and $q$ to be real-valued, so we can take

$$
A=\left(\frac{i}{2} a(t, x), 0, \ldots, 0\right), \quad q=-\frac{1}{2} i \partial_{t} a(t, x)+b(t, x)
$$

then $P$ in (1) is the same as the one above. Then Theorem 4.4 proves unique recovery of $A, q$ up to the gauge transform $A \mapsto A-\mathrm{d} \psi$ with $\psi=0$ on $[0, T] \times \partial \Omega$. Since $A$ is restricted to the class of covector fields with spatial components zero, we must have $\psi=\psi(t)$. However, then $\psi=0$ for $x \in \partial \Omega$ implies $\psi \equiv 0$. Therefore, the logarithmic and the double logarithmic stability estimates in [Bellassoued and Ben Aïcha 2017] for $a$ and for $b$ which are about the DN map can be obtained by Theorem 4.4 combined with the stability estimates in [Begmatov 2001; Salazar 2014]. We can get new uniqueness results however with partial data as in the previous examples. In the Riemannian case studied by Montalto [2014] we can allow an absorption term as well to obtain, up to a gauge transform, stable recovery of a Riemannian simple metric in a generic class, a magnetic field, a potential and an absorption term from the DN map.

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## References

[^1] J. Amer. Math. Soc. 27:4 (2014), 953-981. MR Zbl
[Begmatov 2001] A. K. Begmatov, "A certain inversion problem for the ray transform with incomplete data", Sibirsk. Mat. Zh. 42:3 (2001), 507-514. In Russian; translated in Siberian Math. J. 42:3 (2001), 428-434. MR Zbl
[Belishev 1987] M. I. Belishev, "An approach to multidimensional inverse problems for the wave equation", Dokl. Akad. Nauk SSSR 297:3 (1987), 524-527. In Russian. MR Zbl
[Belishev 2007] M. I. Belishev, "Recent progress in the boundary control method", Inverse Problems 23:5 (2007), R1-R67. MR Zbl
[Belishev and Kurylev 1992] M. I. Belishev and Y. V. Kurylev, "To the reconstruction of a Riemannian manifold via its spectral data (BC-method)", Comm. Partial Differential Equations 17:5-6 (1992), 767-804. MR Zbl
[Bellassoued and Ben Aïcha 2017] M. Bellassoued and I. Ben Aicha, "Stable determination outside a cloaking region of two time-dependent coefficients in an hyperbolic equation from Dirichlet to Neumann map", J. Math. Anal. Appl. 449:1 (2017), 46-76. MR Zbl
[Bellassoued and Dos Santos Ferreira 2011] M. Bellassoued and D. Dos Santos Ferreira, "Stability estimates for the anisotropic wave equation from the Dirichlet-to-Neumann map", Inverse Probl. Imaging 5:4 (2011), 745-773. MR Zbl
[Ben Aïcha 2015] I. Ben Aïcha, "Stability estimate for a hyperbolic inverse problem with time-dependent coefficient", Inverse Problems 31:12 (2015), art. id. 125010. MR Zbl
[Boman and Quinto 1993] J. Boman and E. T. Quinto, "Support theorems for Radon transforms on real analytic line complexes in three-space", Trans. Amer. Math. Soc. 335:2 (1993), 877-890. MR Zbl
[Cooper and Strauss 1984] J. Cooper and W. Strauss, "The leading singularity of a wave reflected by a moving boundary", J. Differential Equations 52:2 (1984), 175-203. MR Zbl
[Dairbekov et al. 2007] N. S. Dairbekov, G. P. Paternain, P. Stefanov, and G. Uhlmann, "The boundary rigidity problem in the presence of a magnetic field", Adv. Math. 216:2 (2007), 535-609. MR Zbl
[Dos Santos Ferreira et al. 2009] D. Dos Santos Ferreira, C. E. Kenig, M. Salo, and G. Uhlmann, "Limiting Carleman weights and anisotropic inverse problems", Invent. Math. 178:1 (2009), 119-171. MR Zbl
[Duistermaat 1996] J. J. Duistermaat, Fourier integral operators, Progress in Mathematics 130, Birkhäuser, Boston, 1996. MR Zbl
[Eskin 2017] G. Eskin, "Inverse problems for general second order hyperbolic equations with time-dependent coefficients", Bull. Math. Sci. 7:2 (2017), 247-307. MR Zbl
[Eskin and Ralston 2010] G. Eskin and J. Ralston, "The determination of moving boundaries for hyperbolic equations", Inverse Problems 26:1 (2010), art. id. 015001. MR Zbl
[Greenleaf and Uhlmann 1989] A. Greenleaf and G. Uhlmann, "Nonlocal inversion formulas for the X-ray transform", Duke Math. J. 58:1 (1989), 205-240. MR Zbl
[Greenleaf and Uhlmann 1990a] A. Greenleaf and G. Uhlmann, "Composition of some singular Fourier integral operators and estimates for restricted X-ray transforms", Ann. Inst. Fourier (Grenoble) 40:2 (1990), 443-466. MR Zbl
[Greenleaf and Uhlmann 1990b] A. Greenleaf and G. Uhlmann, "Microlocal techniques in integral geometry", pp. 121-135 in Integral geometry and tomography (Arcata, CA, 1989), edited by E. Grinberg and E. T. Quinto, Contemp. Math. 113, Amer. Math. Soc., Providence, RI, 1990. MR Zbl
[Hörmander 1985a] L. Hörmander, The analysis of linear partial differential operators, III: Pseudodifferential operators, Grundlehren der Mathematischen Wissenschaften 274, Springer, 1985. MR Zbl
[Hörmander 1985b] L. Hörmander, The analysis of linear partial differential operators, IV: Fourier integral operators, Grundlehren der Mathematischen Wissenschaften 275, Springer, 1985. MR Zbl
[Isakov and Sun 1992] V. Isakov and Z. Q. Sun, "Stability estimates for hyperbolic inverse problems with local boundary data", Inverse Problems 8:2 (1992), 193-206. MR Zbl
[Kian 2016] Y. Kian, "Recovery of time-dependent damping coefficients and potentials appearing in wave equations from partial data", SIAM J. Math. Anal. 48:6 (2016), 4021-4046. MR Zbl
[Kian et al. 2018] Y. Kian, L. Oksanen, E. Soccorsi, and M. Yamamoto, "Global uniqueness in an inverse problem for time fractional diffusion equations", J. Differential Equations 264:2 (2018), 1146-1170. MR Zbl
[Kurylev et al. 2014a] Y. Kurylev, M. Lassas, and G. Uhlmann, "Inverse problems for Lorentzian manifolds and non-linear hyperbolic equations", preprint, 2014. arXiv
[Kurylev et al. 2014b] Y. Kurylev, M. Lassas, and G. Uhlmann, "Inverse problems in spacetime, I: Inverse problems for Einstein equations (extended preprint version)", preprint, 2014. arXiv
[Lassas et al. 2016] M. Lassas, G. Uhlmann, and Y. Wang, "Inverse problems for semilinear wave equations on Lorentzian manifolds", preprint, 2016. arXiv
[Lassas et al. 2017] M. Lassas, L. Oksanen, P. Stefanov, and G. Uhlmann, "On the inverse problem of finding cosmic strings and other topological defects", Comm. Math. Phys. (online publication November 2017).
[Montalto 2014] C. Montalto, "Stable determination of a simple metric, a covector field and a potential from the hyperbolic Dirichlet-to-Neumann map", Comm. Partial Differential Equations 39:1 (2014), 120-145. MR Zbl
[Palais 1960] R. S. Palais, "Extending diffeomorphisms", Proc. Amer. Math. Soc. 11 (1960), 274-277. MR Zbl
[Petrov 1969] A. Z. Petrov, Einstein spaces, Pergamon Press, Oxford, 1969. MR Zbl
[RabieniaHaratbar 2017] S. RabieniaHaratbar, "Support theorem for the light ray transform on Minkoswki spaces", preprint, 2017. To appear in Inverse Probl. Imaging. arXiv
[Ramm and Rakesh 1991] A. G. Ramm and Rakesh, "Property $C$ and an inverse problem for a hyperbolic equation", J. Math. Anal. Appl. 156:1 (1991), 209-219. MR Zbl
[Ramm and Sjöstrand 1991] A. G. Ramm and J. Sjöstrand, "An inverse problem of the wave equation", Math. Z. 206:1 (1991), 119-130. MR Zbl
[Salazar 2013] R. Salazar, "Determination of time-dependent coefficients for a hyperbolic inverse problem", Inverse Problems 29:9 (2013), art. id. 095015. MR Zbl
[Salazar 2014] R. Salazar, "Stability estimate for the relativistic Schrödinger equation with time-dependent vector potentials", Inverse Problems 30:10 (2014), art. id. 105005. MR Zbl
[Stefanov 1989] P. D. Stefanov, "Uniqueness of the multi-dimensional inverse scattering problem for time dependent potentials", Math. Z. 201:4 (1989), 541-559. MR Zbl
[Stefanov 1991] P. D. Stefanov, "Inverse scattering problem for moving obstacles", Math. Z. 207:3 (1991), 461-480. MR Zbl
[Stefanov 2017] P. Stefanov, "Support theorems for the light ray transform on analytic Lorentzian manifolds", Proc. Amer. Math. Soc. 145:3 (2017), 1259-1274. MR Zbl
[Stefanov and Uhlmann 1998] P. Stefanov and G. Uhlmann, "Stability estimates for the hyperbolic Dirichlet to Neumann map in anisotropic media", J. Funct. Anal. 154:2 (1998), 330-358. MR Zbl
[Stefanov and Uhlmann 2005a] P. Stefanov and G. Uhlmann, "Boundary rigidity and stability for generic simple metrics", J. Amer. Math. Soc. 18:4 (2005), 975-1003. MR Zbl
[Stefanov and Uhlmann 2005b] P. Stefanov and G. Uhlmann, "Stable determination of generic simple metrics from the hyperbolic Dirichlet-to-Neumann map", Int. Math. Res. Not. 2005:17 (2005), 1047-1061. MR Zbl
[Stefanov et al. 2016] P. Stefanov, G. Uhlmann, and A. Vasy, "On the stable recovery of a metric from the hyperbolic DN map with incomplete data", Inverse Probl. Imaging 10:4 (2016), 1141-1147. MR Zbl
[Sun 1990] Z. Q. Sun, "On continuous dependence for an inverse initial-boundary value problem for the wave equation", J. Math. Anal. Appl. 150:1 (1990), 188-204. MR Zbl
[Tataru 1995] D. Tataru, "Unique continuation for solutions to PDE's: between Hörmander's theorem and Holmgren's theorem", Comm. Partial Differential Equations 20:5-6 (1995), 855-884. MR Zbl
[Tataru 1999] D. Tataru, "Unique continuation for operators with partially analytic coefficients", J. Math. Pures Appl. (9) 78:5 (1999), 505-521. MR Zbl
[Taylor 1981] M. E. Taylor, Pseudodifferential operators, Princeton Mathematical Series 34, Princeton University Press, 1981. MR Zbl
[Waters 2014] A. Waters, "Stable determination of X-ray transforms of time dependent potentials from partial boundary data", Comm. Partial Differential Equations 39:12 (2014), 2169-2197. MR Zbl

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