## ANALYSIS \& PDE

## Volume 5 <br> No. 5 <br> 2012

Tapio Helin, Matti Lassas and Lauri Oksanen

# AN INVERSE PROBLEM FOR THE WAVE EQUATION WITH ONE MEASUREMENT AND THE PSEUDORANDOM SOURCE 

msp

# AN INVERSE PROBLEM FOR THE WAVE EQUATION WITH ONE MEASUREMENT AND THE PSEUDORANDOM SOURCE 

Tapio Helin, Matti Lassas and Lauri Oksanen

We consider the wave equation $\left(\partial_{t}^{2}-\Delta_{g}\right) u(t, x)=f(t, x)$, in $\mathbb{R}^{n},\left.u\right|_{\mathbb{R}-\times \mathbb{R}^{n}}=0$, where the metric $g=\left(g_{j k}(x)\right)_{j, k=1}^{n}$ is known outside an open and bounded set $M \subset \mathbb{R}^{n}$ with smooth boundary $\partial M$. We define a source as a sum of point sources, $f(t, x)=\sum_{j=1}^{\infty} a_{j} \delta_{x_{j}}(x) \delta(t)$, where the points $x_{j}, j \in \mathbb{Z}_{+}$, form a dense set on $\partial M$. We show that when the weights $a_{j}$ are chosen appropriately, $\left.u\right|_{\mathbb{R} \times \partial M}$ determines the scattering relation on $\partial M$, that is, it determines for all geodesics which pass through $M$ the travel times together with the entering and exit points and directions. The wave $u(t, x)$ contains the singularities produced by all point sources, but when $a_{j}=\lambda^{-\lambda^{j}}$ for some $\lambda>1$, we can trace back the point source that produced a given singularity in the data. This gives us the distance in $\left(\mathbb{R}^{n}, g\right)$ between a source point $x_{j}$ and an arbitrary point $y \in \partial M$. In particular, if $(\bar{M}, g)$ is a simple Riemannian manifold and $g$ is conformally Euclidian in $\bar{M}$, these distances are known to determine the metric $g$ in $M$. In the case when $(\bar{M}, g)$ is nonsimple, we present a more detailed analysis of the wave fronts yielding the scattering relation on $\partial M$.

## 1. Introduction

In this paper we consider an inverse problem for the wave equation

$$
\begin{aligned}
\left(\partial_{t}^{2}-\Delta_{g}\right) u(t, x) & =f(t, x) \quad \text { in }(0, \infty) \times \mathbb{R}^{n}, \\
\left.u\right|_{t=0} & =\left.\partial_{t} u\right|_{t=0}=0,
\end{aligned}
$$

where $\Delta_{g}$ is the Laplace-Beltrami operator corresponding to a Riemannian metric $g(x)=\left[g_{j k}(x)\right]_{j, k=1}^{n}$, that is,

$$
\Delta_{g} u=\sum_{j, k=1}^{n}|g|^{-1 / 2} \frac{\partial}{\partial x^{j}}\left(|g|^{1 / 2} g^{j k} \frac{\partial}{\partial x^{k}} u\right),
$$

where $|g|=\operatorname{det}\left(g_{j k}\right)$ and $\left[g^{j k}\right]_{j, k=1}^{n}=g(x)^{-1}$ is the inverse matrix of $\left[g_{j k}(x)\right]_{j, k=1}^{n}$. We assume that $g_{j k} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ and that there are $c_{1}, c_{2}>0$ such that

$$
\begin{equation*}
c_{1}|\xi|^{2} \leq \sum_{j . k=1}^{n} g_{j k}(x) \xi^{j} \xi^{k} \leq c_{2}|\xi|^{2}, \quad x, \xi \in \mathbb{R}^{n} \tag{1}
\end{equation*}
$$

The authors were partly supported by the Finnish Centre of Excellence in Inverse Problems Research, the Academy of Finland COE 213476. Lassas was also partly supported by the Mathematical Sciences Research Institute. Oksanen was also partly supported by the Finnish Graduate School in Computational Sciences. Helin was also partly supported by Austrian Academy of Sciences.
MSC2010: primary 35R30, 58J32; secondary 35A18.
Keywords: inverse problems, wave equation, single measurement, pseudorandom, Gaussian beams, scattering relation.

Moreover, we assume that the metric $g$ is known outside an open and bounded set $M \subset \mathbb{R}^{n}$ having a $C^{\infty}$-smooth boundary $\partial M$.

We choose the origin of the time axis so that the source $f$ is active at time $t=0$. To ensure compatibility with the initial conditions, we let $T_{0}<0<T$ and define the measurement map $L=L_{g}$,

$$
\begin{equation*}
L: C_{c}^{\infty}\left(T_{0}, T\right) \otimes C_{c}^{\infty}\left(\mathbb{R}^{n}\right) \rightarrow C^{\infty}\left(\left(T_{0}, T\right) \times \partial M\right), \quad L f=\left.u\right|_{\left(T_{0}, T\right) \times \partial M} \tag{2}
\end{equation*}
$$

where $u$ is the solution of the wave equation

$$
\begin{align*}
\left(\partial_{t}^{2}-\Delta_{g}\right) u(t, x) & =f(t, x) \quad \text { in }\left(T_{0}, T\right) \times \mathbb{R}^{n} \\
\left.u\right|_{t=T_{0}} & =\left.\partial_{t} u\right|_{t=T_{0}}=0 . \tag{3}
\end{align*}
$$

Above, $C_{c}^{\infty}\left(T_{0}, T\right)$ denotes the space of smooth functions having compact support in $\left(T_{0}, T\right)$. Its dual space, the space of generalized functions or distributions, is denoted by $\mathscr{D}^{\prime}\left(T_{0}, T\right)$. Moreover, for functions $\phi \in C_{c}^{\infty}\left(T_{0}, T\right)$ and $\psi \in C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$, we denote their pointwise product by $(\phi \otimes \psi)(t, x)=\phi(t) \psi(x)$.

We remark that the assumption (1), together with the finite speed of propagation for the wave equation, implies that the measurement $L f$ does not depend on $g_{j k}(x)$, for $|x|>R$, when $R$ is sufficiently large. Thus we may assume without loss of generality that all the partial derivatives $\partial_{x}^{\alpha} g_{j k}$ are bounded on $\mathbb{R}^{n}$.

Let $x_{j} \in \partial M, j=1,2, \ldots$, be a dense sequence of points in $\partial M$, and let us consider point sources

$$
f_{x_{j}}(t, x):=\delta(t) \delta_{x_{j}}(x), \quad j=1,2, \ldots
$$

In order to study the measurements $L f_{x_{j}}$, we will use the Sobolev spaces (see [Triebel 1978])

$$
\begin{aligned}
H_{p}^{s}\left(\mathbb{R}^{d}\right) & :=\left\{f \in \mathscr{S}^{\prime}\left(\mathbb{R}^{d}\right) ;\|f\|_{H_{p}^{s}\left(\mathbb{R}^{d}\right)}:=\left\|(1-\Delta)^{s / 2} f\right\|_{L^{p}\left(\mathbb{R}^{d}\right)}<+\infty\right\}, \\
\widetilde{H}_{p}^{s}(U) & :=\left\{f \in H_{p}^{s}\left(\mathbb{R}^{d}\right) ; \operatorname{supp} f \subset \bar{U}\right\}, \\
H_{p}^{s}(U) & :=\left\{f \in \mathscr{D}^{\prime}(U) ; f=\left.h\right|_{U} \text { for some } h \in H_{p}^{s}\left(\mathbb{R}^{d}\right)\right\},
\end{aligned}
$$

where $U \subset \mathbb{R}^{d}$ is open and $s \in \mathbb{R}$. When $p=2$ we omit the subscript $p$ in our notation, that is, we write $H^{s}(U)=H_{2}^{s}(U)$, etc. Moreover, we use projective topology on the tensor product $X \otimes Y$ of two Banach spaces $X$ and $Y$, that is, $\|z\|_{X \otimes Y}:=\inf \sum_{j}\left\|x_{j}\right\|_{X}\left\|y_{j}\right\|_{Y}$, where the infimum is taken over all representations $z=\sum_{j} x_{j} \otimes y_{j}$. We also use projective topology on tensor products of locally convex spaces; see, e.g., [Trèves 1967, Definition 43.2]. The measurement $L f_{x_{j}}$ can be defined in the sense of the following lemma.

Lemma 1.1. Let $p \in(1, n /(n-1))$ and let $m \in \mathbb{N}$ satisfy $m>(n+1) / 4$. Then the measurement operator $L$ defined in (2) has a unique continuous extension

$$
L: \tilde{H}^{-1}\left(T_{0}, T\right) \otimes H_{p}^{-1}\left(\mathbb{R}^{n}\right) \rightarrow \mathscr{D}^{\prime}\left(\left(T_{0}, T\right) \times \partial M\right)
$$

We will prove Lemma 1.1 and other results presented in the introduction in Sections 3-6.
In this paper we study a single measurement $L h_{0}$ that simultaneously combines all the measurements $L f_{x_{j}}$ by adding them together with appropriate weights. When the measurements $L f_{x_{j}}$ are summed together, to the authors' knowledge, there are no algorithms that can filter the value of a particular
measurement from the sum. We will ask, however, whether we can find the essential features given by these measurements, like the travel times between points on $\partial M$, so that the metric could be determined under certain geometric conditions. Our main result is that $L h_{0}$ determines the scattering relation $\Sigma_{M, g}$ for the manifold $(\bar{M}, g)$. Here $h_{0}(t, x)$ is an explicit source that we call pseudorandom; see Definition 1 in Section 2.

The scattering relation has been efficiently used to solve several geometric inverse problems [Dairbekov and Uhlmann 2010; Pestov and Uhlmann 2006; Stefanov and Uhlmann 2008; 2009]. To define the scattering relation, let $T M$ denote the tangent space of $M$ and let $\dot{\gamma}$ denote the tangent vector of a smooth curve $\gamma:[a, b] \rightarrow M$. Let $S M=\left\{(x, \xi) \in T M ;\|\xi\|_{g}=1\right\}$ be the unit sphere bundle on $M$ and define

$$
\partial_{ \pm} S M=\left\{(x, \xi) \in S M ; x \in \partial M, \mp(v, \xi)_{g}>0\right\}
$$

where $v$ is the exterior normal vector of $\partial M$. Moreover, let $\tau_{M, g}(x, \xi)$ be the infimum of the set

$$
\left\{t \in(0, \infty] ; \gamma_{x, \xi}(t) \in \partial M\right\}
$$

where $\gamma_{x, \xi}$ denotes the geodesic on $(M, g)$ with initial data $(x, \xi) \in T M$. We write $\tau=\tau_{M, g}$ when the manifold $(M, g)$ is clear from the context. We define the infimum of the empty set to be $+\infty$.

The scattering relation is the map $\Sigma=\Sigma_{M, g}$,

$$
\Sigma: \mathscr{D}(\Sigma) \rightarrow \overline{\partial_{+} S M} \times \mathbb{R}, \quad \mathscr{D}(\Sigma)=\left\{(x, \xi) \in \partial_{-} S M ; \tau(x, \xi)<\infty\right\}
$$

defined by $\Sigma(x, \xi)=\left(\gamma_{x, \xi}(\tau(x, \xi)), \dot{\gamma}_{x, \xi}(\tau(x, \xi)), \tau(x, \xi)\right)$.
Our main result is the following.
Theorem 1.2. Let $M \subset \mathbb{R}^{n}, n \geq 2$ be an open and bounded set having a $C^{\infty}$-smooth boundary. Then there is a generalized function $h_{0}(t, x)$ supported on $\{0\} \times \partial M$ and having the following properties: Assume that $g_{j k}, g_{j k}^{\prime} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ are two Riemannian metric tensors satisfying (1). Moreover, assume that $g_{j k}(x)=g_{j k}^{\prime}(x)$ for $x \in \mathbb{R}^{n} \backslash M$. Let

$$
T>\max \left(\sup _{(x, \xi) \in \partial_{-} S M} \tau_{M, g}(x, \xi), \sup _{(x, \xi) \in \partial_{-} S M} \tau_{M, g^{\prime}}(x, \xi)\right),
$$

and assume that

$$
L_{g} h_{0}=L_{g^{\prime}} h_{0} \quad \text { on }\left(T_{0}, T\right) \times \partial M .
$$

Then the scattering relations $\Sigma_{M, g}$ and $\Sigma_{M, g^{\prime}}$ of Riemannian manifolds $(M, g)$ and $\left(M, g^{\prime}\right)$ are the same. In particular, if $(\bar{M}, g)$ and $\left(\bar{M}, g^{\prime}\right)$ are simple, the restrictions of the distance functions on the boundary satisfy $d_{\bar{M}, g}(x, y)=d_{\bar{M}, g^{\prime}}(x, y)$ for $x, y \in \partial M$.

We remark that if $\sup _{\partial_{-} S M} \tau$ is infinite, then we prove the above result with measurements on an infinite time interval, that is, we prove that the measurement $\left.u\right|_{\left(T_{0}, \infty\right) \times \partial M}$ determines $\mathscr{D}(\Sigma)$ and $\Sigma$.

Recall that a compact Riemannian manifold ( $\bar{M}, g$ ) with boundary is simple if it is simply connected, any geodesic has no conjugate points, and $\partial M$ is strictly convex with respect to the metric $g$. Any two points of a simple manifold can be joined by a unique geodesic.

The key idea of the proof of Theorem 1.2 is to use source $h_{0}(t, x)=\sum_{j=1}^{\infty} a_{j} f_{x_{j}}$. The point source $a_{j_{0}} f_{x_{j_{0}}}$ produces a singularity, which is observed at a point $y \in \mathbb{R}^{n} \backslash M$ at time $t_{0}=d\left(x_{j_{0}}, y\right)$ with a magnitude $a_{j_{0}} \beta\left(x_{j_{0}}, y\right)$, where $\beta$ is an unknown nonvanishing smooth function. Appropriate choice of the weights $a_{j}$ allows us find the index $j_{0}$ by looking at nearby singularities. Indeed, when $x_{j_{k}} \rightarrow x_{j_{0}}$ and $j_{k} \rightarrow \infty$, we see that the asymptotic behavior of the magnitude $a_{j_{k}} \beta\left(x_{j_{k}}, y\right)$ as $k \rightarrow \infty$ will be that of the weights $a_{j_{k}}$. Thus it is possible to factor out $a_{j_{k}}$ in the magnitude and determine $a_{j_{0}}$. This argument is presented in Section 7 and gives us the distances $d\left(x_{j}, y\right)$ in $\left(\mathbb{R}^{n}, g\right)$ for arbitrary point $y \in \mathbb{R}^{n} \backslash M$ and a source point $x_{j}$.

Theorem 1.2 and boundary rigidity results for simple manifolds imply the following:
Corollary 1.3. Let $M \subset \mathbb{R}^{n}$ and let $g_{j k}, g_{j k}^{\prime} \in C^{\infty}\left(\mathbb{R}^{n}\right)$ be two Riemannian metric tensors satisfying the assumptions of Theorem 1.2. Let $(\bar{M}, g)$ and $\left(\bar{M}, g^{\prime}\right)$ be simple Riemannian manifolds. Then:
(i) If $n=2$ and

$$
\begin{equation*}
L_{g} h_{0}=L_{g^{\prime}} h_{0} \quad \text { on }\left(T_{0}, T\right) \times \partial M \tag{4}
\end{equation*}
$$

then there is a diffeomorphism $\Phi: M \rightarrow M$ such that $\left.\Phi\right|_{\partial M}=\mathrm{Id}$ and $g=\Phi_{*} g^{\prime}$.
(ii) For $n \geq 3$, there is $\epsilon=\epsilon_{n, M}>0$ such that if $\left\|g_{j k}-\delta_{j k}\right\|_{C^{2}(M)}<\epsilon_{n},\left\|g_{j k}^{\prime}-\delta_{j k}\right\|_{C^{2}(M)}<\epsilon_{n}$, and (4) holds, then there is a diffeomorphism $\Phi: M \rightarrow M$ such that $\left.\Phi\right|_{\partial M}=\operatorname{Id}$ and $g=\Phi_{*} g^{\prime}$.
(iii) If $g_{j k}(x)=a(x) \delta_{j k}$ and $g_{j k}^{\prime}(x)=a^{\prime}(x) \delta_{j k}$, that is, the metric tensors are conformally Euclidian, and (4) holds, then $g_{j k}(x)=g_{j k}^{\prime}(x)$ for $x \in M$.

Indeed, by Theorem 1.2, case (i) follows from [Pestov and Uhlmann 2005], case (ii) follows from [Burago and Ivanov 2010], and (iii) from [Muhometov 1977; 1981; Muhometov and Romanov 1978].

If Uhlmann's conjecture [2003], that the scattering relation determines the isometry type of nontrapping compact manifolds with nonempty boundary, can be proven, then Corollary 1.3 holds for a more general class of manifolds.

The problem of determining the metric $g$ (possibly up to a diffeomorphism) given the measurement $L h_{0}$ with only one function $h_{0}(t, x)$ is a formally determined inverse problem. Indeed, the formally computed "dimension of the data," that is, the dimension of $\left(T_{0}, T\right) \times \partial M$, is $n$ and coincides with the dimension of the set $M$ on which the unknown functions $g_{j k}(x)$ are defined.

The formally determined inverse problems have been studied in many cases. For instance, the twodimensional Calderón inverse problem [Astala and Päivärinta 2006; Astala et al. 2005; Imanuvilov et al. 2010; Nachman 1996; Sylvester 1990] is formally determined. The same is true for the related inverse problem for the Schrödinger equation in two dimensions [Bukhgeim 2008]. The corresponding inverse problems in dimension $n \geq 3$ - see [Calderón 1980; Kenig et al. 2007; Lassas et al. 2003; Nachman 1988; Sylvester and Uhlmann 1987] and references in [Greenleaf et al. 2009a; 2009b], are overdetermined that is, the dimension of the data is larger than the dimension of the unknown object. Similar classification holds for the elliptic inverse problems on Riemannian manifolds [Guillarmou and Tzou 2010; 2011; Lassas et al. 2003; Lee and Uhlmann 1989; 2001]. Moreover, the boundary rigidity problem [Kurylev et al. 2010; Michel 1981; Muhometov 1977; 1981; Muhometov and Romanov 1978; Romanov 1987; Stefanov and Uhlmann 2005] is formally determined in dimension $n=2$ and overdetermined for $n \geq 3$.

Inverse problems in the time domain related to the Laplace-Beltrami operator $\Delta_{g}$, namely the inverse boundary value problem for the wave, heat, and dynamical Schrödinger equations with Dirichlet-toNeumann as data - see [Anderson et al. 2004; Belishev and Kurylev 1992; Katchalov and Kurylev 1998; Katchalov et al. 2001] — are overdetermined in dimensions $n \geq 2$. However, these problems are equivalent to the inverse boundary spectral problem (see [Katchalov et al. 2004]), and assuming that the eigenvalues are simple, the Dirichlet-to-Neumann map at a generic Dirichlet boundary value determines the boundary spectral data [Lassas 1995; 1998; Ramm 2001]. Thus, under generic conditions on the spectrum and on the boundary value (that is, under the condition that the these data belong in some open and dense set), it is possible to solve a formally determined inverse problem in time domain.

We point out that in this paper we do not impose any generic conditions on the geometry, and we give an explicit construction of the boundary source. The boundary source considered in this paper is based on the idea of imitating a realization of white noise, and due to the many useful properties of the white noise process, we hope that the constructed source may be useful in the study of other inverse problems requiring generic assumptions on the source.

Another formally determined hyperbolic inverse problem, namely measuring Neumann data when the initial data $\left(\left.u\right|_{t=0},\left.\partial_{t} u\right|_{t=0}\right)$ is nonzero and satisfies subharmonicity or positivity conditions, has been studied using Carleman estimates [Bellassoued and Yamamoto 2008; Imanuvilov and Yamamoto 2003; Isakov 2006; Klibanov 1992; Stefanov and Uhlmann 2011]. The present paper is closely related to these studies, but we emphasize that we assume that the initial data for $u$ vanishes.

Moreover, there are two approaches to solving the formally determined hyperbolic inverse problem to determine a potential from a single boundary measurement. The first one uses Carleman estimates analogous to the estimates mentioned above and assumes similar conditions on the initial data [Bukhgeim and Klibanov 1981]. The second one relies on an adaptation of the Gelfand-Levitan method to multidimensional problems [Rakesh and Sacks 2011; Rakesh 2003; 2008; Romanov 2002; Sacks and Symes 1985].

## 2. Pseudorandom source

In this section we define a special source $h_{0}(t, x)$ which we call pseudorandom. The specific assumptions on the amplitudes are explained in Section 7. An important feature of the pseudorandom source is that it is supported only on a single point in time.

Definition 1. Let $x_{j} \in \partial M, j=1,2, \ldots$, be a dense sequence of distinct points in $\partial M$, and let $a_{j} \in \mathbb{R}$, $j=1,2, \ldots$, with $\sum_{j=1}^{\infty}\left|a_{j}\right|<\infty$, be a sequence of distinct numbers.

We define the pseudorandom source on $\left(x_{j}\right)_{j=1}^{\infty} \subset \partial M$ with coefficients $\left(a_{j}\right)_{j=1}^{\infty} \subset \mathbb{R}$ as the following generalized function on $\mathbb{R} \times \mathbb{R}^{n}$ :

$$
h_{0}(t, x):=\sum_{j=1}^{\infty} a_{j} \delta(t) \delta_{x_{j}}(x), \quad(x, t) \in \mathbb{R}^{n+1}
$$

where $\delta(t)$ and $\delta_{x_{j}}(x)$ are Dirac delta distributions on $\mathbb{R}$ and $\mathbb{R}^{n}$, respectively.

It is rather straightforward to show that $h_{0}$ is well-defined. First, it is well known that $\delta(t) \in H^{-1}(\mathbb{R})$ and $\delta_{x_{j}}(x) \in C\left(\mathbb{R}^{n}\right)^{\prime}$. Next, we have $H_{p^{\prime}}^{1}\left(\mathbb{R}^{n}\right) \subset C\left(\mathbb{R}^{n}\right)$ when $1>n / p^{\prime}$, due to [Triebel 1978, Theorem 2.8.1]. According to [ibid., Theorem 2.6.1], the dual space satisfies $\left(H_{p^{\prime}}^{1}\left(\mathbb{R}^{n}\right)\right)^{\prime}=H_{p}^{-1}\left(\mathbb{R}^{n}\right)$ with $1 / p^{\prime}=1-1 / p$, and hence $C\left(\mathbb{R}^{n}\right)^{\prime} \subset H_{p}^{-1}\left(\mathbb{R}^{n}\right)$ for $1<p<n /(n-1)$. Since $\sum_{j=1}^{\infty}\left|a_{j}\right|<\infty$, we have

$$
\sum_{j=1}^{\infty} a_{j} \delta_{x_{j}}(x) \in H_{p}^{-1}\left(\mathbb{R}^{n}\right)
$$

This yields that for any $p \in\left(1, \frac{n}{n-1}\right)$ and $\epsilon>0$, the pseudorandom source $h_{0}$ satisfies

$$
\begin{equation*}
h_{0} \in \widetilde{H}^{-1}(-\epsilon, \epsilon) \otimes \tilde{H}_{p}^{-1}(M) \tag{5}
\end{equation*}
$$

The spatial structure of the pseudorandom source can be motivated by the structure of the white noise. In the 1-dimensional radar imaging models, white noise signals are considered to be optimal sources when imaging a stationary scatterer [Toomay and Hannen 2004]. This is due to the fact that different translations of the white noise signal are uncorrelated. In a similar fashion, we have the following property for the pseudorandom source $h_{0}$ : for each $x_{j_{0}}$ and each sequence $\left(x_{j_{k}}\right)_{k=1}^{\infty}$ converging to $x_{j_{0}}$ and satisfying $x_{j_{k}} \neq x_{j_{0}}$ for all $k \in \mathbb{Z}_{+}$, it holds that $a_{j_{k}} \rightarrow 0$. This property will be crucial in what follows.

A natural strategy to choose the points $x_{j}$ is by random sampling. The term pseudorandom refers to the fact that algorithmic generators of random numbers use, in fact, a deterministic function to produce a sequence of numbers, but the mixingness of the process is such that the user of the algorithm can consider the numbers to be analogous to independent samples of a random variable. In this manner, the pseudorandom source can be seen as an imitation of a realization of a noise process.

Another source of inspiration for us was a rather new measurement paradigm called compressed sensing [Candès et al. 2006; Donoho 2006], where one aims for a sparse reconstruction of a linear problem using a small number of noisy measurements. We point out that by using the pseudorandom source, one can compress the measurements $L f_{x_{j}}$ with point sources $f_{x_{j}}$ into a single measurement $L h_{0}$.

## 3. Measurement map

In this section we prove that the measurement $L h_{0}$ is well-defined. Let us consider the operator $W: f \mapsto u$ mapping $f$ to the solution of (3). We call such an operator the solution operator for (3). First, we note that by [Hörmander 1985, Theorem 23.2.2], the operator $W: f \mapsto u$ extends in a unique way to a continuous linear operator

$$
\begin{equation*}
W: L^{1}\left(\left(T_{0}, T\right) ; H^{s}\left(\mathbb{R}^{n}\right)\right) \rightarrow C\left(\left[T_{0}, T\right] ; H^{s+1}\left(\mathbb{R}^{n}\right)\right), \quad s \in \mathbb{R} \tag{6}
\end{equation*}
$$

Moreover, if $f \in C^{\infty}\left(\left[T_{0}, T\right] \times \mathbb{R}^{n}\right)$ and $\operatorname{supp}(f) \Subset\left(T_{0}, T\right] \times \mathbb{R}^{n}$, that is, $\operatorname{supp}(f)$ is a compact subset of $\left(T_{0}, T\right] \times \mathbb{R}^{n}$, then $W f \in C^{\infty}\left(\left[T_{0}, T\right] \times \mathbb{R}^{n}\right)$.

We will compose the operator $W$ with the one-sided inverse $\mathscr{I}$ of the derivative $\partial_{t}$, which is given by

$$
\mathscr{I} u(t):=\int_{T_{0}}^{t} u\left(t^{\prime}\right) d t^{\prime}, \quad u \in C_{c}^{\infty}\left(T_{0}, T\right)
$$

One sees easily that this operator has a unique continuous linear extension $\mathscr{I}: \widetilde{H}^{-1}\left(T_{0}, T\right) \rightarrow L^{2}\left(T_{0}, T\right)$.
Next we prove Lemma 1.1 formulated in the introduction, that is, we prove that the measurement map $L$ has a unique continuous extension

$$
\begin{equation*}
\tilde{H}^{-1}\left(T_{0}, T\right) \otimes H_{p}^{-1}\left(\mathbb{R}^{n}\right) \rightarrow \mathscr{D}^{\prime}\left(\left(T_{0}, T\right) \times \partial M\right) \tag{7}
\end{equation*}
$$

Proof of Lemma 1.1. For sufficiently large $z \in \mathbb{R}_{+}$, the operator $z-\Delta_{g}$ is an isomorphism between spaces $H^{s+2}\left(\mathbb{R}^{n}\right)$ and $H^{s}\left(\mathbb{R}^{n}\right)$ as well as between spaces $H_{p}^{s+2}\left(\mathbb{R}^{n}\right)$ and $H_{p}^{s}\left(\mathbb{R}^{n}\right)$ for all integers $s$ by [Shubin 1992].

By the definition of $L$, we have that $L=\operatorname{Tr} \circ W$, where $\operatorname{Tr}$ is the trace operator

$$
\operatorname{Tr}(u)=\left.u\right|_{\left(T_{0}, T\right) \times \partial M}, \quad u \in C^{\infty}\left(\left(T_{0}, T\right) \times \mathbb{R}^{n}\right)
$$

Let $f \in C_{c}^{\infty}\left(\left(T_{0}, T\right) \times \mathbb{R}^{n}\right)$. Then the solution $u=W f$ of the wave equation $\left(\partial_{t}^{2}-\Delta_{g}\right) u=f$ can be written in the form

$$
\begin{equation*}
W f=\left(z-\partial_{t}^{2}\right)^{m}\left(z-\Delta_{g}\right)^{-m} W f+\sum_{j=0}^{m-1}\left(z-\partial_{t}^{2}\right)^{j}\left(z-\Delta_{g}\right)^{-1-j} f \tag{8}
\end{equation*}
$$

Now $f=\partial_{t} \mathscr{I} f$, where $\mathscr{I} f$ is $C^{\infty}$-smooth and satisfies $\operatorname{supp}(\mathscr{I} f) \Subset\left(T_{0}, T\right] \times \mathbb{R}^{n}$. By (6), W $\mathscr{f}$ is $C^{\infty}$-smooth and $\partial_{t} W \mathscr{I} f=W \partial_{t} \mathscr{I} f=W f$. Hence

$$
\begin{equation*}
L f=\partial_{t}\left(z-\partial_{t}^{2}\right)^{m} \operatorname{Tr}\left(z-\Delta_{g}\right)^{-m} W \Phi f+\sum_{j=0}^{m-1}\left(z-\partial_{t}^{2}\right)^{j} \operatorname{Tr}\left(z-\Delta_{g}\right)^{-1-j} f \tag{9}
\end{equation*}
$$

Let us next consider terms appearing in (9). First we consider extension of the operator

$$
\begin{align*}
\sum_{k=1}^{N} \phi_{k} \otimes \psi_{k} & \mapsto \sum_{k=1}^{N}\left(z-\partial_{t}^{2}\right)^{j} \operatorname{Tr}\left(z-\Delta_{g}\right)^{-1-j}\left(\phi_{k} \otimes \psi_{k}\right)  \tag{10}\\
& =\sum_{k=1}^{N}\left(\left(z-\partial_{t}^{2}\right)^{j} \phi_{k}\right) \otimes\left(\operatorname{Tr}\left(z-\Delta_{g}\right)^{-1-j} \psi_{k}\right), \quad j=0, \ldots, m-1
\end{align*}
$$

mapping $C_{c}^{\infty}\left(T_{0}, T\right) \otimes C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ to $C^{\infty}\left(\left(T_{0}, T\right) \times \partial M\right)$. By [Triebel 1978, Theorem 4.7.1], the maps

$$
H_{p}^{-1}\left(\mathbb{R}^{n}\right) \xrightarrow{\left(z-\Delta_{g}\right)^{-1-j}} H_{p}^{1}\left(\mathbb{R}^{n}\right) \xrightarrow{\operatorname{Tr}} B_{p, p}^{1-1 / p}(\partial M)
$$

are continuous, where $B_{p, p}^{1-1 / p}(\partial M)$ is the Besov space on $\partial M$. Thus the operator (10) has a continuous extension in spaces (7).

Next, consider extension of the operator

$$
\begin{equation*}
\sum_{k=1}^{N} \phi_{k} \otimes \psi_{k} \mapsto \sum_{k=1}^{N} \partial_{t}\left(z-\partial_{t}^{2}\right)^{m} \operatorname{Tr}\left(z-\Delta_{g}\right)^{-m} W\left(\left(\mathscr{I} \phi_{k}\right) \otimes \psi_{k}\right) \tag{11}
\end{equation*}
$$

mapping $C_{c}^{\infty}\left(T_{0}, T\right) \otimes C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ to $C^{\infty}\left(\left(T_{0}, T\right) \times \partial M\right)$. As $-1-n / p>-1-n$, we have by [Triebel 1978, Theorem 2.8.1] a continuous embedding $H_{p}^{-1}\left(\mathbb{R}^{n}\right) \hookrightarrow H^{-1-n / 2}\left(\mathbb{R}^{n}\right)$. Moreover, the operator
$\Phi: \widetilde{H}^{-1}\left(T_{0}, T\right) \rightarrow L^{2}\left(T_{0}, T\right)$ and the embedding $L^{2}\left(T_{0}, T\right) \otimes H^{-1-n / 2}\left(\mathbb{R}^{n}\right) \hookrightarrow L^{2}\left(\left(T_{0}, T\right) ; H^{-1-n / 2}\left(\mathbb{R}^{n}\right)\right)$ are continuous. Thus, by (6),

$$
W \mathscr{F}: \tilde{H}^{-1}\left(T_{0}, T\right) \otimes H_{p}^{-1}\left(\mathbb{R}^{n}\right) \rightarrow C\left(\left[T_{0}, T\right] ; H^{-n / 2}\left(\mathbb{R}^{n}\right)\right)
$$

is continuous.
As $\left(1-\Delta_{g}\right)^{-m}: C\left(\left[T_{0}, T\right] ; H^{-n / 2}\left(\mathbb{R}^{n}\right)\right) \rightarrow C\left(\left[T_{0}, T\right] ; H^{-n / 2+2 m}\left(\mathbb{R}^{n}\right)\right)$ is continuous and $-\frac{n}{2}+2 m>\frac{1}{2}$, we see that the operator

$$
\operatorname{Tr}\left(1-\Delta_{g}\right)^{-m} W \mathscr{G}: \tilde{H}^{-1}\left(T_{0}, T\right) \otimes H_{p}^{-1}\left(\mathbb{R}^{n}\right) \rightarrow C\left(\left[T_{0}, T\right] ; L^{2}(\partial M)\right)
$$

is continuous.
Combining the above results, we see that the operator (9) has a continuous extension to the spaces (7). As the spaces $C_{c}^{\infty}\left(T_{0}, T\right)$ and $C_{c}^{\infty}\left(\mathbb{R}^{n}\right)$ are dense in $\tilde{H}^{-1}\left(T_{0}, T\right)$ and $H_{p}^{-1}\left(\mathbb{R}^{n}\right)$, respectively, we see that the continuous extension of $L$ is unique.

## 4. Inner product of a solution and a source

Lemma 4.1. Let $f \in C_{c}^{\infty}\left(\left(T_{0}, T\right) \times M\right)$ and $t_{0} \in\left(T_{0}, T\right)$ and let $w \in C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \mathbb{R}^{n}\right)$ satisfy

$$
\left(\partial_{t}^{2}-\Delta_{g}\right) w=0, \quad \text { in }\left(T_{0}, t_{0}\right) \times \mathbb{R}^{n}
$$

Then

$$
\int_{T_{0}}^{t_{0}} \int_{\mathbb{R}^{n}} f(t, x) w(t, x) d t d V(x)=\int_{\mathbb{R}^{n}}\left(\left(\partial_{t} W f\right)\left(t_{0}, x\right) w\left(t_{0}, x\right)-(W f)\left(t_{0}, x\right)\left(\partial_{t} w\right)\left(t_{0}, x\right)\right) d V(x)
$$

where $d V(x)=|g|^{1 / 2} d x$ is the Riemannian volume measure of $\left(\mathbb{R}^{n}, g\right)$ and $W: f \mapsto u$ is the solution operator of the wave equation (3).
Proof. By finite speed of propagation of waves [Ladyzhenskaya 1985, pp. 150-156], $\operatorname{supp}(W f(t))$ is compact in $\mathbb{R}^{n}$. The claim follows by integration by parts:

$$
\begin{gathered}
\int_{\mathbb{R}^{n}}\left(\left(\partial_{t} u\right)\left(t_{0}, x\right) w\left(t_{0}, x\right)-u\left(t_{0}, x\right)\left(\partial_{t} w\right)\left(t_{0}, x\right)\right) d V(x) \\
-\int_{\mathbb{R}^{n}}\left(\left(\partial_{t} u\right)\left(T_{0}, x\right) w\left(T_{0}, x\right)-u\left(T_{0}, x\right)\left(\partial_{t} w\right)\left(T_{0}, x\right)\right) d V(x) \\
=\int_{\left(T_{0}, t_{0}\right) \times \mathbb{R}^{n}}\left(\left(\partial_{t}^{2}-\Delta_{g}\right) u(t, x) w(t, x)-u(t, x)\left(\partial_{t}^{2}-\Delta_{g}\right) w(t, x)\right) d t d V(x) \\
=\int_{\left(T_{0}, t_{0}\right) \times \mathbb{R}^{n}} f(t, x) w(t, x) d t d V(x)
\end{gathered}
$$

Next, we will prove a generalization of the previous lemma for nonsmooth sources $f$. Denote by $B(0, R)=\left\{x \in \mathbb{R}^{n} ;|x|<R\right\}$ the Euclidean ball. The finite speed of propagation for the wave equation yields that there is $R>0$ such that all $f \in C_{c}^{\infty}\left(\left(T_{0}, T\right) \times M\right)$ satisfy $\operatorname{supp}(W f) \Subset\left(T_{0}, T\right] \times B(0, R)$. We define

$$
\begin{equation*}
\Omega:=B(0, R) \backslash \bar{M} . \tag{12}
\end{equation*}
$$

Below, we use the fact (see [Evans 1998, Theorems 7.2.3/6 and 5.6.3/6]) that the operator $W_{\Omega}: h \mapsto v$ mapping $h$ to the solution of the equation

$$
\begin{align*}
\left(\partial_{t}^{2}-\Delta_{g}\right) v(t, x) & =0 \quad \text { in }\left(T_{0}, T\right) \times \Omega, \\
\left.v\right|_{\left(T_{0}, T\right) \times \partial \Omega} & =h,  \tag{13}\\
\left.v\right|_{t=T_{0}} & =0,\left.\quad \partial_{t} v\right|_{t=T_{0}}=0,
\end{align*}
$$

is continuous as a map $W_{\Omega}: C_{c}^{\infty}\left(\left(T_{0}, T\right) \times \partial \Omega\right) \rightarrow C^{\infty}\left(\left[T_{0}, T\right] \times \bar{\Omega}\right)$.
We let $t_{0} \in\left(T_{0}, T\right)$ and write

$$
\begin{equation*}
\Sigma:=\left\{t_{0}\right\} \times \Omega \tag{14}
\end{equation*}
$$

We denote the trace on $\Sigma$ by $\operatorname{Tr}_{\Sigma}$, that is, we define $\left(\operatorname{Tr}_{\Sigma} u\right)(x):=u\left(t_{0}, x\right)$. Let $v=v(z)$ denote the exterior unit normal vector of $\partial M$ at $z$.

Moreover, let $U$ be an open subset (or a submanifold) of $\mathbb{R}^{n}$, and let us denote by $d V$ (or $d S$ ) the Riemannian volume measure of $(U, g)$. We embed the test functions into the spaces of distribution by using the inner product of the space $L^{2}(U ; d V)$, that is, we identify $u \in C_{0}^{\infty}(U)$ with the distribution

$$
\begin{equation*}
\psi \mapsto \int_{U} u(x) \psi(x) d V(x) \tag{15}
\end{equation*}
$$

We will denote the distribution pairing of $u \in \mathscr{D}^{\prime}(U)$ and $\psi \in C_{0}^{\infty}(U)$ by $(u, \psi)_{\mathscr{D}^{\prime}(U)}$ and use analogous notations for other distribution pairings.

Lemma 4.2. Let $t_{0} \in\left(T_{0}, T\right)$ and define $\Sigma$ by (14). Then operators $\operatorname{Tr}_{\Sigma} W_{\Omega}$ and $\operatorname{Tr}_{\Sigma} \partial_{t} W_{\Omega}$ have unique continuous extensions $\mathscr{E}^{\prime}\left(\left(T_{0}, t_{0}\right) \times \partial \Omega\right) \rightarrow \mathscr{D}^{\prime}(\Omega)$.
Proof. Let $v$ satisfy (13). Consider a function $w \in C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \bar{\Omega}\right)$ such that $\left(\partial_{t}-\Delta_{g}\right) w=0$ in $\left(T_{0}, t_{0}\right) \times \Omega$ and $\left.w\right|_{\left(T_{0}, t_{0}\right) \times \partial \Omega}=0$. Then

$$
\begin{aligned}
0 & =\int_{\Omega \times\left(T_{0}, t_{0}\right)}\left(\left(\partial_{t}-\Delta_{g}\right) v\right) w-v\left(\left(\partial_{t}-\Delta_{g}\right) w\right) d V(x) d t \\
& =\left[\int_{\Omega}\left(\left(\partial_{t} v\right) w-v\left(\partial_{t} w\right)\right) d V(x)\right]_{t=T_{0}}^{t=t_{0}}+\int_{\partial \Omega \times\left(T_{0}, t_{0}\right)}\left(\left(\partial_{\nu} v\right) w-v\left(\partial_{\nu} w\right)\right) d S(x) d t \\
& =\left.\int_{\Omega}\left(\left(\partial_{t} v\right) w-v\left(\partial_{t} w\right)\right) d V(x)\right|_{t=t_{0}}-\int_{\partial \Omega \times\left(T_{0}, t_{0}\right)} h\left(\partial_{\nu} w\right) d S(x) d t
\end{aligned}
$$

where $\partial_{\nu}$ is the normal derivative on $\partial \Omega$,
Denote by $W_{1}: f_{1} \mapsto w$ the solution operator of the equation

$$
\begin{aligned}
\left(\partial_{t}-\Delta_{g}\right) w(t, x) & =0 \quad \text { in }\left(T_{0}, t_{0}\right) \times \Omega, \\
\left.w\right|_{\left(T_{0}, t_{0}\right) \times \partial \Omega} & =0, \\
\left.w\right|_{t=t_{0}} & =f_{1},\left.\quad \partial_{t} w\right|_{t=t_{0}}=0 .
\end{aligned}
$$

The operator $W_{1}: C_{c}^{\infty}(\Omega) \rightarrow C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \bar{\Omega}\right)$ is continuous, as can be seen using Theorems 7.2.3/6 and 5.6.3/6 of [Evans 1998]. Hence, the operator

$$
\partial_{\nu} W_{1}: C_{c}^{\infty}(\Omega) \rightarrow C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \partial \Omega\right),\left.\quad f \mapsto \partial_{\nu} W_{1} f\right|_{\partial \Omega}
$$

is continuous. Moreover,

$$
\left(\operatorname{Tr}_{\Sigma} \partial_{t} W_{\Omega} h, f_{1}\right)_{L^{2}(\Omega ; d V)}=\left(h, \partial_{\nu} W_{1} f_{1}\right)_{L^{2}\left(\left(T_{0}, t_{0}\right) \times \partial \Omega ; d t \otimes d S\right)},
$$

where $\partial_{\nu}$ is the normal derivative on $\partial \Omega$. We define the extension of $\operatorname{Tr}_{\Sigma} \partial_{t} W_{\Omega}$ by identifying it with the transpose $\left(\partial_{\nu} W_{1}\right)^{t}: \mathscr{E}^{\prime}\left(\left(T_{0}, t_{0}\right) \times \partial \Omega\right) \rightarrow \mathscr{D}^{\prime}(\Omega)$ of the operator $\partial_{\nu} W_{1}: C_{c}^{\infty}(\Omega) \rightarrow C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \partial \Omega\right)$.

Similarly, we define the extension of $\operatorname{Tr}_{\Sigma} W_{\Omega}$ by the transpose $\left(\partial_{\nu} W_{2}\right)^{t}: \mathscr{E}^{\prime}\left(\left(T_{0}, t_{0}\right) \times \partial \Omega\right) \rightarrow \mathscr{D}^{\prime}(\Omega)$ of $\partial_{\nu} W_{2}: C_{c}^{\infty}(\Omega) \rightarrow C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \partial \Omega\right)$, where $W_{2}: f_{2} \mapsto w$ is the solution operator of the equation

$$
\begin{aligned}
\left(\partial_{t}-\Delta_{g}\right) w(t, x) & =0 \quad \text { in }\left(T_{0}, t_{0}\right) \times \Omega \\
\left.w\right|_{\left(T_{0}, t_{0}\right) \times \partial \Omega} & =0, \\
\left.w\right|_{t=t_{0}} & =0,\left.\quad \partial_{t} w\right|_{t=t_{0}}=-f_{2} .
\end{aligned}
$$

Denote by $d_{\Omega}(x, y), x, y \in \bar{\Omega}$, the distance function of Riemannian manifold $\left(\bar{\Omega},\left.g\right|_{\bar{\Omega}}\right)$. Next we generalize the result of Lemma 4.1 for a larger class of functions.

Lemma 4.3. Let $t_{0} \in(0, T)$ and $\epsilon>0$ satisfy $[-\epsilon, \epsilon] \subset\left(T_{0}, t_{0}\right)$. Define $\Sigma$ by (14). Let

$$
f \in \tilde{H}^{-1}(-\epsilon, \epsilon) \otimes \tilde{H}_{p}^{-1}(M) \quad \text { and } \quad w \in C^{\infty}\left(\left[T_{0}, t_{0}\right] \times \mathbb{R}^{n}\right)
$$

satisfy

$$
\left(\partial_{t}^{2}-\Delta_{g}\right) w=0, \quad \text { in }\left(T_{0}, t_{0}\right) \times \mathbb{R}^{n}
$$

Suppose that $w\left(t_{0}\right), \partial_{t} w\left(t_{0}\right) \in C_{c}^{\infty}(\Omega)$, and let $\chi \in C_{c}^{\infty}\left(T_{0}, t_{0}\right)$ satisfy $\chi=1$ in a neighborhood of $\left[-\epsilon, t_{0}-r\right]$, where

$$
r:=d_{\Omega}\left(\operatorname{supp}\left(w\left(t_{0}\right)\right) \cup \operatorname{supp}\left(\partial_{t} w\left(t_{0}\right)\right), \partial \Omega\right)
$$

Then

$$
\begin{equation*}
(f, w)_{\mathscr{E}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t_{0}\right)\right)}=\left(\operatorname{Tr}_{\Sigma} \partial_{t} W_{\Omega} \chi L f, w\right)_{\mathscr{D}^{\prime}(\Omega)}-\left(\operatorname{Tr}_{\Sigma} W_{\Omega} \chi L f, \partial_{t} w\right)_{\mathscr{D}^{\prime}(\Omega)}, \tag{16}
\end{equation*}
$$

where we have defined $L f=0$ on $\partial B(0, R)$. Here we regard $\Omega$ as a Riemannian manifold $\left(\Omega,\left.g\right|_{\Omega}\right)$.
Proof. We suppose first that $f \in C_{c}^{\infty}((-\epsilon, \epsilon) \times M)$. Recall that $W$ is solution operator of wave equation (3). Then $W f(\cdot, t)=0$ if $t<-\epsilon$, and

$$
L f=\operatorname{Tr}_{\partial \Omega} W f=\chi \operatorname{Tr}_{\partial \Omega} W f, \quad \text { in }\left(T_{0}, t_{0}-r\right) \times \partial \Omega
$$

where $\operatorname{Tr}_{\partial \Omega}$ is the trace on $\left(T_{0}, T\right) \times \partial \Omega$. As $\Omega \cap \bar{M}=\varnothing$, we have that $\left(\partial_{t}^{2}-\Delta_{g}\right) W f=0$ in $\left(T_{0}, T\right) \times \Omega$. By uniqueness of the solution of (13),

$$
W_{\Omega} \chi \operatorname{Tr}_{\partial \Omega} W f=W f, \quad \text { in }\left(T_{0}, t_{0}-r\right) \times \Omega
$$

By finite speed of propagation,

$$
\operatorname{Tr}_{\Sigma} \partial_{t}^{j} W_{\Omega} \chi \operatorname{Tr}_{\partial \Omega} W f=\operatorname{Tr}_{\Sigma} \partial_{t}^{j} W f, \quad j=0,1
$$

on $\left\{t_{0}\right\} \times \operatorname{supp}\left(w\left(t_{0}\right)\right) \cup \operatorname{supp}\left(\partial_{t} w\left(t_{0}\right)\right)$. By Lemma 4.1, (16) holds.

Then the claim follows, as the embeddings

$$
C_{c}^{\infty}(-\epsilon, \epsilon) \hookrightarrow \tilde{H}^{-1}(-\epsilon, \epsilon), \quad C_{c}^{\infty}(M) \hookrightarrow \tilde{H}_{p}^{-1}(M)
$$

are dense and operators $\left(\operatorname{Tr}_{\Sigma} \partial_{t}^{j} W_{\Omega}\right) \chi L: \tilde{H}^{-1}(-\epsilon, \epsilon) \otimes \tilde{H}_{p}^{-1}(M) \rightarrow \mathscr{D}^{\prime}\left(\left(T_{0}, t_{0}\right) \times \partial \Omega\right), j=0,1$, are continuous.

## 5. Gaussian beams

We consider solutions of the wave equation which are known as Gaussian beams [Babich et al. 1985; Babich and Ulin 1981; Ralston 1982]. These solutions have been constructed to analyze the propagation of singularities for the wave equation in the presence of caustics. Here we use Gaussian beams as an auxiliary technical tool to analyze singularities in the measurements.

Definition 2. Let $\epsilon>0, N \in \mathbb{N}$ and let $\gamma$ be a unit speed geodesic on $\left(\mathbb{R}^{n}, g\right)$. A formal Gaussian beam of order $N$ propagating along geodesic $\gamma$ is a function $U_{\epsilon}^{N}$ of form

$$
U_{\epsilon}^{N}(t, x)=\epsilon^{-n / 4} \exp \left\{-(i \epsilon)^{-1} \theta(t, x)\right\} \sum_{m=0}^{N} u_{m}(t, x)(i \epsilon)^{m}, \quad t \in \mathbb{R}, x \in \mathbb{R}^{n}
$$

satisfying the following properties: The phase function $\theta$ and the amplitude functions $u_{m}$, with $m=$ $0,1, \ldots, N$, are complex-valued smooth functions. The phase function $\theta$ satisfies the conditions

$$
\theta(t, \gamma(t))=0, \quad \operatorname{Im} \theta(t, x) \geq C_{0}(t) d(x, \gamma(t))^{2}
$$

where $C_{0}(t)$ is a continuous strictly positive function. The amplitude function $u_{0}$ satisfies $u_{0}(t, \gamma(t)) \neq 0$. Finally, for any compact set $K \subseteq \mathbb{R} \times \mathbb{R}^{n}$, there is a constant $C>0$ such that the inequality

$$
\left|\left(\partial_{t}^{2}-\Delta_{g}\right) U_{\epsilon}^{N}(t, x)\right| \leq C \epsilon^{N-n / 4}
$$

is satisfied uniformly for $(t, x) \in K$.
The construction of a formal Gaussian beam $U_{\epsilon}^{N}(t, x)$ is considered in detail in [Katchalov et al. 2001, Section 2.4]. Next, we recall the construction and pay attention to the properties of Gaussian beams which we need later.

Let us write the geodesic $\gamma$ in the usual coordinates of $\mathbb{R}^{n}$ as $\gamma(t)=\left(\gamma^{1}(t), \ldots, \gamma^{n}(t)\right)$. We construct the phase function $\theta(t, x)$ at each time $t \in \mathbb{R}$ in terms of a finite Taylor expansion in the $x$ variable centered at $\gamma(t)$,

$$
\theta(t, x)=\sum_{|\alpha| \leq N} \frac{\theta_{\alpha}(t)}{\alpha!}(x-\gamma(t))^{\alpha}
$$

where $\theta_{\alpha}$ are smooth functions and $N \in \mathbb{N}$.
Let $e_{j}=\left(\delta_{1 j}, \ldots, \delta_{n j}\right)$ be multi-indices with the value 1 at the $j$-th place. For clarity, we use the notation $p_{j}(t)=\theta_{e_{j}}(t)$ for the first-order coefficients and the notation $H_{j k}(t)=\theta_{\alpha}(t), \alpha=e_{j}+e_{k}$, for the second-order coefficients in the expansion of $\theta$.

The construction of a formal Gaussian beam consists of the following steps.
(1) We define $\theta_{0}(t)=0$ and $p_{j}(t)=\sum_{k=1}^{n} g_{j k}(\gamma(t)) \dot{\gamma}^{k}(t)$, that is, the first-order coefficients $p_{j}(t)$ are the covariant representation of the velocity vector $\dot{\gamma}$.
(2) The symmetric matrix $H(t)=\left[H_{j k}(t)\right]_{j, k=1}^{n}$ of the second-order coefficients is obtained by solving a Riccati equation, or an equivalent system of ordinary differential equations. We write $H(t)=Z(t) Y(t)^{-1}$, where the pair of complex $n \times n$ matrices $(Z(t), Y(t))$ is the solution of the system of ordinary differential equations

$$
\begin{aligned}
\frac{d}{d t} Y(t)=B(t)^{*} Y(t)+C(t) Z(t), & \left.Y\right|_{t=0}=Y^{0} \\
\frac{d}{d t} Z(t)=-D(t) Y(t)-B(t) Z(t), & \left.Z\right|_{t=0}=Z^{0}
\end{aligned}
$$

Here we choose the initial values to be $Z^{0}=i I$ and $Y^{0}=I$, where $I$ is the identity matrix and $i$ is the imaginary unit. The matrices $B(t), C(t)$, and $D(t)$ in $\mathbb{R}^{n \times n}$ have components given by the second derivatives of the Hamiltonian $h(x, p)=\left(\sum_{j, k=1}^{n} g_{j k}(x) p^{j} p^{k}\right)^{1 / 2}$ evaluated in the point $(x, p)=$ $(\gamma(t), p(t))$ :

$$
B_{l}^{j}=\frac{\partial^{2} h}{\partial x^{l} \partial p_{j}} ; \quad C^{j l}=\frac{\partial^{2} h}{\partial p_{j} \partial p_{l}} ; \quad D_{j l}=\frac{\partial^{2} h}{\partial x^{j} \partial x^{l}} .
$$

The fact that the complex matrix $Y(t)$ is invertible for all $t \in \mathbb{R}$ is crucial for the construction, and is discussed in detail in [Katchalov et al. 2001, Section 2.4].
(3) The coefficients $\theta_{\alpha}(t)$ of order $|\alpha|=m \geq 3$ are solved inductively, with respect to $m$. The coefficients $\theta_{\alpha}(t)$ are constructed using the coefficients $\tilde{\theta}_{\alpha}(t)$ defined so that

$$
\sum_{|\alpha|=m} \tilde{\theta}_{\alpha}(t) \tilde{y}^{\alpha}=\sum_{|\alpha|=m} \theta_{\alpha}(t)(x-\gamma(t))^{\alpha}
$$

for all $\tilde{y}=Y^{-1}(t)(x-\gamma(t)), y \in \mathbb{C}^{n}$. We obtain the coefficients $\tilde{\theta}_{\alpha}(t)$ by solving successive linear systems of ordinary differential equations

$$
\frac{d}{d t} \tilde{\theta}_{\alpha}(t)=K_{\alpha}(t), \quad \tilde{\theta}_{\alpha}(0)=0
$$

where $K_{\alpha}(t)$ depend on $\theta_{\beta}(t)$ with $|\beta| \leq m-1$, the matrix $H(t)$, the vector $p(t)$, and the metric $g_{j k}$ and its derivatives at $\gamma(t)$.
(4) When the phase function $\theta(t, x)$ is constructed, the amplitude functions $u_{n}(t, x)$ are solved using the transport equations, or equivalently, the following ordinary differential equations. Let

$$
u_{m}(t, x)=\sum_{|\alpha| \leq N} \tilde{u}_{m, \alpha}(t) \tilde{y}^{\alpha}, \quad \tilde{y}=Y^{-1}(t)(x-\gamma(t))
$$

where the coefficients $\tilde{u}_{m, \alpha}(t)$ are obtained by solving the successive equations

$$
\frac{d}{d t} \tilde{u}_{m, \alpha}(t)+r(t) \tilde{u}_{m, \alpha}(t)=\mathscr{F}_{m, \alpha}(t), \quad \tilde{u}_{m, \alpha}(0)=\delta_{m, 0} \delta_{|\alpha|, 0}
$$

where $r(t)$ and $\mathscr{F}_{m, \alpha}(t)$ depend on $\tilde{u}_{m^{\prime}, \beta}$ with $|\beta| \leq|\alpha|+2$ and $m^{\prime} \leq m-1$, the function $\theta(t, x)$, the metric $g_{j k}$, and their derivatives at $(t, x), x=\gamma(t)$.

By the above construction, we have the following remark.
Remark 1. The phase function $\theta(t, x)$ and the amplitude functions $u_{m}(t, x)$ at time $t=0$ have the form

$$
\theta(0, x)=\sum_{j, k=1}^{n} g_{j k}(y) \eta^{k}\left(x^{j}-y^{j}\right)+i|x-y|^{2}
$$

where $(y, \eta)=(\gamma(0), \dot{\gamma}(0))$ is the initial data of the geodesic $\gamma, u_{0}(0, x)=1$, and $u_{m}(0, x)=0$ for $m>0$. Hence, $U_{\epsilon}^{N}(0, x)$ is dependent on the metric $g_{j k}$ only via $g_{j k}(y)$. Moreover, $\partial_{t} U_{\epsilon}^{N}(0, x)$, although of more complex form, is dependent on the metric $g_{j k}$ only via $\partial^{\alpha} g_{j k}(y)$ for a certain finite collection of multi-indices $\alpha \in \mathbb{N}^{n}$.

If the coefficients of an ordinary differential equation depend smoothly on some parameter, so does the solution [Amann 1990], and thus we see using an induction that the phase function $\theta$ and the amplitude functions $u_{m}$ depend smoothly on the initial data $(y, \eta)=(\gamma(0), \dot{\gamma}(0))$ of the geodesic $\gamma$. In particular, the amplitude function $u_{0}(t, x ; y, \eta)$ satisfies

$$
\begin{equation*}
u_{0} \in C^{\infty}\left(\mathbb{R} \times \mathbb{R}^{n} \times S \mathbb{R}^{n}\right) \tag{17}
\end{equation*}
$$

Thus far we have considered a formal Gaussian beam. By using continuous dependency of the solution of the wave equation on the source term, one obtains the following results [Katchalov et al. 2001]:

Let $\gamma$ be a unit speed geodesic, $N \in \mathbb{N}, \epsilon>0$, and let $U_{\epsilon}^{N}$ be a formal Gaussian beam of order $N$ propagating along geodesic $\gamma$. Let $\chi \in C_{0}^{\infty}\left(\mathbb{R}^{n}\right)$ be a function which is identically one in a neighborhood of $\gamma(0)$ and let $t_{0}>0$ and let $R$ be the radius in (12). Then for $j \in \mathbb{N}$ and $\alpha \in \mathbb{N}^{n}$ satisfying $j+|\alpha|<N-n / 4$, there is $C>0$ such that the solution $w_{\epsilon}$ of the wave equation,

$$
\begin{align*}
\left(\partial_{t}^{2}-\Delta_{g}\right) w_{\epsilon}(t, x) & =0, \quad(t, x) \in\left(T_{0}, t_{0}\right) \times \mathbb{R}^{n} \\
w_{\epsilon}\left(t_{0}, x\right) & =\chi(x) U_{\epsilon}^{N}(0, x)  \tag{18}\\
\partial_{t} w_{\epsilon}\left(t_{0}, x\right) & =-\chi(x) \partial_{t} U_{\epsilon}^{N}(0, x)
\end{align*}
$$

satisfies

$$
\begin{equation*}
\sup _{x \in B(0, R), t \in\left(T_{0}, t_{0}\right)}\left|\partial_{t}^{j} \partial_{x}^{\alpha}\left(w_{\epsilon}\left(t_{0}-t, x\right)-U_{\epsilon}^{N}(t, x)\right)\right| \leq C \epsilon^{N-(j+|\alpha|)-n / 4} \tag{19}
\end{equation*}
$$

We call $w_{\epsilon}$ a Gaussian beam of order $N$ propagating along geodesic $\gamma$ backwards on time interval $\left(T_{0}, t_{0}\right)$.

## 6. Determination of the travel times

Lemma 6.1. Let $w_{\epsilon}$ be a Gaussian beam of order $N \geq 1+n / 4$ propagating along geodesic $\gamma$ backwards on time interval $\left(T_{0}, t_{0}\right)$, that is, let $w_{\epsilon}$ be the solution of (18). Let $h_{0}$ be the pseudorandom source

$$
\begin{equation*}
h_{0}(t, x)=\sum_{j=1}^{\infty} a_{j} \delta(t) \delta_{x_{j}}(x) \tag{20}
\end{equation*}
$$

If $\gamma\left(t_{0}\right) \neq x_{j}$ for all $j=1,2, \ldots$, then

$$
\lim _{\epsilon \rightarrow 0} \epsilon^{n / 4}\left(h_{0}, w_{\epsilon}\right)_{\mathscr{C}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t_{0}\right)\right)}=0
$$

Moreover, if $\gamma\left(t_{0}\right)=x_{j}$, then

$$
\lim _{\epsilon \rightarrow 0} \epsilon^{n / 4}\left(h_{0}, w_{\epsilon}\right)_{\mathscr{E}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t_{0}\right)\right)}=a_{j} u_{0}\left(t_{0}, x_{j}\right)|g|^{1 / 2}\left(x_{j}\right),
$$

where $u_{0}(t, x)$ is the first amplitude function of a formal Gaussian beam propagating along geodesic $\gamma$. We remind the reader that the test functions are embedded in $\mathscr{E}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, T\right)\right)$ using (15). Proof. By (19), we have that

$$
\begin{aligned}
\epsilon^{n / 4}\left(h_{0}, w_{\epsilon}\right)_{\mathscr{E}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t_{0}\right)\right)} & =\epsilon^{n / 4} \sum_{j=1}^{\infty} a_{j} U_{\epsilon}^{N}\left(t_{0}, x_{j}\right)|g|^{1 / 2}\left(x_{j}\right)+O(\epsilon) \\
& =\sum_{j=1}^{\infty} a_{j} u_{0}\left(t_{0}, x_{j}\right) \exp \left\{-(i \epsilon)^{-1} \theta\left(t_{0}, x_{j}\right)\right\}|g|^{1 / 2}\left(x_{j}\right)+O(\epsilon)
\end{aligned}
$$

As $\operatorname{Im} \theta\left(t_{0}, x_{j}\right) \geq C_{0}\left(t_{0}\right) d\left(x_{j}, \gamma\left(t_{0}\right)\right)$, we have that

$$
\left|\exp \left\{-(i \epsilon)^{-1} \theta\left(t_{0}, x_{j}\right)\right\}\right|=O(\epsilon), \quad \text { if } \gamma\left(t_{0}\right) \neq x_{j}
$$

Suppose that $\gamma\left(t_{0}\right)=x_{j}$. Then $\exp \left\{-(i \epsilon)^{-1} \theta\left(t_{0}, x_{j}\right)\right\}=1$ and there is a constant $C>0$ depending on $\gamma$ and $t_{0}$ such that

$$
\begin{aligned}
& \left.\left|\epsilon^{n / 4}\left(h_{0}, w_{\epsilon}\right) \mathscr{\mathscr { C }}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t_{0}\right)\right)-a_{j} u_{0}\left(t_{0}, x_{j}\right)\right| g\right|^{1 / 2}\left(x_{j}\right) \mid \\
& \quad \leq C \sum_{k=1}^{j-1}\left|a_{k}\right|\left|\exp \left\{-(i \epsilon)^{-1} \theta\left(t_{0}, x_{k}\right)\right\}\right|+C \sum_{k=j+1}^{l}\left|a_{k}\right|\left|\exp \left\{-(i \epsilon)^{-1} \theta\left(t_{0}, x_{k}\right)\right\}\right|+C \sum_{l+1}^{\infty}\left|a_{l}\right|+O(\epsilon) .
\end{aligned}
$$

We may first choose large $l \in \mathbb{N}$ and then small $\epsilon>0$ so that the above three sums are arbitrarily small. The case $\gamma\left(t_{0}\right) \neq x_{j}$ for all $j=1,2, \ldots$, is similar.

Next we define an auxiliary function $S\left(y_{0}, \eta_{0}, t_{0}\right)$ which is nonzero if and only if there is $j \in \mathbb{Z}_{+}$such that $\gamma_{y_{0}, \eta_{0}}\left(t_{0}\right)=x_{j}$.

Definition 3. Let $\left(y_{0}, \eta_{0}\right) \in T \mathbb{R}^{n}$ be such that $y_{0} \in \Omega^{\text {int }}$ and $\left\|\eta_{0}\right\|_{g}=1$. We denote by $\gamma\left(t ; y_{0}, \eta_{0}\right)=$ $\gamma_{y_{0}, \eta_{0}}(t)$ the geodesic on $\left(\mathbb{R}^{n}, g\right)$ with $\gamma(0)=y_{0}, \dot{\gamma}(0)=\eta_{0}$. Moreover, let $w_{\epsilon}$ be a Gaussian beam of order $N \geq 1+n / 4$ propagating along $\gamma(t ; y, \eta)$ backwards on time interval $\left(T_{0}, t_{0}\right)$. We define

$$
S\left(y_{0}, \eta_{0}, t_{0}\right):=\lim _{\epsilon \rightarrow 0} \epsilon^{n / 4}\left(h_{0}, w_{\epsilon}\right) \mathscr{E}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t_{0}\right)\right) .
$$

Lemma 6.2. Let $\left(y_{0}, \eta\right) \in S \Omega$ and $t_{0} \in(0, T)$. Then $L h_{0}$, for pseudorandom source $h_{0}$, and $\left(\Omega,\left.g\right|_{\Omega}\right)$, given as a Riemannian manifold, determine $S\left(y_{0}, \eta_{0}, t_{0}\right)$.

Proof. Let $w_{\epsilon}$ be a Gaussian beam of order $N \geq 1+n / 4$ propagating along the geodesic $\gamma\left(\cdot ; y_{0}, \eta_{0}\right)$ backwards on time interval $\left(T_{0}, t_{0}\right)$. We may choose the cut-off function $\chi$ in (18) so that $w_{\epsilon}\left(t_{0}\right)$ and $\partial_{t} w_{\epsilon}\left(t_{0}\right)$ lie in $C_{c}^{\infty}(\Omega)$. As $\left.g\right|_{\Omega}$ is known, we have by Remark 1 that the initial data $w_{\epsilon}\left(t_{0}\right), \partial_{t} w_{\epsilon}\left(t_{0}\right)$ are known. Moreover, operators $\operatorname{Tr}_{\Sigma} \partial_{t}^{j} W_{\Omega}, j=0,1, \Sigma:=\left\{t_{0}\right\} \times \Omega$, are known. After choosing a suitable
cut-off function $\chi$ in Lemma 4.3, we have that the measurement $L h_{0}$ determines the distributional pairing $\left(h_{0}, w_{\epsilon}\right) \mathscr{C}_{\mathscr{C}^{\prime}\left(\mathbb{R}^{n} \times\left(T_{0}, t\right)\right)}$. Hence $S\left(y_{0}, \eta_{0}, t_{0}\right)$ is determined.

The implicit function theorem yields the following remark. Note that $t_{0} \in \mathbb{R}$ in the remark is not necessarily the first intersection time.
Remark 2. Let $\left(y_{0}, \eta_{0}\right) \in S \mathbb{R}^{n}$ and $t_{0} \in \mathbb{R}$ satisfy

$$
\left(\gamma\left(t_{0} ; y_{0}, \eta_{0}\right), \dot{\gamma}\left(t_{0} ; y_{0}, \eta_{0}\right)\right) \in \partial_{ \pm} S M
$$

Then there are neighborhoods $I \subset \mathbb{R}$ and $U \subset S \mathbb{R}^{n}$ of $t_{0}$ and ( $y_{0}, \eta_{0}$ ) and a smooth map $\ell: U \rightarrow I$ such that for $t \in I$ and $(y, \eta) \in U$,

$$
\gamma(t ; y, \eta) \in \begin{cases}M, & \text { for } \pm t< \pm \ell(y, \eta) \\ \partial M, & \text { for } t=\ell(y, \eta) \\ \Omega, & \text { for } \pm t> \pm \ell(y, \eta)\end{cases}
$$

We remind the reader that $\tau(x, \xi),(x, \xi) \in T \mathbb{R}^{n}$, is defined as the first intersection time with $\partial M$ :

$$
\tau\left(y_{0}, \eta_{0}\right):=\inf \left\{t \in(0, \infty] ; \gamma\left(t ; y_{0}, \eta_{0}\right) \in \partial M\right\}
$$

In the following, we use the Sasaki metric on the tangent bundle TM.
Lemma 6.3. The first intersection times $\tau: S \Omega \rightarrow(0, \infty]$ and $\tau: \partial_{-} S M \rightarrow(0, \infty]$ are lower semicontinuous.

Proof. Let us consider $\tau$ on $S \Omega$. Let a sequence $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset S \Omega$ converge to $\left(y_{0}, \eta_{0}\right) \in S \Omega$ as $j \rightarrow \infty$. We write $\gamma_{j}(t):=\gamma\left(t ; y_{j}, \eta_{j}\right)$ and $\tau_{j}:=\tau\left(y_{j}, \eta_{j}\right)$.

We will show next that $\liminf _{j \rightarrow \infty} \tau_{j} \notin\left(0, \tau_{0}\right)$. Let $t \in\left(0, \tau_{0}\right)$. Then $\gamma_{0}(t) \notin \partial M$ and

$$
d\left(\gamma_{0}(t), \partial M\right)>0 .
$$

Let $j \in \mathbb{Z}_{+}$. Suppose for a moment that $\tau_{j}<\infty$. Noting that $\gamma_{j}$ is unit speed and $\gamma_{j}\left(\tau_{j}\right) \in \partial M$, we have

$$
\left|t-\tau_{j}\right| \geq d\left(\gamma_{j}(t), \gamma_{j}\left(\tau_{j}\right)\right) \geq d\left(\gamma_{j}(t), \partial M\right)
$$

If $\tau_{j}=\infty$, then $\left|t-\tau_{j}\right|=\infty>d\left(\gamma_{j}(t), \partial M\right)$.
The convergence $\gamma_{j}(t) \rightarrow \gamma_{0}(t)$, as $j \rightarrow \infty$, implies that for large $j$,

$$
\left|t-\tau_{j}\right| \geq \frac{d\left(\gamma_{0}(t), \partial M\right)}{2}>0
$$

Hence, $\liminf _{j \rightarrow \infty} \tau_{j} \neq t$ for all $t \in\left(0, \tau_{0}\right)$.
Clearly $\liminf _{j \rightarrow \infty} \tau_{j} \geq 0$, and there is $J \in \mathbb{Z}_{+}$such that

$$
\tau_{j} \geq d\left(y_{j}, \partial M\right) \geq \frac{d\left(y_{0}, \partial M\right)}{2}>0, \quad j \geq J
$$

Hence, $\liminf _{j \rightarrow \infty} \tau_{j} \neq 0$ and $\liminf _{j \rightarrow \infty} \tau_{j} \geq \tau_{0}$.
Let us consider $\tau$ on $\partial_{-} S M$. Let a sequence $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset \partial_{-} S M$ converge to $\left(y_{0}, \eta_{0}\right) \in \partial_{-} S M$ as $j \rightarrow \infty$. We write $\gamma_{j}(t):=\gamma\left(t ; y_{j}, \eta_{j}\right)$ and $\tau_{j}:=\tau\left(y_{j}, \eta_{j}\right)$.



Figure 1. On the left, the trajectory of a Gaussian beam propagating along geodesic $\gamma(t):=\gamma\left(t ; y_{j}, \eta_{j}\right)$ backwards on time interval $\left(T_{0}, t_{j}\right)$. If $S\left(y_{j}, \eta_{j}, t_{j}\right) \neq 0$, then there is a point source at $\gamma\left(t_{j}\right)$. On the right, a sequence $\left(y_{j}, \eta_{j}\right) \in S \Omega$ converging to $(x, \xi) \in \partial_{-} S M$ and trajectories of the corresponding geodesics.

Repeating the above argument, we see that $\liminf _{j \rightarrow \infty} \tau_{j} \notin\left(0, \tau_{0}\right)$. Thus it is enough to show that $\liminf _{j \rightarrow \infty} \tau_{j} \neq 0$.

Remark 2 gives neighborhoods $I \subset \mathbb{R}$ and $U \subset S \mathbb{R}^{n}$ of zero and $\left(y_{0}, \eta_{0}\right)$ and a map $\ell: U \rightarrow I$ of boundary intersection times. We write $V:=U \cap \partial_{-} S M$. As $\gamma(0 ; x, \xi) \in \partial M$ for $(x, \xi) \in V$, we have $\ell=0$ in $V$. In particular, $r:=d(\ell(V), \mathbb{R} \backslash I)>0$. For large $j,\left(\gamma_{j}(0), \dot{\gamma}_{j}(0)\right) \in V$, and thus

$$
\gamma_{j}(t) \in M, \quad t \in(0, r) .
$$

Hence, $\tau_{j} \geq r>0$ for large $j$, and $\liminf _{j \rightarrow \infty} \tau_{j} \geq \tau_{0}$.
We easily see the following continuity result for $\tau$.
Lemma 6.4. Let $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset S \Omega$ converge to $(x, \xi) \in \partial_{-} S M$ in the Sasaki metric. Then

$$
\lim _{j \rightarrow \infty} \tau\left(y_{j}, \eta_{j}\right)=0
$$

Theorem 6.5. Let $(x, \xi) \in \partial_{-} S M$, and denote by $J(x, \xi)$ the set of sequences

$$
\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset(0, \infty) \times S \Omega
$$

for which

$$
\lim _{j \rightarrow \infty}\left(y_{j}, \eta_{j}\right)=(x, \xi), \quad \lim _{j \rightarrow \infty} t_{j} \in(0, \infty), \quad S\left(y_{j}, \eta_{j}, t_{j}\right) \neq 0
$$

The function $S: S \Omega \times(0, \infty) \rightarrow \mathbb{C}$ determines $\tau: \partial_{-} S M \rightarrow(0, \infty]$ by the formula

$$
\tau(x, \xi)=\inf \left\{\lim _{j \rightarrow \infty} t_{j} ;\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \in J(x, \xi) \text { for some }\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset S \Omega\right\}
$$

Moreover, if $\tau(x, \xi)<\infty$, then there is a sequence $\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \in J(x, \xi)$ satisfying

$$
\tau(x, \xi)=\lim _{j \rightarrow \infty} t_{j}
$$

Proof. Let $(x, \xi) \in \partial_{-} S M$ and $\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \in J(x, \xi)$. Let us show that $\tau(x, \xi) \leq \lim _{j \rightarrow \infty} t_{j}$. By Lemma 6.4, $\tau_{j}:=\tau\left(y_{j}, \eta_{j}\right) \rightarrow 0$ as $j \rightarrow \infty$. We define

$$
\tilde{y}_{j}:=\gamma\left(\tau_{j} ; y_{j}, \eta_{j}\right), \quad \xi_{j}:=\dot{\gamma}\left(\tau_{j} ; y_{j}, \eta_{j}\right) .
$$

As $S\left(y_{j}, \eta_{j}, t_{j}\right) \neq 0$, we have

$$
\gamma\left(t_{j}-\tau_{j} ; \tilde{y}_{j}, \xi_{j}\right)=\gamma\left(t_{j} ; y_{j}, \eta_{j}\right) \in \partial M .
$$

As $\lim _{j \rightarrow \infty} t_{j}>0$ and $\lim _{j \rightarrow \infty} \tau_{j}=0$, we have $t_{j}-\tau_{j}>0$ for large $j$. Thus $\tau\left(\tilde{y}_{j}, \xi_{j}\right) \leq t_{j}-\tau_{j}$ for large $j$. Moreover,

$$
\lim _{j \rightarrow \infty}\left(\tilde{y}_{j}, \xi_{j}\right)=(\gamma(0 ; x, \xi), \dot{\gamma}(0 ; x, \xi))=(x, \xi)
$$

In particular, $\left(\tilde{y}_{j}, \xi_{j}\right) \in \partial_{-} S M$ for large $j$. Hence, Lemma 6.3 gives

$$
\lim _{j \rightarrow \infty} t_{j}=\lim _{j \rightarrow \infty}\left(t_{j}-\tau_{j}\right) \geq \liminf _{j \rightarrow \infty} \tau\left(\tilde{y}_{j}, \xi_{j}\right) \geq \tau(x, \xi)
$$

In particular, we have proved the claim in the case $\tau(x, \xi)=\infty$.
Let us assume that $\tau(x, \xi)<\infty$. It is enough to show that there is a sequence $\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \in J(x, \xi)$ satisfying $\tau(x, \xi)=\lim _{j \rightarrow \infty} t_{j}$. We write

$$
t_{0}:=\tau(x, \xi), \quad z:=\gamma\left(t_{0} ; x, \xi\right), \quad \zeta:=-\dot{\gamma}\left(t_{0} ; x, \xi\right) .
$$

We have

$$
(x, \xi)=\left(\gamma\left(t_{0} ; z, \zeta\right),-\dot{\gamma}\left(t_{0} ; z, \zeta\right)\right)
$$

As $(x, \xi) \in \partial_{-} S M$, Remark 2 gives neighborhoods $I$ and $U$ of $t_{0}$ and $(z, \zeta)$ and a map $\ell: U \rightarrow I$ of boundary intersection times. After choosing local coordinates around $z$, we may define

$$
\left(y_{j}, \eta_{j}\right):=\left(\gamma\left(t_{j} ; x_{k_{j}}, \zeta\right),-\dot{\gamma}\left(t_{j} ; x_{k_{j}}, \zeta\right)\right)
$$

where $\left(x_{k_{j}}\right)_{j=1}^{\infty} \subset U$ is a subsequence of the dense sequence of source points in (20) satisfying $\lim _{j \rightarrow \infty} x_{k_{j}}=z$ and $\left(t_{j}\right)_{j=1}^{\infty} \subset I$ satisfies

$$
t_{j}>\ell\left(x_{k_{j}}, \zeta\right), \quad \lim _{j \rightarrow \infty} t_{j}=\ell(z, \zeta)=t_{0}
$$

Clearly, $\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \in J(x, \xi)$ and

$$
\lim _{j \rightarrow \infty} t_{j}=t_{0}=\tau(x, \xi)
$$

## 7. Determination of the scattering relation

In the next theorem, we consider the pseudorandom source $h_{0}(t, x)$ with coefficients

$$
a_{j}=\lambda^{-\lambda^{j}}
$$

with some $\lambda>1$, and make computations "modulo an error in $A$ ", where

$$
A=\left\{-\lambda^{j}: j \in \mathbb{N}\right\}
$$

To this end, let $m_{A}(s)$ be the real number $r$ such that $s=r+a$, where $a \in A$ and $r$ has the smallest possible absolute value. In the case when both $r$ and $-r$ satisfy this condition, we choose the positive value.

Lemma 7.1. Let $\left(y_{0}, \eta_{0}\right) \in S \Omega, t_{0} \in(0, T)$, and suppose that $S\left(y_{0}, \eta_{0}, t_{0}\right) \neq 0$. Then there is a sequence $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset S \Omega$ and $\left(t_{j}\right)_{j=1}^{\infty} \subset(0, T)$ such that

$$
\begin{equation*}
\left(y_{j}, \eta_{j}\right) \rightarrow\left(y_{0}, \eta_{0}\right), \quad t_{j} \rightarrow t_{0}, \quad S\left(y_{j}, \eta_{j}, t_{j}\right) \rightarrow 0, \quad \text { as } j \rightarrow \infty, \quad S\left(y_{j}, \eta_{j}, t_{j}\right) \neq 0 \tag{21}
\end{equation*}
$$

Suppose, moreover, that the coefficients of the pseudorandom source $h_{0}$ are $a_{j}=\lambda^{-\lambda^{j}}$. Then for any sequences $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset T \mathbb{R}^{n}$ and $\left(t_{j}\right)_{j=1}^{\infty} \subset(0, T)$ satisfying (21), we have that

$$
\lim _{j \rightarrow \infty} m_{A}\left(\log _{\lambda}\left|S\left(y_{j}, \eta_{j}, t_{j}\right)\right|\right)=\left.\log _{\lambda}\left|u_{0}\left(t_{0}, \gamma\left(t_{0}\right) ; y_{0}, \eta_{0}\right)\right| g\right|^{1 / 2}\left(\gamma\left(t_{0}\right)\right) \mid
$$

where $\gamma(t)=\gamma\left(t ; y_{0}, \eta_{0}\right)$ and $u_{0}$ is defined as in (17).
Proof. We will use the notation

$$
\gamma_{j}(t):=\gamma\left(t ; y_{j}, \eta_{j}\right), \quad z_{j}:=\gamma_{j}\left(t_{j}\right), \quad S_{j}:=S\left(y_{j}, \eta_{j}, t_{j}\right), \quad \beta_{j}:=\left.\left|u_{0}\left(t_{j}, z_{j} ; y_{j}, \eta_{j}\right)\right| g\right|^{1 / 2}\left(z_{j}\right) \mid .
$$

As $S_{0} \neq 0$, we have that $z_{0}=x_{j}$ for some $j=1,2, \ldots$ By continuity of the geodesic flow and density of $\left(x_{j}\right)_{j=1}^{\infty} \subset \partial M$, there exist a subsequence $\left(x_{k_{j}}\right)_{j=1}^{\infty} \subset\left(x_{j}\right)_{j=1}^{\infty}$ and sequences $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset T \mathbb{R}^{n}$ and $\left(t_{j}\right)_{j=1}^{\infty} \subset(0, T)$ such that

$$
x_{k_{j}} \rightarrow z_{0}, \quad\left(y_{j}, \eta_{j}\right) \rightarrow\left(y_{0}, \eta_{0}\right), \quad t_{j} \rightarrow t_{0}, \quad \text { as } j \rightarrow \infty
$$

and $z_{j}=x_{k_{j}} \neq z_{0}$. Then $\left|S_{j}\right|=\left|a_{k_{j}}\right| \beta_{j} \neq 0$. As $x_{k_{j}} \neq z_{0}$ and $x_{k_{j}} \rightarrow z_{0}$, we have that $k_{j} \rightarrow \infty$ and thus $a_{k_{j}} \rightarrow 0$. By (17) and continuity of the geodesic flow, it holds that $\beta_{j} \rightarrow \beta_{0}>0$. Hence $S_{j} \rightarrow 0$.

Next we use the assumption that $a_{j}=\lambda^{-\lambda^{j}}$. Let $\left(\left(y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \subset T \mathbb{R}^{n}$ and $\left(t_{j}\right)_{j=1}^{\infty} \subset(0, T)$ satisfy (21). As $S_{j} \neq 0$, we have that $\left|S_{j}\right|=a_{k_{j}} \beta_{j}$ for some subsequence $\left(a_{k_{j}}\right)_{j=1}^{\infty} \subset\left(a_{j}\right)_{j=1}^{\infty}$. As $S_{j} \rightarrow 0$, we have that $a_{k_{j}} \rightarrow 0$. Moreover, sequence $\left(\log _{2} \beta_{j}\right)_{j=1}^{\infty}$ is bounded. This boundedness, together with $\log _{\lambda} a_{k_{j}} \in A$ and $\log _{\lambda} a_{k_{j}} \rightarrow-\infty$, yields

$$
m_{A}\left(\log _{\lambda} a_{k_{j}}+\log _{\lambda} \beta_{j}\right)=\log _{\lambda} \beta_{j}
$$

for large $j \in \mathbb{N}$. Hence,

$$
\lim _{j \rightarrow \infty} m_{A}\left(\log _{\lambda}\left|S_{j}\right|\right)=\lim _{j \rightarrow \infty} \log _{\lambda} \beta_{j}=\log _{\lambda} \beta_{0}
$$

Theorem 7.2. If the coefficients of the pseudorandom source $h_{0}$ are $a_{j}=\lambda^{-\lambda^{j}}$, then the functions $S: S \Omega \times(0, \infty) \rightarrow \mathbb{C}$ and $\tau: \partial_{-} S M \rightarrow(0, \infty]$ determine $D(\Sigma)$ and

$$
\gamma(\tau(x, \xi) ; x, \xi), \quad(x, \xi) \in D(\Sigma)
$$

Proof. Clearly $\tau$ on $\partial_{-} S M$ determines $D(\Sigma)$. Let $(x, \xi) \in D(\Sigma)$. By Theorem 6.5, we may choose $\left(\left(t_{j} ; y_{j}, \eta_{j}\right)\right)_{j=1}^{\infty} \in J(x, \xi)$ such that $\lim _{j \rightarrow \infty} t_{j}=\tau(x, \xi)$. As $S\left(y_{j}, \eta_{j}, t_{j}\right) \neq 0$, we have $\gamma\left(t_{j} ; y_{j}, \eta_{j}\right)=x_{k_{j}}$ for some subsequence $\left(x_{k_{j}}\right)_{j=1}^{\infty}$ of the sequence of source points. By Lemma 7.1, the function $S$ determines

$$
\frac{\left|S\left(y_{j}, \eta_{j}, t_{j}\right)\right|}{\left.\left|u_{0}\left(t_{j}, x_{k_{j}} ; y_{j}, \eta_{j}\right)\right| g\right|^{1 / 2}\left(x_{k_{j}}\right) \mid}=a_{k_{j}}
$$

As $a_{j}, j \in \mathbb{Z}_{+}$, are disjoint, this determines the index $k_{j}$ and thus also the point $x_{k_{j}}$. Moreover,

$$
\lim _{j \rightarrow \infty} x_{k_{j}}=\lim _{j \rightarrow \infty} \gamma\left(t_{j} ; y_{j}, \eta_{j}\right)=\gamma(\tau(x, \xi) ; x, \xi)
$$

The following result follows from Remark 2.
Lemma 7.3. Let us denote by $X$ either $S \Omega$ or $\partial_{-} S M$. Let $\left(y_{0}, \eta_{0}\right) \in X$ satisfy

$$
\tau\left(y_{0}, \eta_{0}\right)<\infty, \quad \dot{\gamma}\left(\tau\left(y_{0}, \eta_{0}\right) ; y_{0}, \eta_{0}\right) \notin T_{z} \partial M
$$

where $z=\gamma\left(\tau\left(y_{0}, \eta_{0}\right) ; y_{0}, \eta_{0}\right)$. Then there is a neighborhood $V \subset X$ of $\left(y_{0}, \eta_{0}\right)$ such that $\tau=\ell$ in $V$, where $\ell: U \rightarrow I$ is the map of boundary intersection times defined in Remark 2 for neighborhoods $U \subset X$ and $I \subset \mathbb{R}$ of $\left(y_{0}, \eta_{0}\right)$ and $\tau\left(y_{0}, \eta_{0}\right)$. In particular, $\tau$ is smooth in $V$.

Lemma 7.4. The set of $(x, \xi)$ such that $\gamma(\cdot ; x, \xi)$ is transverse to $\partial M$ is open and dense in

$$
\partial S M:=\{(x, \xi) \in S M ; x \in \partial M\}
$$

Proof. As $\partial_{-} S M \cup \partial_{+} S M$ is open and dense in $\partial S M$, it is enough to show that the set of $(x, \xi)$ such that $\gamma(\cdot ; x, \xi)$ is transverse to $\partial M$ is open and dense in $\partial_{ \pm} S M$. By the parametric transversality theorem [Hirsch 1976, Theorem 3.2.7], the claim follows from the fact that the evaluation map

$$
\begin{aligned}
& F^{\mathrm{ev}}: \partial_{ \pm} S M \times \mathbb{R} \rightarrow \mathbb{R}^{n} \\
& F^{\mathrm{ev}}:(x, \xi, t) \mapsto \gamma(t ; x, \xi)
\end{aligned}
$$

is transverse to $\partial M$.
Lemma 7.5. Let $\left(x_{0}, \xi_{0}\right) \in D(\Sigma)$. Then there is a sequence $\left(\left(x_{j}, \xi_{j}\right)\right)_{j=1}^{\infty} \subset D(\Sigma)$ such that $\gamma\left(\cdot ; x_{j}, \xi_{j}\right)$ is transverse to $\partial M$ and

$$
\lim _{j \rightarrow \infty}\left(x_{j}, \xi_{j}\right)=\left(x_{0}, \xi_{0}\right), \quad \lim _{j \rightarrow \infty} \tau\left(x_{j}, \xi_{j}\right)=\tau\left(x_{0}, \xi_{0}\right)
$$

Proof. We write $\tau_{0}:=\tau\left(x_{0}, \xi_{0}\right)$ and

$$
\left(z_{0}, \zeta_{0}\right):=\left(\gamma\left(\tau_{0} ; x_{0}, \xi_{0}\right),-\dot{\gamma}\left(\tau_{0} ; x_{0}, \xi_{0}\right)\right)
$$

Remark 2 gives a map of boundary intersection times $\ell: U \rightarrow I$ for neighborhoods $U \subset S \mathbb{R}^{n}$ and $I \subset \mathbb{R}$ of $\left(z_{0}, \zeta_{0}\right)$ and $\tau_{0}$. By Lemma 7.4, there is a sequence $\left(\left(z_{j}, \zeta_{j}\right)\right)_{j=1}^{\infty} \subset S M \cap U$ converging to $\left(z_{0}, \zeta_{0}\right)$ such that $\gamma\left(\cdot ; z_{j}, \zeta_{j}\right)$ is transverse to $\partial M$.

We define $t_{j}:=\ell\left(z_{j}, \zeta_{j}\right)$ and

$$
\left(x_{j}, \xi_{j}\right):=\left(\gamma\left(t_{j} ; z_{j}, \zeta_{j}\right),-\dot{\gamma}\left(t_{j} ; z_{j}, \zeta_{j}\right)\right) .
$$

Then $\left(x_{j}, \xi_{j}\right) \rightarrow\left(x_{0}, \xi_{0}\right)$ as $j \rightarrow \infty$. In particular, there is $J \geq 1$ such that $\left(x_{j}, \xi_{j}\right) \in \partial_{-} S M$ for $j \geq J$. By Lemma 6.3,

$$
\begin{aligned}
\tau\left(x_{0}, \xi_{0}\right) & \leq \liminf _{j \rightarrow \infty} \tau\left(x_{j}, \xi_{j}\right) \leq \limsup _{j \rightarrow \infty} \tau\left(x_{j}, \xi_{j}\right) \\
& \leq \lim _{j \rightarrow \infty} \ell\left(z_{j}, \zeta_{j}\right)=\ell\left(z_{0}, \zeta_{0}\right)=\tau\left(x_{0}, \xi_{0}\right)
\end{aligned}
$$

Lemma 7.6. Let $\left(x_{0}, \xi_{0}\right) \in D(\Sigma)$ be such that $\gamma\left(\cdot ; x_{0}, \xi_{0}\right)$ is transverse to $\partial M$. Then there is $\left(y_{0}, \eta_{0}\right) \in$ $S \Omega$ lying on the geodesic $\gamma\left(\cdot ; x_{0}, \xi_{0}\right)$ and a neighborhood $V \subset S_{y_{0}} \Omega$ of $\eta_{0}$ such that the following conditions hold.
(C1) The map $\eta \mapsto \tau\left(y_{0}, \eta\right)$ is smooth $V \rightarrow(0, \infty)$.
(C2) The map

$$
\begin{equation*}
(x(\eta), \xi(\eta)):=\left(\gamma\left(\tau\left(y_{0}, \eta\right) ; y_{0}, \eta\right), \dot{\gamma}\left(\tau\left(y_{0}, \eta\right) ; y_{0}, \eta\right)\right) \tag{22}
\end{equation*}
$$

is smooth $V \rightarrow D(\Sigma)$ and $\left(x\left(\eta_{0}\right), \xi\left(\eta_{0}\right)\right)=\left(x_{0}, \xi_{0}\right)$.
(C3) The map

$$
\begin{equation*}
\tilde{\ell}(\eta):=\tau(x(\eta), \xi(\eta))+\tau\left(y_{0}, \eta\right) \tag{23}
\end{equation*}
$$

is smooth $V \rightarrow(0, \infty)$.
(C4) There is a neighborhood $W \subset \partial M$ of $\gamma\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)$ such that

$$
\begin{equation*}
\eta \mapsto \gamma(\tau(x(\eta), \xi(\eta)) ; x(\eta), \xi(\eta)) \tag{24}
\end{equation*}
$$

is a diffeomorphism $V \rightarrow W$.
Proof. We write $\gamma(t):=\gamma\left(t ; x_{0}, \xi_{0}\right)$ and $z_{0}:=\gamma\left(\tau\left(x_{0}, \xi_{0}\right)\right)$. By Remark $2, \gamma(-t) \in \Omega$ for small $t>0$. Moreover, the points that are conjugate to $z_{0}$ along $\gamma$ are discrete on $\gamma$ [Jost 2008].

Thus there is $\tau_{0}>0$ such that

$$
\left(y_{0}, \eta_{0}\right):=\left(\gamma\left(-\tau_{0}\right), \dot{\gamma}\left(-\tau_{0}\right)\right)
$$

is in $S \Omega, y_{0}$ is not conjugate to $z_{0}$ along $\gamma, \tau\left(y_{0}, \eta_{0}\right)=\tau_{0}$, and

$$
\left(\gamma\left(\tau_{0} ; y_{0}, \eta_{0}\right), \dot{\gamma}\left(\tau_{0} ; y_{0}, \eta_{0}\right)\right)=\left(x_{0}, \xi_{0}\right)
$$

By Lemma 7.3, there is a neighborhood $V_{0} \subset S_{y_{0}} \Omega$ of $\eta_{0}$ such that $\eta \mapsto \tau\left(y_{0}, \eta\right)$ is smooth in $V_{0}$. Hence, the function $\eta \mapsto(x(\eta), \xi(\eta))$ maps $\eta_{0}$ to $\left(x_{0}, \xi_{0}\right)$ and is smooth in $V_{0}$. Moreover, this smoothness, transversality of $\gamma\left(\cdot, x_{0}, \xi_{0}\right)$, and Lemma 7.3 imply that there is a neighborhood $V_{1} \subset V_{0}$ of $\eta_{0}$ such that
$(x(\eta), \xi(\eta)) \in \partial_{-} S M$ and $\eta \mapsto \tau(x(\eta), \xi(\eta))$ is smooth $V_{1} \rightarrow(0, \infty)$. In particular, $(x(\eta), \xi(\eta)) \in D(\Sigma)$ for all $\eta \in V_{1}$. We have shown that $\left(y_{0}, \eta_{0}\right)$ and $V_{1}$ satisfy ( C 1$)-(\mathrm{C} 3)$.

We have

$$
\begin{equation*}
\left.\left(\gamma\left(s ; y_{0}, \eta\right), \dot{\gamma}\left(s ; y_{0}, \eta\right)\right)\right|_{s=t+\tau\left(y_{0}, \eta\right)}=(\gamma(t ; x(\eta), \xi(\eta)), \dot{\gamma}(t ; x(\eta), \xi(\eta))) \tag{25}
\end{equation*}
$$

In particular, $\gamma\left(\tilde{\ell}\left(\eta_{0}\right) ; y_{0}, \eta_{0}\right)=z_{0}$ and

$$
\gamma\left(\tilde{\ell}(\eta) ; y_{0}, \eta\right)=\gamma(\tau(x(\eta), \xi(\eta)) ; x(\eta), \xi(\eta)) \in \partial M
$$

Moreover, as $y_{0}$ is not conjugate to $z_{0}$ along $\gamma$, there are neighborhoods $V_{2} \subset V_{1}, I_{0} \subset(0, \infty)$ and $U_{0} \subset \mathbb{R}^{n}$ of $\eta_{0}, \tilde{\ell}\left(\eta_{0}\right)$ and $z_{0}$ such that $(t, \eta) \mapsto \gamma\left(t ; y_{0}, \eta\right)$ is a diffeomorphism $V_{2} \times I_{0} \rightarrow U_{0}$.

There is a neighborhood $V \subset V_{2}$ of $\eta_{0}$ such that $\tilde{\ell}(V) \subset I_{0}$. The graph of $\eta \mapsto \tilde{\ell}(\eta)$ is an $(n-1)$ dimensional submanifold on $V \times I_{0}$. Hence, the diffeomorphism $(t, \eta) \mapsto \gamma\left(t ; y_{0}, \eta\right)$ maps it onto an ( $n-1$ )-dimensional submanifold $W$ of $U_{0}$. Moreover, $z_{0} \in W$ and $W \subset \partial M$. Thus $W$ is a neighborhood of $z_{0}$ in $\partial M$.

Lemma 7.7. Let $\left(x_{0}, \xi_{0}\right) \in D(\Sigma)$ and $\left(y_{0}, \eta_{0}\right) \in S \Omega$ satisfy conditions $(\mathrm{C} 1)-(\mathrm{C} 4)$ of Lemma 7.6 for neighborhoods $V \subset S_{y_{0}} \Omega$ and $W \subset \partial M$ of $\eta_{0}$ and $z_{0}:=\gamma\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)$. We denote by $F: W \rightarrow V$ the inverse map of (24). Then

$$
\begin{equation*}
\left.\operatorname{grad}_{\partial M}(\tilde{\ell} \circ F)\right|_{z=z_{0}}=\dot{\gamma}_{z_{0}}^{\top} \tag{26}
\end{equation*}
$$

where $\tilde{\ell}: V \rightarrow(0, \infty)$ is the function (23) and $\dot{\gamma}_{z_{0}}^{\top}$ is the orthogonal projection of $\dot{\gamma}\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)$ into $T_{z_{0}} \partial M$.

Proof. Let $\sigma:(-\epsilon, \epsilon) \rightarrow W$ be a smooth curve such that $\sigma(0)=z_{0}$. We define

$$
\Gamma:(-\epsilon, \epsilon) \times \mathbb{R} \rightarrow \mathbb{R}^{n}, \quad \Gamma(s, t):=\gamma\left(t ; y_{0}, F(\sigma(s))\right)
$$

We write $\lambda:=\tilde{\ell} \circ F \circ \sigma$ and $\tilde{\ell}_{0}:=\tilde{\ell}\left(\eta_{0}\right)$. By (25),

$$
\begin{aligned}
\Gamma(s, \lambda(s)) & =\left.\gamma(\tau(x(\eta), \xi(\eta)) ; x(\eta), \xi(\eta))\right|_{\eta=F(\sigma(s))}=\sigma(s) \\
\left(\partial_{t} \Gamma\right)\left(0, \tilde{\ell}_{0}\right) & =\dot{\gamma}\left(\tilde{\ell}_{0} ; y_{0}, \eta_{0}\right)=\dot{\gamma}\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right) .
\end{aligned}
$$

Hence

$$
\dot{\sigma}(0)=\left.\partial_{s} \Gamma(s, \lambda(s))\right|_{s=0}=\left(\partial_{s} \Gamma\right)\left(0, \tilde{\ell}_{0}\right)+\left(\partial_{t} \Gamma\right)\left(0, \tilde{\ell}_{0}\right) \lambda^{\prime}(0)
$$

The curve $t \mapsto \Gamma(s, t)$ is a unit speed geodesic for all $s \in(-\epsilon, \epsilon)$. Hence

$$
\begin{align*}
\left(\dot{\sigma}(0), \dot{\gamma}\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)\right)_{g} & =\left.\left(\left(\partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g}+\lambda^{\prime}(0)\left(\partial_{t} \Gamma, \partial_{t} \Gamma\right)_{g}\right)\right|_{s=0, t=\tilde{\ell}_{0}}  \tag{27}\\
& =\left.\left(\partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g}\right|_{s=0, t=\tilde{\ell}_{0}}+\lambda^{\prime}(0)
\end{align*}
$$

We define

$$
\mathscr{L}(s, l):=\int_{0}^{l}\left|\partial_{t} \Gamma(s, t)\right|_{g} d t, \quad(s, l) \in(-\epsilon, \epsilon) \times(0, \infty)
$$

Then $\mathscr{L}(s, l), s \in(-\epsilon, \epsilon)$ is the length of a unit speed geodesic on the interval $[0, l]$. Thus $\mathscr{L}(s, l)=l$ for all $s \in(-\epsilon, \epsilon)$. We may derive an expression for $\left.\partial_{s} \mathscr{L}(s, l)\right|_{s=0}$ as in [Lee 1997, Proposition 6.5]:

$$
\left.\partial_{s} \mathscr{L}(s, l)\right|_{s=0}=\left.\int_{0}^{l}\left(D_{t} \partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g} d t\right|_{s=0}
$$

As $t \mapsto \Gamma(s, t)$ is a geodesic, $D_{t} \partial_{t} \Gamma(s, t)=0$, and thus

$$
\partial_{t}\left(\partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g}=\left(D_{t} \partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g}
$$

Moreover, $\Gamma(s, 0)=y_{0}$ for all $s \in(-\epsilon, \epsilon)$, and thus $\partial_{s} \Gamma(s, 0)=0$. Hence

$$
0=\left.\partial_{s} \mathscr{L}(s, l)\right|_{s=0}=\left.\int_{0}^{l} \partial_{t}\left(\partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g} d t\right|_{s=0}=\left.\left(\partial_{s} \Gamma, \partial_{t} \Gamma\right)_{g}\right|_{s=0, t=l}, \quad l \in(0, \infty) .
$$

By (27), we have

$$
\begin{aligned}
\left(\dot{\sigma}(0), \gamma_{z_{0}}^{\top}\right)_{g} & =\left(\dot{\sigma}(0), \dot{\gamma}\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)\right)_{g} \\
& =\lambda^{\prime}(0)=\left\langle\left. d(\tilde{\ell} \circ F)\right|_{z=z_{0}}, \dot{\sigma}(0)\right\rangle_{T_{z_{0}}^{*} \partial M \times T_{z_{0}} \partial M}=\left(\dot{\sigma}(0),\left.\operatorname{grad}_{\partial M}(\tilde{\ell} \circ F)\right|_{z=z_{0}}\right)_{g}
\end{aligned}
$$

for all smooth curves $\sigma$ in $W$ such that $\sigma(0)=z_{0}$, which proves the claim.
Theorem 7.8. The functions $\tau: \partial_{-} S M \rightarrow(0, \infty]$ and

$$
z: D(\Sigma) \rightarrow \partial M, \quad z(x, \xi):=\gamma(\tau(x, \xi) ; x, \xi)
$$

together with the Riemannian manifold $\left(\Omega,\left.g\right|_{\Omega}\right)$, determine

$$
\dot{\gamma}(\tau(x, \xi) ; x, \xi), \quad(x, \xi) \in D(\Sigma)
$$

Proof. The functions $\tau$ and $z$ on $D(\Sigma)$ determine the set $B$ of points $\left(x_{0}, \xi_{0}\right) \in D(\Sigma)$ such that the conditions (C1)-(C4) of Lemma 7.6 hold for some $\left(y_{0}, \eta_{0}\right) \in S \Omega$.

Let $\left(x_{0}, \xi_{0}\right) \in B$. We write $\zeta_{0}:=\dot{\gamma}\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)$. The map

$$
\eta \mapsto z(x(\eta), \xi(\eta))
$$

determines its local inverse. Hence $\tau$ and $z$ determine the function $F$ of Lemma 7.7, and thus they determine $\dot{\gamma}_{z_{0}}^{\top}$ by the formula (26). As $\zeta_{0}$ is a unit vector,

$$
\zeta_{0}=\dot{\gamma}_{z_{0}}^{\top}+\left(1-\left|\dot{\gamma}_{z_{0}}^{\top}\right|^{2}\right)^{1 / 2} v_{z_{0}}
$$

where $\nu_{z_{0}}$ is the unit exterior normal vector of $\partial M$. Hence $\tau$ and $z$ determine $\zeta_{0}$ for all $\left(x_{0}, \xi_{0}\right) \in B$.
Let $\left(x_{0}, \xi_{0}\right) \in D(\Sigma)$. By Lemmata 7.5 and 7.6 , there is a sequence $\left(\left(x_{j}, \xi_{j}\right)\right)_{j=1}^{\infty} \subset B$ such that

$$
\lim _{j \rightarrow \infty}\left(x_{j}, \xi_{j}\right)=\left(x_{0}, \xi_{0}\right), \quad \lim _{j \rightarrow \infty} \tau\left(x_{j}, \xi_{j}\right)=\tau\left(x_{0}, \xi_{0}\right)
$$

Moreover, the functions $\tau$ and $z$ determine the set of such sequences, and thus they determine

$$
\lim _{j \rightarrow \infty} \dot{\gamma}\left(\tau\left(x_{j}, \xi_{j}\right) ; x_{j}, \xi_{j}\right)=\dot{\gamma}\left(\tau\left(x_{0}, \xi_{0}\right) ; x_{0}, \xi_{0}\right)
$$

Theorems 6.5, 7.2 and 7.8 prove Theorem 1.2 formulated in the introduction.

## References

[Amann 1990] H. Amann, Ordinary differential equations: an introduction to nonlinear analysis, Studies in Mathematics 13, de Gruyter, Berlin, 1990. MR 91e:34001 Zbl 0708.34002
[Anderson et al. 2004] M. Anderson, A. Katsuda, Y. Kurylev, M. Lassas, and M. Taylor, "Boundary regularity for the Ricci equation, geometric convergence, and Gel'fand's inverse boundary problem", Invent. Math. 158:2 (2004), 261-321. MR 2005h:53051 Zbl 1177.35245
[Astala and Päivärinta 2006] K. Astala and L. Päivärinta, "Calderón's inverse conductivity problem in the plane", Ann. of Math. (2) 163:1 (2006), 265-299. MR 2007b:30019 Zbl 1111.35004
[Astala et al. 2005] K. Astala, L. Päivärinta, and M. Lassas, "Calderón's inverse problem for anisotropic conductivity in the plane", Comm. Partial Differential Equations 30:1-3 (2005), 207-224. MR 2005k:35421 Zbl 1129.35483
[Babich and Ulin 1981] V. M. Babich and V. V. Ulin, "The complex space-time ray method and 'quasiphotons"', pp. 5-12 in Математические вопросы теории распространения волн 12, Zap. Nauchn. Sem. Leningrad. Otdel. Mat. Inst. Steklov 117, 1981. In Russian; translated in J. Soviet Math. 24:3 (1984), 269-273. MR 83j:35094 Zbl 0477.35025
[Babich et al. 1985] V. M. Babich, V. S. Buldyrev, and I. A. Molotkov, Пространственно-временной лучевой метод: Линеиные и нелинеиные волны, Leningrad Univ., 1985. MR 88e:35002
[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 94a:58199 Zbl 0812.58094
[Bellassoued and Yamamoto 2008] M. Bellassoued and M. Yamamoto, "Determination of a coefficient in the wave equation with a single measurement", Appl. Anal. 87:8 (2008), 901-920. MR 2009m:35533 Zbl 1149.35404
[Bukhgeim 2008] A. L. Bukhgeim, "Recovering a potential from Cauchy data in the two-dimensional case", J. Inverse Ill-Posed Probl. 16:1 (2008), 19-33. MR 2008m:30049 Zbl 1142.30018
[Bukhgeim and Klibanov 1981] A. L. Bukhgeim and M. V. Klibanov, "Uniqueness in the large of a class of multidimensional inverse problems", Dokl. Akad. Nauk SSSR 260:2 (1981), 269-272. In Russian; translated in Sov. Math. Dokl. 24 (1981), 224-227. MR 83b:35157 Zbl 0497.35082
[Burago and Ivanov 2010] D. Burago and S. Ivanov, "Boundary rigidity and filling volume minimality of metrics close to a flat one", Ann. of Math. (2) 171:2 (2010), 1183-1211. MR 2011d:53079 Zbl 1192.53048
[Calderón 1980] A.-P. Calderón, "On an inverse boundary value problem", pp. 65-73 in Seminar on Numerical Analysis and its Applications to Continuum Physics (Rio de Janeiro, 1980), Soc. Brasil. Mat., Rio de Janeiro, 1980. MR 81k:35160 Zbl 1182.35230
[Candès et al. 2006] E. J. Candès, J. K. Romberg, and T. Tao, "Stable signal recovery from incomplete and inaccurate measurements", Comm. Pure Appl. Math. 59:8 (2006), 1207-1223. MR 2007f:94007 Zbl 1098.94009
[Dairbekov and Uhlmann 2010] N. Dairbekov and G. Uhlmann, "Reconstructing the metric and magnetic field from the scattering relation", Inverse Probl. Imaging 4:3 (2010), 397-409. MR 2011i:53042 Zbl 1202.53042
[Donoho 2006] D. L. Donoho, "Compressed sensing", IEEE Trans. Inform. Theory 52:4 (2006), 1289-1306. MR 2007e:94013
[Evans 1998] L. C. Evans, Partial differential equations, Graduate Studies in Mathematics 19, American Mathematical Society, Providence, RI, 1998. MR 99e:35001 Zbl 0902.35002
[Greenleaf et al. 2009a] A. Greenleaf, Y. Kurylev, M. Lassas, and G. Uhlmann, "Cloaking devices, electromagnetic wormholes, and transformation optics", SIAM Rev. 51:1 (2009), 3-33. MR 2010b:35484 Zbl 1158.78004
[Greenleaf et al. 2009b] A. Greenleaf, Y. Kurylev, M. Lassas, and G. Uhlmann, "Invisibility and inverse problems", Bull. Amer. Math. Soc. (N.S.) 46:1 (2009), 55-97. MR 2010d:35399 Zbl 1159.35074
[Guillarmou and Tzou 2010] C. Guillarmou and L. Tzou, "Calderón inverse problem for the Schrödinger operator on Riemann surfaces", pp. 129-141 in The AMSI-ANU Workshop on Spectral Theory and Harmonic Analysis, edited by A. Hassell et al., Proc. Centre Math. Appl. Austral. Nat. Univ. 44, Austral. Nat. Univ., Canberra, 2010. MR 2011k:58031 Zbl 1231.35302
[Guillarmou and Tzou 2011] C. Guillarmou and L. Tzou, "Calderón inverse problem with partial data on Riemann surfaces", Duke Math. J. 158:1 (2011), 83-120. MR 2012f:35574 Zbl 1222.35212
[Hirsch 1976] M. W. Hirsch, Differential topology, Graduate Texts in Mathematics 33, Springer, New York, 1976. MR 56 \#6669 Zbl 0356.57001
[Hörmander 1985] L. Hörmander, The analysis of linear partial differential operators, III: Pseudodifferential operators, vol. 274, Grundlehren der Math. Wiss., Springer, Berlin, 1985. Reprinted 1994, 2007. MR 87d:35002a Zbl 0601.35001
[Imanuvilov and Yamamoto 2003] O. Y. Imanuvilov and M. Yamamoto, "Determination of a coefficient in an acoustic equation with a single measurement", Inverse Problems 19:1 (2003), 157-171. MR 2004c:35415 Zbl 1020.35117
[Imanuvilov et al. 2010] O. Imanuvilov, G. Uhlmann, and M. Yamamoto, "Partial Cauchy data for general second order operators in two dimensions", preprint, 2010. arXiv 1010.5791
[Isakov 2006] V. Isakov, Inverse problems for partial differential equations, 2nd ed., Applied Mathematical Sciences 127, Springer, New York, 2006. MR 2006h:35279 Zbl 1092.35001
[Jost 2008] J. Jost, Riemannian geometry and geometric analysis, 5th ed., Universitext, Springer, Berlin, 2008. MR 2009g:53036 Zbl 1143.53001
[Katchalov and Kurylev 1998] A. Katchalov and Y. Kurylev, "Multidimensional inverse problem with incomplete boundary spectral data", Comm. Partial Differential Equations 23:1-2 (1998), 55-95. MR 99b:35213 Zbl 0904.65114
[Katchalov et al. 2001] A. Katchalov, Y. Kurylev, and M. Lassas, Inverse boundary spectral problems, Monographs and Surveys in Pure and Applied Mathematics 123, Chapman \& Hall/CRC, Boca Raton, FL, 2001. MR 2003e:58045 Zbl 1037.35098
[Katchalov et al. 2004] A. Katchalov, Y. Kurylev, M. Lassas, and N. Mandache, "Equivalence of time-domain inverse problems and boundary spectral problems", Inverse Problems 20:2 (2004), 419-436. MR 2005d:35270 Zbl 1073.35209
[Kenig et al. 2007] C. E. Kenig, J. Sjöstrand, and G. Uhlmann, "The Calderón problem with partial data", Ann. of Math. (2) 165:2 (2007), 567-591. MR 2008k:35498 Zbl 1127.35079
[Klibanov 1992] M. V. Klibanov, "Inverse problems and Carleman estimates", Inverse Problems 8:4 (1992), 575-596. MR 93h:35210 Zbl 0755.35151
[Kurylev et al. 2010] Y. Kurylev, M. Lassas, and G. Uhlmann, "Rigidity of broken geodesic flow and inverse problems", Amer. J. Math. 132:2 (2010), 529-562. MR 2011g:53179 Zbl 1206.53073
[Ladyzhenskaya 1985] O. A. Ladyzhenskaya, The boundary value problems of mathematical physics, Applied Mathematical Sciences 49, Springer, New York, 1985. MR 87f:35001 Zbl 0588.35003
[Lassas 1995] M. Lassas, "Non-selfadjoint inverse spectral problems and their applications to random bodies", Ann. Acad. Sci. Fenn. Math. Diss. 103 (1995), 108. MR 97d:47058 Zbl 0856.35131
[Lassas 1998] M. Lassas, "Inverse boundary spectral problem for non-selfadjoint Maxwell's equations with incomplete data", Comm. Partial Differential Equations 23:3-4 (1998), 629-648. MR 99c:35174 Zbl 0906.35116
[Lassas and Uhlmann 2001] M. Lassas and G. Uhlmann, "On determining a Riemannian manifold from the Dirichlet-to-Neumann map", Ann. Sci. École Norm. Sup. (4) 34:5 (2001), 771-787. MR 2003e:58037 Zbl 0992.35120
[Lassas et al. 2003] M. Lassas, M. Taylor, and G. Uhlmann, "The Dirichlet-to-Neumann map for complete Riemannian manifolds with boundary", Comm. Anal. Geom. 11:2 (2003), 207-221. MR 2004h:58033 Zbl 1077.58012
[Lee 1997] J. M. Lee, Riemannian manifolds. An introduction to curvature, Graduate Texts in Mathematics 176, Springer, New York, 1997. MR 98d:53001 Zbl 0905.53001
[Lee and Uhlmann 1989] J. M. Lee and G. Uhlmann, "Determining anisotropic real-analytic conductivities by boundary measurements", Comm. Pure Appl. Math. 42:8 (1989), 1097-1112. MR 91a:35166 Zbl 0702.35036
[Michel 1981] R. Michel, "Sur la rigidité imposée par la longueur des géodésiques", Invent. Math. 65:1 (1981), 71-83. MR 83d:58021 Zbl 0471.53030
[Muhometov 1977] R. G. Muhometov, "The reconstruction problem of a two-dimensional Riemannian metric, and integral geometry", Dokl. Akad. Nauk SSSR 232:1 (1977), 32-35. In Russian; translated in Sov. Math. Dokl., 18:1 (1977), 27-31. MR 55 \#4076
[Muhometov 1981] R. G. Muhometov, "On a problem of reconstructing Riemannian metrics", Sibirsk. Mat. Zh. 22:3 (1981), 119-135, 237. MR 82m:53071
[Muhometov and Romanov 1978] R. G. Muhometov and V. G. Romanov, "On the problem of finding an isotropic Riemannian metric in an $n$-dimensional space", Dokl. Akad. Nauk SSSR 243:1 (1978), 41-44. In Russian; translated in Sov. Math. Dokl. 19:6 (1978), 1330-1333. MR 81a:53059
[Nachman 1988] A. I. Nachman, "Reconstructions from boundary measurements", Ann. of Math. (2) 128:3 (1988), 531-576. MR 90i:35283 Zbl 0675.35084
[Nachman 1996] A. I. Nachman, "Global uniqueness for a two-dimensional inverse boundary value problem", Ann. of Math. (2) 143:1 (1996), 71-96. MR 96k:35189 Zbl 0857.35135
[Pestov and Uhlmann 2005] L. Pestov and G. Uhlmann, "Two dimensional compact simple Riemannian manifolds are boundary distance rigid", Ann. of Math. (2) 161:2 (2005), 1093-1110. MR 2006c:53038 Zbl 1076.53044
[Pestov and Uhlmann 2006] L. Pestov and G. Uhlmann, "The scattering relation and the Dirichlet-to-Neumann map", pp. 249-262 in Recent advances in differential equations and mathematical physics, edited by N. Chernov et al., Contemp. Math. 412, Amer. Math. Soc., Providence, RI, 2006. MR 2007k:53055 Zbl 1112.35132
[Rakesh 2003] Rakesh, "An inverse problem for a layered medium with a point source", Inverse Problems 19:3 (2003), 497-506. MR 2004g:35234 Zbl 1033.35141
[Rakesh 2008] Rakesh, "Inverse problems for the wave equation with a single coincident source-receiver pair", Inverse Problems 24:1 (2008), 015012, 16. MR 2009c:35479 Zbl 1151.35104
[Rakesh and Sacks 2011] Rakesh and P. Sacks, "Uniqueness for a hyperbolic inverse problem with angular control on the coefficients", J. Inverse Ill-Posed Probl. 19:1 (2011), 107-126. MR 2012c:35480
[Ralston 1982] J. Ralston, "Gaussian beams and the propagation of singularities", pp. 206-248 in Studies in partial differential equations, edited by W. Littman, MAA Stud. Math. 23, Math. Assoc. America, Washington, DC, 1982. MR 85c:35052 Zbl 0533.35062
[Ramm 2001] A. G. Ramm, "A non-overdetermined inverse problem of finding the potential from the spectral function", Int. J. Differ. Equ. Appl. 3:1 (2001), 15-29. MR 2002f:35223 Zbl 1048.35137
[Romanov 1987] V. G. Romanov, Inverse problems of mathematical physics, VNU Science Press, Utrecht, 1987. MR 88b:35203 Zbl 0576.35001
[Romanov 2002] V. G. Romanov, Investigation methods for inverse problems, Inverse and Ill-posed Problems Series, VSP, Utrecht, 2002. MR 2005a:35281 Zbl 1038.35001
[Sacks and Symes 1985] P. Sacks and W. Symes, "Uniqueness and continuous dependence for a multidimensional hyperbolic inverse problem", Comm. Partial Differential Equations 10:6 (1985), 635-676. MR 87f:35146 Zbl 0574.35078
[Shubin 1992] M. A. Shubin, "Spectral theory of elliptic operators on noncompact manifolds", pp. 35-108 in Méthodes semi-classiques, I (Nantes, 1991), Astérisque 207, 1992. MR 94h:58175 Zbl 0793.5803
[Stefanov and Uhlmann 2005] P. Stefanov and G. Uhlmann, "Boundary rigidity and stability for generic simple metrics", $J$. Amer. Math. Soc. 18:4 (2005), 975-1003. MR 2006h:53031 Zbl 1079.53061
[Stefanov and Uhlmann 2008] P. Stefanov and G. Uhlmann, "Boundary and lens rigidity, tensor tomography and analytic microlocal analysis", pp. 275-293 in Algebraic analysis of differential equations from microlocal analysis to exponential asymptotics (Kyoto, 2005), edited by A. Takashi et al., Springer, Tokyo, 2008. MR 2012c:58046 Zbl 1138.53039
[Stefanov and Uhlmann 2009] P. Stefanov and G. Uhlmann, "Local lens rigidity with incomplete data for a class of non-simple Riemannian manifolds", J. Differential Geom. 82:2 (2009), 383-409. MR 2011d:53081 Zbl 1247.53049
[Stefanov and Uhlmann 2011] P. Stefanov and G. Uhlmann, "Recovery of a source term or a speed with one measurement and applications", preprint, 2011. arXiv 1103.1097
[Sylvester 1990] J. Sylvester, "An anisotropic inverse boundary value problem", Comm. Pure Appl. Math. 43:2 (1990), 201-232. MR 90m:35202 Zbl 0709.35102
[Sylvester and Uhlmann 1987] J. Sylvester and G. Uhlmann, "A global uniqueness theorem for an inverse boundary value problem", Ann. of Math. (2) 125:1 (1987), 153-169. MR 88b:35205 Zbl 0625.35078
[Toomay and Hannen 2004] J. C. Toomay and P. Hannen, Radar principles for the non-specialist, 3rd ed., SciTech Publishing, Herndon, VA, 2004.
[Trèves 1967] F. Trèves, Topological vector spaces, distributions and kernels, Academic Press, New York, 1967. MR 37 \#726 Zbl 0171.10402
[Triebel 1978] H. Triebel, Interpolation theory, function spaces, differential operators, VEB Verlag, Berlin, 1978. MR 80i:46032a Zbl 0387.46033
[Uhlmann 2003] G. Uhlmann, "The Cauchy data and the scattering relation", pp. 263-287 in Geometric methods in inverse problems and PDE control, edited by C. Croke et al., IMA Vol. Math. Appl. 137, Springer, New York, 2003. MR 2006f:58032 Zbl 1061.35175

Received 10 Nov 2010. Accepted 26 May 2011.
TAPIO HELIN: tapio.helin@oeaw.ac.at
Johann Radon Institute for Computational and Applied Mathematics (RICAM), Austrian Academy of Sciences, Altenbergerstrasse 69, A-4040 Linz, Austria
Matti Lassas: Matti.Lassas@helsinki.fi
Department of Mathematics and Statistics, University of Helsinki, P.O. Box 68 (Gustaf Hällströmin katu 2b), FI-00014 Helsinki, Finland

LAURI OKSANEN: lauri.oksanen@helsinki.fi
Department of Mathematics and Statistics, University of Helsinki, P.O. Box 68 (Gustaf Hällströmin katu 2b), FI-00014 Helsinki, Finland

# Analysis \& PDE 

msp.berkeley.edu/apde

## EDITORS

| Editor-IN-CHIEF |  |  |  |
| :---: | :---: | :---: | :---: |
| Maciej Zworski |  |  |  |
|  | University of Berkeley | of California ley, USA |  |
| Board of Editors |  |  |  |
| Michael Aizenman | Princeton University, USA aizenman@math.princeton.edu | Nicolas Burq | Université Paris-Sud 11, France nicolas.burq@math.u-psud.fr |
| Luis A. Caffarelli | University of Texas, USA caffarel@math.utexas.edu | un-Yung Alice Chang | Princeton University, USA chang@math.princeton.edu |
| Michael Christ | University of California, Berkeley, USA mchrist@math.berkeley.edu | Charles Fefferman | Princeton University, USA cf@math.princeton.edu |
| Ursula Hamenstaedt | Universität Bonn, Germany ursula@math.uni-bonn.de | Nigel Higson | Pennsylvania State Univesity, USA higson@math.psu.edu |
| Vaughan Jones | University of California, Berkeley, USA vfr@math.berkeley.edu | Herbert Koch | Universität Bonn, Germany koch@math.uni-bonn.de |
| Izabella Laba | University of British Columbia, Canada ilaba@math.ubc.ca | Gilles Lebeau | Université de Nice Sophia Antipolis, France lebeau@unice.fr |
| László Lempert | Purdue University, USA lempert@math.purdue.edu | Richard B. Melrose | Massachussets Institute of Technology, USA rbm@math.mit.edu |
| Frank Merle | Université de Cergy-Pontoise, France Frank.Merle@u-cergy.fr | William Minicozzi II | Johns Hopkins University, USA minicozz@math.jhu.edu |
| Werner Müller | Universität Bonn, Germany mueller@math.uni-bonn.de | Yuval Peres | University of California, Berkeley, USA peres@stat.berkeley.edu |
| Gilles Pisier | Texas A\&M University, and Paris 6 pisier@math.tamu.edu | Tristan Rivière | ETH, Switzerland riviere@math.ethz.ch |
| Igor Rodnianski | Princeton University, USA irod@math.princeton.edu | Wilhelm Schlag | University of Chicago, USA schlag@math.uchicago.edu |
| Sylvia Serfaty | New York University, USA serfaty@cims.nyu.edu | Yum-Tong Siu | Harvard University, USA siu@math.harvard.edu |
| Terence Tao | University of California, Los Angeles, USA tao@math.ucla.edu | A Michael E. Taylor | Univ. of North Carolina, Chapel Hill, USA met@math.unc.edu |
| Gunther Uhlmann | University of Washington, USA gunther@math.washington.edu | András Vasy | Stanford University, USA andras@math.stanford.edu |
| Dan Virgil Voiculescu | University of California, Berkeley, USA dvv@math.berkeley.edu | Steven Zelditch | Northwestern University, USA zelditch@math.northwestern.edu |

## PRODUCTION

production@msp.org

$$
\text { Silvio Levy, Scientific Editor } \quad \text { Sheila Newbery, Senior Production Editor }
$$

## See inside back cover or msp.berkeley.edu/apde for submission instructions

The subscription price for 2012 is US $\$ 140 /$ year for the electronic version, and $\$ 240 /$ year for print and electronic. Subscriptions, requests for back issues from the last three years and changes of subscribers address should be sent to Mathematical Sciences Publishers, Department of Mathematics, University of California, Berkeley, CA 94720-3840, USA.

Analysis \& PDE, at Mathematical Sciences Publishers, Department of Mathematics, University of California, Berkeley, CA 94720-3840 is published continuously online. Periodical rate postage paid at Berkeley, CA 94704, and additional mailing offices.

APDE peer review and production are managed by EditFLOw ${ }^{\mathrm{TM}}$ from Mathematical Sciences Publishers.

## PUBLISHED BY

E. mathematical sciences publishers
http://msp.org/
A NON-PROFIT CORPORATION
Typeset in LATEX
Copyright ©2012 by Mathematical Sciences Publishers

## ANALYSIS \& PDE

Volume 5 No. 52012
An inverse problem for the wave equation with one measurement and the pseudorandom ..... 887sourceTapio Helin, Matti Lassas and Lauri Oksanen
Two-dimensional nonlinear Schrödinger equation with random radial data ..... 913
Yu Deng
Schrödinger operators and the distribution of resonances in sectors ..... 961
Tanya J. Christiansen
Weighted maximal regularity estimates and solvability of nonsmooth elliptic systems, II ..... 983
Pascal Auscher and Andreas Rosén
The two-phase Stefan problem: regularization near Lipschitz initial data by phase dynamics ..... 1063
Sunhi Choi and Inwon Kim
$C^{\infty}$ spectral rigidity of the ellipse ..... 1105
Hamid Hezari and Steve Zelditch
A natural lower bound for the size of nodal sets ..... 1133
Hamid Hezari and Christopher D. Sogge
Effective integrable dynamics for a certain nonlinear wave equation ..... 1139
Patrick Gérard and Sandrine Grellier
Nonlinear Schrödinger equation and frequency saturation ..... 1157
Rémi Carles

