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INVERSE SCATTERING FOR A RANDOM POTENTIAL

PEDRO CARO*, TAPIO HELIN[◦] AND MATTI LASSAS[◦]

ABSTRACT. In this paper we consider an inverse problem for the n -dimensional random Schrödinger equation $(\Delta - q + k^2)u = 0$. We study the scattering of plane waves in the presence of a potential q which is assumed to be a Gaussian random function such that its covariance is described by a pseudodifferential operator. Our main result is as follows: given the backscattered far field, obtained from a single realization of the random potential q , we uniquely determine the principal symbol of the covariance operator of q . Especially, for $n = 3$ this result is obtained for the full non-linear inverse backscattering problem. Finally, we present a physical scaling regime where the method is of practical importance.

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1. INTRODUCTION

In inverse scattering theory one aims at determining an electric potential q in $\mathbb{R}^n$ with $n \geq 2$ from measurements describing how this scatters certain incoming waves. In many applications, the scatterer is so rough and vastly complicated that there is an apparent lack of systematic patterns in its micro-scale structure. In these situations, the potential is assumed to be created by a physical random process and the goal is not any more to recover the full potential but to determine some parameters or functions describing properties of its micro-structure. In this paper, we are interested in reconstructing statistical properties, more precisely, *the local strength* of a potential q .

In the usual mathematical approach to the inverse scattering theory, one considers the scattering problem

$$(1) \quad \begin{cases} (\Delta - q(x) + k^2)u(x) = 0 & \text{in } \mathbb{R}^n \\ u(x) = e^{ik\theta \cdot x} + u_{sc}(x) \\ u_{sc}(x) \text{ satisfies the Sommerfeld radiation condition,} \end{cases}$$

where the *incident wave* is assumed to be the plane wave $e^{ik\theta \cdot x}$ and the *scattered* and *total* waves are denoted by u_{sc} and u , respectively. The scattered wave satisfies the following asymptotic expansion

$$u_{sc}(x) = c_n k^{\frac{n-1}{2}} |x|^{-\frac{n-1}{2}} e^{ik|x|} u^\infty \left(k, \theta, \frac{x}{|x|} \right) + o \left(|x|^{-\frac{n-1}{2}} \right),$$

where u^∞ is known as the *far-field pattern* of u_{sc} . In this context, the inverse backscattering problem aims at answering the question:

(Q) *Given the backscattered far-field pattern $u^\infty(k, \theta, -\theta)$ for multiple values of $k > 0$ and $\theta \in \mathbb{S}^{n-1}$, what kind of information of q can be recovered?*

The deterministic inverse back-scattering problem—which asks whether a potential $q(x)$ can be uniquely determined from its backscattered far-field pattern $u^\infty(k, \theta, -\theta)$ —is a long-standing open problem. At this moment, the problem has been solved only under assumptions on controlled angular regularity of the potential (see [45]). We discuss below the literature about this problem. In this paper we consider a related stochastic inverse problem where the statistical parameter functions of the potential q are determined from the observations.

Since we are interested in the situations where the scatterer presents a random behaviour, we need to rephrase our approach to the inverse scattering theory. In order to do so, we assume the potential $\omega \in \Omega \mapsto q(x, \omega)$ to be a generalized random function in a probability space $(\Omega, \mathcal{H}, \mathbb{P})$. This makes the far-field pattern be random as well, which means that it changes with each realization $q(x, \omega)$. However, our approach consists of assuming the backscattering data $u^\infty(k, \theta, -\theta)$ with $k > 0$ and $\theta \in \mathbb{S}^{n-1}$ to be generated by a single realization $q(x, \omega_0)$ for certain $\omega_0 \in \Omega$. Then, the inverse backscattering problem in this context asks to determine the parameters characterizing the probability law of q from the backscattering data.

As we previously advanced, we reconstruct in this paper the local strength of the potential, which is one of the parameters describing the probability law of q . In order to provide an interpretation of this parameter, we will need to make some general assumptions on q . Firstly,

we assume q to be a generalized Gaussian field supported in a bounded domain D and its expected potential $\mathbb{E}q$ to be a smooth function. Additionally, we assume the covariance function $K_q(x, y)$ to be smooth out of the diagonal, which means that the long distance interactions depends smoothly on their locations; we also assume the average roughness (or smoothness) of q to remain unchanged for every sub-domain of D . However, we allow the size of this roughness to change in different sub-domains of D . The local strength of the potential measures or controls these different sizes. These assumptions¹ can be rigorously introduced, assuming that the covariance operator C_q is a classical pseudodifferential operator (see for example [31]) of order $-m$ with $m > n - 1$, and such that, C_q has

$$(2) \quad \sigma(x, \xi) = \frac{\mu(x)}{|\xi|^m}$$

as a principal symbol, with μ a smooth non-negative function supported on D —called the *local strength* of the potential. As we will see in Definition 2.1, this is to say that q is a Gaussian *microlocally isotropic* random field. Eventually, assuming q as above with μ unknown, our goal will be to reconstruct μ from the backscattering data. As we will see in the section 2.2, μ yields valuable control on the oscillations of q : where μ is large, the rough oscillations of q are most likely large as well.

In order to cover a broad spectrum of well-known random field models, we also include the possibility of realizations of q being generalized functions almost surely (a.s. for short). In fact, whenever $n - 1 < m \leq n$, we can only ensure that q belongs to a Sobolev space with negative smoothness index almost surely. This consideration requires a carefully analysis² of the forward problem with compactly supported potentials in the Sobolev spaces $L^p_{-s}(\mathbb{R}^n) = W^{-s,p}(\mathbb{R}^n)$ with $0 < s \leq 1/2$ and $n/s \leq p < \infty$. Inspired by the works [1, 36], we provide new insights to the classical scattering theory for rough potentials.

A microlocally isotropic Gaussian random field q of order $-m$ in $D \subset \mathbb{R}^n$ can be written in the form $q = (C_q)^{1/2}W$, where W is a white noise. We will later see that $q \in L^p_{-s}(\mathbb{R}^n)$ almost surely for any $1 < p < \infty$ and $-s < (m - n)/2$. The local strength μ determines the roughest component of q in the sense that if $C_{\tilde{q}}$ is a properly classical pseudodifferential operator of order $-m$ having the same principal symbol (2) as the operator C_q , then $\tilde{q} = (C_{\tilde{q}})^{1/2}W$ is also a microlocally isotropic Gaussian random field of order $-m$ such that $\tilde{q} - q \in L^p_{-s+1}(\mathbb{R}^n)$ almost surely for any $1 < p < \infty$ and $-s < (m - n)/2$, that is, $\tilde{q} - q$ is one degree smoother than q .

In applications the measurement data is often obtained as an average of signals at multiple frequencies. Also, in many standard references in the literature in applied sciences, see e.g. [33, 39, 43, 50], one considers the effective equations for the expectations and covariances of the scattered waves. This means that one considers the averages of waves that are generated by many independent samples of the scatterers. This approach can be poorly justified if the scatterer changes slowly during the measurements or is independent of time. In this paper,

¹The random model is discussed in detail in the section 2.1.

²The readers who are expert on uniqueness for the Calderón problem for non-regular conductivities will notice the connection with the work [16] of Brown and Torres.

the data is assumed to be a weighted average of far-field patterns at a given separation $\tau > 0$:

$$(3) \quad M(\tau, \theta) = \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m u^\infty(k, \theta, -\theta) \overline{u^\infty(k + \tau, \theta, -\theta)} dk,$$

where $\theta \in \mathbb{S}^{n-1}$. As we pointed out before, the far-field pattern is random and consequently our data is random as well. Since we want to show that the data generated by a single realization of q allows us to reconstruct the local strength μ —which is non-random, we will need to prove that the randomness averages out at the limit. Actually, we will prove that, for $n = m = 3$, there exists a known constant $c > 0$ such that

$$(4) \quad M(\tau, \theta) = c \hat{\mu}(2\tau\theta) \quad \text{almost surely.}$$

Formula (4) means that the measurement function $M(\tau, \theta)$, computed from the measured far-field patterns, does not depend (with probability one) on the realization of the random potential q . Such measurement functions that are independent of the realization of the random media are said to be *statistically stable*, see [15, 23]. The study of statistically stable measurement functions have turned to be very useful in particular in the study of inverse source problems in random medium background [4, 11, 12, 13, 14, 15, 23].

Despite the non-linear terms generated by the 2nd order and multiple scattering, it is interesting to note the linear relation between the data and the local strength μ in (2). This suggests that, whenever (4) holds, the local strength of the Born approximation of q equals the local strength of the full potential q .

Let us now formulate the main theorem of this paper.

Theorem 1.1. *Let q be a Gaussian microlocally isotropic random field of order -3 in $D \subset \mathbb{R}^3$. Then, the measurement data $\{M(\tau_j, \theta_j) : j \in \mathbb{N}\}$, with $\{(\tau_j, \theta_j) : j \in \mathbb{N}\}$ any dense subset of $\mathbb{R}_+ \times \mathbb{S}^2$, determines the local strength μ almost surely.*

Theorem 1.1 can be interpreted as follows: We consider a complicated potential $q(x, \omega_0)$ that is assumed to be created, before the measurements are made, by a random process, that is, the potential $q(x, \omega_0)$ is a single realization of the Gaussian microlocally isotropic random $q(x, \omega)$. We show the measurement data $\{M(\tau_j, \theta_j) : j \in \mathbb{N}\}$, obtained from this single realization of the process, determine with probability one the principal symbol of the covariance operator of the random process $q(x, \omega)$.

In order to prove Theorem 1.1, we explore the limit in (3) by separating the effects of different orders of scattering in the Born series

$$u^\infty(k, \theta, -\theta) = \sum_{j=1}^{\infty} u_j^\infty(k, \theta, -\theta),$$

where u_j^∞ describes the far field of the j -th order scattering. Under the assumptions of Theorem 1.1, it can be rigorously shown that the interactions coming from the simple backscattering, i.e.,

$$M_1(\tau, \theta) = \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)} dk$$

coincides with $M(\tau, \theta)$ almost surely, whereas the contribution from 2nd order and multiple backscattering becomes negligible at the limit:

$$\lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m u_j^\infty(k, \theta, -\theta) \overline{u_l^\infty(k + \tau, \theta, -\theta)} dk = 0 \quad \text{for } j + l \geq 3.$$

In particular, there exists a known constant $c > 0$ such that

$$M_1(\tau, \theta) = c \hat{\mu}(2\tau\theta) \quad \text{almost surely.}$$

Our work also has direct implications to the problem in any dimension if the interaction of 2nd order and multiple backscattering can be a priori neglected.

Theorem 1.2. *Let q be a Gaussian microlocally isotropic random field of order $-m$ in $D \subset \mathbb{R}^n$ with $m > n - 1$ and $n \geq 2$. Then, the data $\{M_1(\tau_j, \theta_j) : j \in \mathbb{N}\}$, with $\{(\tau_j, \theta_j) : j \in \mathbb{N}\}$ any dense subset of $\mathbb{R}_+ \times \mathbb{S}^{n-1}$, determines the local strength μ almost surely.*

This second theorem suggests that for general dimensions $n \geq 2$, and $m > n - 1$, the measurement data $\{M(\tau_j, \theta_j) : j \in \mathbb{N}\}$, with $\{(\tau_j, \theta_j) : j \in \mathbb{N}\}$ any dense subset of $\mathbb{R}_+ \times \mathbb{S}^{n-1}$, determines the local strength of the Born approximation of q almost surely.

In the literature on scattering, one often makes use of different physical scaling regimes in order to estimate the size of relevant mathematical objects and to design effective reconstruction methods. Our hope is that the theoretical framework we set up here can produce interesting stable algorithms for the inversion of the probabilistic backscattering problem. In this spirit, we have included a brief analysis in the appendix A, where we consider scaling regimes such that the analogue of Theorem 1.1 holds for $n = m$ with $n \geq 2$.

Our work follows the line of the previous papers [38] by Lassas, Päiväranta and Sakasman, and [28] by Helin, Lassas and Päiväranta. In [38], a similar problem was solved in $\mathbb{R}^2$ for a backscattering problem with point sources in an open and bounded set, but assuming the knowledge of the full scattered wave. The present paper improves this setting by studying scattering of plane waves and assuming only knowledge of the far-field pattern of the backscattered wave. Moreover, the results are generalized to arbitrary dimension. Although, our work draws inspiration of this paper, these two aspects require more sophisticated techniques in several parts. Later, in [28] Helin *et al* considered backscattering from random Robin boundary condition in half-space geometry of $\mathbb{R}^3$.

The literature on the deterministic inverse backscattering problem is considerably wide. For uniqueness results in generic (i.e., dense and open sets) class of potentials, see [21, 22]. Uniqueness of the problem for potentials with controlled angular regularity has been proved in [45]. Earlier partial results for the inverse backscattering problem for Schrödinger equation has been obtained in [32, 40, 41, 44, 54, 56, 63]. Approximative or numerical reconstructions have been studied in [9, 32]. The recovery of singularities of the potential from backscattering data is analyzed in [8, 26, 42, 46, 47, 49, 51, 52, 53]. Other references on inverse backscattering for a time-harmonic Schrödinger equation are [10, 37, 59]. The backscattering problem has also been studied in the framework of acoustic scattering (see [55, 61, 62]) and Maxwell equations (see [60]). For a concise treatment of classical inverse scattering, we refer to [17].

The wave and particle propagation in heterogeneous media has been extensively studied. Often, heterogeneous media is not known precisely and is modelled as a realization of random

media with known statistics. Mathematical theory being developed typically relies on multi-scale analysis or homogenization with the aim of capturing the effective properties of the propagation. We refer to the works in [2, 33] for various perspectives on wave propagation (whether classical or quantum) in random media. Let us also mention the papers on random Schrödinger models [3, 27], where the potential model involving slowly decaying correlations corresponds closely to the random potential model in the present paper. Notice that our work does not involve assumptions on scaling regimes nor any approximations. However, as mentioned above, we have included Appendix A discussing our method from the perspective of multi-scale analysis.

Recently, inverse problems related to imaging of random media have received wide attention [4, 5, 6, 7, 11, 12, 13, 14, 19, 20, 24]. The key feature of time reversal in a randomly inhomogeneous media is that it leads to focusing resolution that is much better than in a homogeneous media. This phenomenon is called super-resolution and appears due to multipathing caused by the random media [15]. Similar to our work, the back-propagated fields are self-averaging and the imaging method is statistically stable, i.e., independent of the realization of the random media

This paper is organized as follows. In the section 2, we describe in detail our stochastic model for a random potential q and the implications it has for the regularity of q . As discussed above, these regularity considerations require to develop the theory of the forward problem for non-regular potentials. This is studied in the section 3. The inverse problem is then covered in the section 4. The effects of first, second and higher order scattering for zero-mean potentials are studied separately in the sections 4.1, 4.2 and 4.3, respectively. In the section 4.4 we proof Theorem 1.1 for non-zero-mean potentials. Finally, in Appendix A we consider the physical scaling regimes where our method could be numerically effective and, afterwards, in Appendix B give some basic results regarding Gaussian distributions.

2. RANDOM POTENTIAL

2.1. Microlocally isotropic random field. In order to provide a precise mathematical description of the random potential to be considered, let $(\Omega, \mathcal{H}, \mathbb{P})$ be a complete probability space. Since we are interested in the properties of an object with a complicated micro-structure, we start by assuming that q is a *generalized random function*³. Below, the generalized function u defines a linear and continuous function $u : C_0^\infty(\mathbb{R}^n; \mathbb{R}) \rightarrow \mathbb{R}$. For $\phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ we denote $u(\phi) = \langle u, \phi \rangle$ and for the set of generalized functions (or distributions) we use the notation $\mathcal{D}'(\mathbb{R}^n; \mathbb{R})$. We also recall that any function $u \in L_{loc}^1(\mathbb{R}^n)$ defines a generalized function given by $\langle u, \phi \rangle = \int_{\mathbb{R}^n} u(x)\phi(x) dx$.

The assumption that q is a generalized random function means that q is a mapping defined on Ω such that, for every $\omega \in \Omega$, the realization $q(\omega)$ is a linear real valued functional on $C_0^\infty(\mathbb{R}^n; \mathbb{R})$ —the space of smooth real-valued functions with compact support in $\mathbb{R}^n$ —with $n \geq 2$ and the function

$$\omega \in \Omega \longmapsto \langle q(\omega), \phi \rangle \in \mathbb{R}$$

is a random variable for all $\phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$. Moreover we assume that, for every compact $K \subset \mathbb{R}^n$, there exists a non-negative random variable $C : (\Omega, \mathcal{H}) \rightarrow \mathbb{R}_+$ with $\mathbb{E}C^2 < \infty$ and

³For properties of generalized random functions, see [48].

$N \in \mathbb{N}$ such that for $\mathbb{P}$ -a.e. $\omega \in \Omega$ we have

$$(5) \quad |\langle q(\omega), \phi \rangle| \leq C(\omega) \sum_{|\alpha| \leq N} \sup_{x \in \mathbb{R}^n} |\partial^\alpha \phi(x)|$$

for all $\phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ with compact support $\text{supp } \phi \subset K$. Note that $q(\omega) \in \mathcal{D}'(\mathbb{R}^n; \mathbb{R})$, which denotes the space of real distributions in $\mathbb{R}^n$. Obviously, $q(\omega)$ can be extended to the space of smooth (complex-valued) functions with compact support as

$$\langle q(\omega), \phi \rangle = \langle q(\omega), \text{Re } \phi \rangle + i \langle q(\omega), \text{Im } \phi \rangle, \quad \phi \in C_0^\infty(\mathbb{R}^n).$$

A generalized random function is said to be *Gaussian* if the random variable

$$(6) \quad r_1 \langle q, \phi_1 \rangle + \cdots + r_l \langle q, \phi_l \rangle$$

has a Gaussian distribution for every $r_1, \dots, r_l \in \mathbb{R}$, $\phi_1, \dots, \phi_l \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ and $l \in \mathbb{N} \setminus \{0\}$. We say it is compactly supported if there exists a bounded domain D in $\mathbb{R}^n$ such that $\text{supp } q \subset D$ almost surely. Note that the probability law of a generalized Gaussian field q is determined by

$$\mathbb{E}q : \phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R}) \mapsto \mathbb{E}\langle q, \phi \rangle \in \mathbb{R}$$

$$\text{Cov } q : (\phi_1, \phi_2) \in C_0^\infty(\mathbb{R}^n; \mathbb{R})^2 \mapsto \text{Cov}(\langle q, \phi_1 \rangle, \langle q, \phi_2 \rangle) \in \mathbb{R},$$

where $\mathbb{E}\langle q, \phi \rangle$ denotes the expected value of $\langle q, \phi \rangle$ and

$$\text{Cov}(\langle q, \phi_1 \rangle, \langle q, \phi_2 \rangle) = \mathbb{E}((\langle q, \phi_1 \rangle - \mathbb{E}\langle q, \phi_1 \rangle)(\langle q, \phi_2 \rangle - \mathbb{E}\langle q, \phi_2 \rangle))$$

denotes the covariance of $\langle q, \phi_1 \rangle$ and $\langle q, \phi_2 \rangle$. Note that $\mathbb{E}q \in \mathcal{D}'(\mathbb{R}^n; \mathbb{R})$. The *covariance operator* $C_q : \phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R}) \mapsto C_q \phi \in \mathcal{D}'(\mathbb{R}^n; \mathbb{R})$ is defined as

$$\langle C_q \phi, \psi \rangle = \text{Cov}(\langle q, \phi \rangle, \langle q, \psi \rangle).$$

Since C_q is continuous, by the Schwartz kernel theorem, there exists a unique $K_q \in \mathcal{D}'(\mathbb{R}^n \times \mathbb{R}^n; \mathbb{R})$, usually called the *covariance function*, such that

$$\langle K_q, \psi \otimes \phi \rangle = \langle C_q \phi, \psi \rangle = (\text{Cov } q)(\phi, \psi)$$

for all $\phi, \psi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ (or more generally, for all $\phi, \psi \in C_0^\infty(\mathbb{R}^n)$). In particular,

$$(7) \quad K_q = \mathbb{E}((q - \mathbb{E}q) \otimes (q - \mathbb{E}q)).$$

It is often convenient to write

$$K_q(x, y) = \mathbb{E}((q(x) - \mathbb{E}q(x))(q(y) - \mathbb{E}q(y)))$$

and

$$\langle K_q, \psi \otimes \phi \rangle = \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} K_q(x, y) \psi(x) \phi(y) dx dy.$$

Definition 2.1. A generalized function q on $\mathbb{R}^n$ is called *microlocally isotropic* of order $-m$ in D , if the following conditions hold:

- (1) $\mathbb{E}q$ is smooth,
- (2) q is supported in D a.s.,
- (3) the covariance operator C_q is a classical pseudo differential operator of order $-m$ with $n - 1 < m \leq n + 1$ and
- (4) C_q has a principal symbol of the form $\mu(x)|\xi|^{-m}$ with $\mu \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$, $\text{supp } \mu \subset D$ and $\mu(x) \geq 0$ for all $x \in \mathbb{R}^n$.

In consequence, if q is microlocally isotropic, we have $\text{supp}(\mathbb{E}q) \subset D$ and

$$C_q \phi(x) = \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} e^{i(x-y) \cdot \xi} c_q(x, \xi) \phi(y) dy d\xi$$

for a classical symbol $c_q \in \mathcal{S}^{-m}(\mathbb{R}^n \times \mathbb{R}^n)$. Moreover, there exists a classical symbol $a \in \mathcal{S}^{-m-1}(\mathbb{R}^n \times \mathbb{R}^n)$ such that

$$(8) \quad a(x, \xi) = c_q(x, \xi) - \mu(x)|\xi|^{-m} \text{ for } x \in \mathbb{R}^n \text{ and } |\xi| \geq 1.$$

The covariance function and the symbol of C_q are connected via the following identity

$$(9) \quad K_q(x, y) = \frac{1}{(2\pi)^{n/2}} \mathcal{F}^{-1}(c_q(x, \cdot))(x - y),$$

where $\mathcal{F}^{-1}$ denotes the inverse Fourier transform. Here, for an integrable function f , we define $(\mathcal{F}f)(\xi) = (2\pi)^{-n/2} \int_{\mathbb{R}^n} e^{-i\xi \cdot x} f(x) dx$ and frequently abbreviate $\widehat{f} = (\mathcal{F}f)$. In our particular case, $\text{supp } K_q \subset D \times D$ and

$$(10) \quad \int_{\mathbb{R}^n} \left(\int_{\mathbb{R}^n} K_q(x, y) e^{-i\xi \cdot (x-y)} dy \right) \phi(x) dx = \int_D c_q(x, \xi) \phi(x) dx$$

$$(11) \quad = \int_{\mathbb{R}^n} \mu(x) |\xi|^{-m} \phi(x) dx + \int_D a(x, \xi) \phi(x) dx$$

for every $\phi \in C^\infty(\mathbb{R}^n)$ and all $|\xi| \geq 1$.

Let us illustrate this definition with a brief example.

Example 2.2. Let W stand for the generalized Gaussian white noise. That is, W is a generalized Gaussian field with $\mathbb{E}W = 0$ and its covariance operator satisfying

$$\mathbb{E}(\langle W, \phi \rangle \langle W, \psi \rangle) = \int_{\mathbb{R}^n} \phi(x) \psi(x) dx$$

for every $\phi, \psi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$. It is well known that $W \in H_{loc}^{-n/2-\epsilon}(\mathbb{R}^n)$ a.s. for any $\epsilon > 0$. Our example of a microlocally isotropic random potential of order $-m$ in D is given by

$$q = \sqrt{\mu}(\mathbf{I} - \Delta)^{-m/4} W + q_0$$

with μ and m as in Definition 2.1, and q_0 a smooth real-valued function in $\mathbb{R}^n$ with support in D . Thus, q is a generalized Gaussian field with covariance operator

$$C_q = M_{\sqrt{\mu}}(\mathbf{I} - \Delta)^{-m/2} M_{\sqrt{\mu}}$$

where $M_{\sqrt{\mu}}\phi(x) = \sqrt{\mu(x)}\phi(x)$. Its covariance function is

$$K_q(x, y) = \sqrt{\mu(x)}\sqrt{\mu(y)}G_m(x - y)$$

with $G_m \in L^1(\mathbb{R}^n)$ (see page 132 in [57]) such that

$$\widehat{G}_m(\xi) = \frac{1}{(2\pi)^{n/2}} \frac{1}{(1 + |\xi|^2)^{-m/2}}.$$

Finally, since C_q has a principal symbol of the form $\mu(x)|\xi|^{-m}$, we see that q is a microlocally isotropic random potential of order $-m$ in D .

Example 2.3. *Following the Example 1 in [38] let us define the multidimensional fractional Brownian motion in $\mathbb{R}^n$ for the Hurst index H as the centered Gaussian process $X_H(z)$ indexed by $z \in \mathbb{R}^n$ with following properties:*

$$\begin{aligned} \mathbb{E}|X_H(z_1) - X_H(z_2)|^2 &= |z_1 - z_2|^{2H} \quad \text{for all } z_1, z_2 \in \mathbb{R}^n \\ X(z_0) &= 0 \quad \text{and} \\ \text{the paths } z &\mapsto X_H(z) \text{ are a.s. continuous.} \end{aligned}$$

The existence and basic properties of X_H are well-known [34]. Let us define the potential q by setting

$$q(z, \omega) = \sqrt{\mu}(z)X_H(z, \omega)$$

for some $\mu \in C_0^\infty(\mathbb{R}^n)$ and index $H > 0$. It follows that the principal symbol of C_q is of the form $\mu(z)|\xi|^{-n-2H}$.

FIGURE 1. (Color in the online version of the paper) A realization of the random potential $q(x)$ in Example 2.3 was generated on $\mathbb{R}^2$ with the Hurst index $H = 0.25$ and the local strength μ . A realization $q(x) = q(x, \omega_0)$ is shown in the figure on the left and in the middle. The figure on the right is the local strength $\mu(x)$ of the random field $q(x)$, i.e., the principal symbol of the covariance operator C_q . Notice that the function μ is large on the areas where the realization of the random field $q(x, \omega_0)$ has large local fluctuations.

2.2. Regularity of the potential. We will use the potential Sobolev spaces $L_s^p(\mathbb{R}^n)$ to determine the regularity of the realizations of q . Let us recall the definition and a basic property of these spaces.

Let $\mathcal{J}_s$ denote the Bessel potential $\mathcal{J}_s = (I - \Delta)^{-s/2}$, i.e., for any Schwartz function $f \in \mathcal{S}(\mathbb{R}^n)$ it holds that

$$\widehat{\mathcal{J}_s f}(\xi) = (1 + |\xi|^2)^{-s/2} \widehat{f}(\xi)$$

for all $\xi \in \mathbb{R}^n$. The Bessel potential extends to temperate distributions and, in particular, can be applied to almost every realization of q . The potential Sobolev space $L_s^p(\mathbb{R}^n)$ with $1 \leq p \leq \infty$ and $s \in \mathbb{R}$ is defined by the set of $f = \mathcal{J}_s g$ such that $g \in L^p(\mathbb{R}^n)$ and it is endowed with the norm

$$\|f\|_{L_s^p} = \|g\|_{L^p}.$$

When $p = 2$, the space $L_s^2(\mathbb{R}^n)$ is commonly denoted by $H^s(\mathbb{R}^n)$. When $k \in \mathbb{N}$ and $1 < p < \infty$, the space $L_k^p(\mathbb{R}^n)$ can be identified with usual Sobolev spaces $W^{k,p}(\mathbb{R}^n)$. We also use the

notation $W^{k,2}(\mathbb{R}^n) = H^k(\mathbb{R}^n)$. Note that, since $\mathcal{J}_t : L_s^p(\mathbb{R}^n) \longrightarrow L_{s+t}^p(\mathbb{R}^n)$ is an isometric isomorphism, $f \in L_s^p(\mathbb{R}^n)$ if and only if $\mathcal{J}_t f \in L_{s+t}^p(\mathbb{R}^n)$.

Now we can establish the regularity of the realizations of q .

Proposition 2.4. *Let q be a microlocally isotropic random potential of order $-m$ in D . Then, $q \in L_s^p(\mathbb{R}^n)$ almost surely for any $1 < p < \infty$ and $s < (m - n)/2$.*

Corollary 2.5. *Let q be a microlocally isotropic random potential of order $-m$ in D with $n < m \leq n + 1$. Then, $q \in C^{0,\alpha}(\mathbb{R}^n)$ almost surely for any $0 < \alpha < (m - n)/2$.*

The corollary follows from the proposition using the Sobolev embeddings and the fact that, if $m > n$ and $0 < \alpha < (m - n)/2$, there exist $s < (m - n)/2$ close to $(m - n)/2$ and $p > n$ very large such that $\alpha = s - n/p < (m - n)/2$.

Let us prove Proposition 2.4 in the case $\mathbb{E}q = 0$. The general case follows by our assumption of $\mathbb{E}q$ being smooth. For $\varepsilon \in (0, 1]$, define

$$(12) \quad f_\varepsilon(\omega, x) = \langle \mathcal{J}_{-s}q(\omega), \varphi_\varepsilon(x - \bullet) \rangle$$

with $\varphi_\varepsilon(x) = \varepsilon^{-n} \varphi(x/\varepsilon)$, $\varphi \in C_0^\infty(\mathbb{R}^n; [0, 1])$ and $\int_{\mathbb{R}^n} \varphi(x) dx = 1$. Note that $f_\varepsilon(\omega)$ tends to $\mathcal{J}_{-s}q(\omega)$ in $\mathcal{D}'(\mathbb{R}^n)$ as ε goes to 0 for $\mathbb{P}$ -almost every $\omega \in \Omega$.

Lemma 2.6. *For every $p \in [1, \infty)$ there exists $C = C(p)$ such that*

$$(\mathbb{E}|f_\varepsilon(x)|^p)^{1/p} \leq C(\mathbb{E}|f_\varepsilon(x)|^2)^{1/2},$$

where $f_\varepsilon(x) = f_\varepsilon(\bullet, x)$.

Proof. If $p \in [1, 2)$, the lemma follows applying Hölder inequality. The case $p = 2$ is obvious. So it only remains the case $p \in (2, \infty)$. For every $p \in (2, \infty)$, there is an only $j \in \mathbb{N} \setminus \{0\}$ such that $2j < p \leq 2(j + 1)$, and, by Hölder inequality, we see that

$$(13) \quad (\mathbb{E}|f_\varepsilon(x)|^p)^{1/p} \leq (\mathbb{E}|f_\varepsilon(x)|^{2(j+1)})^{\frac{1}{2(j+1)}}.$$

By the triangle inequality, we see that

$$(\mathbb{E}|f_\varepsilon(x)|^{2(j+1)})^{1/(j+1)} \leq \left(\mathbb{E}(\operatorname{Re} f_\varepsilon(x))^{2(j+1)} \right)^{1/(j+1)} + \left(\mathbb{E}(\operatorname{Im} f_\varepsilon(x))^{2(j+1)} \right)^{1/(j+1)}.$$

Note that

$$\operatorname{Re} f_\varepsilon(\omega, x) = \langle q(\omega), \operatorname{Re} \mathcal{J}_{-s} \varphi_\varepsilon(x - \bullet) \rangle, \quad \operatorname{Im} f_\varepsilon(\omega, x) = \langle q(\omega), \operatorname{Im} \mathcal{J}_{-s} \varphi_\varepsilon(x - \bullet) \rangle;$$

and

$$(14) \quad \mathbb{E} \operatorname{Re} f_\varepsilon(x) = \mathbb{E} \operatorname{Im} f_\varepsilon(x) = 0$$

since $\mathbb{E}q = 0$. Using that $\operatorname{Re} f_\varepsilon(x)$ and $\operatorname{Im} f_\varepsilon(x)$ are Gaussian for every $x \in \mathbb{R}^n$, the identity (14) and Lemma B.1, we can check that

$$(\mathbb{E}|f_\varepsilon(x)|^{2(j+1)})^{1/(j+1)} \leq C \left(\mathbb{E}(\operatorname{Re} f_\varepsilon(x))^2 + \mathbb{E}(\operatorname{Im} f_\varepsilon(x))^2 \right) = C \mathbb{E}|f_\varepsilon(x)|^2$$

for $C = C(j)$. Finally, plugging the previous inequality into (13), we get the estimate in the lemma. $\square$

It is an immediate consequence of the previous lemma that

$$(15) \quad \mathbb{E} \|f_\varepsilon\|_{L^p(K)}^p = \int_K \mathbb{E} |f_\varepsilon(x)|^p dx \leq C \int_K (\mathbb{E} |f_\varepsilon(x)|^2)^{p/2} dx$$

for every compact $K \subset \mathbb{R}^n$. On the other hand, we have the identity

$$\begin{aligned} \mathbb{E} |f_\varepsilon(x)|^2 &= \mathbb{E} (\langle q, \mathcal{J}_{-s}(\varphi_\varepsilon(x - \cdot)) \rangle \langle q, \overline{\mathcal{J}_{-s}(\varphi_\varepsilon(x - \cdot))} \rangle) \\ &= \langle C_q \mathcal{J}_{-s}(\varphi_\varepsilon(x - \cdot)), \overline{\mathcal{J}_{-s}(\varphi_\varepsilon(x - \cdot))} \rangle \\ &= \langle \mathcal{J}_{-s} C_q \mathcal{J}_{-s}(\varphi_\varepsilon(x - \cdot)), \varphi_\varepsilon(x - \cdot) \rangle. \end{aligned}$$

Note that $\mathcal{J}_{-s} C_q \mathcal{J}_{-s}$ is pseudo-differential operator of order $-m + 2s$ with a symbol $\tilde{c} \in \mathcal{S}^{-m+2s}(\mathbb{R}^n \times \mathbb{R}^n)$ and Schwartz kernel

$$\tilde{K}(y, z) = \frac{1}{(2\pi)^{n/2}} \mathcal{F}^{-1}(\tilde{c}(y, \cdot))(y - z).$$

For $s < (m - n)/2$, we can check that

$$\mathbb{E} |f_\varepsilon(x)|^2 = \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} \left(\int_{\mathbb{R}^n} e^{i\varepsilon(z-y) \cdot \xi} \tilde{c}(x - \varepsilon y, \xi) d\xi \right) \varphi(y) \varphi(z) dy dz$$

and consequently that

$$(16) \quad \mathbb{E} |f_\varepsilon(x)|^2 \leq C \int_{\mathbb{R}^n} (1 + |\xi|)^{-m+2s} d\xi.$$

Therefore, by (16) and (15) we have that, for every compact K in $\mathbb{R}^n$, there exists a constant C independent of ε such that

$$(17) \quad \mathbb{E} \|f_\varepsilon\|_{L^p(K)}^p \leq C,$$

for all $\varepsilon \in (0, 1]$.

Lemma 2.7. *For $1 \leq p < \infty$ and $s < (m - n)/2$ with $n - 1 < m \leq n + 1$, we have that*

$$\mathcal{J}_{-s} q \in L_{\text{loc}}^p(\mathbb{R}^n)$$

almost surely.

Proof. The case $p = 1$ follows from the case $1 < p < \infty$ by Hölder's inequality. The case $1 < p < \infty$ is a consequence of (17). Indeed, let K be an arbitrary compact set in $\mathbb{R}^n$. The bound (17) implies, by the Banach-Alaoglu theorem and since $L^p(\Omega \times K)$ with $1 < p < \infty$ is reflexive, that there exist $f \in L^p(\Omega \times K)$ and a vanishing sequence $\{\varepsilon_j\}_{j=1}^\infty$, such that, f_{ε_j} converges weakly to f as j goes to infinity. In particular, we have

$$(18) \quad \lim_{j \rightarrow \infty} \int_{\Omega} \int_{\mathbb{R}^n} f_{\varepsilon_j}(\omega, x) \phi(x) \psi(\omega) dx d\mathbb{P}(\omega) = \int_{\Omega} \int_{\mathbb{R}^n} f(\omega, x) \phi(x) \psi(\omega) dx d\mathbb{P}(\omega)$$

for every $\phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ such that $\text{supp}(\phi) \subset K$ and every simple function $\psi : \Omega \rightarrow \mathbb{R}$.

Next, we show that

$$(19) \quad \lim_{\varepsilon \rightarrow 0} \int_{\Omega} \langle f_\varepsilon(\omega), \phi \rangle \psi(\omega) d\mathbb{P}(\omega) = \int_{\Omega} \langle \mathcal{J}_{-s} q(\omega), \phi \rangle \psi(\omega) d\mathbb{P}(\omega)$$

for every $\phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ such that $\text{supp}(\phi) \subset K$ and every simple function $\psi : \Omega \rightarrow \mathbb{R}$. The limit (19) holds by the dominate convergence theorem since $f_\varepsilon(\omega)$ tends to $\mathcal{J}_{-s} q(\omega)$ in

$\mathcal{D}'(\mathbb{R}^n)$ as ε goes to 0 almost surely and $|\langle f_\varepsilon(\omega), \phi \rangle|$ is bounded, for all ε , by an integrable function in Ω . To check this last point, note that

$$\langle f_\varepsilon(\omega), \phi \rangle = \left\langle q(\omega), \chi \int_{\mathbb{R}^n} \varphi_\varepsilon(x) \mathcal{J}_{-s} \phi(x + \bullet) dx \right\rangle$$

with $\chi \in C_0^\infty(\mathbb{R}^n)$ such that $\chi(x) = 1$ for all $x \in \overline{D}$, and by (5),

$$|\langle f_\varepsilon(\omega), \phi \rangle| \leq C(\omega) \sum_{|\alpha| \leq N} \sup_{y \in \mathbb{R}^n} \left| \partial^\alpha (\chi(y) \int_{\mathbb{R}^n} \varphi_\varepsilon(x) \mathcal{J}_{-s} \phi(x + y) dx) \right|,$$

where the second term on the right-hand side is bounded by a constant independent of ε and $C = C(\omega)$ is integrable in Ω .

Therefore, by the density of the simple functions and the limits (18) and (19), we have that for $\mathbb{P}$ -almost every $\omega \in \Omega$ we have

$$\langle \mathcal{J}_{-s} q(\omega), \phi \rangle = \int_{\mathbb{R}^n} f(\omega, x) \phi(x) dx$$

for any $\phi \in C_0^\infty(\mathbb{R}^n; \mathbb{R})$ such that $\text{supp}(\phi) \subset K$. This concludes the proof of the lemma. $\square$

In order to conclude the regularity of the realizations of q , choose $\chi \in C_0^\infty(\mathbb{R}^n)$ such that $\chi(x) = 1$ for all $x \in \overline{D}$ and write

$$q(\omega) = \mathcal{J}_s(\chi \mathcal{J}_{-s} q(\omega)) + \mathcal{J}_s((1 - \chi) \mathcal{J}_{-s} q(\omega)).$$

The first term belongs to $L_s^p(\mathbb{R}^n)$ a.s. by the previous lemma while the second one belongs to $L_t^p(\mathbb{R}^n)$ a.s. for every $t \in \mathbb{R}$ and every $1 < p < \infty$ because the supports of $1 - \chi$ and $q(\omega)$ are disjoint. This ends the proof of Proposition 2.4.

2.3. Covariance function. We conclude this section providing a more detailed description of the covariance function. This will be proved in the next proposition and used at the end of the section 4.2.

Proposition 2.8. *Let q be a microlocally isotropic random potential of order $-m$ in D . The covariance function K_q has the following form:*

- (a) *If $n < m \leq n + 1$, there exists a compactly supported function $F_{1+\alpha} \in C^{1,\alpha}(\mathbb{R}^n \times \mathbb{R}^n)$ for $0 < \alpha < m - n$ such that*

$$K_q(x, y) = c_{n,m} \mu(x) |x - y|^{m-n} + F_{1+\alpha}(x, y).$$

where $c_{n,m}$ is a constant depending on n and m .

- (b) *If $n - 1 < m \leq n$, there exists a compactly supported function $F_\alpha \in C^{0,\alpha}(\mathbb{R}^n \times \mathbb{R}^n)$ for $0 < \alpha < m - (n - 1)$ such that*

$$K_q(x, y) = \begin{cases} c_{n,m} \mu(x) |x - y|^{-(n-m)} + F_\alpha(x, y), & \text{if } n - 1 < m < n, \\ c_{n,m} \mu(x) \log |x - y| + F_\alpha(x, y), & \text{if } m = n. \end{cases}$$

where $c_{n,m}$ is a constant depending on n and m .

Proof. As a consequence of identities (8) and (9), for a radially symmetric $\psi \in C_0^\infty(\mathbb{R}^n)$ such that $\psi(\xi) = 1$ for $|\xi| \leq 1$, we have

$$(20) \quad \begin{aligned} (2\pi)^{n/2} \langle K_q(x, x - \cdot), \phi \rangle &= \int_{\mathbb{R}^n} \mu(x)(1 - \psi(\xi))|\xi|^{-m} (\mathcal{F}^{-1}\phi)(\xi) d\xi \\ &+ \int_{\mathbb{R}^n} \mathcal{F}^{-1}((1 - \psi)a(x, \cdot))(y)\phi(y) dy + \int_{\mathbb{R}^n} \mathcal{F}^{-1}(\psi c_q(x, \cdot))(y)\phi(y) dy. \end{aligned}$$

Note that the function $v(x, y) = \mathcal{F}^{-1}(\psi c_q(x, \cdot))(y)$ satisfies $v \in \mathcal{S}(\mathbb{R}^n \times \mathbb{R}^n)$ with $\text{supp}(v) \subset D \times \mathbb{R}^n$. On the other hand, $w(x, y) = \mathcal{F}^{-1}((1 - \psi)a(x, \cdot))(y)$ is smooth and compactly supported in D with respect to the variable x . Moreover, using the Hausdorff–Young inequality, we see that for fixed $x \in D$ we have $w(x, \cdot) \in L_s^p(\mathbb{R}^n)$ for $2 \leq p \leq \infty$ and $s > 0$ such that $s - n/p < m - (n - 1)$. Furthermore, for every $0 < \alpha \leq 1$ and $k = 0, 1$ such that $k + \alpha < m - (n - 1)$ there exists $s > 0$ and $p > n$ with $k + \alpha = s - n/p$. Therefore, by the Sobolev embedding theorem, we have $w(x, \cdot) \in C^{k, \alpha}(\mathbb{R}^n)$.

If $n < m \leq n + 1$, the first term on the right-hand side of (20) coincides with

$$\mu(x) \int_{\mathbb{R}^n} \mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m})(y)\phi(y) dy$$

where $\mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m}) \in L_s^p(\mathbb{R}^n)$ with $2 \leq p \leq \infty$ and $s > 0$ such that $s - n/p < m - n$. Again, by the Sobolev embedding theorem, we can ensure that $\mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m}) \in C^{0, \alpha}(\mathbb{R}^n)$ for $0 < \alpha < m - n$.

Let us now prove that

$$\mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m})(y) = c_{n, m} |y|^{m-n} + g(y)$$

for some $g \in C^\infty(\mathbb{R}^n)$. The previous identity follows from the identity

$$(21) \quad \begin{aligned} \mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m})(\lambda z) &= \lambda^{m-n} \mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m})(z) \\ &+ \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{i\lambda z \cdot \xi} |\xi|^{-m} (\psi(\lambda \xi) - \psi(\xi)) d\xi \end{aligned}$$

with $z = y/|y|$, $\lambda = |y|$ and

$$c_{n, m} = \mathcal{F}^{-1}((1 - \psi)|\cdot|^{-m})(y/|y|).$$

Recall that ψ is radially symmetric and hence $c_{n, m}$ is independent of y . The identity (21) is, in turn, an immediate consequence of

$$(1 - \psi(\xi))|\xi|^{-m} = (1 - \psi(\lambda \xi))|\xi|^{-m} + (\psi(\lambda \xi) - \psi(\xi))|\xi|^{-m}.$$

Thus, in the case $n < m \leq n + 1$ the identity (20) becomes

$$(2\pi)^{n/2} \langle K_q(x, x - y), \phi(y) \rangle = \int_{\mathbb{R}^n} (c_{n, m} \mu(x) |y|^{m-n} + g(x, y)) \phi(y) dy$$

for some $g \in C^{1, \alpha}(\mathbb{R}^n \times \mathbb{R}^n)$ with compact support with respect to the variable x . Since K_q is supported on $D \times D$, we can just introduce an appropriate compactly supported smooth function in the previous identity and write the identity

$$K_q(x, y) = c_{n, m} \mu(x) |x - y|^{m-n} + F_{1+\alpha}(x, y)$$

with $c_{n, m}$ and $F_{1+\alpha}$ as in the (a). This concludes the proof of statement (a).

Consider the case (b) with $n - 1 < m < n$. Let us first record two useful identities. First, we can write

$$(22) \quad |\xi|^{-m} = \frac{2^{-m/2}}{\Gamma(m/2)} \int_0^\infty t^{m/2} e^{-t|\xi|^2/2} \frac{dt}{t},$$

where Γ stands for

$$\Gamma(\beta) = \int_0^\infty t^\beta e^{-t} \frac{dt}{t}.$$

Second, the basic properties of the Fourier transform yield the identity

$$(23) \quad \int_{\mathbb{R}^n} e^{-t|\xi|^2/2} (1 - \psi(\xi)) \mathcal{F}^{-1}\phi(\xi) d\xi = t^{-n/2} \int_{\mathbb{R}^n} e^{-|y|^2/(2t)} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy.$$

Using (22) and (23) on the first term on the right-hand side of (20) we have

$$(24) \quad \begin{aligned} & \int_{\mathbb{R}^n} \mu(x) (1 - \psi(\xi)) |\xi|^{-m} \mathcal{F}^{-1}\phi(\xi) d\xi \\ &= \frac{2^{-m/2}}{\Gamma(m/2)} \mu(x) \int_0^\infty t^{m/2-n/2} \int_{\mathbb{R}^n} e^{-|y|^2/(2t)} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy \frac{dt}{t} \\ &= \frac{2^{-m+n/2}}{\Gamma(m/2)} \Gamma(n/2 - m/2) \mu(x) \int_{\mathbb{R}^n} |y|^{m-n} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy. \end{aligned}$$

In the last identity we have used the definition of Γ . Since

$$\int_{\mathbb{R}^n} |y|^{m-n} \mathcal{F}(\psi\mathcal{F}^{-1}\phi)(y) dy = \int_{\mathbb{R}^n} \left(\frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} |y|^{m-n} \mathcal{F}^{-1}\psi(z - y) dy \right) \phi(z) dz,$$

we can deduce from (24) and (20), proceeding as before, that

$$K_q(x, y) = c_{n,m} \mu(x) |x - y|^{m-n} + F_\alpha(x, y)$$

with $c_{n,m}$ and F_α as in the (b).

The case $m = n$ follows from a limit argument, making $m < n$ goes to n . To do so, note that

$$\Gamma(n/2 - m/2) = \frac{\Gamma(n/2 - m/2 + 1)}{n/2 - m/2},$$

which is a consequence of the definition of Γ , and rewrite (24) as

$$\begin{aligned} & \int_{\mathbb{R}^n} \mu(x) (1 - \psi(\xi)) |\xi|^{-m} \mathcal{F}^{-1}\phi(\xi) d\xi \\ &= \frac{2^{-m+n/2}}{\Gamma(m/2)} \Gamma(n/2 - m/2 + 1) \mu(x) \int_{\mathbb{R}^n} \frac{|y|^{m-n} - 1}{n/2 - m/2} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy \\ &+ \frac{2^{-m+n/2}}{\Gamma(m/2)} \Gamma(n/2 - m/2) \mu(x) \int_{\mathbb{R}^n} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy. \end{aligned}$$

On the other hand, since $(|y|^{-(n-m)} - 1)/(n - m) \rightarrow \log |y|^{-1}$ as $m \rightarrow n$, we can use the dominate convergence theorem to pass to the limit and obtain

$$\begin{aligned} \int_{\mathbb{R}^n} \mu(x)(1 - \psi(\xi))|\xi|^{-n} \mathcal{F}^{-1}\phi(\xi) d\xi &= c_{n,m} \mu(x) \int_{\mathbb{R}^n} \log |y|^{-1} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy \\ &\quad + \tilde{c}_{n,m} \mu(x) \int_{\mathbb{R}^n} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy \end{aligned}$$

for $c_{n,m}$ and $\tilde{c}_{n,m}$ positive constants. Checking that the identities

$$\begin{aligned} \int_{\mathbb{R}^n} \log |y| \mathcal{F}(\psi \mathcal{F}^{-1}\phi)(y) dy &= \int_{\mathbb{R}^n} \left(\frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} \log |y| \mathcal{F}^{-1}\psi(z - y) dy \right) \phi(z) dz, \\ \int_{\mathbb{R}^n} \mathcal{F}((1 - \psi)\mathcal{F}^{-1}\phi)(y) dy &= \int_{\mathbb{R}^n} \left(1 - \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} \mathcal{F}^{-1}\psi(z - y) dy \right) \phi(z) dz \end{aligned}$$

hold and proceeding as before we have

$$K_q(x, y) = c_{n,m} \mu(x) \log |x - y| + F_\alpha(x, y)$$

with $c_{n,m}$ and F_α as in the (c), which ends the proof of this proposition. $\square$

3. DIRECT SCATTERING FOR A ROUGH POTENTIAL

In the previous section we established that the realizations of our random potential field model can be rough (Proposition 2.4). Therefore, we need to show that the scattering problem in (1) is well-defined. To achieve this, we leave the randomness aside for a while and consider the deterministic scattering problem

$$(25) \quad \begin{cases} (\Delta + k^2 - V)u = 0 \\ u(x) = e^{ik\theta \cdot x} + u_{\text{sc}}(x) \\ \frac{x}{|x|} \cdot \nabla u_{\text{sc}}(x) - ik u_{\text{sc}}(x) = o(|x|^{-\frac{n-1}{2}}) \text{ as } |x| \rightarrow \infty \end{cases}$$

for a rough potential V . Later, our plan is to apply these results to the random scattering scenario for pointwise values $q = q(\omega)$. The condition satisfied by u_{sc} at infinity is usually referred as the *outgoing Sommerfeld radiation condition* (SRC for short). Here we assume V to be in $L^p_{-s}(\mathbb{R}^n)$ with $0 < s \leq 1/2$ and $n/s \leq p < \infty$ and to have support in a bounded domain D in $\mathbb{R}^n$. Note that the limiting case $V \in L^\infty(\mathbb{R}^n)$ (corresponding to $s \rightarrow 0$) is classical.

We will construct the scattered wave to be $u_{\text{sc}} = \sum_{j=1}^{\infty} u_j$ with

$$\begin{cases} (\Delta + k^2)u_j = V u_{j-1} & j \in \mathbb{N} \setminus \{0\}, \\ u_0(x) = e^{ik\theta \cdot x}. \end{cases}$$

Since the function $\xi \mapsto (-|\xi|^2 + k^2)^{-1}$ does not define a temperate distribution, we will introduce a temperate distribution $\Psi_k^\pm$ as the limit of

$$\xi \mapsto \frac{1}{-|\xi|^2 + k^2 \pm i\varepsilon}$$

as $\varepsilon > 0$ vanishes. This limit exists and it is given by

$$\langle \Psi_k^\pm, \phi \rangle = \text{p.v.} \int_{\mathbb{R}^n} \frac{\phi(\xi)}{-|\xi|^2 + k^2} d\xi \mp i \frac{\pi}{2k} \int_{|\xi|=k} \phi(\xi) d\sigma_k(\xi)$$

where the principal value $\text{p.v.} \int_{\mathbb{R}^n}$ stands for $\lim_{\epsilon \rightarrow 0} \int_{0 < \epsilon \leq |k^2 - |\xi|^2|}$ and $d\sigma_k$ denotes the volume form on $\{\xi \in \mathbb{R}^n : |\xi| = k\}$. Now, we let $\Phi_k^\pm$ denote the inverse Fourier transform of $\Psi_k^\pm$ and $\mathcal{R}_k^\pm$ be defined by

$$(26) \quad \widehat{\mathcal{R}_k^\pm f} = \Psi_k^\pm \widehat{f}$$

for f in the Schwartz class, hence,

$$\mathcal{R}_k^\pm f = \frac{1}{(2\pi)^{n/2}} \Phi_k^\pm * f.$$

Thus, we set $u_j = \mathcal{R}_k^\pm(Vu_{j-1}) = (\mathcal{R}_k^\pm \circ V)^j u_0$, where $\mathcal{R}_k^\pm \circ V$ denotes the *resolvent operator* $\mathcal{R}_k^\pm$ composed with the operator *multiplication-by* V .

On the other hand, it is known that $\Phi_k^\pm$ satisfies

$$(27) \quad \frac{x}{|x|} \cdot \nabla \Phi_k^\pm(x) \mp ik \Phi_k^\pm(x) = o(|x|^{-\frac{n-1}{2}}) \text{ as } |x| \rightarrow \infty.$$

Therefore, the scattered wave, which has to satisfy the SRC, will be constructed with $\mathcal{R}_k^+$ as

$$(28) \quad u_{\text{sc}} = \sum_{j=1}^{\infty} (\mathcal{R}_k^+ \circ V)^j u_0$$

provided this infinite sum makes sense.

In the following, we prove that the construction (28) is well-defined, by proving boundedness properties for the operators $\mathcal{R}_k^\pm$ and multiplication-by V in the weighted spaces $H_k^{s,\delta}(\mathbb{R}^n)$. The space $H_k^{s,\delta}(\mathbb{R}^n)$ is constructed in the following manner. Let $L^{2,\delta}(\mathbb{R}^n)$, with $\delta \in \mathbb{R}$, denote the equivalence class of measurable functions f in $\mathbb{R}^n$ such that

$$\int_{\mathbb{R}^n} (1 + |x|^2)^\delta |f(x)|^2 dx < \infty,$$

and be endowed with the norm

$$\|f\|_{L^{2,\delta}} = \left(\int_{\mathbb{R}^n} (1 + |x|^2)^\delta |f(x)|^2 dx \right)^{1/2}.$$

If $\delta = 0$, the space $L^{2,0}(\mathbb{R}^n)$ is just $L^2(\mathbb{R}^n)$. Let $H_k^{s,\delta}(\mathbb{R}^n)$, with δ and s in $\mathbb{R}$, be the space of $f = (k^2 - \Delta)^{-s/2} g$ such that $g \in L^{2,\delta}(\mathbb{R}^n)$, where $(k^2 - \Delta)^{-s/2}$ is defined as the multiplier with symbol $(k^2 + |\xi|^2)^{-s/2}$. Let this space be endowed with the norm

$$\|f\|_{H_k^{s,\delta}} = \|g\|_{L^{2,\delta}}.$$

When $\delta = 0$, the space $H_k^{s,0}(\mathbb{R}^n)$ will be denoted by $H_k^s(\mathbb{R}^n)$.

The boundedness properties of the operators $\mathcal{R}_k^\pm$ and multiplication-by V needed for our purposes, will be studied in the sections 3.1 and 3.2 respectively, and stated here as follows:

Theorem 3.1. *The operator $\mathcal{R}_k^\pm$ is bounded from $H_k^{-s,\delta}(\mathbb{R}^n)$ to $H_k^{s,-\delta}(\mathbb{R}^n)$ with $0 \leq s \leq 1/2$ and $\delta > 1/2$ and satisfies the inequality*

$$\|\mathcal{R}_k^\pm f\|_{H_k^{s,-\delta}} \lesssim k^{-(1-2s)} \|f\|_{H_k^{-s,\delta}}.$$

Proposition 3.2. *If V belong either $L_{-s}^p(\mathbb{R}^n)$ with $0 < s \leq 1/2$ and $n/s \leq p < \infty$, then the operator multiplication-by V is bounded from $H_k^{s,-\delta}(\mathbb{R}^n)$ to $H_k^{-s,\delta}(\mathbb{R}^n)$ and satisfies*

$$\|Vf\|_{H_k^{-s,\delta}} = o\left(\|f\|_{H_k^{s,-\delta}}\right),$$

expressed with the little o of Landau.

As a consequence of these statements, there exists $k_0 > 0$, depending only on $\|V\|_{L_{-s}^p}$, n , s and δ , such that the operator $\mathcal{R}_k^\pm \circ V$ maps $H_k^{s,-\delta}(\mathbb{R}^n)$ into itself with a norm strictly less than 1 for $k \geq k_0$. Then, the sum (28) converges in $H_k^{s,-\delta}(\mathbb{R}^n)$, and u_{sc} can be constructed by the infinite sum.

Remark 3.3. *Notice carefully that in the probabilistic problem setting the threshold wavelength $k_0 = k_0(\omega)$ becomes random due to the dependence on the term $\|q(\omega)\|_{L_{-s}^p}$. However, we clearly have $k_0(\omega) < \infty$ almost surely and as our reconstruction method is based on a single realization of the potential, the randomness poses no problems.*

Finally, by construction, the scattered wave satisfies the Lippmann–Schwinger equation

$$u_{\text{sc}}(x) = \int_{\mathbb{R}^n} \Phi_k^+(x-y)V(y)(e^{ik\theta \cdot y} + u_{\text{sc}}(y))dy \quad \text{in } H_k^{s,-\delta}(\mathbb{R}^n).$$

By the asymptotic behaviour of Φ_k^+ and the fact that u_{sc} solves $(\Delta + k^2)u_{\text{sc}} = 0$ in the exterior of a ball containing D , we have that $u_{\text{sc}}(x)$ is asymptotically equivalent, as $|x|$ grows, to

$$k^{(n-1)/2} \frac{e^{ik|x|}}{|x|^{(n-1)/2}} u^\infty(k, \theta, x/|x|),$$

where $u^\infty(k, \theta, x/|x|)$ is the *far-field pattern* and can be expressed as

$$u^\infty(k, \theta, x/|x|) = c_n \int_{\mathbb{R}^n} e^{-ik \frac{x}{|x|} \cdot y} V(y)(e^{ik\theta \cdot y} + u_{\text{sc}}(y)) dy$$

with c_n a constant only depending on n . Furthermore, we can conclude that the u_{sc} satisfies the SRC.

3.1. Resolvent estimates. Start by noting that whenever $f \in \mathcal{S}(\mathbb{R}^n)$, the identity

$$\langle \Psi_k^\pm \widehat{f}, \phi \rangle = \text{p.v.} \int_{\mathbb{R}^n} \frac{\widehat{f}(\xi)\phi(\xi)}{-|\xi|^2 + k^2} d\xi \mp i \frac{\pi}{2k} \int_{|\xi|=k} \widehat{f}(\xi)\phi(\xi) d\sigma_k(\xi)$$

holds for every bounded smooth function ϕ . Thus, according to (26), we have that

$$\mathcal{R}_k^\pm f(x) = \frac{1}{(2\pi)^{n/2}} \left(\text{p.v.} \int_{\mathbb{R}^n} \frac{e^{ix \cdot \xi} \widehat{f}(\xi)}{-|\xi|^2 + k^2} d\xi \mp i \frac{\pi}{2k} \int_{|\xi|=k} e^{ix \cdot \xi} \widehat{f}(\xi) d\sigma_k(\xi) \right)$$

for every $f \in \mathcal{S}(\mathbb{R}^n)$. For convenience, let us write

$$\mathcal{P}_k f(x) = \text{p.v.} \int_{\mathbb{R}^n} \frac{e^{ix \cdot \xi} \widehat{f}(\xi)}{-|\xi|^2 + k^2} d\xi, \quad \mathcal{Q}_k f(x) = \frac{1}{k} \int_{|\xi|=k} e^{ix \cdot \xi} \widehat{f}(\xi) d\sigma_k(\xi).$$

The main goal of this section is to prove the inequalities

$$(29) \quad \|\mathcal{P}_k f\|_{H_k^{1,-\delta}} \lesssim \|f\|_{L^{2,\delta}},$$

$$(30) \quad \|\mathcal{Q}_k f\|_{H_k^{1,-\delta}} \lesssim \|f\|_{L^{2,\delta}}$$

for every $f \in \mathcal{S}(\mathbb{R}^n)$, since Theorem 3.1 is a consequence of them. Indeed, these inequalities imply that

$$\begin{aligned} \|\mathcal{R}_k^\pm f\|_{H_k^{1-s,-\delta}} &= \|\mathcal{R}_k^\pm (k^2 - \Delta)^{-s/2} f\|_{H_k^{1,-\delta}} \\ &\lesssim \|(k^2 - \Delta)^{-s/2} f\|_{L^{2,\delta}} = \|f\|_{H_k^{-s,\delta}}. \end{aligned}$$

On the other hand, checking that

$$(31) \quad \|\mathcal{R}_k^\pm f\|_{H_k^{s,-\delta}} \lesssim k^{-(1-2s)} \|\mathcal{R}_k^\pm f\|_{H_k^{1-s,-\delta}},$$

we conclude the estimate

$$\|\mathcal{R}_k^\pm f\|_{H_k^{s,-\delta}} \lesssim k^{-(1-2s)} \|f\|_{H_k^{-s,\delta}},$$

and consequently the theorem. Note that (31) is a simple consequence of the following proposition.

Proposition 3.4. *Let s, t, δ and k be real numbers such that $s \leq t$. Then,*

$$\|f\|_{H_k^{s,\delta}} \lesssim k^{-(t-s)} \|f\|_{H_k^{t,\delta}}$$

for all $f \in \mathcal{S}(\mathbb{R}^n)$.

Proof. Firstly note that it is enough to prove that for $s, \delta, k \in \mathbb{R}$ with $s \geq 0$ the estimate

$$(32) \quad \|f\|_{H_k^{-s,\delta}} \lesssim k^{-s} \|f\|_{L^{2,\delta}}$$

holds for all $f \in \mathcal{S}(\mathbb{R}^n)$.

In order to prove (32), it will be convenient to introduce some notation. Set

$$D_0 = \{x \in \mathbb{R}^n : |x| \leq 1\},$$

$$D_j = \{x \in \mathbb{R}^n : 2^{j-1} < |x| \leq 2^j\} \quad j \in \mathbb{N} \setminus \{0\}$$

and consider χ_0 a smooth function with values in $[0, 1]$ such that $\text{supp } \chi_0 \subset D_0$ and $\chi_0(x) = 1$ for all $|x| \leq 1/2$. Let χ_j with $j \in \mathbb{N} \setminus \{0\}$ denote $\chi_j(x) = \chi_0(x/2^{j+1}) - \chi_0(x/2^j)$ and note that $\text{supp } \chi_j \subset \{2^{j-1} \leq |x| \leq 2^{j+1}\}$ for $j \geq 1$. By construction, $\sum_{j \in \mathbb{N}} \chi_j(x) = 1$ for every $x \in \mathbb{R}^n$.

Using this notation, we can write

$$\|f\|_{H_k^{-s,\delta}}^2 \lesssim \sum_{j \in \mathbb{N}} \left(\sum_{l \in \mathbb{N}} 2^{j\delta} \|(k^2 - \Delta)^{-s/2} (\chi_l f)\|_{L^2(D_j)} \right)^2.$$

The sums on the right hand side will be studied separately according to

$$\sum_{j \in \mathbb{N}} \left(\sum_{l \in \mathbb{N}} \cdots \right)^2 \sim \sum_{j \in \mathbb{N}} \sum_{|l-j| \leq 2} (\cdots)^2 + \sum_{j \in \mathbb{N}} \left(\sum_{|l-j| > 2} \cdots \right)^2.$$

Start by the case where, the support of χ_l and D_j intersect each other:

$$\sum_{j \in \mathbb{N}} \sum_{|l-j| \leq 2} 2^{2j\delta} \|(k^2 - \Delta)^{-s/2} (\chi_l f)\|_{L^2(D_j)}^2 \lesssim \sum_{l \in \mathbb{N}} 2^{2l\delta} \|(k^2 - \Delta)^{-s/2} (\chi_l f)\|_{L^2}^2$$

Using now Plancherel's identity, the right hand side of the previous inequality can be bounded by

$$k^{-2s} \sum_{l \in \mathbb{N}} 2^{2l\delta} \int_{\mathbb{R}^n} |\chi_l(x) f(x)|^2 dx \lesssim k^{-2s} \|f\|_{L^{2,\delta}}^2.$$

Consider now the case where the support of χ_l and D_j are far from each other. The norm $\|(k^2 - \Delta)^{-s/2} (\chi_l f)\|_{L^2(D_j)}$ will be estimated by duality:

$$(33) \quad \frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} e^{i(x-y) \cdot \xi} (k^2 + |\xi|^2)^{-s/2} f_l(y) dy d\xi g(x) dx$$

where g is any smooth function compactly supported in D_j and f_l denotes for simplicity $\chi_l f$. Let ϕ be a real-valued smooth function defined on $[0, \infty)$ with compact $\{0 \leq t \leq 1/4\}$ and such that $\phi(t) = 1$ for $0 \leq t \leq 1/8$. Then, (33) is equal to

$$\frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} e^{i(x-y) \cdot \xi} \left[1 - \phi\left(\frac{|x-y|}{|2^j - 2^l|}\right) \right] (k^2 + |\xi|^2)^{-s/2} f_l(y) dy d\xi g(x) dx.$$

This holds because if $x \in D_j$ and $y \in \text{supp } \chi_l$ with $|j - l| > 2$, then $|x - y| \geq 2^{\max(j,l)}/4 > |2^j - 2^l|/4$. Furthermore, using the identity

$$-\frac{\Delta_\xi e^{i(x-y) \cdot \xi}}{|x-y|^2} = e^{i(x-y) \cdot \xi}, \quad x \neq y$$

(N times) and integrating by parts ($2N$ times) in ξ , we see that (33) equals

$$-\frac{1}{(2\pi)^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} e^{i(x-y) \cdot \xi} \frac{1}{|x-y|^{2N}} \left[1 - \phi\left(\frac{|x-y|}{|2^j - 2^l|}\right) \right] \Delta_\xi^N (k^2 + |\xi|^2)^{-s/2} f_l(y) dy d\xi g(x) dx.$$

For $N > n/2$, we have that

$$\left| \int_{\mathbb{R}^n} e^{i(x-y) \cdot \xi} \Delta_\xi^N (k^2 + |\xi|^2)^{-s/2} d\xi \right| \lesssim k^{-s},$$

so we can apply Fubini to integrate first in ξ and then apply the Cauchy-Schwarz. Consequently, we have

$$\left| \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} \int_{\mathbb{R}^n} e^{i(x-y) \cdot \xi} (k^2 + |\xi|^2)^{-s/2} f_l(y) dy d\xi g(x) dx \right| \lesssim k^{-s} \|\Phi * f_l\|_{L^2} \|g\|_{L^2}$$

with

$$\Phi(x) = \frac{1}{|x|^{2N}} \left[1 - \phi\left(\frac{|x|}{|2^j - 2^l|}\right) \right].$$

This implies, by duality and Young's inequality, that

$$(34) \quad \sum_{j \in \mathbb{N}} \left(\sum_{|l-j|>2} 2^{j\delta} \|(k^2 - \Delta)^{-s/2} (\chi_l f)\|_{L^2(D_j)} \right)^2 \lesssim k^{-2s} \sum_{j \in \mathbb{N}} \left(\sum_{|l-j|>2} 2^{j\delta} \|\Phi\|_{L^1} \|\chi_l f\|_{L^2} \right)^2$$

where

$$\|\Phi\|_{L^1} = \frac{1}{|2^j - 2^l|^{2N-n}} \int_{\mathbb{R}^n} \frac{1}{|x|^{2N}} [1 - \phi(|x|)] dx.$$

The right-hand side of (34) can be bounded as follows

$$k^{-2s} \sum_{j \in \mathbb{N}} \left(\sum_{|l-j|>2} \frac{2^{(j-l)\delta}}{|2^j - 2^l|^{2N-n}} 2^{l\delta} \|\chi_l f\|_{L^2} \right)^2 \lesssim k^{-2s} \sum_{j \in \mathbb{N}} \left(\sum_{|l-j|>2} \frac{1}{2^{|j-l|(2N-n-|\delta|)}} 2^{l\delta} \|\chi_l f\|_{L^2} \right)^2.$$

To check the previous inequality, it may be convenient to notice that $|2^j - 2^l| > 2^{\max(j,l)-1}$ for $|l-j| > 0$. By Young's inequality for convolutions, the right-hand side can be bounded by

$$k^{-2s} \left(\sum_{j \in \mathbb{N}} 2^{-j(2N-n-|\delta|)} \right)^2 \sum_{l \in \mathbb{N}} 2^{2l\delta} \|\chi_l f\|_{L^2}^2.$$

Choosing $N > (n + |\delta|)/2$ we have that the right-hand side of (34) is bounded by

$$k^{-2s} \sum_{l \in \mathbb{N}} 2^{2l\delta} \|\chi_l f\|_{L^2}^2 \lesssim k^{-2s} \|f\|_{L^{2,\delta}}^2.$$

This ends the proof of this lemma. $\square$

We turn our attention to estimates (29) and (30):

Proof of the inequality (29). Let D_j with $j \in \mathbb{N}$ be as in the proof of Proposition 3.4 and bound

$$(35) \quad \|\mathcal{P}_k f\|_{H_k^{1,-\delta}} \lesssim \left(\sum_{j \in \mathbb{N}} 2^{-(2\delta-1)j} \right)^{1/2} \sup_{j \in \mathbb{N}} (2^{-j/2} \|(k^2 - \Delta)^{1/2} \mathcal{P}_k f\|_{L^2(D_j)}).$$

As we see below, the inequality (29) will be a consequence of the following lemma.

Lemma 3.5. *Let λ be a positive constant and $\mathcal{P}$ be given by*

$$\mathcal{P}f(x) = \text{p.v.} \int_{\mathbb{R}^n} \frac{e^{ix \cdot \xi} \widehat{f}(\xi)}{1 - |\xi|^2} d\xi.$$

Then,

$$\int_{|x|<R} |(1 - \Delta)^{1/2} \mathcal{P}f(x)|^2 dx \lesssim R\lambda \|f\|_{L^2}^2$$

for all $R > 0$ and $f \in \mathcal{S}(\mathbb{R}^n)$ such that $\text{supp } f \subset \{\lambda/2 \leq |x| \leq 2\lambda\}$.

Lemma 3.5 is a slight modification of Lemma 2.4 in [36] due to Kenig, Ponce and Vega.

Let us show that Lemma 3.5 implies (29). From the inequality (35) we have

$$(36) \quad \|\mathcal{P}_k f\|_{H_k^{1,-\delta}} \lesssim \sup_{R>0} (R^{-1/2} \|(k^2 - \Delta)^{1/2} \mathcal{P}_k f\|_{L^2(|x|\leq R)})$$

$$(37) \quad = k^{-n/2-1/2} \sup_{R>0} (R^{-1/2} \|(1 - \Delta)^{1/2} \mathcal{P}S_k f\|_{L^2(|x|\leq R)}),$$

where $S_k f(x) = f(x/k)$. Let χ be a smooth function with values in $[0, 1]$ such that $\text{supp } \chi \subset D_0$ and $\chi(x) = 1$ for all $|x| \leq 1/2$. Let χ_j with $j \in \mathbb{Z}$ denote $\chi_j(x) = \chi(x/2^{j+1}) - \chi(x/2^j)$ and note that $\text{supp } \chi_j \subset \{2^{j-1} \leq |x| \leq 2^{j+1}\}$. By construction, $\sum_{j \in \mathbb{Z}} \chi_j(x) = 1$ for every $x \in \mathbb{R}^n \setminus \{0\}$. Then, by Lemma 3.5, we have that

$$\begin{aligned} \|\mathcal{P}_k f\|_{H_k^{1,-\delta}} &\lesssim k^{-n/2-1/2} \sum_{j \in \mathbb{Z}} \sup_{R>0} (R^{-1/2} \|(1-\Delta)^{1/2} \mathcal{P} S_k(\chi_j f)\|_{L^2(|x|<R)}) \\ &\lesssim k^{-n/2-1/2} \sum_{j \in \mathbb{Z}} k^{1/2} 2^{j/2} \|S_k(\chi_j f)\|_{L^2} \\ &= \sum_{j \in \mathbb{Z}} 2^{j/2} \|\chi_j f\|_{L^2}. \end{aligned}$$

This last term can be manipulate to obtain inequality (29):

$$\begin{aligned} \sum_{j \in \mathbb{Z}} 2^{j/2} \|\chi_j f\|_{L^2} &\leq \sum_{j<0} 2^{j/2} \|f\|_{L^2(D_0)} + \sum_{j \in \mathbb{N}} 2^{j(1-2\delta)/2} 2^{j\delta} \|\chi_j f\|_{L^2} \\ &\lesssim \|f\|_{L^2(D_0)} + \left(\sum_{j \in \mathbb{N}} 2^{j(1-2\delta)} \right)^{1/2} \left(\sum_{j \in \mathbb{N}} 2^{2j\delta} \|\chi_j f\|_{L^2}^2 \right)^{1/2} \\ &\lesssim \|f\|_{L^{2,\delta}}. \end{aligned}$$

In order to prove (29), the only ingredient to be checked is Lemma 3.5.

Proof of Lemma 3.5. Start by writing

$$(I - \Delta)^{1/2} \mathcal{P} f = \mathcal{K} f + \mathcal{L} f,$$

where the operators $\mathcal{K}$ and $\mathcal{L}$ are given by

$$\begin{aligned} \mathcal{K} f(x) &= \text{p.v.} \int_{\mathbb{R}^n} e^{ix \cdot \xi} \frac{\phi(1 - |\xi|^2)(1 + |\xi|^2)^{1/2}}{1 - |\xi|^2} \widehat{f}(\xi) d\xi, \\ \mathcal{L} f(x) &= (2\pi)^{n/2} \mathcal{J}_1 m(D) f, \end{aligned}$$

where ϕ is a smooth function defined on $\mathbb{R}$ and taking values in $[0, 1]$ such that $\text{supp } \phi \subset [-1/(8n), 1/(8n)]$ and $\phi(t) = 1$ for all $|t| \leq 1/(16n)$, $\mathcal{J}_1$ is the Bessel potential defined as in the section 2.2, and $m(D)$ is the multiplier with symbol

$$m(\xi) = \frac{[1 - \phi(1 - |\xi|^2)](1 + |\xi|^2)}{1 - |\xi|^2}.$$

By Hölder's inequality with $1/2 = 1/p + 1/(2n)$ we have

$$\|\mathcal{L} f\|_{L^2(|x|<R)}^2 \lesssim R \|\mathcal{L} f\|_{L^p}^2.$$

By Sobolev embeddings⁴, the right-hand side is bounded by a multiple constant of

$$(38) \quad R \|m(D) f\|_{L^{p'}}^2$$

⁴This is nothing but Lemma 2 in §3.2 of chapter V in [57] with $\alpha = 1$, Young's inequality for functions convolved with finite measures and Theorem 1 in §1.2 of chapter V in [57] with $\alpha = 1$

with $1/p + 1/p' = 1$. By the Mikhlin–Hörmander multiplier theorem and Hölder’s inequality with $1/p' = 1/2 + 1/(2n)$, we have that (38) is bounded by a constant multiple of

$$R \|f\|_{L^{p'}}^2 \lesssim R\lambda \|f\|_{L^2}^2.$$

Therefore, we have

$$\|\mathcal{L}f\|_{L^2(|x|<R)}^2 \lesssim R\lambda \|f\|_{L^2}^2.$$

To finish the proof of Lemma 3.5, it only remains to prove the corresponding estimate for the singular part $\mathcal{K}$. To this end, we introduce a partition of unity $\{\psi_1, \dots, \psi_{2n}\}$ of $\Gamma = \{\xi \in \mathbb{R}^n : 1/2 < |\xi| < 3/2\}$ subordinated to $\{\Gamma_1, \dots, \Gamma_{2n}\}$, where

$$\Gamma_{2l-1} = \{\xi \in \Gamma : \xi_l > 1/(2\sqrt{n})\}, \quad \Gamma_{2l} = \{\xi \in \Gamma : \xi_l < -1/(2\sqrt{n})\}.$$

This partition of unity can be assumed to satisfy:

- (i) $\psi_2(x) = \psi_1(I_1x)$ with I_1 the reflection given by the matrix $[-e_1|e_2| \dots |e_n]$ and $e_1, \dots, e_n$ the elements of canonical base of $\mathbb{R}^n$;
- (ii) for $l = 2, \dots, n$,

$$\psi_{2l-1}(x) = \psi_1(I_lx), \quad \psi_{2l}(x) = \psi_2(I_lx)$$

with I_l the rotation given the matrix $[-e_l|e_2| \dots |e_{l-1}|e_1|e_{l+1}| \dots |e_n]$.

Thus, $\mathcal{K}f = \sum_{j=1}^{2n} \mathcal{K}_j f$ with

$$\mathcal{K}_j f(x) = \text{p.v.} \int_{\mathbb{R}^n} e^{ix \cdot \xi} \frac{\phi(1 - |\xi|^2)(1 + |\xi|^2)^{1/2}}{1 - |\xi|^2} \psi_j(\xi) \widehat{f}(\xi) d\xi,$$

and, in order to prove the lemma, it is enough to prove that

$$(39) \quad \|\mathcal{K}_1 f\|_{L^2(|x|<R)}^2 \lesssim R\lambda \|f\|_{L^2}^2$$

for all $f \in \mathcal{S}(\mathbb{R}^n)$ with $\text{supp } f \subset \{\lambda/2 \leq |x| \leq 2\lambda\}$.

In order to prove (39), we first write

$$\Phi(\xi) = \frac{\phi(1 - |\xi|^2)(1 + |\xi|^2)^{1/2}}{1 - |\xi|^2} \psi_1(\xi).$$

Note that ϕ has been chosen in such a way that it satisfies

$$\text{supp } \Phi \subset \{\xi = (\xi_1, \xi') \in \mathbb{R} \times \mathbb{R}^{n-1} : 1 - |\xi'|^2 > 1/(8n)\}$$

and for what follows we write

$$\Phi(\xi) = \frac{\Psi(\xi)}{(1 - |\xi'|^2)^{1/2} - \xi_1} \quad \text{with} \quad \Psi(\xi) = \frac{\phi(1 - |\xi|^2)(1 + |\xi|^2)^{1/2}}{(1 - |\xi'|^2)^{1/2} + \xi_1} \psi_1(\xi).$$

This also allows us to use the expression

$$\mathcal{K}_1 f(x) = \lim_{\epsilon \rightarrow 0} \int_{0 < \epsilon < |(1 - |\xi'|^2)^{1/2} - \xi_1|} e^{ix \cdot \xi} \Phi(\xi) \widehat{f}(\xi) d\xi.$$

Let $\mathcal{F}'$ denote the Fourier transform on the variable x' with dual variable ξ' . It follows that

$$\mathcal{K}_1 f(x) = \int_{\mathbb{R}^{n-1}} e^{ix' \cdot \xi'} \int_{\mathbb{R}} a(x_1, y_1, \xi') \mathcal{F}' f(y_1, \xi') dy_1 d\xi'$$

with

$$a(x_1, y_1, \xi') = \text{p.v.} \int_{\mathbb{R}} e^{i(x_1 - y_1)\xi_1} \frac{\Psi(\xi)}{(1 - |\xi'|^2)^{1/2} - \xi_1} d\xi_1$$

where the principal value is understood as $\lim_{\epsilon \rightarrow 0} \int_{0 < \epsilon < |(1 - |\xi'|^2)^{1/2} - \xi_1|}$. By Plancherel's identity we obtain

$$\begin{aligned} \|\mathcal{K}_1 f\|_{L^2(|x| < R)}^2 &\leq \int_{|x_1| < R} \int_{\mathbb{R}^{n-1}} |\mathcal{K}_1 f(x_1, x')|^2 dx' dx_1 \\ &\sim \int_{|x_1| < R} \int_{\mathbb{R}^{n-1}} \left| \int_{\mathbb{R}} a(x_1, y_1, \xi') \mathcal{F}' f(y_1, \xi') dy_1 \right|^2 d\xi' dx_1. \end{aligned}$$

Furthermore, since $\text{supp } f \subset \{\lambda/2 \leq |x| \leq 2\lambda\}$, we have by applying the Cauchy–Schwarz inequality, and then Plancherel's identity, that

$$\begin{aligned} \|\mathcal{K}_1 f\|_{L^2(|x| < R)}^2 &\lesssim \lambda \int_{|x_1| < R} \int_{\mathbb{R}^{n-1}} \sup_{\lambda/2 \leq |y_1| \leq 2\lambda} |a(x_1, y_1, \xi')|^2 \int_{\mathbb{R}} |\mathcal{F}' f(y_1, \xi')|^2 dy_1 d\xi' dx_1 \\ &\lesssim \lambda R \sup_{|x_1| < R} \sup_{\lambda/2 \leq |y_1| \leq 2\lambda} |a(x_1, y_1, \xi')|^2 \|f\|_{L^2}^2. \end{aligned}$$

Therefore, to conclude the proof of the inequality (39), it is enough to show that

$$(40) \quad \sup_{|x_1| < R} \sup_{\lambda/2 \leq |y_1| \leq 2\lambda} |a(x_1, y_1, \xi')| \lesssim 1.$$

Now recall the following identity

$$\frac{1}{(2\pi)^{1/2}} \text{p.v.} \int_{\mathbb{R}} \frac{e^{-ist}}{t} \varphi(t) dt = -i \frac{1}{2} \int_{\mathbb{R}} \text{sign}(s - t) \widehat{\varphi}(t) dt.$$

Consequently, $a(x_1, y_1, \xi')$ is a multiple constant of

$$e^{i(x_1 - y_1)(1 - |\xi'|^2)^{1/2}} \int_{\mathbb{R}} \text{sign}(x_1 - y_1 - t) \int_{\mathbb{R}} e^{-its} \Psi((1 - |\xi'|^2)^{1/2} - s, \xi') ds dt,$$

which can be written as a convolution changing variables according to $\xi_1 = (1 - |\xi'|^2)^{1/2} - s$:

$$\int_{\mathbb{R}} e^{i(x_1 - y_1 - t)(1 - |\xi'|^2)^{1/2}} \text{sign}(x_1 - y_1 - t) \int_{\mathbb{R}} e^{it\xi_1} \Psi(\xi) d\xi_1 dt.$$

Finally, by Young's inequality and then the non-stationary phase principle, we have

$$|a(x_1, y_1, \xi')| \lesssim \int_{\mathbb{R}} \left| \int_{\mathbb{R}} e^{it\xi_1} \Psi(\xi) d\xi_1 \right| dt \lesssim 1.$$

This proves (40), so that (39) holds and the proof of Lemma 3.5 is over. $\square$

Proof of the inequality (30). In the same way we obtained (36), we see that

$$\|\mathcal{Q}_k f\|_{H_k^{1, -\delta}} \lesssim \sup_{R > 0} (R^{-1/2} \|(k^2 - \Delta)^{1/2} \mathcal{Q}_k f\|_{L^2(|x| < R)}).$$

Note that $(k^2 - \Delta)^{1/2} \mathcal{Q}_k f(x)$ is a constant multiple of

$$\int_{|\xi|=k} e^{ix \cdot \xi} \widehat{f}(\xi) d\sigma_k(\xi) = \frac{1}{k} \int_{|\eta|=1} e^{ikx \cdot \eta} \widehat{S_k f}(\eta) d\sigma(\eta),$$

where $d\sigma$ is the volume form on $\{|\eta| = 1\}$ with $S_k f$ as in (37). Note that the previous expression is essentially the Fourier transform of the measure $\widehat{S_k f}(\eta) d\sigma(\eta)$ supported on $\{|\eta| = 1\}$. Because of the same considerations as in (37), we have that

$$\|\mathcal{Q}_k f\|_{H_k^{1,-\delta}} \lesssim k^{-n/2-1/2} \sup_{R>0} \left(R^{-1/2} \left| \int_{|\eta|=1} e^{ix \cdot \eta} \widehat{S_k f}(\eta) d\sigma(\eta) \right|_{L^2(|x|<R)} \right).$$

By Theorem 7.1.26 in [29], we can bound the right hand side of the previous inequality and obtain

$$\|\mathcal{Q}_k f\|_{H_k^{1,-\delta}} \lesssim k^{-n/2-1/2} \left(\int_{|\eta|=1} |\widehat{S_k f}(\eta)|^2 d\sigma(\eta) \right)^{1/2}.$$

By an appropriate duality argument (in the spirit of Theorem 14.1.1 in [30]) applied to Theorem 7.1.26 in [29] we have that

$$\left(\int_{|\eta|=1} |\widehat{S_k f}(\eta)|^2 d\sigma(\eta) \right)^{1/2} \lesssim \sum_{j \in \mathbb{Z}} k^{1/2} 2^{j/2} \left(\int_{k2^{j-1} < |x| \leq k2^j} |S_k f(x)|^2 dx \right)^{1/2}.$$

By performing a rescaling in k , we have

$$\|\mathcal{Q}_k f\|_{H_k^{1,-\delta}} \lesssim k^{-n/2-1/2} k^{1/2} k^{n/2} \sum_{j \in \mathbb{Z}} 2^{j/2} \left(\int_{2^{j-1} < |x| \leq 2^j} |f(x)|^2 dx \right)^{1/2}.$$

The right-hand side can be obviously manipulated as follows to obtain inequality (30)

$$\begin{aligned} \sum_{j \in \mathbb{Z}} 2^{j/2} \left(\int_{2^{j-1} < |x| \leq 2^j} |f(x)|^2 dx \right)^{1/2} &\lesssim \sum_{j \leq 0} 2^{j/2} \left(\int_{|x| \leq 1} |f(x)|^2 dx \right)^{1/2} \\ &\quad + \left(\sum_{j \geq 1} 2^{j(1-2\delta)} \right)^{1/2} \left(\sum_{j \geq 1} 2^{2j\delta} \int_{D_j} |f(x)|^2 dx \right)^{1/2} \\ &\lesssim \|f\|_{L^{2,\delta}} \end{aligned}$$

with D_j as in the proof of Proposition 3.4.

3.2. Boundedness of the multiplication operator. Consider any two $f, g \in \mathcal{S}(\mathbb{R}^n)$, the multiplication by V is defined by

$$\langle Vf, g \rangle = \langle V, fg \rangle,$$

where the brackets denote the corresponding dualities. Let χ_D be a compactly supported smooth function with values in $[0, 1]$ such that $\chi(x) = 1$ for all $x \in D$. Since $V \in L_{-s}^p(\mathbb{R}^n)$ with support in D , there exists $W \in L^p(\mathbb{R}^n)$ such that $V = (I - \Delta)^{s/2} W$, and consequently

$$\langle Vf, g \rangle = \langle W, (I - \Delta)^{s/2} (f_D g_D) \rangle$$

with $f_D = \chi_D f$ and $g_D = \chi_D g$. Let ϕ_ε be as in (12) and set $W^\sharp = \phi_\varepsilon * W$ and $W^\flat = W - W^\sharp$. For $s - n/p < t < s$ we have that

$$\langle Vf, g \rangle = \langle (I - \Delta)^{t/2} W^\sharp, (I - \Delta)^{(s-t)/2} (f_D g_D) \rangle + \langle W^\flat, (I - \Delta)^{s/2} (f_D g_D) \rangle$$

and, by Hölder's inequality,

$$(41) \quad |\langle Vf, g \rangle| \leq \|(I - \Delta)^{t/2} W^\sharp\|_{L^q} \|(I - \Delta)^{(s-t)/2} (f_D g_D)\|_{L^{q'}} + \|W^\flat\|_{L^p} \|(I - \Delta)^{s/2} (f_D g_D)\|_{L^{p'}}$$

with $q = n/(s - t)$ and q' and p' the dual exponents of q and p , respectively. Since $q > p$, we have by Young's inequality that

$$\|(I - \Delta)^{t/2} W^\sharp\|_{L^q} \lesssim \varepsilon^{-t+n/q-n/p} \|W\|_{L^p}.$$

Furthermore, we will prove in Lemma 3.6 that

$$(42) \quad \|(I - \Delta)^{(s-t)/2} (fg)\|_{L^{q'}} \lesssim \|f_D\|_{H^{s-t}} \|g_D\|_{H^{s-t}}$$

$$(43) \quad \|(I - \Delta)^{s/2} (fg)\|_{L^{p'}} \lesssim \|f_D\|_{H^s} \|g_D\|_{H^s},$$

hence, for $k \geq 1$, we have

$$(44) \quad \begin{aligned} |\langle Vf, g \rangle| &\lesssim \varepsilon^{-t+n/q-n/p} \|W\|_{L^p} \|f_D\|_{H_k^{s-t}} \|g_D\|_{H_k^{s-t}} + \|W^b\|_{L^p} \|f_D\|_{H_k^s} \|g_D\|_{H_k^s} \\ &\lesssim \varepsilon^{-t+n/q-n/p} \|W\|_{L^p} \|f\|_{H_k^{s-t,-\delta}} \|g\|_{H_k^{s-t,-\delta}} + \|W^b\|_{L^p} \|f\|_{H_k^{s,-\delta}} \|g\|_{H_k^{s,-\delta}}. \end{aligned}$$

The estimate (44) will be justified by Lemma 3.7. Finally, by Proposition 3.4, we have

$$|\langle Vf, g \rangle| \lesssim (\varepsilon^{-t+n/q-n/p} k^{-2t} \|W\|_{L^p} + \|W^b\|_{L^p}) \|f\|_{H_k^{s,-\delta}} \|g\|_{H_k^{s,-\delta}}.$$

Choosing $\varepsilon = k^{-1/2}$, a simple duality argument show that the estimate stated in Proposition 3.2 holds.

In order to end the proof of this proposition, we need to show inequalities (42), (43) and (44). Note that they follow from the next lemmas:

Lemma 3.6. *Let $s > 0$ and $p \in (1, \infty)$ satisfy $p \geq n/s$. Then,*

$$\|fg\|_{L_s^{p'}} \lesssim \|f\|_{H^s} \|g\|_{H^s}$$

for all $f, g \in \mathcal{S}(\mathbb{R}^n)$ with $1/p + 1/p' = 1$.

Proof. The Kato–Ponce inequality (see [25] and the original reference [35]) ensures that

$$\|fg\|_{L_s^{p'}} \lesssim \|f\|_{H^s} \|g\|_{L^r} + \|f\|_{L^r} \|g\|_{H^s}$$

for $1/p' = 1/2 + 1/r$, which is equivalent to $1/2 = 1/p + 1/r$. By the Sobolev embeddings, we have that

$$\|fg\|_{L_s^{p'}} \lesssim \|f\|_{H^s} \|g\|_{H^t} + \|f\|_{H^t} \|g\|_{H^s}$$

with $t - n/2 = -n/r$. Since $t = n/2 - n/r = n/p \leq s$, the estimate claimed in the statement follows immediately from the last one. $\square$

Lemma 3.7. *Let s and δ be positive constants and $\phi \in \mathcal{S}(\mathbb{R}^n)$. Then,*

$$\|\phi f\|_{H_k^s} \lesssim \|f\|_{H_k^{s,-\delta}}$$

for all $f \in \mathcal{S}(\mathbb{R}^n)$.

Proof. Note that it is enough to prove the lemma for $\delta = 2N$ with $N \in \mathbb{N} \setminus \{0\}$. In this case, we just need to show that

$$(45) \quad \|T_N g\|_{H_k^s} \lesssim \|g\|_{L^2}$$

for all $g \in \mathcal{S}(\mathbb{R}^n)$, where

$$T_N g(x) = \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{ix \cdot \xi} a_N(x, \xi) \widehat{g}(\xi) d\xi$$

and

$$a_N(x, \xi) = \phi(x) e^{-ix \cdot \xi} (I - \Delta_\xi)^N (e^{ix \cdot \xi} (k^2 + |\xi|^2)^{-s/2}).$$

In order to check this claim, it suffices to test the inequality (45) for the functions

$$g = (1 + |\cdot|^2)^{-N} (k^2 - \Delta)^{s/2} f$$

with f any function in $\mathcal{S}(\mathbb{R}^n)$, and note that

$$T_N((1 + |\cdot|^2)^{-N} (k^2 - \Delta)^{s/2} f) = \phi f.$$

The inequality (45) follows from two general results for pseudodifferential operators. To apply them, we first observe that a_N is a smooth function in $\mathbb{R}^n \times \mathbb{R}^n$ and

$$|\partial_x^\beta \partial_\xi^\alpha a_N(x, \xi)| \leq A_{\alpha, \beta, N} (k + |\xi|)^{-s - |\alpha|}$$

for all multi-indices α and β . Then, by symbolic calculus for pseudodifferential operators (for example Theorem 2 of §3 in chapter VI of [58]) we see that

$$(46) \quad \|T_N g\|_{H_k^s} \leq \|S_0 g\|_{L^2} + \|S_{-1} g\|_{L^2},$$

where

$$S_j g(x) = \frac{1}{(2\pi)^{n/2}} \int_{\mathbb{R}^n} e^{ix \cdot \xi} b_j(x, \xi) \widehat{g}(\xi) d\xi$$

with

$$b_0(x, \xi) = \phi(x) e^{-ix \cdot \xi} (k^2 + |\xi|^2)^{s/2} (1 - \Delta_\xi)^N (e^{ix \cdot \xi} (k^2 + |\xi|^2)^{-s/2}),$$

and $|\partial_x^\beta \partial_\xi^\alpha b_{-1}(x, \xi)| \leq A_{\alpha, \beta, N} (k + |\xi|)^{-1 - |\alpha|}$. By the L^2 boundedness of pseudodifferential operators (for example Theorem 1 of §3 in chapter VI of [58]), we have that the right-hand side of (46) can be bounded as

$$\|S_0 g\|_{L^2} + \|S_{-1} \circ (k^2 - \Delta)^{1/2} ((k^2 - \Delta)^{-1/2} g)\|_{L^2} \lesssim \|g\|_{L^2} + \|g\|_{H_k^{-1}} \lesssim \|g\|_{L^2}.$$

This ends the proof of this lemma. $\square$

4. RECONSTRUCTION OF THE LOCAL STRENGTH

Let us recall that our aim is to reconstruct μ , the local strength of q , from one single realization of the measurement data

$$M(\tau, \theta) = \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m u^\infty(k, \theta, -\theta) \overline{u^\infty(k + \tau, \theta, -\theta)} dk$$

given for some set of $\tau \geq 0$ and $\theta \in \mathbb{S}^{n-1}$. Recall that, by Proposition 2.4, we have $q \in L_s^p(\mathbb{R}^n)$ with $1 < p < \infty$ and $s < (m - n)/2$ almost surely. In the section 3, we have studied the direct scattering theory for a potential under slightly more general regularity assumptions, and shown that the backscattering far-field pattern

$$(47) \quad u^\infty(k, \theta, -\theta) = c_n \int_{\mathbb{R}^n} e^{ik\theta \cdot y} q(y) (e^{ik\theta \cdot y} + u_{\text{sc}}(y)) dy$$

can be expressed with the Born series (28) of the scattered wave for $k \geq k_0(\omega)$ almost surely, where k_0 is the threshold wavelength in Remark 3.3. Here c_n is a constant depending only on n .

Below, we give proofs to Theorems 1.1 and 1.2. Notice that for the full non-linear inverse scattering problem in Theorem 1.1, we restrict to the case $n = m = 3$. This condition will be needed in sections 4.2 and 4.3.

In a nutshell, the idea behind the reconstruction of μ consists of writing u^∞ as the Born series

$$u^\infty(k, \theta, -\theta) = \sum_{j \in \mathbb{N} \setminus \{0\}} u_j^\infty(k, \theta, -\theta)$$

for $k \geq k_0(\omega)$ and $\theta \in \mathbb{S}^{n-1}$ in such a way that the higher order terms

$$(48) \quad \frac{1}{K} \int_K^{2K} k^m u_j^\infty(k, \theta, -\theta) \overline{u_l^\infty(k + \tau, \theta, -\theta)} dk$$

with $j + l \geq 3$ are negligible in comparison with the realization of

$$(49) \quad \frac{1}{K} \int_K^{2K} k^m u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)} dk.$$

Next, one proceeds by proving that the limit of (49) as $K \rightarrow \infty$ provides enough information to reconstruct μ , when given at multiple values of τ and θ . Recall that according to (28) we have

$$(50) \quad u_j^\infty(k, \theta, -\theta) = c_n \int_{\mathbb{R}^n} e^{ik\theta \cdot y} q(y) (\mathcal{R}_k^+ \circ q)^{j-1} u_0(y) dy$$

where $u_0(x) = e^{ik\theta \cdot x}$ and c_n is the same as in (47). With the expression (50) the connection of measurement data and the statistics of q becomes apparent.

For the sake of clarity, we first reconstruct the local strength assuming $\mathbb{E}q = 0$, and then in the section 4.4 we consider the more general case where $\mathbb{E}q$ is smooth and supported in D .

4.1. Single backscattering. In order to prove that the limit of (49) coincides almost surely with a deterministic function, we need to apply suitable ergodicity arguments. The following theorem (available in [18, p. 94]) provides a useful condition:

Theorem 4.1. *Let X_t with $t \geq 0$ be a real-valued stochastic process with continuous paths and zero-mean $\mathbb{E}X_t = 0$. Assume that for some positive constants c, ϵ the condition*

$$|\mathbb{E}(X_t X_{t+r})| \leq c(1+r)^{-\epsilon}$$

holds for all $t \geq 0$ and $r \geq 0$. Then,

$$\lim_{T \rightarrow \infty} \frac{1}{T} \int_0^T X_t dt = 0 \quad \text{almost surely.}$$

Note that under the same assumptions of this theorem,

$$\lim_{T \rightarrow \infty} \frac{1}{T} \int_T^{2T} X_t dt = 0 \quad \text{almost surely.}$$

The single backscattering $u_1^\infty(k, \theta, -\theta)$ in the formula (49) is described by a constant multiple of

$$\langle q, e^{i2k\theta \cdot y} \rangle = U_k + iV_k,$$

where U_k and V_k denote the real and imaginary parts. We can rewrite the product in (49) as a constant multiple of

$$(51) \quad 2\langle q, e^{i2k\theta \cdot y} \rangle \overline{\langle q, e^{i2(k+\tau)\theta \cdot y} \rangle} = (1+i)(U_k^2 + U_{k+\tau}^2 + V_k^2 + V_{k+\tau}^2) - (U_k - U_{k+\tau})^2 - (V_k - V_{k+\tau})^2 - i(U_k + V_{k+\tau})^2 - i(V_k - U_{k+\tau})^2.$$

Let W_k denote any of random variables

$$(52) \quad U_k, U_{k+\tau}, V_k, V_{k+\tau}, U_k - U_{k+\tau}, V_k - V_{k+\tau}, U_k + V_{k+\tau}, V_k - U_{k+\tau}.$$

Using Theorem 4.1, we will prove that

$$(53) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m (W_k^2 - \mathbb{E}W_k^2) dk = 0$$

almost surely, and hence we deduce by (51) that

$$(54) \quad \begin{aligned} \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)} dk \\ = \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m \mathbb{E}(u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}) dk \end{aligned}$$

almost surely.

According to Theorem 4.1, identity (53) holds if there exists an $\varepsilon > 0$ such that

$$(55) \quad |\mathbb{E}(k^m (W_k^2 - \mathbb{E}W_k^2)(k+r)^m (W_{k+r}^2 - \mathbb{E}W_{k+r}^2))| \lesssim (1+r)^{-\varepsilon}$$

for all $k \geq 1$ and $r \geq 0$. In order to verify condition (55), we observe that (W_k, W_{k+r}) is always a centred Gaussian random vector (which holds because of (6) and $\mathbb{E}q = 0$), and we use the following lemma.

Lemma 4.2. *Let X and Y be two random variables such that the pair (X, Y) is a Gaussian random vector. If $\mathbb{E}X = \mathbb{E}Y = 0$, then*

$$\mathbb{E}((X^2 - \mathbb{E}X^2)(Y^2 - \mathbb{E}Y^2)) = 2(\mathbb{E}(XY))^2.$$

Proof. Note that one can assume that the variances be $\mathbb{E}X^2 = \mathbb{E}Y^2 = 1$. Thus, one is reduced to prove

$$(56) \quad \mathbb{E}((X^2 - 1)(Y^2 - 1)) = 2(\mathbb{E}(XY))^2$$

for X and Y satisfying $\mathbb{E}X^2 = \mathbb{E}Y^2 = 1$. In order to prove (56), it is enough to show that it holds for a Gaussian vector (X, Y') having the same probability law as (X, Y) . Indeed,

$$\begin{aligned} \mathbb{E}((X^2 - 1)(Y^2 - 1)) &= \int_{\mathbb{R}^2} (x^2 - 1)(y^2 - 1) d\mathbb{P}_{(X,Y)}(x, y) \\ &= \int_{\mathbb{R}^2} (x^2 - 1)(y^2 - 1) d\mathbb{P}_{(X,Y')}(x, y) = \mathbb{E}((X^2 - 1)((Y')^2 - 1)) \end{aligned}$$

and

$$\mathbb{E}(XY) = \int_{\mathbb{R}^2} xy d\mathbb{P}_{(X,Y)}(x, y) = \int_{\mathbb{R}^2} xy d\mathbb{P}_{(X,Y')}(x, y) = \mathbb{E}(XY').$$

Let X' be a Gaussian random variable with mean 0, variance 1 and independent of X . Consider $Y' = \cos \alpha X + \sin \alpha X'$ with $\cos \alpha = \mathbb{E}(XY)$, which is possible since $|\mathbb{E}(XY)| \leq 1$.

Note that Y' is a Gaussian random variable and the pair (X, Y') is a Gaussian random vector. Moreover, since

$$\begin{pmatrix} \mathbb{E}X \\ \mathbb{E}Y \end{pmatrix} = \begin{pmatrix} \mathbb{E}X \\ \mathbb{E}Y' \end{pmatrix}, \quad \begin{pmatrix} \mathbb{E}X^2 & \mathbb{E}(XY) \\ \mathbb{E}(XY) & \mathbb{E}Y^2 \end{pmatrix} = \begin{pmatrix} \mathbb{E}X^2 & \mathbb{E}(XY') \\ \mathbb{E}(XY') & \mathbb{E}(Y')^2 \end{pmatrix},$$

the Gaussian vector (X, Y) and (X, Y') are equally distributed. Therefore, it only remains to show that (56) holds for (X, Y') , but this is a simple computation that can be verified using that $\mathbb{E}(XY') = \cos \alpha$,

$$\begin{aligned} \mathbb{E}((X^2 - \mathbb{E}X^2)((Y')^2 - \mathbb{E}(Y')^2)) &= (\cos \alpha)^2 \mathbb{E}X^4 - (\cos \alpha)^2 + (\sin \alpha)^2 \mathbb{E}(X^2(X')^2) - (\sin \alpha)^2 \\ &\quad + 2 \cos \alpha \sin \alpha \mathbb{E}(X^3 X') \end{aligned}$$

and

$$\mathbb{E}X^4 = 3, \quad \mathbb{E}(X^2(X')^2) = \mathbb{E}X^2 \mathbb{E}(X')^2 = 1, \quad \mathbb{E}(X^3 X') = \mathbb{E}X^3 \mathbb{E}X' = 0.$$

This concludes the proof. $\square$

Using Lemma 4.2, the identity (53) holds if there exist constants $c > 0$, which may depend on τ , and $\varepsilon > 0$ such that

$$(57) \quad |\mathbb{E}(k^{m/2}(k+r)^{m/2}W_k W_{k+r})| \leq c(1+r)^{-\varepsilon}.$$

Note that the inequality (57) is a consequence of the following proposition and therefore (53) holds and consequently (54).

Proposition 4.3. *Suppose the potential satisfies $\mathbb{E}q = 0$. Then, there exists a known constant $c_{n,m}$, depending on n and m , such that*

$$(58) \quad \mathbb{E}(u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}) = c_{n,m} k^{-m} \widehat{\mu}(2\tau\theta) + \mathcal{O}(k^{-m-1})$$

for $k \geq 1/2$ and $\tau \geq 0$. Moreover, for all $k_1, k_2 > 0$, we have that

$$(59) \quad |\mathbb{E}(U_{k_1} U_{k_2})| \lesssim k_1^{-m} (1 + |k_1 - k_2|)^{-N}$$

$$(60) \quad |\mathbb{E}(V_{k_1} V_{k_2})| \lesssim k_1^{-m} (1 + |k_1 - k_2|)^{-N}$$

$$(61) \quad |\mathbb{E}(U_{k_1} V_{k_2})| \lesssim k_1^{-m} (1 + |k_1 - k_2|)^{-N}$$

for all $N \in \mathbb{N}$.

Proof. Note that

$$(62) \quad \mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q, e^{i2k_2\theta \cdot x} \rangle}) = \int_{\mathbb{R}^n} \left(\int_{\mathbb{R}^n} K_q(x, y) e^{-i2k_1\theta \cdot (x-y)} dy \right) e^{-i2(k_2-k_1)\theta \cdot x} dx.$$

By (11), we have that

$$\begin{aligned} &\mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q, e^{i2k_2\theta \cdot x} \rangle}) \\ &= \int_{\mathbb{R}^n} \mu(x) |2k_1|^{-m} e^{-i2(k_2-k_1)\theta \cdot x} dx + \int_D a(x, 2k_1\theta) e^{-i2(k_2-k_1)\theta \cdot x} dx \end{aligned}$$

for $k_1 \geq 1/2$, which implies (58).

On the other hand, by (62) and (10), we have

$$\mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q, e^{i2k_2\theta \cdot x} \rangle}) = \int_{\mathbb{R}^n} c_q(x, 2k_1\theta) \chi(x) e^{-i2(k_2-k_1)\theta \cdot x} dx$$

with $\chi \in C_0^\infty(\mathbb{R}^n)$ such that $\chi(x) = 1$ for all $x \in D$, which implies, by the non-stationary phase principle, that

$$(63) \quad |\mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q, e^{i2k_2\theta \cdot x} \rangle})| \lesssim (1 + k_1)^{-m} (1 + |k_1 - k_2|)^{-N}$$

for all $N \in \mathbb{N}$. By the same kind of considerations, one can prove that

$$\mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \langle q, e^{i2k_2\theta \cdot x} \rangle) = \int_{\mathbb{R}^n} c_q(x, 2k_1\theta) \chi(x) e^{i2(k_1+k_2)\theta \cdot x} dx$$

with $\chi \in C_0^\infty(\mathbb{R}^n)$ such that $\chi(x) = 1$ for all $x \in D$, which implies, again by the non-stationary phase principle, that

$$(64) \quad |\mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \langle q, e^{i2k_2\theta \cdot x} \rangle)| \lesssim (1 + k_1)^{-m} (1 + k_1 + k_2)^{-N}$$

for all $N \in \mathbb{N}$. Finally, the estimates (59), (60) and (61) follow from (63) and (64). $\square$

As a consequence of the identities (54) and (58), the information of the local strength of q provided by the single backscattering can be recorded as follows.

Corollary 4.4. *Suppose $\mathbb{E}q = 0$. The Fourier transform of the local strength of q can be recovered from the single backscattering as follows:*

$$(65) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)} dk = c_{n,m} \widehat{\mu}(2\tau\theta)$$

almost surely for any fixed $\tau \geq 0$ and $\theta \in \mathbb{S}^{n-1}$.

Proof of Theorem 1.2 in the case $\mathbb{E}q = 0$. Corollary 4.4 connects the measurement data to the Fourier transform of μ at a point $2\tau\theta$. We can now proceed by repeating the same measurement at a countable dense set $\{(\tau_j, \theta_j)\}_{j \in \mathbb{N}} \subset \mathbb{R}_+ \times \mathbb{S}^2$. The almost sure convergence takes place simultaneously on this countable dense set. Since μ is smooth with compact support, it follows that $\widehat{\mu}$ is in the Schwartz class, in particular is continuous. Then, we can recover $\widehat{\mu}$ everywhere from the countable dense set extending by continuity. Finally, the inverse formula of the Fourier transform provides μ . $\square$

The following consequence of identity (51) and Proposition 4.3 is not needed for Corollary 4.4. However, it will allow us to study the convergence speed of integral (65) in appendix A.

Corollary 4.5. *Suppose $\mathbb{E}q = 0$ and let $Z : \Omega \times \mathbb{R}_+ \rightarrow \mathbb{C}$ be the random process defined by*

$$Z(k) = u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}.$$

It follows that

$$(66) \quad |\mathbb{E}[(Z(k) - \mathbb{E}Z(k))(Z(k + r) - \mathbb{E}Z(k + r))]| \lesssim (1 + \min(r, |r - \tau|))^{-N}$$

$$(67) \quad |\mathbb{E}[(Z(k) - \mathbb{E}Z(k)) \overline{(Z(k + r) - \mathbb{E}Z(k + r))})]| \lesssim (1 + \min(r, |r - \tau|))^{-N}$$

for any N .

Proof. We observe that using identity (51) one can write the expectation in (67) (as well as (66)) as a product of sums of random variables appearing in formula (52). Multiplying the terms and applying triangle inequality yields a bound to the right hand side of (67) given as a sum of terms

$$|\mathbb{E}[(W_k^2 - \mathbb{E}W_k^2)((W'_{k+r})^2 - \mathbb{E}(W'_{k+r})^2)]| = 2(\mathbb{E}W_k W'_{k+r})^2,$$

where W_k and W'_{k+r} represent random variables given in (52) and we used Lemma 4.2 for the identity. Finally, the result is obtained by applying Proposition 4.3 and the fact that the cross-covariance is computed at all pairs of frequencies $k_1 \in \{k, k + \tau\}$ and $k_2 \in \{k + r, k + r + \tau\}$. Same reasoning applies to inequality (67). $\square$

4.2. 2nd order backscattering. We now consider the interactions between the single and 2nd order backscattering for the case $n = m = 3$ and show that they are negligible. Recall that in this particular case $q \in L^p_{-s}(\mathbb{R}^3)$ for $1 < p < \infty$ and $s > 0$ almost surely. According to (50), the 2nd order backscattering $u_2^\infty(k, \theta, -\theta)$ is described by a constant multiple of

$$(68) \quad \int_{\mathbb{R}^3} \int_{\mathbb{R}^3} e^{ik\theta \cdot (x+y)} q(x)q(y)\Phi_k^+(x-y) dx dy,$$

where Φ_k^+ is the fundamental solution which satisfies the outgoing Sommerfeld radiation condition. Note that the integrals in (68) have to be understood as a distributional pairing. The interaction to be considered now correspond to the terms (48) with $j + l = 3$ and $j = l = 2$ and the goal is to show that they are negligible, more precisely to prove the following statement:

Proposition 4.6. *Let us assume that $\mathbb{E}q = 0$. We have that, for every $\tau \geq 0$ and every $\theta \in \mathbb{S}^2$,*

$$\lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 u_j^\infty(k, \theta, -\theta) \overline{u_l^\infty(k + \tau, \theta, -\theta)} dk = 0$$

almost surely whenever $j + l = 3$ or $j = l = 2$.

Let us prove this proposition. By the Cauchy-Schwarz inequality and changing variables, the modulus of (48) can be bounded by

$$\left(\frac{1}{K} \int_K^{2K} k^3 |u_j^\infty(k, \theta, -\theta)|^2 dk \right)^{1/2} \left(\frac{1}{K} \int_{K+\tau}^{2K+\tau} k^3 |u_l^\infty(k, \theta, -\theta)|^2 dk \right)^{1/2}.$$

Thus, in order to study the interactions between the single and 2nd order scattering is enough to consider

$$\frac{1}{K} \int_K^{2K} k^3 |u_j^\infty(k, \theta, -\theta)|^2 dk$$

for any $K \geq 1$ with $j = 1, 2$. After identity (54) and (63), we know that

$$(69) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 |u_1^\infty(k, \theta, -\theta)|^2 dk < \infty$$

almost surely. Therefore, in order to prove Proposition 4.6, it is enough to check that

$$(70) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 |u_2^\infty(k, \theta, -\theta)|^2 dk = 0$$

almost surely or alternatively,

$$(71) \quad \lim_{K \rightarrow \infty} \frac{1}{K-1} \int_1^K k^3 |u_2^\infty(k, \theta, -\theta)|^2 dk = 0$$

almost surely. Note that

$$\frac{1}{K-1} \int_1^K k^3 |u_2^\infty(k, \theta, -\theta)|^2 dk = \int_1^\infty \frac{k \mathbf{1}_{[1, K]}(k)}{K-1} k^2 |u_2^\infty(k, \theta, -\theta)|^2 dk,$$

where $\mathbf{1}_{[1, K]}$ denotes the characteristic function of the interval $[1, K]$. Since $k \mathbf{1}_{[1, K]}(k)/(K-1)$ converges point-wise to zero as K goes to infinity, we have by the dominate convergence theorem that, if

$$\int_1^\infty k^2 |u_2^\infty(k, \theta, -\theta)|^2 dk < \infty,$$

then (71) holds. Obviously, by the continuity of the function $\theta \in \mathbb{S}^2 \mapsto u_2^\infty(k, \theta, -\theta)$, it will be enough to show that

$$(72) \quad \mathbb{E} \int_{\mathbb{S}^2} \int_1^\infty k^2 |u_2^\infty(k, \theta, -\theta)|^2 dk d\sigma(\theta) < \infty,$$

where $d\sigma$ denotes the volume form on $\mathbb{S}^2$. In order to prove (72), we will show that

$$(73) \quad \limsup_{\varepsilon \rightarrow 0} \mathbb{E} \int_{\mathbb{S}^2} \int_1^\infty k^2 |v_\varepsilon(k, \theta)|^2 dk d\sigma(\theta) < \infty$$

with

$$v_\varepsilon(k, \theta) = \int_{\mathbb{R}^3} \int_{\mathbb{R}^3} e^{ik\theta \cdot (x+y)} q_\varepsilon(x) q_\varepsilon(y) \Phi_k^+(x-y) dx dy,$$

where $q_\varepsilon(\omega, x) = \langle q(\omega), \varphi_\varepsilon(x - \cdot) \rangle$ and φ_ε is as in (12). Then, as a consequence of Fatou's lemma, we see that (72) holds.

Let us prove (73). Since in dimension $n = 3$, $\Phi_k^+(x)$ is given by a constant multiple of $e^{ik|x|}/|x|$, the approximation of the 2nd order backscattering $v_\varepsilon(k, \theta)$ is

$$(74) \quad \int_{\mathbb{R}^3} \int_{\mathbb{R}^3} e^{ik[\theta \cdot (x+y) + |x-y|]} \frac{q_\varepsilon(x) q_\varepsilon(y)}{|x-y|} dx dy.$$

Changing variables, (74) becomes

$$\int_{\mathbb{R}^3} \int_{\mathbb{S}^2} \int_0^\infty e^{ik(\theta \cdot z + \rho)} q_\varepsilon\left(\frac{z + \rho\omega}{2}\right) q_\varepsilon\left(\frac{z - \rho\omega}{2}\right) \rho d\rho d\sigma(\omega) dz.$$

Denoting

$$(75) \quad f_\varepsilon(z, \rho) = \rho \mathbf{1}_{[0, \infty)}(\rho) \int_{\mathbb{S}^2} q_\varepsilon\left(\frac{z + \rho\omega}{2}\right) q_\varepsilon\left(\frac{z - \rho\omega}{2}\right) d\sigma(\omega),$$

we have that

$$v_\varepsilon(k, \theta) = \int_{\mathbb{R}^3} \int_{\mathbb{R}} e^{ik(\theta \cdot z + \rho)} f_\varepsilon(z, \rho) d\rho dz.$$

Note that there exists an R which depends on D such that

$$(76) \quad |z| + |\rho| = \left| \frac{z + \rho\omega}{2} + \frac{z - \rho\omega}{2} \right| + \left| \frac{z + \rho\omega}{2} - \frac{z - \rho\omega}{2} \right| \leq |z + \rho\omega| + |z - \rho\omega| \leq R$$

whenever

$$\frac{z + \rho\omega}{2}, \frac{z - \rho\omega}{2} \in \bigcup_{\varepsilon \in (0,1]} \text{supp } q_\varepsilon.$$

Note that

$$(77) \quad \int_{\mathbb{S}^2} k^2 |v_\varepsilon(k, \theta)|^2 d\sigma(\theta) \sim \int_{\mathbb{S}^2} \left| \int_{\mathbb{R}^3} e^{ik\theta \cdot z} T f_\varepsilon(z, k) dz \right|^2 k^2 d\sigma(\theta)$$

with

$$T f_\varepsilon(z, k) = \int_{\mathbb{R}} e^{ik\rho} f_\varepsilon(z, \rho) d\rho.$$

Lemma 4.7. *The right hand-side of (77) can be bounded by above as follows:*

$$\int_{\mathbb{S}^2} \left| \int_{\mathbb{R}^3} e^{ik\theta \cdot z} T f_\varepsilon(z, k) dz \right|^2 k^2 d\sigma(\theta) \lesssim \int_{\mathbb{R}^3} |T f_\varepsilon(z, k)|^2 dz$$

almost surely, where the implicit constant depends on D .

Proof. Start by noting that the term to be estimated by above can be rewritten as

$$(78) \quad \int_{S_k} \left| \int_{\mathbb{R}^3} e^{i\theta \cdot z} T f_\varepsilon(z, k) dz \right|^2 d\sigma_k(\theta)$$

with $S_k = \{x \in \mathbb{R}^3 : |x| = k\}$ and $d\sigma_k$ denoting its volume form. The term (78) is equivalent to the square of the L^2 norm of the Fourier transform of $T f_\varepsilon(\cdot, k)$ restricted to S_k . We will estimate it by duality: let g be a smooth function on S_k , then

$$\int_{S_k} g(\theta) \int_{\mathbb{R}^3} e^{i\theta \cdot z} T f_\varepsilon(z, k) dz d\sigma_k(\theta) = \int_{\mathbb{R}^3} \int_{S_k} e^{i\theta \cdot z} g(\theta) d\sigma_k(\theta) T f_\varepsilon(z, k) dz.$$

Since $\text{supp } T f_\varepsilon(\cdot, k) \subset \{|z| < R\}$ with R as in (76), we have by Cauchy-Schwarz that

$$(79) \quad \begin{aligned} & \left| \int_{S_k} g(\theta) \int_{\mathbb{R}^3} e^{i\theta \cdot z} T f_\varepsilon(z, k) dz d\sigma_k(\theta) \right| \\ & \leq \left\| \int_{S_k} e^{i\theta \cdot z} g(\theta) d\sigma_k(\theta) \right\|_{L^2(|z| < R)} \|T f_\varepsilon(\cdot, k)\|_{L^2} \end{aligned}$$

almost surely. Consider $\{\chi_j : j = 1, \dots, 6\}$ a partition of unity of S_k subordinated to the sets

$$\Gamma_{2l} = \{x \in S_k : 2\sqrt{3}x_l > k\}, \quad \Gamma_{2l-1} = \{x \in S_k : 2\sqrt{3}x_l < -k\}$$

for $l = 1, 2, 3$. Then,

$$(80) \quad \int_{S_k} e^{i\theta \cdot z} g(\theta) d\sigma_k(\theta) = \sum_{j=1}^6 \int_{\mathbb{R}^2} e^{i\vartheta_j(y) \cdot z} g_j(\vartheta_j(y)) \frac{k}{\sqrt{k^2 - |y|^2}} dy,$$

where $g_j = \chi_j g$ and

$$\begin{aligned} \vartheta_1(y) &= (-\sqrt{k^2 - |y|^2}, y_1, y_2), & \vartheta_2(y) &= (\sqrt{k^2 - |y|^2}, y_1, y_2), \\ \vartheta_3(y) &= (y_1, -\sqrt{k^2 - |y|^2}, y_2), & \vartheta_4(y) &= (y_1, \sqrt{k^2 - |y|^2}, y_2), \\ \vartheta_5(y) &= (y_1, y_2, -\sqrt{k^2 - |y|^2}), & \vartheta_6(y) &= (y_1, y_2, \sqrt{k^2 - |y|^2}). \end{aligned}$$

Note that every term on the sum of (80) is a multiple of the two-dimensional Fourier transform of

$$y \mapsto e^{i(-1)^{j(l)} z_l \sqrt{k^2 - |y|^2}} g_{j(l)}(\vartheta_{j(l)}(y)) \frac{k}{\sqrt{k^2 - |y|^2}}$$

evaluated at $\widehat{z}_l$, where $\widehat{z}_1 = (z_2, z_3)$, $\widehat{z}_2 = (z_1, z_3)$ and $\widehat{z}_3 = (z_1, z_2)$, and $j(l)$ stand for $2l - 1$ or $2l$ with $l = 1, 2, 3$.

Letting $d\widehat{z}_1, d\widehat{z}_2$ and $d\widehat{z}_3$ denote $dz_2 dz_3, dz_1 dz_3$ and $dz_1 dz_2$ respectively, we have, by the Plancherel identity in $\mathbb{R}^2$, that

$$(81) \quad \int_{\mathbb{R}^2} \left| \int_{\mathbb{R}^2} e^{i\vartheta_{j(l)}(y) \cdot z} g_{j(l)}(\vartheta_{j(l)}(y)) \frac{k}{\sqrt{k^2 - |y|^2}} dy \right|^2 d\widehat{z}_l \lesssim \int_{\mathbb{R}^2} |g_{j(l)}(\vartheta_{j(l)}(y))|^2 \frac{k}{\sqrt{k^2 - |y|^2}} dy \\ \leq \int_{S_k} |g(\theta)|^2 d\sigma_k(\theta),$$

since

$$\frac{k}{\sqrt{k^2 - |y|^2}} < 2\sqrt{3} \quad \text{for } y \in \text{supp } g_{j(l)} \circ \vartheta_{j(l)}.$$

Therefore, from (80) and (81), we conclude that

$$(82) \quad \int_{|z| < R} \left| \int_{S_k} e^{i\theta \cdot z} g(\theta) d\sigma_k(\theta) \right|^2 dz \lesssim R \int_{S_k} |g(\theta)|^2 d\sigma_k(\theta).$$

Finally, by duality, we can ensure that (78) is almost surely bounded by

$$\|Tf_\varepsilon(\cdot, k)\|_{L^2}^2$$

with a constant which depends on R . Therefore, the lemma is proven. $\square$

After (77) and Lemma 4.7 we obtain

$$\int_{\mathbb{S}^2} \int_1^\infty k^2 |v_\varepsilon(k, \theta)|^2 dk d\sigma(\theta) \lesssim \int_{\mathbb{R}^3} \int_1^\infty |Tf_\varepsilon(z, k)|^2 dk dz$$

almost surely with an implicit constant depending on the domain $D \subset \mathbb{R}^3$. Note that Tf_ε is a constant multiple of the inverse Fourier transform of f_ε in the variable ρ , so we have that

$$(83) \quad \int_{\mathbb{S}^2} \int_1^\infty k^2 |v_\varepsilon(k, \theta)|^2 dk d\sigma(\theta) \lesssim \int_{\mathbb{R}^3} \int_{\mathbb{R}} |f_\varepsilon(z, \rho)|^2 d\rho dz$$

almost surely by the Plancherel identity. By (83) and (76), we can conclude that

$$(84) \quad \limsup_{\varepsilon \rightarrow 0} \mathbb{E} \int_{\mathbb{S}^2} \int_1^\infty k^2 |v_\varepsilon(k, \theta)|^2 dk d\sigma(\theta) \lesssim \lim_{\varepsilon \rightarrow 0} \int_{|z| < R} \int_{|\rho| < R} \mathbb{E} |f_\varepsilon(z, \rho)|^2 d\rho dz,$$

provided that the limit on the right-hand side exists. We now show that the limit exists and this equals

$$\begin{aligned} \int_{|z|<R} \int_{0<\rho<R} \rho^2 & \left[\int_{\mathbb{S}^2} \int_{\mathbb{S}^2} K_q\left(\frac{z+\rho\omega}{2}, \frac{z-\rho\omega}{2}\right) K_q\left(\frac{z+\rho\theta}{2}, \frac{z-\rho\theta}{2}\right) d\sigma(\omega) d\sigma(\theta) \right. \\ & + \int_{\mathbb{S}^2} \int_{\mathbb{S}^2} K_q\left(\frac{z+\rho\omega}{2}, \frac{z+\rho\theta}{2}\right) K_q\left(\frac{z-\rho\omega}{2}, \frac{z-\rho\theta}{2}\right) d\sigma(\omega) d\sigma(\theta) \\ & \left. + \int_{\mathbb{S}^2} \int_{\mathbb{S}^2} K_q\left(\frac{z+\rho\omega}{2}, \frac{z-\rho\theta}{2}\right) K_q\left(\frac{z-\rho\omega}{2}, \frac{z+\rho\theta}{2}\right) d\sigma(\omega) d\sigma(\theta) \right] d\rho dz. \end{aligned}$$

Before proving this claim, note that this already ensures that (73) holds, since $K_q(x, y) = c_{n,m} \mu(x) \log|x-y| + F_\alpha(x, y)$, according to the Proposition 2.8, and hence K_q is integrable over $\mathbb{S}^2 \times \mathbb{S}^2$.

Finally, we show that the limit of the right-hand side of (84) exists and we compute it. Start by noting that (75) makes $\mathbb{E}|f_\varepsilon(z, \rho)|^2$ be equal to

$$\rho^2 \mathbf{1}_{[0,\infty)}(\rho) \int_{\mathbb{S}^2} \int_{\mathbb{S}^2} \mathbb{E}\left[q_\varepsilon\left(\frac{z+\rho\omega}{2}\right) q_\varepsilon\left(\frac{z-\rho\omega}{2}\right) q_\varepsilon\left(\frac{z+\rho\theta}{2}\right) q_\varepsilon\left(\frac{z-\rho\theta}{2}\right)\right] d\sigma(\omega) d\sigma(\theta).$$

By the Isserlis' theorem, this equals

$$\begin{aligned} \rho^2 \mathbf{1}_{[0,\infty)}(\rho) & \left[\int_{\mathbb{S}^2} \int_{\mathbb{S}^2} \mathbb{E}\left[q_\varepsilon\left(\frac{z+\rho\omega}{2}\right) q_\varepsilon\left(\frac{z-\rho\omega}{2}\right)\right] \mathbb{E}\left[q_\varepsilon\left(\frac{z+\rho\theta}{2}\right) q_\varepsilon\left(\frac{z-\rho\theta}{2}\right)\right] d\sigma(\omega) d\sigma(\theta) \right. \\ & + \int_{\mathbb{S}^2} \int_{\mathbb{S}^2} \mathbb{E}\left[q_\varepsilon\left(\frac{z+\rho\omega}{2}\right) q_\varepsilon\left(\frac{z+\rho\theta}{2}\right)\right] \mathbb{E}\left[q_\varepsilon\left(\frac{z-\rho\omega}{2}\right) q_\varepsilon\left(\frac{z-\rho\theta}{2}\right)\right] d\sigma(\omega) d\sigma(\theta) \\ & \left. + \int_{\mathbb{S}^2} \int_{\mathbb{S}^2} \mathbb{E}\left[q_\varepsilon\left(\frac{z+\rho\omega}{2}\right) q_\varepsilon\left(\frac{z-\rho\theta}{2}\right)\right] \mathbb{E}\left[q_\varepsilon\left(\frac{z-\rho\omega}{2}\right) q_\varepsilon\left(\frac{z+\rho\theta}{2}\right)\right] d\sigma(\omega) d\sigma(\theta) \right]. \end{aligned}$$

Recalling (7), it is a simple observation to note that $\mathbb{E}[q_\varepsilon(x)q_\varepsilon(y)]$ converges to $K_q(x, y)$ pointwise as ε vanishes. Moreover, using again that $K_q(x, y) = c_{n,m} \mu(x) \log|x-y| + F_\alpha(x, y)$, we can check that

$$|\mathbb{E}[q_\varepsilon(x)q_\varepsilon(y)]| \lesssim |\log|x-y|| \mathbf{1}_{\{|x-y|<1\}}(x, y) + 1$$

assuming $\text{supp } \varphi \subset \{|x| \leq 1/4\}$, which can always be assumed. Hence, by applying the dominated convergence theorem to the integral we have that the limit of the right hand-side of (84) exists and is the one claimed above.

To sum up, we have shown that (73) holds and consequently,

$$\lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 |u_2^\infty(k, \theta, -\theta)|^2 dk = 0$$

almost surely for every $\theta \in \mathbb{S}^2$. This ends the proof of Proposition 4.6.

4.3. Multiple backscattering. Again, we only consider the effects of multiple scattering under the assumptions $m = n = 3$. In that case, the realizations of q are almost surely in

$L^p_{-s}(\mathbb{R}^3)$ for $1 < p < \infty$ and $s > 0$. According to the previous decomposition on single and 2nd order backscattering, we can write

$$\begin{aligned} & \frac{1}{K} \int_K^{2K} k^3 u^\infty(k, \theta, -\theta) \overline{u^\infty(k + \tau, \theta, -\theta)} dk \\ &= \sum_{1 \leq j \leq 2} \sum_{1 \leq l \leq 2} \frac{1}{K} \int_K^{2K} k^3 u_j^\infty(k, \theta, -\theta) \overline{u_l^\infty(k + \tau, \theta, -\theta)} dk \\ &+ \frac{1}{K} \int_K^{2K} k^3 \left(u^\infty(k, \theta, -\theta) - \sum_{1 \leq j \leq 2} u_j^\infty(k, \theta, -\theta) \right) \overline{u^\infty(k + \tau, \theta, -\theta)} dk \\ &+ \sum_{1 \leq j \leq 2} \frac{1}{K} \int_K^{2K} k^3 u_j^\infty(k, \theta, -\theta) \left(\overline{u^\infty(k + \tau, \theta, -\theta)} - \sum_{1 \leq l \leq 2} \overline{u_l^\infty(k + \tau, \theta, -\theta)} \right) dk. \end{aligned}$$

The first term on the right-hand side of the previous inequality corresponds to the single and 2nd order backscattering terms studied in the sections 4.1 and 4.2. The other two terms describe the multiple backscattering and will be shown here that they are negligible, that is, they vanish as K grows. Applying the Cauchy–Schwarz inequality and changing variable as in the proof of Proposition 4.6, we only need to check that

$$(85) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 |u^\infty(k, \theta, -\theta)|^2 dk < \infty,$$

$$(86) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 \left| u^\infty(k, \theta, -\theta) - \sum_{1 \leq l \leq 2} u_l^\infty(k, \theta, -\theta) \right|^2 dk = 0$$

hold almost surely. Note that (85) follows from (86), (69) and (70). Thus, it is enough to show that (86) holds. This is a straight consequence of the next lemma:

Lemma 4.8. *Suppose that $\mathbb{E}q = 0$. We have that*

$$\sup_{\theta \in \mathbb{S}^2} \left| u^\infty(k, \theta, -\theta) - \sum_{1 \leq l \leq 2} u_l^\infty(k, \theta, -\theta) \right| = o(k^{-2(1-3s)})$$

almost surely.

Proof. By (47) and (50), and then Lemma 3.6, we have that

$$\begin{aligned} \left| u^\infty(k, \theta, -\theta) - \sum_{1 \leq l \leq 2} u_l^\infty(k, \theta, -\theta) \right| &\lesssim \left| \int_{\mathbb{R}^3} e^{ik\theta \cdot y} q(y) (u_{\text{sc}}(y) - u_1(y)) dy \right| \\ &\lesssim \|q\|_{L^p_{-s}} \|e^{ik\theta \cdot y} \chi\|_{H^s} \|\chi(u_{\text{sc}} - u_1)\|_{H^s}, \end{aligned}$$

where χ is a smooth function with compact support such that $\chi(x) = 1$ for every x in a ball containing D . A direct computation shows that $\|e^{ik\theta \cdot y} \chi\|_{H^s} = \mathcal{O}(k^s)$. On the other hand,

$$\|\chi(u_{\text{sc}} - u_1)\|_{H^s} \leq \|\chi(u_{\text{sc}} - u_1)\|_{H_k^s} \lesssim \|u_{\text{sc}} - u_1\|_{H_k^{s, -\delta}}$$

by Lemma 3.7. By Theorem 3.1 and Proposition 3.2, we know that the operator $\mathcal{R}_k^+ \circ q$ maps $H_k^{s, -\delta}(\mathbb{R}^3)$ into itself, for $k \geq k_0(\omega)$ almost surely, with a norm

$$\|\mathcal{R}_k^+ \circ q\|_{H_k^{s, -\delta} \rightarrow H_k^{s, -\delta}} = o(k^{-(1-2s)}).$$

Therefore, $u_{\text{sc}} - u_1 = \sum_{j>1} (\mathcal{R}_k^+ \circ q)^j u_0$ can be bounded as follows:

$$\|u_{\text{sc}} - u_1\|_{H_k^{s,-\delta}} \leq \sum_{j>1} \|\mathcal{R}_k^+ \circ q\|_{H_k^{s,-\delta} \rightarrow H_k^{s,-\delta}}^j \|\chi u_0\|_{H_k^{s,-\delta}} = o(k^s k^{-2(1-2s)})$$

almost surely. This concludes the proof of this lemma. $\square$

Proof of Theorem 1.1 in the case $\mathbb{E}q = 0$. Following the discussion in the beginning of this section we notice that we finally possess all necessary tools to prove the main theorem in the case $\mathbb{E}q = 0$. Having established the well-posedness of the forward problem and the measurement in the sections 2 and 3, we have shown in Corollary 4.4 that the first order contribution $M_1(\tau, \theta)$ coincides almost surely (and up to a multiplicative constant) with $\hat{\mu}(2\tau\theta)$ at fixed $\tau \geq 0$ and $\theta \in \mathbb{S}^2$. In Proposition 4.6 and this section, we have proven that in $\mathbb{R}^3$ the contribution from the second and higher order scattering vanishes in $M(\tau, \theta, -\theta)$. We can now proceed similarly to the proof of Theorem 1.2 and repeat the measurement at a countable dense set $\{(\tau_j, \theta_j)\}_{j \in \mathbb{N}} \subset \mathbb{R}_+ \times \mathbb{S}^2$. Finally, continuation from a dense set yields the result. $\square$

4.4. Non-zero-mean potentials. This section is devoted to extend the proof in sections 4.1, 4.2 and 4.3 to the case of non-zero-mean potentials. We proceed pointing out the places where some changes have to be made.

4.4.1. Single backscattering. As in the section 4.1, the first goal is to show that the identity (54) holds when $\mathbb{E}q \neq 0$. As in the zero-mean case, this will be a consequence of the fact that

$$(87) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m (W_k^2 - \mathbb{E}W_k^2) dk = 0$$

almost surely for W_k as in (52). The difference now is that $\mathbb{E}W_k \neq 0$, so we write

$$W_k^2 - \mathbb{E}W_k^2 = Z_k^2 - \mathbb{E}Z_k^2 + 2\mathbb{E}W_k Z_k$$

with $Z_k = W_k - \mathbb{E}W_k$. We will prove that (87) holds showing that

$$(88) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m (Z_k^2 - \mathbb{E}Z_k^2) dk = 0 \quad \text{and}$$

$$(89) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^m \mathbb{E}W_k Z_k dk = 0$$

almost surely. In order to check (88), we use Theorem 4.1 verifying that there exists an $\varepsilon > 0$ such that

$$(90) \quad |\mathbb{E}(k^m (Z_k^2 - \mathbb{E}Z_k^2) (k+r)^m (Z_{k+r}^2 - \mathbb{E}Z_{k+r}^2))| \lesssim (1+r)^{-\varepsilon}$$

for all $k \geq 1$ and $r \geq 0$. To do so, we observe that (Z_k, Z_{k+r}) is a centred Gaussian random vector and use Lemma 4.2. Thus, instead of condition (90) we just need to see that there exists $\varepsilon > 0$ such that

$$(91) \quad |\mathbb{E}(k^{m/2} (k+r)^{m/2} Z_k Z_{k+r})| \lesssim (1+r)^{-\varepsilon}.$$

The inequality (91) follows from (93), (94) and (95) in the next proposition, and therefore (88) holds.

Proposition 4.9. *Let q be the potential given by Definition 2.1. There exists a known constant $c_{n,m}$, depending on n and m , such that*

$$(92) \quad \mathbb{E}(u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}) = c_{n,m} k^{-m} \widehat{\mu}(2\tau\theta) + \mathcal{O}(k^{-m-1})$$

for $k \geq 1/2$ and $\tau \geq 0$. Moreover, for all $k_1, k_2 > 0$, we have that

$$(93) \quad |\mathbb{E}((U_{k_1} - \mathbb{E}U_{k_1})(U_{k_2} - \mathbb{E}U_{k_2}))| \lesssim k_1^{-m}(1 + |k_1 - k_2|)^{-N}$$

$$(94) \quad |\mathbb{E}((V_{k_1} - \mathbb{E}V_{k_1})(V_{k_2} - \mathbb{E}V_{k_2}))| \lesssim k_1^{-m}(1 + |k_1 - k_2|)^{-N}$$

$$(95) \quad |\mathbb{E}((U_{k_1} - \mathbb{E}U_{k_1})(V_{k_2} - \mathbb{E}V_{k_2}))| \lesssim k_1^{-m}(1 + |k_1 - k_2|)^{-N}$$

for all $N \in \mathbb{N}$.

Proof. Note that

$$(96) \quad \begin{aligned} & \mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle}) \\ &= \int_{\mathbb{R}^n} \left(\int_{\mathbb{R}^n} K_q(x, y) e^{-i2k_1\theta \cdot (x-y)} dy \right) e^{-i2(k_2-k_1)\theta \cdot x} dx. \end{aligned}$$

By (11), we have that

$$(97) \quad \begin{aligned} & \mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle}) \\ &= \int_{\mathbb{R}^n} \mu(x) |2k_1|^{-m} e^{-i2(k_2-k_1)\theta \cdot x} dx + \int_D a(x, 2k_1\theta) e^{-i2(k_2-k_1)\theta \cdot x} dx \end{aligned}$$

for $k_1 \geq 1/2$. Furthermore,

$$(98) \quad \begin{aligned} & \mathbb{E}(\langle q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q, e^{i2k_2\theta \cdot x} \rangle}) \\ &= \mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle}) + \langle \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle} \end{aligned}$$

with $\mathbb{E}q$ smooth and compactly supported. By (97) and the non-stationary phase principle, we have that (92) holds.

On the other hand, by (96) and (10), we have

$$\mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle}) = \int_{\mathbb{R}^n} c_q(x, 2k_1\theta) \chi(x) e^{-i2(k_2-k_1)\theta \cdot x} dx$$

with $\chi \in C_0^\infty(\mathbb{R}^n)$ such that $\chi(x) = 1$ for all $x \in D$, which implies, by the non-stationary phase principle, that

$$(99) \quad |\mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle})| \lesssim (1 + k_1)^{-m}(1 + |k_1 - k_2|)^{-N}$$

for all $N \in \mathbb{N}$. By the same kind of considerations, one can prove that

$$\mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle}) = \int_{\mathbb{R}^n} c_q(x, 2k_1\theta) \chi(x) e^{i2(k_1+k_2)\theta \cdot x} dx$$

with $\chi \in C_0^\infty(\mathbb{R}^n)$ such that $\chi(x) = 1$ for all $x \in D$, which implies, again by the non-stationary phase principle, that

$$(100) \quad |\mathbb{E}(\langle q - \mathbb{E}q, e^{i2k_1\theta \cdot y} \rangle \overline{\langle q - \mathbb{E}q, e^{i2k_2\theta \cdot x} \rangle})| \lesssim (1 + k_1)^{-m}(1 + k_1 + k_2)^{-N}$$

for all $N \in \mathbb{N}$. Finally, the estimates (93), (94) and (95) follow from (99) and (100). $\square$

Eventually, we prove that (89) holds. To do so, we check that the quasi-orthogonality condition in Theorem 4.1 holds: there exists $\varepsilon > 0$ such that

$$(101) \quad |\mathbb{E}(k^m \mathbb{E} W_k Z_k (k+r)^m \mathbb{E} W_{k+r} Z_{k+r})| \lesssim (1+r)^{-\varepsilon}.$$

Indeed, by (91) we have

$$|\mathbb{E}(k^m \mathbb{E} W_k Z_k (k+r)^m \mathbb{E} W_{k+r} Z_{k+r})| \lesssim (1+r)^{-\varepsilon} k^{m/2} (k+r)^{m/2} \mathbb{E} W_k \mathbb{E} W_{k+r}.$$

Since $\mathbb{E}q$ is smooth and compactly supported, we have by the non-stationary phase principle that

$$k^{m/2} (k+r)^{m/2} \mathbb{E} W_k \mathbb{E} W_{k+r} \lesssim 1,$$

and consequently (101). Then, by Theorem 4.1, we can ensure that (89) holds.

Summarizing, (88) and (89) imply that (87) holds, and so does (54). On the other hand, by (92) we can conclude that Corollary 4.4 also holds in the case that $\mathbb{E}q$ is smooth and compactly supported in D . Finally, Theorem 1.2 follows in the case $\mathbb{E}q \neq 0$ by the same density argument performed in the section 4.1.

4.4.2. 2nd order backscattering. The goal of this section is to prove that Proposition 4.6 holds when $\mathbb{E}q$ is smooth and has support in D . As in the case $\mathbb{E}q = 0$, it is enough to show that

$$(102) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 |u_1^\infty(k, \theta, -\theta)|^2 dk < \infty \quad \text{and}$$

$$(103) \quad \lim_{K \rightarrow \infty} \frac{1}{K} \int_K^{2K} k^3 |u_2^\infty(k, \theta, -\theta)|^2 dk = 0$$

hold almost surely. The finiteness of the limit in (102) is a consequence of (98), (99) and the non-stationary phase principle. However, showing that (103), or alternatively

$$\lim_{K \rightarrow \infty} \frac{1}{K-1} \int_1^K k^3 |u_2^\infty(k, \theta, -\theta)|^2 dk = 0,$$

holds almost surely requires a more subtle argument. Fortunately, this is exactly the same as in the case $\mathbb{E}q = 0$, and it reduces to prove that

$$(104) \quad \lim_{\varepsilon \rightarrow 0} \int_{|z| < R} \int_{|\rho| < R} \mathbb{E} |f_\varepsilon(z, \rho)|^2 d\rho dz$$

exists and is finite, with $f_\varepsilon(z, \rho)$ as in (75). It is convenient to write

$$\begin{aligned} f_\varepsilon(z, \rho) = & \rho \mathbf{1}_{[0, \infty)}(\rho) \int_{\mathbb{S}^2} (q_\varepsilon - \mathbb{E} q_\varepsilon) \left(\frac{z + \rho \omega}{2} \right) (q_\varepsilon - \mathbb{E} q_\varepsilon) \left(\frac{z - \rho \omega}{2} \right) d\sigma(\omega) \\ & + 2\rho \mathbf{1}_{[0, \infty)}(\rho) \int_{\mathbb{S}^2} (q_\varepsilon - \mathbb{E} q_\varepsilon) \left(\frac{z + \rho \omega}{2} \right) \mathbb{E} q_\varepsilon \left(\frac{z - \rho \omega}{2} \right) d\sigma(\omega) \\ & + \rho \mathbf{1}_{[0, \infty)}(\rho) \int_{\mathbb{S}^2} \mathbb{E} q_\varepsilon \left(\frac{z + \rho \omega}{2} \right) \mathbb{E} q_\varepsilon \left(\frac{z - \rho \omega}{2} \right) d\sigma(\omega). \end{aligned}$$

Thus,

$$\mathbb{E} |f_\varepsilon(z, \rho)|^2 = \rho^2 \mathbf{1}_{[0, \infty)}(\rho) \int_{\mathbb{S}^2} \int_{\mathbb{S}^2} \sum_{j=1}^4 I_j(z, \rho, \omega, \theta) d\sigma(\omega) d\sigma(\theta)$$

with

$$\begin{aligned}
I_4 &= \mathbb{E} \left[(q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\omega}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z - \rho\omega}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\theta}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z - \rho\theta}{2} \right) \right], \\
I_3 &= 4\mathbb{E} \left[(q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\omega}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z - \rho\omega}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\theta}{2} \right) \right] \mathbb{E}q_\varepsilon \left(\frac{z - \rho\theta}{2} \right), \\
I_2 &= 2\mathbb{E} \left[(q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\omega}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z - \rho\omega}{2} \right) \right] \mathbb{E}q_\varepsilon \left(\frac{z + \rho\theta}{2} \right) \mathbb{E}q_\varepsilon \left(\frac{z - \rho\theta}{2} \right) \\
&\quad + 4\mathbb{E} \left[(q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\omega}{2} \right) (q_\varepsilon - \mathbb{E}q_\varepsilon) \left(\frac{z + \rho\theta}{2} \right) \right] \mathbb{E}q_\varepsilon \left(\frac{z - \rho\omega}{2} \right) \mathbb{E}q_\varepsilon \left(\frac{z - \rho\theta}{2} \right), \\
I_1 &= \mathbb{E}q_\varepsilon \left(\frac{z + \rho\omega}{2} \right) \mathbb{E}q_\varepsilon \left(\frac{z - \rho\omega}{2} \right) \mathbb{E}q_\varepsilon \left(\frac{z + \rho\theta}{2} \right) \mathbb{E}q_\varepsilon \left(\frac{z - \rho\theta}{2} \right).
\end{aligned}$$

Applying Isserlis' theorem on the terms I_3 and I_4 , we see that $I_3 = 0$ and I_4 can be studied in the same way as we did at the end of the section 4.2. Finally, the terms with I_2 and I_1 give no problem when showing the existence and finiteness of (104) since $\mathbb{E}[(q_\varepsilon - \mathbb{E}q_\varepsilon)(x)(q_\varepsilon - \mathbb{E}q_\varepsilon)(y)]$ converges to $K_q(x, y)$ pointwise as ε vanishes and

$$|\mathbb{E}[(q_\varepsilon - \mathbb{E}q_\varepsilon)(x)(q_\varepsilon - \mathbb{E}q_\varepsilon)(y)]| \lesssim 1 + |\log |x - y|| \mathbf{1}_{\{|x - y| < 1\}}(x, y).$$

4.4.3. Multiple backscattering. As we argued in the section 4.3, we just needed to show that (85) and (86) hold almost surely. Note that (85) follows from (86), (102) and (103). In turn, (86) follows from Lemma 4.8. As discussed in the proof of Theorem 1.1 for the case $\mathbb{E}q = 0$, we can conclude now this theorem for non-zero-mean potentials.

APPENDIX A. SCALING REGIMES

A.1. Two-scale model of a non-smooth scatterer. In this paper we have set up a theoretical framework for the inversion of the probabilistic backscattering problem that hopefully helps to inspire stable imaging algorithms. Here, we briefly consider two practical quantities—the correlation length and the characteristic size of the potential—and how their multi-scale analysis can be used to estimate the statistics of error in the reconstruction in terms of these quantities. Unlike in the previous sections, *our analysis is not rigorous in this appendix* as we apply several approximations often used in multi-scale analysis in physics literature. In this appendix our aim is to formulate our previous results in the terminology used in multi-scale analysis. Also, we will also show how our microlocal techniques can be applied also in much more general setting than is done in the main text of the paper when one adds several scales of orders in the scattering model and makes certain approximations. Indeed, using those approximations, we can consider the multi-scale analysis in arbitrary dimension n and for general potential models.

Before we formulate the results, we make a few remarks. First, in any practical measurement device the maximum frequency K in the data is bounded, i.e., we can never have infinite precision. In order to estimate how wide frequency band is needed to achieve reasonable accuracy we need to understand effective scales of frequency correlations in the backscattered far field. Second, in physical applications one often has some prior information about the correlations within the random scatterer. This information is typically described via so-called correlation length of the random media. Third, if the potential q has a form $q(x) = \epsilon Q(x)$

and the random field Q is roughly of a constant order, we say that the parameter $\epsilon > 0$ is the *characteristic size* of q . The effect of higher order scattering in the Schrödinger equation

$$(\Delta + k^2 + \epsilon Q(x))u(x) = 0$$

becomes negligible (problem is approximately linear) if the characteristic size ϵ is small compared to other parameters appearing in the system.

In what follows, we employ scaling regimes that are characterized by three quantities described above: the (effective) length of the *frequency band* K , the *correlation length* ℓ and the *characteristic size* of the potential ϵ . Moreover, we consider the following scaling regime with respect to the relative sizes of K, ℓ and ϵ : for the reference wavelength K_0 , distance of propagation L_0 and correlation ℓ_0 scales we have

$$(105) \quad K \gg K_0 \max \left(\left(\frac{\ell}{\ell_0} \right)^{-\beta_1}, \left(\frac{L}{L_0} \right)^n \right) \quad \text{and} \quad \epsilon \ll \left(\frac{K}{K_0} \right)^{-\beta_2},$$

with some $\beta_1 > 1$ and $\beta_2 > m/2 - 1$, where $L = \text{diam}(D)$ is the diameter of the domain where the potential q is supported and m is the order of the covariance operator of Q . For convenience, we assume below that the reference scales K_0 and ℓ_0 are of constant order can be neglected from the error estimate analysis.

In the scaling regime given in (105) we show below that the measurement data with a finite frequency band,

$$(106) \quad M_K(\theta, \tau, \omega) = \frac{1}{K} \int_K^{2K} k^m u^\infty(k, \theta, -\theta) \overline{u^\infty(k + \tau, \theta, -\theta)} dk,$$

is close to the ideal data $M(\theta, \tau, \omega)$ in (3). Therefore, the local strength $\mu(x)$ can be approximately estimated from $M_K(\theta, \tau, \omega)$.

At this point, notice carefully that in the specific case $n = m = 3$ the higher order scattering can be analysed rigorously (as in the section 4) and we need no prior assumptions regarding the characteristic size ϵ , that is, the latter inequality in (105) is not needed. In the same spirit, the estimates regarding the second order scattering (and therefore for the full non-linear scattering) can be improved by techniques used in the section 4.2. We emphasize that in the case $n = m = 3$ the main results of the paper apply for general potentials and do not require the form assumed below.

Let us now define the two-scale model that we analyse in detail.

Definition A.1 (Two-scale model). *Consider a microlocally isotropic random field*

$$(107) \quad q(x) = \epsilon \sqrt{\mu(x)} Q \left(\frac{x}{\ell} \right),$$

where $\mu \in C_0^\infty(D)$ is the local strength function and Q is a stationary zero-mean microlocally isotropic field of order m such that

$$\sigma_p(\xi) = |\xi|^{-m} + a(\xi),$$

where $a \in \mathcal{S}^{-m-1}(\mathbb{R}^n)$, the correlation function $K_Q(z) = \mathbb{E}(Q(x)Q(x-z))$ satisfies $K_Q \in L^1(\mathbb{R}^n)$ and Q has a correlation length of constant order.

Notice that realizations of the random field Q are not compactly supported whereas the realizations of q are due to the compact support of μ . For orders $m > n - 1$, such that $m \neq n$, we know according to Proposition 2.8 and since Q stationary that

$$\mathbb{E}(Q(x)Q(y)) = K_Q(x - y) = c|x - y|^{m-n} + F_\alpha(x - y),$$

where $c \in \mathbb{R}$ and F_α is smooth (for case $m = n$, the leading term is logarithmic). An example of a Gaussian random process with such asymptotics is given by covariance function

$$K_Q(x - y) = \exp(-|x - y|^{m-n})$$

for $m > n$. For more discussion of random processes of type (107), consider Example 2.2 and examples given in [28, 38].

Remark A.2. In the case $m > n$ the correlation length can be defined by

$$(108) \quad L_Q(x) = \left| \frac{\int_{\mathbb{R}^n} \mathbb{E}(Q(x)Q(x+y))dy}{\mathbb{E}(Q(x)^2)} \right|^{1/n}.$$

for a random field Q such that $\mathbb{E}(Q(x)Q(x+\cdot)) \in L^1(\mathbb{R}^n)$.

Since we have

$$K_q(x, y) = \epsilon^2 \sqrt{\mu(x)\mu(y)} K_Q\left(\frac{x-y}{\ell}\right)$$

and $K_q(x, x) = \epsilon^2 \mu(x) K_Q(0)$, we find that

$$L_q(x) = \left| \frac{\int_{\mathbb{R}^n} \sqrt{\mu(x+y)} K_Q\left(\frac{y}{\ell}\right) dy}{\sqrt{\mu(x)} K_Q(0)} \right|^{\frac{1}{n}} = \ell \left| \frac{\sqrt{\mu(x)} \int_{\mathbb{R}^n} K_Q(y) dy + \mathcal{O}(\ell)}{\sqrt{\mu(x)} K_Q(0)} \right|^{\frac{1}{n}} = \ell L_Q(x) + \mathcal{O}(\ell^2).$$

Therefore, the correlation length of process Q_ℓ is of order ℓ .

Consider the covariance operator C_{Q_ℓ} of random field $Q(x/\ell)$. Its symbol c_{Q_ℓ} satisfies (in the sense of generalized functions)

$$\begin{aligned} c_{Q_\ell}(x, \xi) &= \int_{\mathbb{R}^n} K_Q\left(\frac{x}{\ell}, \frac{y}{\ell}\right) \exp(-i\xi \cdot (x - y)) dy \\ &= \ell^n \int_{\mathbb{R}^n} K_Q\left(\frac{x}{\ell}, y'\right) \exp\left(-i\ell\xi \cdot \left(\frac{x}{\ell} - y'\right)\right) dy' \\ &= \ell^n c_Q(\ell\xi), \end{aligned}$$

since c_Q is independent on the spatial variable. Consequently, due to the additional multiplication by $\epsilon\sqrt{\mu(x)}$ the symbol of C_q is of the form

$$(109) \quad c_q(x, \xi) = \epsilon^2 \ell^n (\mu(x) |\ell\xi|^{-m} + \tilde{a}(x, \ell\xi)),$$

where $\tilde{a} \in \mathcal{S}^{-m-1}(\mathbb{R}^n \times \mathbb{R}^n)$ and $\tilde{a}$ has a compact support with respect to x .

A.2. Inverse scattering with the two-scale model. Next we analyze the statistical properties of the measurement data with a finite frequency band, $M_K(\theta, \tau, \omega)$, given in (106). Let us start by considering the correlations for first order scattering from potential model described in Definition A.1. By utilizing the asymptotic of the symbol c_q in (109), we obtain

$$\begin{aligned} \mathbb{E}(u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}) &= \epsilon^2 \ell^n \int_{\mathbb{R}^n} (\mu(x) |2k\ell|^{-m} + \tilde{a}(x, 2k\ell\theta)) e^{-i2\tau\theta \cdot x} dx \\ (110) \qquad \qquad \qquad &= \epsilon^2 \ell^{n-m} (2k)^{-m} \left(\hat{\mu}(2\tau\theta) + L\mathcal{O}\left(\frac{1}{k\ell}\right) \right), \end{aligned}$$

where the second term appears due to the fact that $\tilde{a}$ is supported in $D \times \mathbb{R}^2$. This identity corresponds to the formula (58) and is valid in our scaling regime in equation (105).

Using similar technique as in Proposition 4.3 and utilizing non-stationary phase principle (more precisely, see equations (63) and (64)), we obtain inequality

$$(111) \qquad \left| \mathbb{E}(u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}) \right| \lesssim \epsilon^2 \ell^{n-m} L^n (1+k)^{-m} (1+\tau)^{-N},$$

where $L = \text{diam}(D)$ and $N > 0$ is arbitrary and the implicit constant depends on N . The bounding constant in non-stationary phase principle also depends on the domain (and the phase, which is here trivial), which we have explicitly included on the right-hand side. Similar deduction also yields the bound

$$(112) \qquad |\mathbb{E}(u_1^\infty(k, \theta, -\theta) u_1^\infty(k + \tau, \theta, -\theta))| \lesssim \epsilon^2 \ell^{n-m} L^n (1+k)^{-m} (1+k+\tau)^{-N}.$$

Now let us turn our attention to the convergence of the data (106), i.e., the mean integral over the frequency band and, for a brief moment, consider a real-valued random process $X(k)$ as in Theorem 4.1. It is well-known that the variance of the mean integral

$$Y(K) = \frac{1}{K} \int_K^{2K} (X(k) - \mathbb{E}X(k)) dk$$

decays at rate $\mathcal{O}(1/K)$, if the covariance of X decays at a polynomial rate

$$(113) \qquad \mathbb{E}(X(k) - \mathbb{E}X(k))(X(k + \tau) - \mathbb{E}X(k + \tau)) \lesssim (1 + \tau)^{-N}$$

for $N > 2$. Similar arguments carry over to complex valued random processes, when the modulus of $Y(K)$ and decay of the complex covariance are considered.

Notice that due to Corollary 4.5 the decay rates in equations (111) and (112) imply similar decay of covariance for the product

$$Z(k) = k^m u_1^\infty(k, \theta, -\theta) \overline{u_1^\infty(k + \tau, \theta, -\theta)}.$$

The first order backscattered term u_1^∞ is a linear mapping of the potential and we find that

$$(114) \qquad M_K^{(1)}(\theta, \tau, \omega) := \frac{1}{K} \int_K^{2K} Z(k) dk = c_{n,m} \epsilon^2 (\hat{\mu}(2\tau\theta) + N_1(\omega)),$$

where $c_{n,m} \in \mathbb{R}$, μ is the local strength of the random field q and the error term $N_1(\omega)$ can be factorized in random and deterministic parts

$$N_1(\omega) = E(\omega) + F.$$

The error term $E : \Omega \rightarrow \mathbb{C}$ represents the random deviation from expected value

$$E(\omega) = \frac{1}{\epsilon^2}(M_K - \mathbb{E}M_K)$$

and is a zero-mean random variable satisfying

$$\mathbb{E}|E|^2 \lesssim \frac{\ell^{n-m}L^n}{K}.$$

The deterministic part F represents approximation

$$(115) \quad |F| = \frac{1}{\epsilon^2} \left| \mathbb{E}M_K - M(\theta, \tau, \omega) \right| \lesssim \ell^{n-m}L^n \cdot \frac{\log K}{K\ell}$$

due to the Proposition 4.3. In the particular case of $n = m$, the error in the first order backscattered data is roughly of order (standard deviation of the random error and deterministic error)

$$(116) \quad \mathbb{E}|N_1|^2 \lesssim \mathbb{E}|E|^2 + F^2 \lesssim \frac{L^n}{K} + L^{2n} \left(\frac{\log K}{K\ell} \right)^2.$$

Summarizing, the inequality (116) gives an estimate the how well the first order scattering term $M_K^{(1)}(\theta, \tau, \omega)$ in (114) estimates the local strength μ of the random field q . Next we consider the effect of the higher order scattering.

A.3. Effects from the higher order scattering. Our analysis in the section 4 regarding the full non-linear backscattering assumes the specific case $n = m = 3$. Therefore, under general conditions and, in particular, general dimension the data $M_K(\theta, \tau, \omega)$ from full backscattered field does not necessarily converge when K increases. Below we show that in the scaling regime (105) the smallness of the potential reduces the effect of the higher-order scattering, which is rather straightforward to quantify given the techniques in the proof of Lemma 4.8. However, in addition, we wish to quantify the effect of correlation length, which requires more care.

Let us first consider the effect of correlation length in the model A.1 and state the following well-known fact: the dilation operator $D_\ell : f \mapsto f(\frac{\cdot}{\ell})$ is bounded in $L_{-s}^p(\mathbb{R}^n)$ and its norm scales according to

$$(117) \quad \|D_\ell f\|_{L_{-s}^p} \lesssim \ell^{\frac{n}{p}} \|f\|_{L_{-s}^p}$$

for $\ell \leq 1$ and $s \geq 0$. Same inequality holds as well for $\ell \geq 1$ and $s \leq 0$. The inequality can be shown to hold by first noting that the standard norm of $L_s^p(\mathbb{R}^n)$ and $\|\cdot\|_{L^p(\mathbb{R}^n)} + \|(-\Delta)^s \cdot\|_{L^p(\mathbb{R}^n)}$ are equivalent for $s \geq 0$. Thereafter, one has directly $(-\Delta)^s(D_\ell u) = \ell^{-2s}D_\ell(\Delta u)$ for any $s \in \mathbb{N}$ and by utilizing interpolation and duality argument, the result follows.

Let $\Psi \in C_0^\infty(\mathbb{R}^n)$ be a smooth function such that $\text{supp}(\Psi) \subset [-1, 1]^n$. Moreover, we assume that functions $\Psi_{\vec{j}} := \Psi(\cdot - \vec{j})$ define a partition of unity, i.e.

$$(118) \quad \sum_{\vec{j} \in \mathbb{Z}^n} \Psi_{\vec{j}}(x) = 1 \quad \text{for all } x \in \mathbb{R}^n.$$

Next, let $R(\ell) \in \mathbb{N}$ be the (smallest) number of translations in one dimension that are needed to cover the dilated domain $\frac{1}{\ell} \cdot D$, i.e. we assume

$$(119) \quad \sum_{|\vec{j}| \leq R(\ell)} \Psi_{\vec{j}}(x) = 1 \quad \text{for all } x \text{ such that } \ell x \in D.$$

Clearly, we have that $R(1) \lesssim \text{diam}(D)$ and $R(\ell) \lesssim \frac{\text{diam}(D)}{\ell}$.

In order to quantify the effect of ℓ in the Born series of the data, we need a scaling estimate of q in the norm of $L^p_{-s}(\mathbb{R}^n)$ as well as an operator norm of $f \mapsto qf$. Since the local strength $\mu \in C_0^\infty(D)$ is bounded, it follows by (117) that

$$(120) \quad \left\| \sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) \right\|_{L^p_{-s}} \lesssim \ell^{\frac{n}{p}} \left\| \sqrt{\mu}(\ell \cdot) Q \right\|_{L^p_{-s}} \lesssim \ell^{\frac{n}{p}} \sum_{|\vec{j}| \leq R(\ell)} \left\| \Psi_{\vec{j}} Q \right\|_{L^p_{-s}} \\ \lesssim \ell^{\frac{n}{p} - n} \frac{L^n}{R(\ell)^n} \sum_{|\vec{j}| \leq R(\ell)} \left\| \Psi_{\vec{j}} Q \right\|_{L^p_{-s}} = G_1(\omega) L \ell^{\left(\frac{1}{p} - 1\right)n}.$$

Above, the random bound

$$G_1(\omega) = \frac{1}{R(\ell)^n} \sum_{|\vec{j}| \leq R(\ell)} \left\| \Psi_{\vec{j}} Q \right\|_{L^p_{-s}}$$

is bounded almost surely, since the expectation $\mathbb{E}G_1$ is bounded. Furthermore, $\mathbb{E}G_1$ is independent of ℓ , since random variables $\left\| \Psi_{\vec{j}} Q \right\|_{L^p_{-s}}$ are identically distributed (although not independent).

Considering the multiplication operator $f \mapsto qf$, we reproduce the ideas of the section 3.2 by replacing V with $\sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right)$. By dilation, boundedness of μ and identity (118), we can directly record the following identity

$$(121) \quad \left| \left\langle \sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) f, g \right\rangle \right| = \ell^n \left| \left\langle \sqrt{\mu}(\ell \cdot) Q, f(\ell \cdot) g(\ell \cdot) \right\rangle \right| \lesssim \ell^n \left| \left\langle \sum_{|\vec{j}| \leq R(\ell)} \Psi_{\vec{j}} Q, f(\ell \cdot) g(\ell \cdot) \right\rangle \right|.$$

We write

$$W_{\vec{j}} = (I - \Delta)^{-s/2} \left(\Psi_{\vec{j}} Q \right) \quad \text{and} \quad W = \sum_{|\vec{j}| \leq R(\ell)} W_{\vec{j}},$$

and define $W^\sharp$ and $W^\flat$ according to the section 3.2. Following the idea in inequality (41) and applying the estimate (121), decomposition (119) and dilation scaling (117), we obtain

$$\left| \left\langle \sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) f, g \right\rangle \right| \\ \leq \ell^n \left(\|(I - \Delta)^{t/2} W^\sharp\|_{L^r} \|(I - \Delta)^{(s-t)/2} (f(\ell \cdot) g(\ell \cdot))\|_{L^{r'}} + \|W^\flat\|_{L^p} \|(I - \Delta)^{s/2} (f(\ell \cdot) g(\ell \cdot))\|_{L^{p'}} \right) \\ \lesssim \ell^n \left(\ell^{\frac{n}{r'}} \|(I - \Delta)^{t/2} W^\sharp\|_{L^r} \|(I - \Delta)^{(s-t)/2} (fg)\|_{L^{r'}} + \ell^{\frac{n}{p'}} \|W^\flat\|_{L^p} \|(I - \Delta)^{s/2} (fg)\|_{L^{p'}} \right),$$

where $s - \frac{n}{p} < t < s$, $r = \frac{n}{s-t}$ and p and r are Hölder conjugates of p' and r' , respectively. Since $\frac{n}{p'} < \frac{n}{r'}$, we have that for small $\ell < 1$ the term $\ell^{n/p'}$ dominates and

$$\begin{aligned} \left| \langle \sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) f, g \rangle \right| &\lesssim \ell^{n(2-\frac{1}{p})} \sum_{|\vec{j}| \leq R(\ell)} (C(k) \|W_{\vec{j}}\|_{L^p} + \|W_{\vec{j}}^b\|_{L^p}) \|f\|_{H_k^{s,-\delta}} \|g\|_{H_k^{s,-\delta}} \\ &\lesssim L^n \ell^{n(1-\frac{1}{p})} \frac{1}{R(\ell)^n} \sum_{|\vec{j}| \leq R(\ell)} (C(k) \|W_{\vec{j}}\|_{L^p} + \|W_{\vec{j}}^b\|_{L^p}) \|f\|_{H_k^{s,-\delta}} \|g\|_{H_k^{s,-\delta}} \\ &\simeq L^n \ell^{n(1-\frac{1}{p})} G_2(\omega, \ell) \|f\|_{H_k^{s,-\delta}} \|g\|_{H_k^{s,-\delta}}, \end{aligned}$$

where

$$G_2(\omega, \ell) = \frac{1}{R(\ell)^n} \sum_{|\vec{j}| \leq R(\ell)} (C(k) \|W_{\vec{j}}\|_{L^p} + \|W_{\vec{j}}^b\|_{L^p})$$

is almost surely finite, since $\mathbb{E}G_2 < \infty$ and $C(k)$ decays to zero as k increases. In consequence, we obtain

$$(122) \quad \left\| \sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) f \right\|_{H_k^{-s,\delta}} = L^n \ell^{n(1-\frac{1}{p})} G_2(\omega, \ell) \|f\|_{H_k^{s,-\delta}}.$$

We are finally ready to analyse the effect from higher order scattering in the measurement data. Let us write

$$u_R^\infty(k, \theta, -\theta) = \sum_{j=2}^{\infty} u_j^\infty(k, \theta, -\theta).$$

Similarly to the proof of Lemma 4.8 we can apply the norm bounds obtained in equations (120) and (122) and further in the section 3 to arrive at

$$\begin{aligned} \sup_{\theta \in \mathbb{S}^2} |u_R^\infty(k, \theta, -\theta)| &\leq \epsilon \left\| \sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) \right\|_{L_{-s}^p} \|\exp(ik\theta \cdot y)\chi\|_{H^s} \|\chi u_0\|_{H_k^{s,-\delta}} \\ &\quad \times \sum_{j>0} \epsilon^j \left\| \mathcal{R}_k^+ \circ \left(\sqrt{\mu} Q \left(\frac{\cdot}{\ell} \right) \right) \right\|_{H_k^{s,-\delta} \rightarrow H_k^{s,-\delta}}^j \\ &\lesssim G_1(\omega) L^n \epsilon \ell^{-(1-\frac{1}{p})n} k^{2s} \sum_{j>0} \left(G_2(\omega, \ell) L^n \epsilon \ell^{(1-\frac{1}{p})n} k^{-(1-2s)} \right)^j \\ &\leq G_3(\omega, \ell) L^{2n} \epsilon^2 \frac{k^{-1+4s}}{1 - G_2(\omega, \ell) L^n \epsilon \ell^{(1-\frac{1}{p})n} k^{-1+2s}} \\ (123) \quad &\leq G_4(\omega, \ell) L^{2n} \epsilon^2 k^{-1+4s}, \end{aligned}$$

where expectations of random coefficients G_j , $j = 1, \dots, 4$, are uniformly bounded by a constant independent of ℓ and $s > 0$ is arbitrarily small. Therefore, for k large enough so that $1 - G_2(\omega, \ell) L^n \epsilon \ell^{(1-\frac{1}{p})n} k^{-1+2s} > \frac{1}{2}$, we have an upper bound

$$\begin{aligned} \frac{1}{K} \int_K^{2K} k^m |u_R^\infty(k, \theta, -\theta)|^2 dk &\lesssim \frac{G_4(\omega, \ell) L^{4n} \epsilon^4}{K} \int_K^{2K} k^m k^{-2+\delta} dk \\ &\simeq G_4(\omega, \ell) L^{4n} \epsilon^4 K^{m-2+\delta} \end{aligned}$$

for some small $\delta = 4s > 0$.

By combining the two approximations and formulas (114) and (116), we see using the Cauchy–Schwarz inequality that the finite frequency band $M_K(\theta, \tau, \omega)$ given in (106) satisfies

$$(124) \quad M_K(\theta, \tau, \omega) = \epsilon^2(\widehat{\mu}(2\tau\theta) + N_1(\omega) + N_2(\omega))$$

for

$$N_2(\omega) = G_4(\omega, \ell) L^{4n} \epsilon^2 K^{m-2+\delta}.$$

In fact, since G_4 can be considered to be of constant order we have that the full error in $M_K(\theta, \tau, \omega)$ (when scaled by ϵ^2) in the case $n = m$ is of order

$$(125) \quad \mathbb{E}|N_1 + N_2|^2 \lesssim \frac{L^n}{K} + L^{2n} \left(\frac{\log K}{K\ell} \right)^2 + L^{4n} \epsilon^2 K^{n-2+\delta}.$$

where $\delta > 0$ is arbitrarily small. Thus we can summarise the main result of this appendix: Under the assumptions on scaling regimes (105), the finite frequency band measurement $M_K(\theta, \tau, \omega)$, given in (106), determines the Fourier transform of local strength $\mu(x)$ by the formula

$$(126) \quad \widehat{\mu}(2\tau\theta) = \epsilon^{-2} M_K(\theta, \tau, \omega) + \mathcal{E}, \quad \mathcal{E} = -\epsilon^{-2}(N_1(\omega) + N_2(\omega))$$

where the error $\mathcal{E}$ can be estimated using the formula (125). In the scaling regime (105) the error term $\mathcal{E}$ is negligible with high probability. In consequence, we have means to approximate the Fourier transform of μ at the given frequency 2τ . Note that as $\mu \in C_0^\infty(D)$, the Fourier transform of μ decays rapidly, and therefore the formula (126) and the error estimate (125) are useful for small values of τ , e.g. $|\tau| \ll K(L_0/L)^n$.

APPENDIX B. RANDOM VARIABLES WITH GAUSSIAN PROBABILITY LAWS

Let X be a random variable, we say that it has a Gaussian law with mean μ and variance σ^2 , if its law $\mathbb{P}_X$ satisfies:

$$\mathbb{P}_X((-\infty, a]) = \frac{1}{\sqrt{2\pi\sigma^2}} \int_{-\infty}^a e^{-\frac{(x-\mu)^2}{2\sigma^2}} dx, \quad \forall a \in \mathbb{R}.$$

Lemma B.1. *Let X be a Gaussian random variable with mean 0. Then, there exists a constant $c_k > 0$ such that*

$$(\mathbb{E}X^{2k})^{\frac{1}{2k}} = c_k (\mathbb{E}X^2)^{\frac{1}{2}}, \quad \forall k \in \mathbb{N} \setminus \{0\}.$$

Proof. By homogeneity, one can assume that the variance be $\mathbb{E}X^2 = 1$. Integrating by parts,

$$\mathbb{E}X^{2k} = \int_{\mathbb{R}} x^{2k} e^{-\frac{x^2}{2}} dx = - \int_{\mathbb{R}} x^{2k-1} \frac{d}{dx} \left(e^{-\frac{x^2}{2}} \right) dx = (2k-1) \int_{\mathbb{R}} x^{2(k-1)} e^{-\frac{x^2}{2}} dx.$$

Iterating the process we see that $\mathbb{E}X^{2k} = c'_k$ for some constant which only depends on k . This concludes the proof of this lemma. $\square$

Proposition B.2. *The linear combination of two independent Gaussian random variables has a Gaussian law.*

Given two random variables X_1 and X_2 , the pair $X = (X_1, X_2)$ is said to be a Gaussian random vector if $r_1 X_1 + r_2 X_2$ has a Gaussian law for any $r_1, r_2 \in \mathbb{R}$. If X is a Gaussian random vector, its law is determined by the vector $(\mathbb{E}X_1, \mathbb{E}X_2)$ and the matrix

$$\begin{pmatrix} \text{Cov}(X_1, X_1) & \text{Cov}(X_1, X_2) \\ \text{Cov}(X_2, X_1) & \text{Cov}(X_2, X_2) \end{pmatrix}.$$

Remark B.3. (i) If $X = (X_1, X_2)$ is a Gaussian random vector, then X_1 and X_2 are Gaussian random variables.

(ii) If X_1 and X_2 are independent Gaussian random variables, then $X = (X_1, X_2)$ is a Gaussian random vector.

(iii) The pair $X = (X_1, X_2)$ of two Gaussian random variables X_1 and X_2 is not in general a Gaussian vector. To see this, is enough to consider a Gaussian random variable X with mean 0 and variance 1 and the random variable X_c defined as $X_c(\omega) = X(\omega)$ if $X(\omega) > c$ and $X_c(\omega) = -X(\omega)$ if $X(\omega) \leq c$. One can check that X_c has the same law than X . However, $X - X_c$ does not have a Gaussian law since $\mathbb{P}_{X-X_c}(\{0\}) > 0$.

(iv) Two Gaussian random variables X_1 and X_2 may be uncorrelated (i.e. $\text{Cov}(X_1, X_2) = 0$) and not be independent. This can be verified using the same X and X_c of (iii), which are not independent. To do so, note that the function $c \mapsto \mathbb{E}(XX_c)$ is continuous and

$$\lim_{c \rightarrow \mp\infty} \mathbb{E}(XX_c) = \pm 1.$$

Thus, there exists a c_0 such that $\mathbb{E}(XX_{c_0}) = 0$ and the random variables X and X_{c_0} are uncorrelated.

(v) Let X_1, X_2 and X'_2 be three Gaussian random variables equally distributed and assume that $\mathbb{E}(X_1 X_2) = \mathbb{E}(X_1 X'_2)$. In general, it is not true that (X_1, X_2) and (X_1, X'_2) are equally distributed. We show that using the example in (iii). Let X' be a Gaussian random variable with mean 0 and variance 1 and assume it to be independent of X . Then, the vectors (X, X_{c_0}) and (X, X') do not have the same law, since (X, X') is Gaussian and (X, X_{c_0}) is not.

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