

THE OBSERVATIONAL LIMIT OF WAVE PACKETS WITH NOISY MEASUREMENTS

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ABSTRACT. The authors consider the problem of recovering an observable from certain measurements containing random errors. The observable is given by a pseudodifferential operator while the random errors are generated by a Gaussian white noise. The authors show how wave packets can be used to partially recover the observable from the measurements almost surely. Furthermore, they point out the limitation of wave packets to recover the remaining part of the observable, and show how the errors hide the signal coming from the observable. The recovery results are based on an ergodicity property of the errors produced by wave packets.

1. INTRODUCTION AND MAIN RESULTS

We study the problem of recovering an *observable* P from certain *measurements* \mathcal{N}_P that contain some random errors. In our case, the observable P is described by a pseudodifferential operator

$$(1) \quad Pf(x) = \frac{1}{(2\pi)^{d/2}} \int_{\mathbb{R}^d} e^{ix \cdot \xi} a(x, \xi) \widehat{f}(\xi) d\xi,$$

with a classical symbol a of order $m \in \mathbb{R}$. Here $d \geq 1$ and \widehat{f} denotes the Fourier transform of f defined as $\widehat{f}(\xi) = 1/(2\pi)^{d/2} \int_{\mathbb{R}^d} e^{-i\xi \cdot x} f(x) dx$. The measurements \mathcal{N}_P consists of the ideal data or *signal* yielded by P plus a centered complex Gaussian variable modeling the random error. More precisely, for arbitrary *states* f and g , the measurements are given by

$$\mathcal{N}_P(f, g) = \int_{\mathbb{R}^d} \bar{f} P g + \mathcal{E}(\bar{f}, g),$$

where $\mathcal{E}(\bar{f}, g)$ is a complex Gaussian with zero mean, and variance depending on the states f and g . For the purpose of this introduction, we think of the random error given as follows —a more complete description can be found in the

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section 3. Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space so that there exists a countable family $\{X_\alpha : \alpha \in \mathbb{N}^2\}$ of independent complex Gaussian variables satisfying

$$(2) \quad \mathbb{E}X_\alpha = 0, \quad \mathbb{E}(\overline{X_\alpha}X_\alpha) = 1.$$

Then, the error $\mathcal{E} = \mathcal{E}_\beta$, contained in the measurements $\mathcal{N}_P = \mathcal{N}_{\beta,P}$, can be expressed as

$$(3) \quad \mathcal{E}_\beta(f, g) = \sum_{\alpha \in \mathbb{N}^2} (f|e_{\alpha_1})_\beta (g|e_{\alpha_2})_\beta X_\alpha,$$

where $\alpha = (\alpha_1, \alpha_2)$, $\{e_n : n \in \mathbb{N}\}$ is an orthonormal basis of the Sobolev space $H^\beta(\mathbb{R}^d)$ with $\beta \in \mathbb{R}$ and $(\cdot|\cdot)_\beta$ denotes its inner product,

$$(4) \quad (f|g)_\beta = \int_{\mathbb{R}^d} (1 + |\xi|^2)^\beta \widehat{f}(\xi) \overline{\widehat{g}(\xi)} d\xi.$$

As usual, $\|f\|_\beta^2 = (f|f)_\beta$. For convenience, if $\beta = 0$, we simplify the notation so that the inner product is denoted by $(\cdot|\cdot)$ and the norm by $\|\cdot\|$. In the section 3 we shall show that

$$(5) \quad \mathbb{E} \mathcal{E}_\beta(f, g) = 0 \quad \text{and} \quad \mathbb{E} |\mathcal{E}_\beta(f, g)|^2 = \|f\|_\beta^2 \|g\|_\beta^2.$$

Note that different values of β might correspond to different variances of the error \mathcal{E}_β . This can be interpreted as follows. If $\beta = 0$ the oscillations of the states f and g do not affect the size of the error, only their masses influence its size; while for $\beta > 0$ the noise can be increased in the cases where the states f and g present many oscillations. When $\beta < 0$ the oscillations of the states might reduce the size of the noise.

For simplicity, we only consider observables P so that their symbols accept an expansion in terms of homogeneous functions. More precisely, we consider an observable P so that its symbol a of order m satisfies that

$$(6) \quad a \sim \sum_{j=1}^{\infty} a_j,$$

with a_j being a classical symbol of order $m_j \in \mathbb{R}$, for $m_j < m_{j-1} < \dots < m_1 = m$, which is *homogeneous in the variable ξ* , that is,

$$(7) \quad a_j(x, t\xi) = t^{m_j} a_j(x, \xi), \quad \text{whenever } |\xi| \geq 1/2 \text{ and } t \geq 1.$$

Note that such an expansion is not unique —one can choose any smooth extension of the $a_j(x, \xi)$ for $|\xi| \leq 1/2$. In the section 2, we provide more details about classical symbols and the exact meaning of the expansion (6). It can be convenient for the reader to note that differential operators of any order with smooth coefficients are observable satisfying these properties.

Note that recovering the symbol a of a pseudodifferential operator P is an easy task: if one extends the operator P to the tempered distributions in \mathbb{R}^d , one can verify that *plane waves* $e^{ix \cdot \xi}$ yield the symbol, $e^{-ix \cdot \xi} P(e^{ix \cdot \xi}) = a(x, \xi)$. In the

case of a differential operator $P = \sum_{|\alpha| \leq m} a_\alpha(x) (-i)^{|\alpha|} \partial_x^\alpha$, this is straightforward $e^{-ix \cdot \xi} P(e^{ix \cdot \xi}) = \sum_{|\alpha| \leq m} a_\alpha(x) \xi^\alpha$. This suggests that, to recover the observable in the presence of a random noise, one needs to approximate plane waves with functions in $H^\beta(\mathbb{R}^d)$. The problem is that this cannot be achieved with sequences that are uniformly bounded in $H^\beta(\mathbb{R}^d)$ since no Sobolev space in \mathbb{R}^d contains the plane waves. This seems a relevant point since, if the sequence is not bounded, by (5), it looks unlikely that we can filter the noise. Fortunately, there is a way around this at least when the symbol has an homogeneous expansion as in (6) and (7). Indeed, it is enough to introduce a family of states $\{f_t : t \geq 1\}$ that concentrate around a point in \mathbb{R}^d as t grows, and oscillate at a higher order than the rate of concentration, namely, wave packets.

At this point, we are ready to state the main results of this article. The first theorem consists of several reconstruction formulas that will be used recursively to obtain the expansion of the observable P up to some extent. Let us first state this first theorem and then we explain how to use these formulas to obtain the expansion.

Theorem 1. *Let P be an observable whose symbol a of order m accepts an homogeneous expansion as in (6), and Q_j denote the pseudodifferential operators with symbols a_j . Let $\beta \in \mathbb{R}$ be the value in the error \mathcal{E}_β and assume that $m > 2\beta - 1/2$. Set*

$$k_\beta = \min\{j \in \mathbb{N} : m_{j+1} \leq 2\beta - 1/2\} \text{ and } j_\beta = \min\{j \in \mathbb{N} \cup \{0\} : m_{j+1} \leq 2\beta\},$$

and note that $k_\beta \geq 1$ while $0 \leq j_\beta \leq k_\beta$. For every $x_0, \xi_0 \in \mathbb{R}^d$ with $|\xi_0| = 1$, $\lambda > 1$ and $t \geq 1$, set the wave packets

$$f_{t,\lambda}(x) = t^{d/2} \chi(t(x - x_0)) e^{it^\lambda(x - x_0) \cdot \xi_0},$$

where χ is a function in the Schwartz class satisfying (10) in the section 2. Assuming that we know β and the orders $\{m_1 \dots m_{k_\beta}\}$, we have:

- (a) Whenever $j_\beta \geq 1$, set $\lambda_j = \max(1/(m_j - m_{j+1}), 2)$ if $1 \leq j < j_\beta$, and $\lambda_{j_\beta} \geq \max(1/(m_{j_\beta} - 2\beta), 2)$. Then,

$$\lim_{N \rightarrow \infty} N^{-\lambda_j m_j} [\mathcal{N}_{\beta, P}(f_{N, \lambda_j}, f_{N, \lambda_j}) - \sum_{0 < k < j} (f_{N, \lambda_j} | Q_k f_{N, \lambda_j})] = a_j(x_0, \xi_0)$$

almost surely for every $j \in \{1, \dots, j_\beta\}$.

- (b) Set $\lambda_j = \max(1/(m_j - m_{j+1}), 2)$ if $j_\beta < j < k_\beta$, and $\lambda_{k_\beta} > \max(1/(m_{k_\beta} - 2\beta + 1/2), 2)$. Then,

$$\lim_{N \rightarrow \infty} \frac{1}{N} \int_N^{2N} t^{-\lambda_j m_j} [\mathcal{N}_{\beta, P}(f_{t, \lambda_j}, f_{t, \lambda_j}) - \sum_{0 < k < j} (f_{t, \lambda_j} | Q_k f_{t, \lambda_j})] dt = a_j(x_0, \xi_0)$$

almost surely whenever $0 \leq j_\beta < j \leq k_\beta$.

To explain how to use this formulas to obtain the expansion, assume for instance that the observable P is so that $m > m_{j_\beta} > 2\beta \geq m_{k_\beta} > 2\beta - 1/2$. Choosing λ_1 as described in (a), we compute a_1 from

$$\lim_{N \rightarrow \infty} N^{-\lambda_1 m_1} [\mathcal{N}_{\beta, P}(f_{N, \lambda_1}, f_{N, \lambda_1})] = a_1(x_0, \xi_0),$$

which is a limit that converges almost surely. Combining this with the homogeneity property (7) one essentially obtains $a_1(x, \xi)$ for all $x \in \mathbb{R}$ and $|\xi| \geq 1/2$, to define the values of $a_1(x, \xi)$ when $|\xi| \leq 1/2$, one can choose any smooth extension. Once we know a_1 , we construct Q_1 and compute a_2 from

$$\lim_{N \rightarrow \infty} N^{-\lambda_2 m_2} [\mathcal{N}_{\beta, P}(f_{N, \lambda_2}, f_{N, \lambda_2}) - (f_{N, \lambda_2} | Q_1 f_{N, \lambda_2})] = a_2(x_0, \xi_0),$$

whose convergence is almost surely. Iterating the process —still applying (a)— we obtain a_j with $j \leq j_\beta$. Note that the choices of $\lambda_1, \dots, \lambda_{j_\beta}$ are only based on the a priori knowledge of β and $\{m_1, \dots, m_{j_\beta}\}$. In order to compute $a_{j_\beta+1}$, we use the limit in (b):

$$\lim_{N \rightarrow \infty} \frac{1}{N} \int_N^{2N} t^{-\lambda_{j_\beta+1} m_{j_\beta+1}} [\mathcal{N}_{\beta, P}(f_{t, \lambda_{j_\beta+1}}, f_{t, \lambda_{j_\beta+1}}) - \sum_{k=1}^{j_\beta} (f_{t, \lambda_{j_\beta+1}} | Q_k f_{t, \lambda_{j_\beta+1}})] dt,$$

which again converges almost surely to $a_{j_\beta+1}(x_0, \xi_0)$. Iterating the process — applying now (b)— we obtain a_j for $j_\beta < j \leq k_\beta$. Again, the choices of λ_j with $j \in \{j_\beta + 1, \dots, k_\beta\}$ are only based on the a priori knowledge of β and the orders $\{m_{j_\beta+1}, \dots, m_{k_\beta}\}$. This completes the iteration to construct a_1, \dots, a_{k_β} in the situation where the observable P satisfies $m > m_{j_\beta} > 2\beta \geq m_{k_\beta} > 2\beta - 1/2$. Note that, in the case where $j_\beta = 0$, we would only use (b) to recover a_1, \dots, a_{k_β} in the expansion.

It is worth to pay attention to the case where the observable P is a differential operator of order m , in that case it is possible to recover the full P from $\mathcal{N}_{\beta, P}$ with $\beta < 1/4$. This means that, even with an error \mathcal{E}_β that gets amplified with the oscillations of the states, one can obtain the full observable. Furthermore, in the case that P is elliptic —the leading order term of its symbol satisfies that $\sum_{|\alpha|=m} a_\alpha(x) \xi^\alpha \neq 0$ for all $\xi \in \mathbb{R}^d \setminus \{0\}$, one can construct a suitable parametrix from the measurements $\mathcal{N}_{\beta, P}$ with $\beta < 1/4$.

Theorem 1 shows that wave packets are appropriate states to obtain part of the expansion of the observable P . Furthermore, as we will show in the [section 2](#), they are suitable to recover the full expansion in the case of ideal data, that is, in the absence of \mathcal{E}_β . However, in the presence of noise, we cannot recover a_j without uncertainty for any $j > k_\beta$ since the variance of the errors becomes so large that the signal produced by the observable is lost in the noise. This claim is actually the statement of the second theorem of this paper. For convenience, introduce the

following notation

$$(8) \quad \mathcal{N}_{\beta, P_j}(f, g) = \mathcal{N}_{\beta, P}(f, g) - \sum_{0 < k < j} (f | Q_k g) \text{ for } j \in \mathbb{N},$$

where Q_j is as in [Theorem 1](#).

Theorem 2. *Let $\beta \in \mathbb{R}$ be given.*

(a) *When $m_j \leq 2\beta$ there is no $\lambda \in (1, \infty)$ such that the sequence*

$$\{N^{-\lambda m_j} \mathcal{N}_{\beta, P_j}(f_{N, \lambda}, f_{N, \lambda}) : N \in \mathbb{N}\}$$

converges to $a_j(x_0, \xi_0)$ in probability. In fact, if $m_j < 2\beta$, then for every $\lambda \in (1, \infty)$ and every $c > 0$,

$$\lim_{t \rightarrow \infty} \mathbb{P} \left\{ \left| t^{-\lambda m_j} \mathcal{N}_{\beta, P_j}(f_{t, \lambda}, f_{t, \lambda}) - a_j(x_0, \xi_0) \right| > c \right\} = 1.$$

(b) *When $m_j \leq 2\beta - 1/2$ there is no $\lambda \in (1, \infty)$ such that the sequence*

$$\left\{ \frac{1}{N} \int_N^{2N} t^{-\lambda m_j} \mathcal{N}_{\beta, P_j}(f_{t, \lambda}, f_{t, \lambda}) dt : N \in \mathbb{N} \right\}$$

converges to $a_j(x_0, \xi_0)$ in probability. Actually, for every $\lambda \in (1, \infty)$ and every $c > 0$,

$$\lim_{T \rightarrow \infty} \mathbb{P} \left\{ \left| \frac{1}{T} \int_T^{2T} t^{-\lambda m_j} \mathcal{N}_{\beta, P_j}(f_{t, \lambda}, f_{t, \lambda}) dt - a_j(x_0, \xi_0) \right| > c \right\} = 1.$$

Finally, we state the third main result of our paper, which provides the rate of convergence in probability of the limits in [Theorem 1](#).

Theorem 3. *Under the same assumptions of [Theorem 1](#) and adopting the notation in (8), we have:*

(a) *Whenever $j_\beta \geq 1$, set $\lambda_j = \max(1/(m_j - m_{j+1}), 2)$ if $1 \leq j < j_\beta$, and $\lambda_{j_\beta} > \max(1/(m_{j_\beta} - 2\beta), 2)$. Then, for all $\varepsilon > 0$ and $\delta \in (0, 1]$, there exists N_0 such that for all $N \geq N_0$*

$$\mathbb{P} \left\{ \left| N^{-\lambda m_j} \mathcal{N}_{\beta, P_j}(f_{N, \lambda}, f_{N, \lambda}) - a_j(x_0, \xi_0) \right| \leq \varepsilon \right\} \geq 1 - \delta.$$

Here $N_0 = C \max(1/\varepsilon, (\log 1/\delta)^{1/\theta})$ with $\theta > 0$ and $C \geq 1$ explicit and only depending on $d, \{m_1, \dots, m_{k_\beta}\}$ and β .

(b) *Set $\lambda_j = \max(1/(m_j - m_{j+1}), 2)$ if $j_\beta < j < k_\beta$, and $\lambda_{k_\beta} > \max(1/(m_{k_\beta} - 2\beta + 1/2), 2)$. Then, for all $\varepsilon > 0$ and $\delta \in (0, 1]$, there exists N_0 such that for all $N \geq N_0$*

$$\mathbb{P} \left\{ \left| \frac{1}{N} \int_N^{2N} t^{-\lambda m_j} \mathcal{N}_{\beta, P_j}(f_{t, \lambda}, f_{t, \lambda}) dt - a_j(x_0, \xi_0) \right| \leq \varepsilon \right\} \geq 1 - \delta.$$

Here $N_0 = C \max(1/\varepsilon, (\log 1/\delta)^{1/\theta})$ with $\theta > 0$ and $C \geq 1$ explicit and only depending on $d, \{m_1, \dots, m_{k_\beta}\}$ and β .

The proof of the part (b) of [Theorem 1](#) is based on an ergodicity property for $t^{-\lambda m} \mathcal{E}_\beta(\overline{f_{t,\lambda}}, f_{t,\lambda})$ (see the [lemma 3.2](#)). To show the ergodicity, we have to compute the covariance of

$$\frac{1}{T} \int_T^{2T} t^{-\lambda m} \mathcal{E}_\beta(\overline{f_t}, f_t) dt,$$

which is reduced to analyse an oscillatory integral and isolate appropriately the stationary points.

The problem we study in this article and our results are of interest by themselves. Think for example of the fact that one can construct a parametrix for an elliptic differential operator in the presence of noise. The problem also has connections with quantum mechanics, where the observables are self-adjoint operators whose symbols are classical observables (functions in the phase space). Furthermore, if a unit vector ψ in a Hilbert space describes the state of a quantum system, the measurement $\langle \psi, P\psi \rangle$ corresponds to the expected value of the classical observable a —here P is the quantum observable and a its symbol. As the reader can guess, we have borrowed the terms observables, measurements, states, plane waves and wave packets from quantum mechanics.

Despite these interesting connections, our motivation originates in the Calderón problem. This inverse problem consists in recovering the conductivity of a medium from suitable electrical superficial data (voltage and current on the surface of the medium). In technical terms, the conductivity is described by a positive function γ defined in a bounded domain $D \subset \mathbb{R}^d$ with $d \geq 2$. The boundary data is modeled by the so-called Dirichlet-to-Neumann map Λ_γ , which is an operator that maps the voltage f on ∂D , the boundary of D , to the current through the boundary $\Lambda_\gamma f = \gamma \partial_\nu u|_{\partial D}$. Here u is the electric potential generated by f , that is, the solution of the Dirichlet problem

$$\begin{cases} \nabla \cdot (\gamma \nabla u) = 0 \text{ in } D, \\ u|_{\partial D} = f, \end{cases}$$

and ∂_ν denotes the partial derivative in the direction given by the exterior unit normal vector on ∂D . Thus, the Calderón problem consists in deciding if γ is uniquely determined by Λ_γ , and if it is so, reconstructing the conductivity from the boundary data. This inverse problem has received a lot of attention since Calderón formulated it in [\[5\]](#). Some important contributions on this problem are [\[14\]](#), [\[13\]](#), [\[12\]](#), [\[2\]](#), [\[9\]](#) and [\[8\]](#).

Since the Dirichlet-to-Neumann map contains physical data about the voltage and the current on the boundary, one assumes that this information has been obtained after some measurement procedure. In applications, these measurements will contain errors. In [\[6\]](#), the authors assumed that the available data had been corrupted in the process of measurements so that they only have access to

$$\mathcal{N}(f, g) = \int_{\partial D} \Lambda_\gamma f g + \mathcal{E}_0(f, g),$$

where the Hilbert space involved in the definition of the error \mathcal{E}_0 is $L^2(\partial D)$ —which explains the subindex 0. The authors of [6] showed that from the measurements \mathcal{N} one can reconstruct γ and $\partial_\nu \gamma$ on ∂D almost surely. Their approach is based on a method proposed by Brown [3], and later developed in collaboration with Salo [4], to reconstruct the unknown on the boundary assuming full knowledge of ideal data. This method uses a family of solutions whose boundary data remind to wave packets. The formula, given in [6], to reconstruct the conductivity on the boundary from \mathcal{N} is a limit in the spirit of part (a) in Theorem 1. However, to reconstruct its normal derivative on the boundary, the formula requires an average in the parameter of the family, as the formula in (b) of Theorem 1.

Let us clarify in detail the connection of Theorem 1 with the results in [6]. First, we notice that, if D and γ are smooth, the map Λ_γ can be locally identified with a pseudodifferential operator of order $m = 1$ whose symbol admits an expansion as (6) with $m_j = 2 - j$. In this case each a_j contains information of $\partial_\nu^{j-1} \gamma$ on the boundary (this was proved by Sylvester and Uhlmann [15]). Thus, for $\beta = 0$, $m_1 = 1$ is in the assumptions of Theorem 1 (a), so no average is needed to reconstruct a_1 , which contains information of γ on the boundary of D —as shown in [6]. The second order $m_2 = 0$ satisfies the conditions of Theorem 1 (b) for $\beta = 0$, so the formula to reconstruct a_2 (containing information on $\partial_\nu \gamma$ on ∂D) requires an average —as in [6]. On the other hand, Theorem 2 shows that the method adopted in [6] does not allow to recover $\partial_\nu^k \gamma$ on the boundary for $k > 1$ with $\beta = 0$. This point could suggest that is not possible to reconstruct γ in D from \mathcal{N} almost surely. However, making the level of the noise \mathcal{N} small, Abraham and Nickl [1] have been able to provide an algorithm to reconstruct γ in D .

Finally, our results can also be connected to those in [7], where the authors study the question of data corruption in electromagnetism and elasticity. In electrodynamics, the boundary data is given by either the admittance map or the impedance map. These maps can be identified with pseudodifferential operators of order $m = 0$ whose symbols can be expanded as in (6), containing information of the unknown on the boundary (see [11] and [10]). Thus, according to our results, for $\beta = 0$ we would require making an average to obtain the leading order term of the expansion, while for $\beta = -1$ we would not. This actually coincides with the formulas given by the authors in [7]. In this regard, the results of our paper establish a framework that helps to obtain a better understanding of the results in [6] and [7].

The contents of the next sections are the following. In the section 2, we recall the basic definition of classical symbols and show how to use waves packets to recover the full expansion of the observable. In the section 3, we define rigorously the random error and prove the main results of this paper, meaning, Theorem 1, Theorem 2 and Theorem 3.

Throughout the paper we write $a \lesssim b$ or equivalently $b \gtrsim a$, when a and b are positive constants and there exists $C > 0$ so that $a \leq Cb$. We refer to C as the

implicit constant and will only depends on d , $\{m_j : j \in \mathbb{N}\}$ and β . Additionally, if $a \lesssim b$ and $b \lesssim a$, we write $a \simeq b$.

2. RECOVERING AN OBSERVABLE FROM IDEAL DATA

In this section we provide an algorithm to recover the expansion $\sum_{j=1}^{\infty} a_j$ of the symbol of an observable P when the available data has not been corrupted by measurement errors. The advantage of our method is that it can be used to filter the noise up to some extent, as described in [Theorem 1](#).

Let x_0 and ξ_0 belong to \mathbb{R}^d with $|\xi_0| = 1$ and take a smooth function χ , such that its Fourier transform $\widehat{\chi}$ has compact support in B_1 , the open ball of radius 1. Consider the *wave packets*

$$(9) \quad f_t(x) = t^{d/2} \chi(t(x - x_0)) e^{it^\lambda(x-x_0) \cdot \xi_0},$$

where $t \geq 1$ and $1 < \lambda < \infty$. To reduce the notation in this section, we write f_t instead of $f_{t,\lambda}$. In particular, we choose $\widehat{\chi}$ as a smooth cutoff satisfying

$$(10) \quad b \mathbf{1}_{B_{1/2}}(\xi) \leq \widehat{\chi}(\xi) \leq b \mathbf{1}_{B_1}(\xi) \quad \forall \xi \in \mathbb{R}^d \quad \text{and} \quad \|\chi\| = \|\widehat{\chi}\| = 1,$$

where $b > 0$, and $\mathbf{1}_{B_{1/2}}$ and $\mathbf{1}_{B_1}$ denote the characteristic functions of the balls of radius 1/2 and 1, respectively. One can check that the wave packet f_t satisfies

$$(11) \quad \|f_t\|_\beta \simeq t^{\lambda\beta}.$$

Heuristically the wave packets f_t will look increasingly like a plane wave close to the point x_0 . This is the reason to introduce the parameter $\lambda > 1$, since it guarantees that the frequency of the wave increases more quickly than the support of $\chi(t(\cdot - x_0))$ shrinks around x_0 .

Before providing the algorithm to recover the expansion of the observable P , let us recall a few definitions. We say that a smooth function $a = a(x, \xi)$ in the *phase space* $\mathbb{R}^d \times \mathbb{R}^d$ is a *classical symbol* if it satisfies that

$$|\partial_x^\alpha \partial_\xi^\beta a(x, \xi)| \leq C_{\alpha,\beta} \langle \xi \rangle^{m-|\beta|},$$

where α and β are multi-indices and $\langle \xi \rangle = (1 + |\xi|^2)^{1/2}$. Additionally, if $\{m_j : j \in \mathbb{N}\}$ is a strictly decreasing sequence of real numbers so that $\lim_{j \rightarrow \infty} m_j = -\infty$, and a_j is a classical symbol of order m_j , we say that $\sum_{j=1}^{\infty} a_j$ expands the symbol a of order $m = m_1$, and we write $a \sim \sum_{j=1}^{\infty} a_j$, if

$$a - \sum_{j=1}^k a_j \text{ with } k \in \mathbb{N}$$

is a classical symbol of order m_{k+1} . Notice that, as remarked in the introduction, the expansion of a is non unique, by the definition, any a_j can be modified by adding to it a lower order symbol. On the other hand, if the homogeneity condition (7) is added, then $a_j(x, \xi)$ is uniquely determined for $|\xi| \geq 1/2$ —there are no lower

order symbols with the same homogeneity of a_j — but it still can be modified freely for $|\xi| \leq 1/2$.

We are now ready to describe the algorithm to recover the homogeneous expansion of a given observable. However, for simplicity, we start by assuming that the observable P has a symbol of the form $a + b$ with a and b being, respectively, classical symbols of order m and n such that, a is homogeneous in the variable ξ , and $n < m$. The objective is to show that the homogeneous part of the symbol can be recovered from knowledge of $(f_t|Pf_t)$, where f_t is the family of wave packets given in (9). We want to prove that

$$(12) \quad \lim_{t \rightarrow \infty} t^{-\lambda m} (f_t|Pf_t) = a(x_0, \xi_0),$$

with $1 < \lambda < \infty$.

We begin with a simple lemma that will be used to estimate the part of the observable P corresponding to b .

Lemma 2.1. *Let R be a pseudodifferential operator of order $n \in \mathbb{R}$. Then*

$$|(f_t|Rf_t)| \lesssim t^{\lambda n}.$$

Proof. Let b denote the symbol of R . Note that we have

$$(f_t|R(f_t)) = \frac{1}{(2\pi)^{d/2}} \int_{\mathbb{R}^d} \overline{f_t(x)} \left(\int_{\mathbb{R}^d} e^{ix \cdot \xi} b(x, \xi) \widehat{f_t}(\xi) d\xi \right) dx.$$

By the basic properties of the Fourier transform and the form of (9) we have that

$$(13) \quad \widehat{f_t}(\xi) = t^{-d/2} e^{-i\xi \cdot x_0} \widehat{\chi}((\xi - t^\lambda \xi_0)/t).$$

After the changes of variables

$$(14) \quad \eta = (\xi - t^\lambda \xi_0)/t \quad \text{and} \quad y = t(x - x_0),$$

we can ensure that

$$(15) \quad |(f_t|R(f_t))| \lesssim \int_{\mathbb{R}^d} |\chi(y)| \int_{\mathbb{R}^d} |b(x_0 + t^{-1}y, t^\lambda \xi_0 + t\eta)| |\widehat{\chi}(\eta)| d\eta dy.$$

Since b a classical symbol of order n , we know that $|b(x, \xi)| \lesssim \langle \xi \rangle^n$ for all $x, \xi \in \mathbb{R}^d$, in particular

$$|b(x_0 + t^{-1}y, t^\lambda \xi_0 + t\eta)| \lesssim \langle t^\lambda \xi_0 + t\eta \rangle^n.$$

Note now that whenever $\eta \in \text{supp } \widehat{\chi} \subset B_1$, we have that $\langle t^\lambda \xi_0 + t\eta \rangle^n \simeq t^{\lambda n}$. This equivalence together with (15) finishes the proof of this lemma. \square

We now prove that the limit (12) holds.

Proposition 2.2. *Let P be a pseudodifferential operator with a symbol of the form $a + b$ with a and b classical symbols of order m and $n < m$, respectively. Additionally, assume that a is homogeneous in the variable ξ . Consider the wave packets $\{f_t : t \geq 1\}$ defined in (9). Then,*

$$|t^{-\lambda m} (f_t|Pf_t) - a(x_0, \xi_0)| = O(t^{-1}) + O(t^{\lambda(n-m)}) + O(t^{1-\lambda}),$$

and consequently

$$\lim_{t \rightarrow \infty} t^{-\lambda m} (f_t | P f_t) = a(x_0, \xi_0).$$

Proof. Let Q and R denote the pseudodifferential operators with symbols a and b , respectively. Thus, $P = Q + R$. By the [lemma 2.1](#) we have that

$$t^{-\lambda m} (f_t | R f_t) = O(t^{\lambda(n-m)}).$$

Hence, it is enough to show that

$$(16) \quad t^{-\lambda m} (f_t | Q f_t) = t^{-\lambda m} (f_t | a(\cdot, t^\lambda \xi_0) f_t) + O(t^{1-\lambda}),$$

since, using that $a(\cdot, t^\lambda \xi_0) = t^{\lambda m} a(\cdot, \xi_0)$, one has

$$\begin{aligned} |t^{-\lambda m} (f_t | a(\cdot, t^\lambda \xi_0) f_t) - a(x_0, \xi_0)| &= \left| \int_{\mathbb{R}^d} \overline{f_t(x)} a(x, \xi_0) f_t(x) dx - a(x_0, \xi_0) \right| \\ &\leq \int_{\mathbb{R}^d} t^d |\chi(t(x - x_0))|^2 |a(x, \xi_0) - a(x_0, \xi_0)| dx = O(t^{-1}), \end{aligned}$$

using that $a(\cdot, \xi_0)$ is a Lipschitz function with uniform bound for $|\xi_0| = 1$. Thus, we just need to prove (16). Note that using the Leibniz rule this is immediate if Q is a differential operator. Since Q has a as symbol, we can write

$$\begin{aligned} (f_t | Q f_t) - (f_t | a(\cdot, t^\lambda \xi_0) f_t) \\ = \frac{1}{(2\pi)^{d/2}} \int_{\mathbb{R}^d} \overline{f_t(x)} \left(\int_{\mathbb{R}^d} e^{ix \cdot \xi} (a(x, \xi) - a(x, t^\lambda \xi_0)) \widehat{f_t}(\xi) d\xi \right) dx \end{aligned}$$

using the inversion formula of the Fourier transform. We now use the form of $\widehat{f_t}$ given in (13) and the changes of variables (14). Thus,

$$(17) \quad \begin{aligned} &|(f_t | Q f_t) - (f_t | a(\cdot, t^\lambda \xi_0) f_t)| \\ &\leq \int_{\mathbb{R}^d} |\chi(y)| \int_{\mathbb{R}^d} |a(x_0 + t^{-1}y, t^\lambda \xi_0 + t\eta) - a(x_0 + t^{-1}y, t^\lambda \xi_0)| |\widehat{\chi}(\eta)| d\eta dy. \end{aligned}$$

Now, since $\widehat{\chi}$ is supported in the ball of radius 1, we can use that for $|\eta| \leq 1$, one has

$$\begin{aligned} t^{-\lambda m} |a(x_0 + t^{-1}y, t^\lambda \xi_0 + t\eta) - a(x_0 + t^{-1}y, t^\lambda \xi_0)| \\ = |a(x_0 + t^{-1}y, \xi_0 + t^{1-\lambda}\eta) - a(x_0 + t^{-1}y, \xi_0)| \\ \leq t^{1-\lambda} \sup_{\zeta \in B_1} |\nabla_\xi a(x_0 + t^{-1}y, \xi_0 + \zeta)| = O(t^{1-\lambda}). \end{aligned}$$

Note that we have used the homogeneity property to get the second line. Multiplying (17) by $t^{-\lambda m}$ and using the previous inequality we obtain the desired result. \square

Finally, we use the [proposition 2.2](#) to derive the algorithm to recover the full expansion of the observable P from the ideal data. Recall that its symbol satisfies

$$a \sim \sum_{j=1}^{\infty} a_j$$

with a_j of order m_j and homogeneous in the variable ξ . For convenience, set $P_1 = P$ and note that its symbol can be written as $a = a_1 + b_1$ with a_1 homogeneous of order $m_1 = m$ and $b_1 = a - a_1$ of order $m_2 < m_1$. By the [proposition 2.2](#),

$$(18) \quad |t^{-\lambda m_1}(f_t|P_1 f_t) - a_1(x_0, \xi_0)| = O(t^{-1}) + O(t^{\lambda(m_2 - m_1)}) + O(t^{1-\lambda}),$$

and consequently

$$(19) \quad \lim_{t \rightarrow \infty} t^{-\lambda m_1}(f_t|P f_t) = a_1(x_0, \xi_0).$$

Next, let Q_1 denote the pseudodifferential operator with symbol a_1 , and set $P_2 = P - Q_1$. The symbol of P_2 can be written as $a_2 + b_2$ with a_2 homogeneous of order m_2 and $b_2 = a - a_1 - a_2$ of order $m_3 < m_2$. Again, by the [proposition 2.2](#) we have that

$$|t^{-\lambda m_2}(f_t|P_2 f_t) - a_2(x_0, \xi_0)| = O(t^{-1}) + O(t^{\lambda(m_3 - m_2)}) + O(t^{1-\lambda}),$$

and consequently

$$\lim_{t \rightarrow \infty} t^{-\lambda m_2} [(f_t|P f_t) - (f_t|Q_1 f_t)] = a_2(x_0, \xi_0).$$

Now, assume that we have already recovered a_1, \dots, a_{j-1} and let Q_1, \dots, Q_{j-1} denote the pseudodifferential operators with symbols a_1, \dots, a_{j-1} , respectively. Then, set

$$(20) \quad P_j = P - \sum_{k=1}^{j-1} Q_k, \text{ for } j \in \mathbb{N} \text{ with } j > 1.$$

Note that the symbol of P_j can be written as $a_j + b_j$ with a_j homogeneous of order m_j and $b_j = a - \sum_{k=1}^j a_k$ of order $m_{j+1} < m_j$. By the [proposition 2.2](#) we have that

$$(21) \quad |t^{-\lambda m_j}(f_t|P_j f_t) - a_j(x_0, \xi_0)| = O(t^{-1}) + O(t^{\lambda(m_{j+1} - m_j)}) + O(t^{1-\lambda}),$$

and consequently

$$(22) \quad \lim_{t \rightarrow \infty} t^{-\lambda m_j} [(f_t|P f_t) - \sum_{k=1}^{j-1} (f_t|Q_k f_t)] = a_j(x_0, \xi_0).$$

Note that this algorithm only requires knowledge of P and the different orders $\{m_j : j \in \mathbb{N}\}$ of the homogeneous expansion of its symbol.

3. FILTERING THE NOISE IN THE MEASUREMENTS

In this section we show how to filter out the random noise to recover the observable from the measurements with full certainty. Unfortunately, —as described in [Theorem 2](#)— with our method it is impossible to recover the full expansion of the symbol of P since, for very low order terms of the expansion, the variance of the error becomes so large that the signal produced by the observable is lost in the noise. We also show in this section how this phenomenon works.

Before addressing these questions, we recall some properties about white noise, and define the random errors contained in our measurements. Let $(\Omega, \mathcal{F}, \mathbb{P})$ be a probability space and \mathcal{H} be a Hilbert space. A linear map $\mathbb{W} : \mathcal{H} \rightarrow L^2(\Omega, \mathcal{F}, \mathbb{P})$ is said to be a *complex Gaussian white noise* if, for every $f \in \mathcal{H}$, $\mathbb{W}f$ is a centered complex Gaussian variable and

$$(23) \quad \mathbb{E}(\overline{\mathbb{W}f}\mathbb{W}g) = (f|g)_{\mathcal{H}}, \text{ for all } f, g \in \mathcal{H}.$$

Here $(\cdot|\cdot)_{\mathcal{H}}$ denotes the inner product of \mathcal{H} . If \mathcal{H} is a separable Hilbert space and $(\Omega, \mathcal{F}, \mathbb{P})$ is a probability space so that there exists a sequence $\{X_n : n \in \mathbb{N}\}$ of independent complex Gaussian so that

$$\mathbb{E}X_n = 0, \quad \mathbb{E}(\overline{X_n}X_n) = 1,$$

then, there always exists a complex Gaussian white noise for $(\Omega, \mathcal{F}, \mathbb{P})$ and \mathcal{H} . It is enough to consider an orthonormal basis $\{e_n : n \in \mathbb{N}\}$ of \mathcal{H} , and define $\mathbb{W}f = \sum_{n \in \mathbb{N}} (e_n|f)_{\mathcal{H}} X_n$, where the convergence takes place in $L^2(\Omega, \mathcal{F}, \mathbb{P})$. On the other hand, if \mathbb{W} is a complex Gaussian white noise with \mathcal{H} separable, we set $X_n = \mathbb{W}e_n$, and then $\mathbb{W}f = \sum_{n \in \mathbb{N}} (e_n|f)_{\mathcal{H}} X_n$. By definition X_n is a centered complex Gaussian variable and since \mathbb{W} is an isometry, $\mathbb{E}(\overline{X_n}X_n) = 1$. Additionally, X_1, \dots, X_n, \dots are independent, since they are uncorrelated and every finite linear combination of any of them is a centered Gaussian variable.

With these properties at hand, we define now the *error* \mathcal{E}_β from a complex Gaussian white noise in the Hilbert space $\mathcal{H} = H^\beta(\mathbb{R}^d) \otimes H^\beta(\mathbb{R}^d)$ as follows:

$$(24) \quad \mathcal{E}_\beta(f, g) = \mathbb{W}(f \otimes g).$$

Recall that $(f_1 \otimes g_1 | f_2 \otimes g_2)_{\mathcal{H}} = (f_1 | f_2)_\beta (g_1 | g_2)_\beta$. This explains why the error \mathcal{E}_β takes the form in [\(3\)](#), with $X_\alpha = \mathbb{W}(e_{\alpha_1} \otimes e_{\alpha_2})$ and $\{e_n : n \in \mathbb{N}\}$ an orthonormal basis of $H^\beta(\mathbb{R}^d)$. In [\[6\]](#) and [\[7\]](#), the authors chose the form of [\(3\)](#) to define the error. Here we have chosen a more intrinsic way of defining it. After these comments, it is straightforward to check [\(5\)](#): $\mathbb{E}\mathcal{E}_\beta(f, g) = \mathbb{E}\mathbb{W}(f \otimes g) = 0$, and

$$\mathbb{E}|\mathcal{E}_\beta(f, g)|^2 = \mathbb{E}|\mathbb{W}(f \otimes g)|^2 = (f \otimes g | f \otimes g)_{\mathcal{H}} = (f|f)_\beta (g|g)_\beta = \|f\|_\beta^2 \|g\|_\beta^2.$$

Note that filtering out the noise, consists essentially in proving that the error $\mathcal{E}_\beta(\overline{f_t}, f_t)$ goes to zero almost surely when t tends to infinity. This property can be seen from its variance $\mathbb{E}|\mathcal{E}_\beta(\overline{f_t}, f_t)|^2$, and this is the reason to state here the following lemma.

Lemma 3.1. *Let (X, Σ, μ) be a measure space, and let $\{f_n : n \in \mathbb{N}\}$ be a sequence in $L^p(X, \Sigma, \mu)$ with $1 \leq p < \infty$ converging to f in $L^p(X, \Sigma, \mu)$. Assume there exists a sequence $\{b_n : n \in \mathbb{N}\}$ of positive real numbers whose limit vanishes, and*

$$(25) \quad \sum_{n=1}^{\infty} \frac{1}{b_n^p} \int_X |f_n - f|^p d\mu < \infty.$$

Then, $f_n(x)$ goes to $f(x)$ as n tends to ∞ for almost every $x \in X$.

This fact is well known, but for completeness, we include a proof.

Proof. Set $E_n = \{x \in X : |f_n(x) - f(x)| > b_n\}$. Using Chebyshev's inequality we see that

$$\mu \left(\bigcap_{n=1}^{\infty} \bigcup_{k=n}^{\infty} E_k \right) \leq \sum_{k=n}^{\infty} \mu(E_k) \leq \sum_{k=n}^{\infty} \frac{1}{b_n^p} \int_X |f_n - f|^p d\mu,$$

which, after (25), yields that

$$\mu \left(\bigcap_{n=1}^{\infty} \bigcup_{k=n}^{\infty} E_k \right) = 0.$$

On the other hand, if $x \notin \bigcap_{n=1}^{\infty} \bigcup_{k=n}^{\infty} E_k$, it means that there exists an n_x such that $|f_n(x) - f(x)| \leq b_n$ for all $n \geq n_x$. Thus, $f_n(x)$ tends to $f(x)$ for almost every $x \in X$. \square

Now, we are in the position to prove the theorems stated in the introduction of this paper. We start by showing when it is possible to use wave packets to filter the noise.

Proof of Theorem 1. Start by noting that, whenever $1 < \lambda < \infty$, $\beta \in \mathbb{R}$ and $j \in \mathbb{N}$, we have by (5) and (11) that

$$(26) \quad \mathbb{E} |t^{-\lambda m_j} \mathcal{E}_\beta(\bar{f}_t, f_t)|^2 \simeq t^{-2\lambda(m_j - 2\beta)},$$

where the implicit constant only depends on d , λ , m_j and β . Then, we have, by (18) and (26), that

$$\mathbb{E} |t^{-\lambda_1 m_1} \mathcal{N}_{\beta, P}(f_{t, \lambda_1}, f_{t, \lambda_1}) - a_1(x_0, \xi_0)|^2 \lesssim t^{-2}.$$

While for $j \in \{2, \dots, j_\beta\}$, we have, by (21) and (26), that

$$\mathbb{E} \left| t^{-\lambda_j m_j} \left[\mathcal{N}_{\beta, P}(f_{t, \lambda_j}, f_{t, \lambda_j}) - \sum_{k=1}^{j-1} (f_{t, \lambda_j} | Q_k f_{t, \lambda_j}) \right] - a_j(x_0, \xi_0) \right|^2 \lesssim t^{-2}.$$

Thus, to prove the part (a) of Theorem 1, we apply the lemma 3.1 with $b_n = n^{-1/2+\varepsilon}$ with $\varepsilon < 1/2$.

Next we continue with the part (b) of Theorem 1. Note that for $j > j_\beta$ the variance of the error does not decay any more, it could even grow —see (26). Thus, inspired by the strong law of large numbers or, more generally, basic ergodicity

results, we perform an average on the parameter t of the family of wave packets. This motivates the next lemma, which is a refinement of [6, Lemma 3.5].

Lemma 3.2. *Consider $1 < \lambda < \infty$ and $m, \beta \in \mathbb{R}$. Assume also that $x_0, \xi_0 \in \mathbb{R}^d$ with $|\xi_0| = 1$ in the definition of the wave packets f_t . Then, for $T > 2^{1/(\lambda-1)}$*

$$\mathbb{E} \left| \frac{1}{T} \int_T^{2T} t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t) dt \right|^2 \simeq T^{2\lambda(2\beta - \frac{1}{2} - m) + 1}.$$

where the implicit constants in the equivalence depend on d, m, λ and β .

Proof. For convenience, set $Q_T = [T, 2T] \times [T, 2T]$. Note that by (24) and (23)

$$\mathbb{E} [\mathcal{E}_\beta(\bar{f}_t, f_t) \overline{\mathcal{E}_\beta(\bar{f}_s, f_s)}] = |(f_t|f_s)_\beta|^2,$$

and consequently,

$$(27) \quad \mathbb{E} \left| \frac{1}{T} \int_T^{2T} t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t) dt \right|^2 = \frac{1}{T^2} \int_{Q_T} t^{-\lambda m} s^{-\lambda m} |(f_t|f_s)_\beta|^2 d(t, s) \simeq T^{-2\lambda m - 2} \int_{Q_T} |(f_t|f_s)_\beta|^2 d(t, s).$$

We claim that

$$(28) \quad (f_t|f_s)_\beta \simeq T^{2\lambda\beta} (f_t|f_s) \text{ for all } (t, s) \in Q_T.$$

To see this we plug in $\hat{f}_t(\xi) = t^{-d/2} e^{-i\xi \cdot x_0} \hat{\chi}((\xi - t^\lambda \xi_0)/t)$ to (4) and recall that $\hat{\chi}$ is non-negative and supported in the ball B_1 . The change of variables $\eta = (\xi - s^\lambda \xi_0)/s$ yields

$$(29) \quad |(f_t|f_s)_\beta| = t^{-d/2} s^{-d/2} \int_{\mathbb{R}^d} \hat{\chi} \left(\frac{\xi - t^\lambda \xi_0}{t} \right) \hat{\chi} \left(\frac{\xi - s^\lambda \xi_0}{s} \right) \langle \xi \rangle^{2\beta} d\xi \\ = t^{-d/2} s^{d/2} \int_{B_1} \hat{\chi} \left(\frac{s}{t} \eta + \frac{s^\lambda - t^\lambda}{t} \xi_0 \right) \hat{\chi}(\eta) \langle s\eta + s^\lambda \xi_0 \rangle^{2\beta} d\eta \simeq s^{2\lambda\beta} |(f_t|f_s)|.$$

To obtain the last line we have used that $1/2s^\lambda \leq \langle s\eta + s^\lambda \xi_0 \rangle$ for $|\eta| \leq 1$ and $s \geq 2^{1/(\lambda-1)}$, and consequently that $\langle s\eta + s^\lambda \xi_0 \rangle \simeq s^\lambda$. This proves the claim. Then, by (27) and (28), to finish the proof of the lemma it is enough to show that

$$(30) \quad \int_{Q_T} |(f_t|f_s)|^2 d(t, s) \simeq T^{3-\lambda}.$$

We begin with the upper bound. Since

$$(31) \quad (f_t|f_s) = t^{d/2} s^{d/2} \int_{\mathbb{R}^d} e^{i(t^\lambda - s^\lambda)x \cdot \xi_0} \chi(t(x - x_0)) \chi(s(x - x_0)) dx,$$

to exploit the cancellations produced by the exponential factor, let L be the formally self-adjoint operator given by $Lf = t + s - i\frac{t-s}{|t-s|}\xi_0 \cdot \nabla f$. Then, since $|\xi_0| = 1$, we have that

$$\begin{aligned} (f_t|f_s) &= \frac{t^{d/2}s^{d/2}}{t+s+|t^\lambda-s^\lambda|} \int_{\mathbb{R}^d} L(e^{i(t^\lambda-s^\lambda)x \cdot \xi_0}) \chi(t(x-x_0)) \chi(s(x-x_0)) \, dx \\ &= \frac{t^{d/2}s^{d/2}}{t+s+|t^\lambda-s^\lambda|} \int_{\mathbb{R}^d} e^{i(t^\lambda-s^\lambda)x \cdot \xi_0} L(\chi(t(x-x_0)) \chi(s(x-x_0))) \, dx, \end{aligned}$$

and using that

$$\begin{aligned} |L(\chi(t(x-x_0)) \chi(s(x-x_0)))| &\leq (t+s)|\chi(t(x-x_0)) \chi(s(x-x_0))| \\ &\quad + t|\nabla \chi(t(x-x_0)) \chi(s(x-x_0))| + s|\chi(t(x-x_0)) \nabla \chi(s(x-x_0))|, \end{aligned}$$

leads, since $t+s \simeq T$ whenever $(t,s) \in Q_T$, to

$$|(f_t|f_s)| \lesssim \frac{t+s}{t+s+|t^\lambda-s^\lambda|} \|\chi\|_1^2 \lesssim \frac{1}{1+T^{-1}|t^\lambda-s^\lambda|}.$$

Moreover, whenever $(t,s) \in Q_T$, we have

$$|(f_t|f_s)|^2 \lesssim T^{1-\lambda} \frac{s^{\lambda-1}}{(1+T^{-1}|t^\lambda-s^\lambda|)^2}.$$

Now, we make the change of variables $u = t$, $v = T^{-1}(s^\lambda - t^\lambda)$ where the volume form is $du \, dv = \lambda T^{-1} s^{\lambda-1} d(t,s)$, and therefore

$$\begin{aligned} (32) \quad \int_{Q_T} |(f_t|f_s)|^2 \, d(t,s) &\lesssim T^{1-\lambda} \int_{Q_T} \frac{s^{\lambda-1}}{(1+T^{-1}|t^\lambda-s^\lambda|)^2} \, d(t,s) \\ &\lesssim T^{2-\lambda} \int_T^{2T} du \int_{-\infty}^{\infty} \frac{1}{(1+|v|)^2} \, dv \simeq T^{3-\lambda}. \end{aligned}$$

This finishes the proof of the upper bound, next we prove the lower bound. Using the second identity in (29) with $\beta = 0$, we have

$$\begin{aligned} |(f_t|f_s)| &= t^{-d/2} s^{d/2} \int_{B_1} \widehat{\chi} \left(\frac{s}{t} \eta + \frac{s^\lambda - t^\lambda}{t} \xi_0 \right) \widehat{\chi}(\eta) \, d\eta \\ &\geq b t^{d/2} s^{-d/2} \int_{B_{1/2}} \widehat{\chi} \left(\frac{s}{t} \eta + \frac{s^\lambda - t^\lambda}{t} \xi_0 \right) \, d\eta, \end{aligned}$$

where to get the last line we have used that by (10), $\widehat{\chi}$ is nonnegative and $\widehat{\chi}(\eta) = b$ in $B_{1/2}$. This also implies that the last integral can be bounded below by b times the measure of the set $B_{1/2} \cap B_{s,t}$ where $B_{s,t}$ is the ball of radius $r = t/(2s)$ and center $c = (t^\lambda - s^\lambda)/s \xi_0$. Whenever $(t,s) \in Q_T$, we have that $r \geq 1/4$ and $|c| \leq T^{-1}|t^\lambda - s^\lambda|$. If additionally, $(t,s) \in Q_T$ is so that $|t^\lambda - s^\lambda| \leq T/4$, then $|c| \leq 1/4$. This implies that actually we have at least $|B_{1/2} \cap B_{s,t}| \geq |B_{1/4}|$ since

the previous estimates for r and c imply that the intersection must contain a ball of radius $1/4$. Therefore, we have

$$|(f_t|f_s)| \geq b^2 t^{d/2} s^{-d/2} |B_{1/2} \cap B_{s,t}| \geq b^2 2^{-d/2} |B_{1/4}|,$$

in the set $D = \{(t, s) \in Q_T : |t^\lambda - s^\lambda| \leq T/4\}$, and consequently that

$$\int_{Q_T} |(f_t|f_s)|^2 d(t, s) \gtrsim \int_D d(t, s) \simeq T^{1-\lambda} \int_D s^{\lambda-1} d(t, s) \simeq T^{2-\lambda} \int_T^{2T} du \int_{-1/4}^{1/4} dv,$$

by the same change of variables used in (32). This proves the lower bound in (30), which ends the proof of the lemma 3.2. \square

By (21) and the lemma 3.2, we have that

$$\mathbb{E} \left| \frac{1}{T} \int_T^{2T} t^{-\lambda_j m_j} [\mathcal{N}_{\beta, P}(f_{t, \lambda_j}, \bar{f}_{t, \lambda_j}) - \sum_{0 < k < j} (f_{t, \lambda_j} | Q_k f_{t, \lambda_j})] dt - a_j(x_0, \xi_0) \right|^2 \lesssim T^{-1-2\epsilon_j}$$

for certain $\epsilon_j > 0$. Thus, the part (b) of Theorem 1 follows from the lemma 3.1 with $b_n = n^{-\epsilon_j/2}$. This finishes the proof of Theorem 1. \square

Next, we show why wave packets are not suitable states to recover the lower order terms of the expansion of the observable.

Proof of Theorem 2. In the view of (8) and the limits (19) and (22), one realizes that the relevant quantities to prove (a) and (b) are the random variables given by the error and the averaged error produced by the wave packets:

$$(33) \quad X_t = t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t), \quad Y_T = \frac{1}{T} \int_T^{2T} t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t) dt.$$

Here we simplify the notation as we did in the section 2, we write f_t instead of $f_{t, \lambda}$ for the states.

The variables X_t and Y_T are complex Gaussian with zero mean. The variable X_t is complex Gaussian because $\mathbb{W}(\bar{f}_t \otimes f_t)$ is a complex Gaussian variable. To see that Y_T is complex Gaussian, we should understand the integral in dt as a limit of finite linear combinations $\{\mathcal{E}_\beta(\bar{f}_{t_1}, f_{t_1}), \dots, \mathcal{E}_\beta(\bar{f}_{t_n}, f_{t_n})\}$ for every $n \in \mathbb{N}$. These linear combinations are actually complex Gaussian variables by the definition and the linearity of \mathbb{W} . Eventually, the limit of complex Gaussian variables is complex Gaussian, so consequently Y_T is a complex Gaussian variable.

Recall that if Z is a complex Gaussian variable with zero mean we have that

$$(34) \quad \mathbb{P}\{|Z| > c\} = \frac{1}{2\pi |\det K|^{1/2}} \int_{|z| > c} \exp(-z \cdot K^{-1} z / 2) dz,$$

where K is the covariance matrix of Z given by

$$K = \frac{1}{2} \begin{bmatrix} \mathbb{E}|Z|^2 + \Re(\mathbb{E}Z^2) & \Im(\mathbb{E}Z^2) \\ \Im(\mathbb{E}Z^2) & \mathbb{E}|Z|^2 - \Re(\mathbb{E}Z^2) \end{bmatrix}.$$

In the special case where $\mathbb{E}|Z|^2 = \sigma^2$ and $\mathbb{E}Z^2 = 0$, one can check —by computing explicitly the integral (34)— that its distribution is generated by

$$(35) \quad \mathbb{P}\{|Z| > c\} = e^{-c^2/\sigma^2}.$$

We will see in the next lemma that $\mathbb{E}X_t^2 = \mathbb{E}Y_T^2 = 0$, and therefore, the distributions of X_t and Y_T are generated in the same way.

Lemma 3.3. *Consider $1 < \lambda < \infty$ and $m, \beta \in \mathbb{R}$. Assume also that $x_0, \xi_0 \in \mathbb{R}^d$ with $|\xi_0| = 1$ in the definition of the wave packets f_t . Then, for any $t, T > 1$ we have*

$$(36) \quad \mathbb{E} \left(t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t) \right)^2 = 0.$$

and

$$(37) \quad \mathbb{E} \left(\frac{1}{T} \int_T^{2T} t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t) dt \right)^2 = 0.$$

Proof. We start with the proof of (37). By (24) and (23)

$$(38) \quad \begin{aligned} \mathbb{E} \left(\mathcal{E}_\beta(\bar{f}_t, f_t) \mathcal{E}_\beta(\bar{f}_s, f_s) \right) &= \mathbb{E} \left(\mathbb{W}(\bar{f}_t \otimes f_t) \mathbb{W}(\bar{f}_s \otimes f_s) \right) \\ &= (f_t \otimes \bar{f}_t | \bar{f}_s \otimes f_s)_{\mathcal{H}} = (f_t | \bar{f}_s)_\beta (\bar{f}_t | f_s)_\beta. \end{aligned}$$

On the one hand, using the notation of the lemma 3.2

$$(39) \quad \begin{aligned} \mathbb{E} \left(\frac{1}{T} \int_T^{2T} t^{-\lambda m} \mathcal{E}_\beta(\bar{f}_t, f_t) dt \right)^2 \\ &= \frac{1}{T^2} \int_{Q_T} t^{-\lambda m} s^{-\lambda m} \mathbb{E} \left(\mathcal{E}_\beta(\bar{f}_t, f_t) \mathcal{E}_\beta(\bar{f}_s, f_s) \right) d(t, s) \\ &= \frac{1}{T^2} \int_{Q_T} t^{-\lambda m} s^{-\lambda m} (\bar{f}_t | f_s)_\beta (f_t | \bar{f}_s)_\beta d(t, s), \end{aligned}$$

Furthermore, note that

$$(40) \quad (\bar{f}_t | f_s)_\beta = t^{-d/2} s^{-d/2} \int_{\mathbb{R}^d} \widehat{\chi} \left(\frac{-\xi - t^\lambda \xi_0}{t} \right) \widehat{\chi} \left(\frac{\xi - s^\lambda \xi_0}{s} \right) \langle \xi \rangle^{2\beta} d\xi = 0,$$

since the supports of the two smooth cut-off functions do not overlap for any $s, t > 1$, $|\xi_0| = 1$, and $1 < \lambda < \infty$. Indeed, the function $\xi \mapsto \widehat{\chi}(-(\xi + t^\lambda \xi_0)/t)$ is supported in a ball of radius t centered at $-t^\lambda \xi_0$ while the function $\xi \mapsto \widehat{\chi}((\xi - s^\lambda \xi_0)/s)$ is supported in a ball of radius s and center $s^\lambda \xi_0$. By (39), this proves (37). Finally, (36) follows directly from (38) and (40) with $s = t$. \square

From (26), we know that

$$\mathbb{E}|X_t|^2 \simeq t^{2\lambda(2\beta-m)},$$

hence, using this in (35), there exists $C > 1$ so that

$$(41) \quad e^{-c^2 C t^{2\lambda(m-2\beta)}} \leq \mathbb{P}\{|X_t| > c\} \leq e^{-c^2 C^{-1} t^{2\lambda(m-2\beta)}}.$$

Thus, if $m < 2\beta$, the previous quantity tends to 1 as t goes to infinity. On the other hand, if $m = 2\beta$, the previous quantity also grows as c becomes smaller. These two facts imply (a). If we turn now our attention to Y_T , we have by the lemma 3.2 that

$$\mathbb{E}|Y_T|^2 \simeq T^{2\lambda(2\beta-\frac{1}{2}-m)+1}.$$

Once again, there exists $C > 1$ so that

$$(42) \quad e^{-c^2 C T^{-2\lambda(2\beta-\frac{1}{2}-m)-1}} \leq \mathbb{P}\{|Y_T| > c\} \leq e^{-c^2 C^{-1} T^{-2\lambda(2\beta-\frac{1}{2}-m)-1}}.$$

Thus, if $m \leq 2\beta - 1/2$, the quantity above tends to 1 as T goes to infinity, which proves (b). This concludes the proof of Theorem 2. \square

Finally, we prove Theorem 3, which provide the rate of convergence in probability of the limits of Theorem 1.

Proof of Theorem 3. Start by proving (a). Using the notation of P_j given in (20), we have, by (21) and the choice of λ_j , that

$$|N^{-\lambda_j m_j} \mathcal{N}_{\beta, P_j}(f_{N, \lambda}, f_{N, \lambda}) - a_j(x_0, \xi_0)| \leq O(N^{-1}) + |N^{-\lambda_j m_j} \mathcal{E}_\beta(\overline{f_{N, \lambda_j}}, f_{N, \lambda_j})|.$$

There exists C' such that $N \geq C'/\varepsilon$,

$$\begin{aligned} \mathbb{P}\{|N^{-\lambda_j m_j} \mathcal{N}_{\beta, P_j}(f_{N, \lambda}, f_{N, \lambda}) - a_j(x_0, \xi_0)| \leq \varepsilon\} &\geq \mathbb{P}\{|N^{-\lambda_j m_j} \mathcal{E}_\beta(\overline{f_{N, \lambda_j}}, f_{N, \lambda_j})| \leq \varepsilon/2\} \\ &= 1 - \mathbb{P}\{|N^{-\lambda_j m_j} \mathcal{E}_\beta(\overline{f_{N, \lambda_j}}, f_{N, \lambda_j})| > \varepsilon/2\}. \end{aligned}$$

By (33) and (41), we know that

$$\mathbb{P}\{|N^{-\lambda_j m_j} \mathcal{E}_\beta(\overline{f_{N, \lambda_j}}, f_{N, \lambda_j})| > \varepsilon/2\} \leq e^{-(\varepsilon^2/4C)N^{2\lambda_j(m_j-2\beta)}}.$$

From this point a simple computation allows to find N_0 . This proves (a).

Next we turn our attention to (b). Using the notation of P_j given in (20), we have, by (21) and the choice of λ_j , that

$$\begin{aligned} \left| \frac{1}{N} \int_N^{2N} t^{-\lambda_j m_j} \mathcal{N}_{\beta, P_j}(f_{t, \lambda}, f_{t, \lambda}) dt - a_j(x_0, \xi_0) \right| \\ = O(N^{-1}) + \left| \frac{1}{N} \int_N^{2N} t^{-\lambda_j m_j} \mathcal{E}_\beta(\overline{f_{t, \lambda}}, f_{t, \lambda}) dt \right|. \end{aligned}$$

Arguing as in (a), we have that

$$\begin{aligned} \mathbb{P}\left\{ \left| \frac{1}{N} \int_N^{2N} t^{-\lambda_j m_j} \mathcal{N}_{\beta, P_j}(f_{t, \lambda}, f_{t, \lambda}) dt - a_j(x_0, \xi_0) \right| \leq \varepsilon \right\} \\ \geq 1 - \mathbb{P}\left\{ \left| \frac{1}{N} \int_N^{2N} t^{-\lambda_j m_j} \mathcal{E}_\beta(\overline{f_{t, \lambda}}, f_{t, \lambda}) dt \right| > \varepsilon/2 \right\}. \end{aligned}$$

By (33) and (42), we know that

$$\mathbb{P} \left\{ \left| \frac{1}{N} \int_N^{2N} t^{-\lambda m_j} \mathcal{E}_\beta(\overline{f_{t,\lambda}}, f_{t,\lambda}) dt \right| > \varepsilon/2 \right\} \leq e^{-(\varepsilon^2/4C)N^{-2\lambda(2\beta-\frac{1}{2}-m)-1}}.$$

Obtaining now N_0 is a simple computation, which proves (b). This concludes the proof of [Theorem 3](#). \square

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