# A NONLINEAR SCHRÖDINGER EQUATION WITH FRACTIONAL NOISE 

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#### Abstract

We study a stochastic Schrödinger equation with a quadratic nonlinearity and a space-time fractional perturbation, in space dimension $d \leq 3$. When the Hurst index is large enough, we prove local well-posedness of the problem using classical arguments. However, for a small Hurst index, even the interpretation of the equation needs some care. In this case, a renormalization procedure must come into the picture, leading to a Wick-type interpretation of the model. Our fixed-point argument then involves some specific regularization properties of the Schrödinger group, which allows us to cope with the strong irregularity of the solution.


## 1. Introduction and main results

1.1. General introduction. In this paper we study the following $d$-dimensional stochastic Schrödinger equation with a quadratic nonlinearity and a space-time fractional perturbation:

$$
\left\{\begin{array}{l}
\imath \partial_{t} u-\Delta u=\rho^{2}|u|^{2}+\dot{B}, \quad t \in[0, T], x \in \mathbb{R}^{d}  \tag{1.1}\\
u_{0}=\phi
\end{array}\right.
$$

where $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is a smooth cut-off function in space and $\dot{B}$ stands for the derivative of a space-time fractional Brownian motion of Hurst index $H=\left(H_{0}, \ldots, H_{d}\right) \in(0,1)^{d+1}$.

We first show that, when $2 H_{0}+\sum_{i=1}^{d} H_{i}>d+1$ (the so-called regular case), the interpretation and local well-posedness of (1.1) can be derived from quite direct arguments, based on a firstorder expansion and the use of Strichartz inequalities.

The equation behaves less favorably when $2 H_{0}+\sum_{i=1}^{d} H_{i} \leq d+1$ (the irregular or rough case). In this situation, we first need a Wick-type renormalization procedure in order to interpret the model. The fixed-point argument then relies on the smoothing properties of the Schrödinger equation, and in particular on its local regularization effect.

We can loosely sum up our results as follows (see the subsequent Sections 1.2, 1.4 for precise statements).
Theorem 1.1. Assume that $1 \leq d \leq 3$ and set $\alpha_{H}:=\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right)-(d+1)$. The following picture holds true:
(i) Case $\alpha_{H}>0$. The equation (1.1) is almost surely locally well-posed in $H^{\beta}\left(\mathbb{R}^{d}\right)$ for some $\beta>0$.
(ii) Case $\alpha_{H} \leq 0$. There exists $\alpha_{d}<0$ such that if $\alpha_{H}>\alpha_{d}$ then the equation (1.1) can be interpreted in the Wick (renormalized) sense and it is almost surely locally well-posed in $H^{-\beta}\left(\mathbb{R}^{d}\right)$ for some $\beta>0$.

We refer to Definition 1.4 and Theorem 1.10 for precise statements in the regular case $(i)$, and to Definition 1.8 and Theorem 1.11 in the rough case ( $i i$ ) (see in particular the condition (H2') for the exact value of $\alpha_{d}$, depending on $d$ ). To our knowledge, and although the statement is still restricted to values of $\alpha_{H}$ close to 0 , Theorem 1.1 (ii) is the first result in the context of

[^0]nonlinear Schrödinger equations where both renormalization arguments and local regularization properties are used to control an irregular noise (in Sobolev spaces of negative order).

The stochastic Schrödinger equation is a widely studied model in the SPDE literature. Just as for stochastic heat or wave equations, the stochastic Schrödinger model admits numerous possible variants and is known to be the source of many challenging questions, whose treatment can only be achieved through the sophisticated combination of PDE tools with probabilistic analysis.

The study of nonlinear stochastic Schrödinger models includes for instance the consideration of a multiplicative noise, that is (see e.g. [10, 13, 16, 17])

$$
\imath \partial_{t} u-\Delta u=u|u|^{2 \sigma}+u \cdot \dot{\xi}, \quad t \in[0, T], x \in \mathbb{T}^{d} \text { or } x \in \mathbb{R}^{d}, \sigma>0,
$$

or the consideration of a random dispersion, that is (see e.g. [14])

$$
\imath \partial_{t} u-\Delta u \cdot \dot{\xi}=u|u|^{2 \sigma}, \quad t \in[0, T], x \in \mathbb{T}^{d} \text { or } x \in \mathbb{R}^{d}, \sigma>0
$$

In these equations, $\dot{\xi}$ stands for a random noise which, in most situations, is taken as white in time.

In the additive-noise framework (which our model (1.1) belongs to), the first mathematical study of nonlinear stochastic Schrödinger equations essentially traces back to a series of works by De Bouard and Debussche (see e.g. [11, 12]), where the authors considered the general dynamics

$$
\begin{equation*}
\imath \partial_{t} u+\Delta u=|u|^{2 \sigma} u+\dot{\xi}, \quad t \in[0, T], x \in \mathbb{R}^{d} \tag{1.2}
\end{equation*}
$$

for suitable values of $\sigma>0$ (depending on $d$ ), and with $\dot{\xi}$ a white noise in time admitting a high regularity in space. These results have been recently extended to noises $\dot{\xi}$ less regular in space, first for the equation on the one-dimensional torus $\mathbb{T}$ with $\sigma=1$ (see [23]), and then for the equation on the whole Euclidean space $\mathbb{R}^{d}$ with suitable values of $\sigma>0$ (see [34, 36]).

In spite of these substantial advances, it can be observed that the analysis of equation (1.2) is so far limited to situations where $\dot{\xi}$ is a white noise in time.

In the present paper, we propose to investigate a new direction in this field, by considering the model of a nonlinear Schrödinger equation with quadratic perturbation and additive forcing term given by a space-time fractional noise. To be more specific, the dynamics we will focus on can be described as follows:

$$
\left\{\begin{array}{l}
\imath \partial_{t} u-\Delta u=\rho^{2}|u|^{2}+\dot{B}, \quad t \in[0, T], x \in \mathbb{R}^{d},  \tag{1.3}\\
u_{0}=\phi,
\end{array}\right.
$$

where

- $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is a smooth compactly-supported function of the space variable, allowing us to bring the analysis back to compact domains,
- $\phi$ is a deterministic initial condition, the regularity of which will be specified later on,
- one has $\dot{B}:=\partial_{t} \partial_{x_{1}} \cdots \partial_{x_{d}} B$, that is $\dot{B}$ is the space-time derivative (in the sense of distributions) of $B$, with $B$ a fractional sheet of Hurst index $H=\left(H_{0}, H_{1}, \ldots, H_{d}\right) \in(0,1)^{d+1}$ (see Definition 2.1 below for details).

The consideration of such a fractional noise $\dot{B}$ represents the main specificity of our analysis, and the main novelty with respect to the above-reported (white-noise) literature. The other new contribution brought by our analysis will consist in the possible treatment of a distributionalvalued solution $u$ for (1.3), which enables us to cover more irregular noises in the model (see the description of the so-called rough case in Sections 1.3 .2 and 1.4.2).

In fact, our aim in the sequel will be to offer as much flexibility as possible regarding the choice of the Hurst index $H \in(0,1)^{d+1}$ that governs $\dot{B}$. Thus, for a (hopefully) large range of such indexes, we intend to - at least - prove local well-posedness of equation (1.3). Note that this objective was already at the core of the investigations of the first author in [20, 21] for a quadratic fractional wave equation.

Before we go further, let us recall that over the last decade, tremendous developments have been observed in the field of singular stochastic PDEs. This progress has been especially prominent in the parabolic (SPDE) setting, with the introduction of the theory of regularity structures [29] or the paracontrolled approach [26]. Among other contributions, those theories provide a convenient framework towards renormalization procedures, thus paving the way to a rigorous treatment of many long-standing problems. If one focuses on additive-noise models only (in the vein of (1.3) , let us quote for instance, among a flourishing literature, the work of Catellier and Chouk [3] about the stochastic quantization equation on the three-dimensional torus

$$
\begin{equation*}
\partial_{t} u-\Delta u=-u^{3}+\xi, \quad t \in[0, T], x \in \mathbb{T}^{3} \tag{1.4}
\end{equation*}
$$

which extends the pioneering results of Da Prato and Debussche [9] for the two-dimensional equation (let us also mention [22] about the consideration of a quadratic nonlinearity - similar to ours - in this parabolic setting, for $d=3$ ). The regularity-structure approach has also proven to be highly effective (see [15, 30]) for the study of the parabolic sine-Gordon model

$$
\begin{equation*}
\partial_{t} u+\frac{1}{2}(1-\Delta) u+\sin (\beta u)=\xi, \quad t \in[0, T], x \in \mathbb{T}^{2} \tag{1.5}
\end{equation*}
$$

with $\xi$ a space-time white noise (see also [37, 38] for recent hyperbolic versions of the model).
Unfortunately, the application of those new groundbreaking approaches beyond the parabolic setting has proved to be very limited so far (this holds true for both regularity structures and paracontrolled theories). To our knowledge, the only attempt to extend such a strategy to a non-parabolic setting is due to Gubinelli, Koch and Oh in their recent work [28], dealing with a stochastic wave model. In particular, we are not aware of any similar extension to a singular stochastic Schrödinger equation.

As regards the deterministic Schrödinger equation with polynomial nonlinearities, its wellposedness in positive Sobolev spaces was established long ago using Strichartz estimates (see [25, 5] and also the monography [4]). More recent developments, applying in particular to NLS with quadratic nonlinearities, also led to well-posedness results for the model in negative Sobolev spaces, thanks to subtle bilinear estimates in the so-called Bourgain spaces (see [7, 41, 2, 31], and also Remark 4.7 below).

With this background in mind, let us now go back to the analysis of equation (1.3). The starting point of the study will be the mild formulation of (1.3), that is the equation

$$
\begin{equation*}
u_{t}=S_{t} \phi-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|u_{\tau}\right|^{2}\right) d \tau+\varrho_{t} \tag{1.6}
\end{equation*}
$$

where $S$ stands for the Schrödinger group, and where we have set

$$
\begin{equation*}
q_{t}:=-\imath \int_{0}^{t} S_{t-\tau}\left(\dot{B}_{\tau}\right) d \tau \tag{1.7}
\end{equation*}
$$

Note that ? can also be seen as the solution of the following "linear" counterpart of (1.3):

$$
\left\{\begin{array}{l}
\imath \partial_{t}!-\Delta!=\dot{B}, \quad t \in[0, T], \quad x \in \mathbb{R}^{d}  \tag{1.8}\\
\quad(0, .)=0
\end{array}\right.
$$

For this reason, we will henceforth refer to ? as the linear solution of the problem.

A first essential part of our work will be to give a precise meaning to both definition (1.7) and equation 1.6). In order to initiate this analysis, let us proceed with a standard transformation of the problem, by considering the equation satisfied by the process $v:=u-?$, that is the equation (which is still formal at this point)

$$
\begin{align*}
v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath & \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho_{\mathrm{O}_{\tau}}\right)\right) d \tau \\
& -\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot\left(\overline{\rho_{\mathrm{O}}} \boldsymbol{)}\right)\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left|\rho_{\mathrm{I}_{\tau}}\right|^{2}\right) d \tau \tag{1.9}
\end{align*}
$$

The main idea behind this transition from $u$ to $v$ is that $?$ is expected to behave as some "firstorder expansion" of $u$. In other words, due to the specific properties of the group $S$ (which we will detail later on), we expect the process $v$ to be more regular than $u$ and $g$, making equation 1.9 ) more tractable than 1.6 . Following this idea, we will focus our subsequent investigations on equation 1.9 .

Let us now elaborate on the successive steps of our reasoning, and introduce our main results. Note that these steps are overall similar to those developed in [20, 21, 27] for the corresponding quadratic wave model. Nevertheless, when going into the technical details, some new fundamental difficulties arise in the analysis of the Schrödinger model, as we will try to highlight it in the presentation below.

### 1.2. Step 1: Study of the linear equation.

Recall that the noise $\dot{B}$ involved in $\sqrt{1.3}$ is defined as the derivative of a fractional sheet $B$, which is a non-differentiable process (in the standard sense). Consequently, just as the white noise $\xi$ in $(1.4)-(1.5), \dot{B}$ can only be understood as some random negative-order distribution, and thus the interpretation of the convolution in (1.7) requires some clarification.

To do so, let us start from a smooth approximation $\left(B_{n}\right)_{n \geq 0}$ of $B$, that is, for each fixed $n \geq 0$, $(t, x) \mapsto B_{n}(t, x)$ is a.s. smooth, and $B_{n}(t, x) \xrightarrow{n \rightarrow \infty} B(t, x)$ for every $(t, x) \in[0, T] \times \mathbb{R}^{d}$ (the choice of such an approximation process will be specified in Section 2 below). Then consider the corresponding sequence of linear solutions, that is the sequence $i_{n}$ of solutions to the equation

$$
\left\{\begin{array}{l}
\imath \partial_{t} \varrho_{n}-\Delta \varrho_{n}=\dot{B}_{n}, \quad t \in[0, T], x \in \mathbb{R}^{d}  \tag{1.10}\\
\varrho_{n}(0, .)=0
\end{array}\right.
$$

where $\dot{B}_{n}:=\partial_{t} \partial_{x_{1}} \cdots \partial_{x_{d}} B_{n}$. For every fixed $n \geq 0$, the smoothness of $B_{n}$ (and accordingly the smoothness of $\dot{B}_{n}$ ) makes the analysis of $(1.10)$ considerably easier than the one of $(1.8)$, and readily allows us to define a unique Gaussian solution process $\left\{{ }^{\circ}(t, x), t \in[0, T], x \in \mathbb{R}^{d}\right\}$ (see Definition 2.3).

The solution ? of the rough equation (1.8) is then interpreted through the following convergence result, which can be seen as our first main contribution:

Proposition 1.2. Let $d \geq 1$ and fix $\left(H_{0}, \ldots, H_{d}\right) \in(0,1)^{d+1}$. Let $\left(B_{n}\right)_{n \geq 0}$ be the sequence of smooth processes defined by formula (2.2), and let ${ }_{n}$ be the solution of $(1.10$ associated with $B_{n}$.

Then, for every test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ (i.e., smooth and compactly-supported), the sequence $\left(\chi \cap_{n}\right)_{n \geq 0}$ converges almost surely in the space $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$, for all $2 \leq p \leq \infty$ and

$$
\begin{equation*}
\alpha>d+1-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right) . \tag{1.11}
\end{equation*}
$$

We denote the limit of this sequence by $\chi$ ?.

The proof of this convergence result relies on the stochastic properties of $\dot{B}$, and will be developed in Section $[2.2$. As the reader might expect it, the resulting regularity property (i.e., the fact that $\chi$ ! $\in \mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$, for every $\alpha$ satisfying (1.11) will be of crucial importance in the analysis of (1.9).

Using a standard patching argument, the limit elements $\left\{\chi\right.$., $\left.\chi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)\right\}$ provided by Proposition 1.2 can then be merged into a single locally-controlled distribution? (see Proposition 2.10), which we will refer to in the sequel.

Remark 1.3. We can compare the above regularity restriction (1.11) for 9 with the corresponding result of [20] in the wave setting, that is when replacing $\imath \partial_{t} ?-\Delta$ ? with $\partial_{t}^{2} ?-\Delta$ ? in (1.8). In the latter situation, and according to [20, Proposition 1.2], one must have

$$
\begin{equation*}
\alpha_{\text {wave }}>d-\frac{1}{2}-\left(H_{0}+\sum_{i=1}^{d} H_{i}\right) . \tag{1.12}
\end{equation*}
$$

Observe in particular that the Hurst indexes are involved through the combination $H_{0}+\sum_{i=1}^{d} H_{i}$ in (1.12), in contrast with the combination $2 \times H_{0}+\sum_{i=1}^{d} H_{i}$ in 1.11). This difference is in fact a direct echo of the behaviour of the underlying linear operators: in the Schrödinger case, time variable somehow "counts twice" with respect to space.

Besides, although such a property cannot be found in the existing literature, we could show along the same pattern that in the heat situation (that is with $\partial_{t}!-\Delta$ ? instead of $\partial_{t}$ ! $Q$ ? ? in (1.8), the restriction for $\alpha$ becomes

$$
\alpha_{\text {heat }}>d-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right),
$$

which, compared to (1.11), reflects the stronger regularization properties of the heat kernel.

### 1.3. Step 2: Interpretation of the main equation.

Now equipped with a proper definition of $?$ (as well as a sharp control on its regularity), we can go back to our interpretation issue for the main equation (1.6) (or equivalently (1.9). In fact, according to the result of Proposition 1.2, two cases need to be distinguished.

### 1.3.1. The regular case.

In the sequel, we will call regular case the situation where

$$
\begin{equation*}
2 H_{0}+\sum_{i=1}^{d} H_{i}>d+1 . \tag{H1}
\end{equation*}
$$

In this case, we can pick $\alpha<0$ such that condition $(1.11)$ is satisfied, and therefore, thanks to Proposition 1.2 , we are allowed to consider every element $\chi$ ? $\left(\chi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)\right)$ as a function of both time and space variables (in an almost sure way). As a result, the square element $|\rho!|^{2}$ involved in (1.9) simply makes sense as a standard pointwise product of functions. This immediately leads us to the following direct interpretation of the equation:

Definition 1.4. Fix $d \geq 1$ and assume that condition (H1) is satisfied. Recall that for every test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}, \chi$ ? is the process defined in Proposition 1.2.

Then we call a solution (on $[0, T]$ ) of equation (1.3) any stochastic process $(u(t, x))_{t \in[0, T], x \in \mathbb{R}^{d}}$ such that, almost surely, the process $v:=u-$ ! is a solution of the mild equation

$$
\begin{aligned}
& v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho_{\cdot}{ }^{\circ}\right)\right) d \tau \\
&-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot\left(\overline{\rho_{\cdot} \cdot}\right)\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left|\rho_{\mathrm{O}_{\tau}}\right|^{2}\right) d \tau, \quad t \in[0, T]
\end{aligned}
$$

### 1.3.2. The rough case.

Let us now turn to the second situation, where

$$
\begin{equation*}
2 H_{0}+\sum_{i=1}^{d} H_{i} \leq d+1 \tag{1.13}
\end{equation*}
$$

In this case, we can no longer pick $\alpha<0$ such that condition 1.11 is satisfied, and so, referring to Proposition 1.2, the element $\rho$ ! involved in 1.9 must be considered as a function of time with values in a space of negative-order distribution. We will call this situation the rough case, to signify that we are here working with very irregular processes.

Naturally, the fact that $\rho_{0_{\tau}}$ must be seen as a distribution (for every fixed $\tau$ ) makes the interpretation of $\left|\rho_{{ }^{\circ} \tau}\right|^{2}$ in (1.9) unclear.

This problem can be emphasized through a regularization approach. Namely, let us go back to the sequence of approximated linear solutions $\left(I_{n}\right)_{n \geq 0}$ satisfying (1.10). For every fixed $n \geq 0$, $i_{n}$ is (almost surely) a function of time and space, and so, for every $(t, x) \in(0, T] \times \mathbb{R}^{d}$, we can compute the moment $\mathbb{E}\left[\left|Q_{n}(t, x)\right|^{2}\right]$. The following asymptotic result (which will be proved in Section 2.4) rules out any possibility to extend the pointwise interpretation to the limit element $|?|^{2}$ :

Proposition 1.5. Fix $d \geq 1$ and assume that $2 H_{0}+\sum_{i=1}^{d} H_{i} \leq d+1$. Then the quantity $\mathbb{E}\left[\left|冃_{n}(t, x)\right|^{2}\right]$ does not depend on $x$. Denoting $\sigma_{n}(t):=\mathbb{E}\left[\left|१_{n}(t, x)\right|^{2}\right]$, the following asymptotic equivalence property holds true: for every $(t, x) \in(0, T] \times \mathbb{R}^{d}$,

$$
\sigma_{n}(t) \underset{n \rightarrow \infty}{\sim} \begin{cases}c_{H}^{1} t n & \text { if } 2 H_{0}+\sum_{i=1}^{d} H_{i}=d+1  \tag{1.14}\\ c_{H}^{2} t 2^{2 n\left(d+1-\left[2 H_{0}+\sum_{i=1}^{d} H_{i}\right]\right)} & \text { if } 2 H_{0}+\sum_{i=1}^{d} H_{i}<d+1\end{cases}
$$

for some constants $c_{H}^{1}, c_{H}^{2}>0$.
A natural way to circumvent this divergence issue and to offer a possible interpretation of $|q|^{2}$ is to turn to a renormalization procedure. In fact, it will here be sufficient to consider the most elementary of those methods, namely the Wick renormalization. Thus, for all fixed $n \geq 0$, $(t, x) \in[0, T] \times \mathbb{R}^{d}$, we set

$$
\begin{equation*}
\wp_{n}(t, x):=\left|१_{n}(t, x)\right|^{2}-\sigma_{n}(t) \quad \text { where } \sigma_{n}(t):=\mathbb{E}\left[\left|१_{n}(t, x)\right|^{2}\right] . \tag{1.15}
\end{equation*}
$$

Our second main contribution now reads as follows:
Proposition 1.6. Let $d \geq 1$ and consider $\left(H_{0}, H_{1}, \ldots, H_{d}\right) \in(0,1)^{d+1}$ such that

$$
\begin{equation*}
d+\frac{3}{4}<2 H_{0}+\sum_{i=1}^{d} H_{i} \leq d+1 \tag{H2}
\end{equation*}
$$

Then, for every test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}$, the sequence $\left(\chi^{2} \propto_{\rho_{n}}\right)_{n \geq 1}$ (defined by 1.15 ) converges almost surely in the space $\mathcal{C}\left([0, T] ; \mathcal{W}^{-2 \alpha, p}\left(\mathbb{R}^{d}\right)\right)$, for all $2 \leq p \leq \infty$ and $\alpha>0$ satisfying (1.11).

We denote the limit of this sequence by $\chi^{2}$ o.

Just as the proof of Proposition 1.2, the proof of Proposition 1.6 will strongly lean on the stochastic properties of $\dot{B}$. It will be the topic of Section 2.3 below.

Remark 1.7. The restriction $2 H_{0}+\sum_{i=1}^{d} H_{i}>d+\frac{3}{4}$ in ( $\mathbf{H 2}$ ) (which will stem from our technical computations) can be compared with the restriction $\sum_{i=0}^{d} H_{i}>d-\frac{3}{4}$ in [20, Proposition 1.4] for the quadratic wave model. We suspect that this condition might not be optimal, that is, we can certainly extend the construction of $\chi^{2} \propto \rho$ below this threshold. However, condition (H2) will prove to be sufficient for our purpose, as it can be seen from the comparison with the more restrictive assumption (H2') in our main theorem (see also Remark 1.12).

With the above constructions in hand, the following Wick interpretation of the equation naturally arises:
Definition 1.8. Fix $d \geq 1$ and assume that condition (H2) is satisfied. Recall that for every test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}, \chi!$ and $\chi^{2} \otimes$ are the processes defined in Proposition 1.2 and Proposition 1.6. respectively.

In this setting, we call a Wick solution (on $[0, T]$ ) of equation (1.3) any stochastic process $(u(t, x))_{t \in[0, T], x \in \mathbb{R}^{d}}$ such that, almost surely, the process $v:=u-\hat{?}$ is a solution of the mild equation

$$
\begin{align*}
& v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho \rho_{\tau}\right)\right) d \tau \\
&-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot\left(\overline{\rho \rho_{\tau}}\right)\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2} \otimes_{\rho_{\tau}}\right) d \tau, \quad t \in[0, T] \tag{1.16}
\end{align*}
$$

Remark 1.9. As the reader may have noticed, our tree notation $\cap$ and $\& \rho$ for the main stochastic processes follows the convention used in [29]: the circle $\circ$ therein stands for the random noise $\dot{B}$, while the line $!$ represents the Duhamel integral operator $\mathcal{I}=\left(i \partial_{t}-\Delta\right)^{-1}$.

### 1.4. Step 3: Local wellposedness results.

At this stage of the procedure, the stochastic part of our analysis can be considered as complete, and our aim now is to solve equation (1.3) (understood along Definition 1.4 or Definition (1.8) by means of a deterministic fixed-point argument.

As in the previous section, and for the sake of clarity, let us separate the regular and rough situations in the presentation of our results.

### 1.4.1. The regular case.

Let us first handle the regular case, where condition (H1) is satisfied and Definition 1.4 of a solution prevails. By relying on the most standard estimates associated with the Schrödinger group $S$ (the so-called Strichartz inequalities, summed up in Lemma 3.3 below), we are here able to establish the following result:

Theorem 1.10 (Local well-posedness under (H1)). Assume that $1 \leq d \leq 4$ and that condition (H1) is satisfied. Let $\beta$ be such that $0<\beta<2 H_{0}+\sum_{i=1}^{d} H_{i}-(d+1)$, and consider the pair $(p, q)$ given by the formulas

$$
p=\frac{12}{d-\beta}, q=\frac{6 d}{2 d+\beta} .
$$

Assume finally that $\phi \in H^{\beta}\left(\mathbb{R}^{d}\right)$. In this setting, the following assertions hold true:
(i) Almost surely, there exists a time $T_{0}>0$ such that equation 1.3) admits a unique solution $u$ (in the sense of Definition 1.4) in the set

$$
\mathcal{S}_{T_{0}}:=ף+X^{\beta}\left(T_{0}\right), \quad \text { where } X^{\beta}\left(T_{0}\right):=\mathcal{C}\left(\left[0, T_{0}\right] ; H^{\beta}\left(\mathbb{R}^{d}\right)\right) \cap L^{p}\left(\left[0, T_{0}\right] ; \mathcal{W}^{\beta, q}\left(\mathbb{R}^{d}\right)\right) .
$$

(ii) For every $n \geq 1$, let $u_{n}$ denote the smooth solution of (1.3), that is $u_{n}$ is the solution (in the sense of Definition (1.4) associated with $\rho \mathrm{I}_{n}$. Then, for every

$$
0<\beta<2 H_{0}+\sum_{i=1}^{d} H_{i}-(d+1)
$$

and for every test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}$, the sequence $\left(\chi u_{n}\right)_{n \geq 1}$ converges almost surely in $\mathcal{C}\left(\left[0, T_{0}\right] ; H^{\beta}\left(\mathbb{R}^{d}\right)\right)$ to $\chi u$, where $u$ is the solution exhibited in item $(i)$.

Let us be slightly more specific about the convergence property in item (ii). In fact, what we will establish in the sequel (see Theorem 3.4 below) is that for every element $\boldsymbol{\Psi}^{\mathbf{1}}$ in a suitable space of functions, the equation obtained by replacing $\rho$ ? with $\Psi^{1}$ in 1.9) admits a unique solution $v$, on some time interval depending only on $\boldsymbol{\Psi}^{\mathbf{1}}$. Besides, this solution is a continuous function of $\boldsymbol{\Psi}^{\mathbf{1}}$ (see Proposition 3.5). Item (i) in the above statement is then an application of these general results to $\boldsymbol{\Psi}^{\mathbf{1}}=\rho!$, while item (ii) corresponds to taking $\boldsymbol{\Psi}^{\mathbf{1}}=\rho_{\mathrm{i}}$, which provides us with the desired time $T_{0}>0$.

### 1.4.2. The rough case.

Let us conclude this presentation of our main results by considering the wellposedness issue for the equation in the rough case. To be more specific, we assume from now on that the assumptions in $\overline{\mathbf{H 2}}$ are satisfied, so that the two processes $\rho$ ? and $\rho^{2} Q_{\rho}$ are well defined and the equation can be understood in the sense of Definition 1.8. In other words, we now focus on the analysis of equation (1.16).

In order to describe the major technical difficulty arising in this case, recall first that under assumption $\mathbf{H 2}$ ) and following Proposition 1.2 , the element $\rho_{\mathfrak{l}_{\tau}}$ must here be treated as a distribution of negative Sobolev regularity (for every fixed time $\tau$ ). Consequently, the term $\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho \rho_{\tau}\right)\left(\right.$ or $\left.\left(\rho v_{\tau}\right) \cdot\left(\overline{\rho \rho_{\tau}}\right)\right)$ in 1.16) can only be understood as the product of a distribution, namely $\rho_{\mathrm{I}_{\tau}}$, with a function, that is $\rho \bar{v}_{\tau}$. Such a product is known to obey the following simple rule (see Lemma 4.1 for a precise statement): if $\rho_{\tau_{\tau}}$ is of Sobolev regularity $-\alpha$ (with $\alpha>0$ ), then $\rho \bar{v}_{\tau}$ must be a function of Sobolev regularity $\beta$ with $\beta>\alpha$, and in this case $\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho \rho_{\tau}\right)$ is indeed well defined as a distribution of negative order $-\alpha$.

Going back to equation 1.16, our only hope to settle a fixed-point argument thus lies in the possibility to control the terms

$$
\int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho \rho_{\tau}\right)\right) d \tau, \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot\left(\rho \overline{\mathfrak{l}}_{\tau}\right)\right) d \tau
$$

as functions of Sobolev regularity $\beta>\alpha>0$. Otherwise stated, we (morally) need convolution with $S$ to produce a regularization effect and allow the transition from $H^{-\alpha}\left(\mathbb{R}^{d}\right)$ to $H^{\beta}\left(\mathbb{R}^{d}\right)$.

Unfortunately, such a regularization property, which corresponds to a well-controlled phenomenon in the heat or in the wave situation, is not as standard in the Schrödinger setting. In particular, the most classical estimates on $S$ (the so-called Strichartz inequalities, which we have already mentioned in the regular situation) cannot provide us with the desired smoothing effect (see Lemma 3.3 for more details).

To overcome this obstacle, we propose to turn to more specific local regularization properties, similar to those exhibited by Constantin and Saut in [8]. It indeed appears that if we only focus on the local regularity of the distributions at stake (meaning that a cut-off function is inserted within the usual Sobolev norm, see (1.19p), then a small gain can be expected from the
convolution with $S$. This will be the topic of our intermediate Lemma 4.3, which can be seen as an extension of the main result in [8].

For this technical property to be implemented here, an additional condition must be imposed on the function $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ in (1.3) (or in 1.16)): namely, we need $\rho$ to be of the form

$$
\rho\left(x_{1}, \ldots, x_{d}\right)=\rho_{1}\left(x_{1}\right) \cdots \rho_{d}\left(x_{d}\right)
$$

for smooth compactly-supported functions $\rho_{1}, \ldots, \rho_{d}$ on $\mathbb{R}$.
Note also that for stability reasons (toward a fixed-point argument), the consideration of a local Sobolev norm is of course not without consequences: it will have to be counterbalanced by a commutator-type estimate, that is a suitable control on switching $\rho$ with the fractional Laplacian in Sobolev norms, which will be the purpose of Lemma 4.4

With the above elements in mind, we are finally in a position to state our main result in the rough situation (the spaces involved in this statement will all be introduced into details in the subsequent Section 1.5):
Theorem 1.11 (Local well-posedness under $\left(\overline{\mathbf{H 2}}{ }^{\prime}\right)$ ). Assume that $1 \leq d \leq 3$ and that the cut-off function $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ in $(1.3)$ is of the form $\left(\mathbf{F}_{\rho}\right)$. Besides, assume that

$$
-\alpha_{d}<2 H_{0}+\sum_{i=1}^{d} H_{i}-(d+1) \leq 0, \quad \text { where } \alpha_{d}:=\left\{\begin{array}{ll}
\frac{3}{20} & \text { if } d=1 \\
\frac{1}{10} & \text { if } d=2 \\
\frac{1}{24} & \text { if } d=3
\end{array} .\right.
$$

Fix $\alpha>0$ such that $d+1-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right)<\alpha<\alpha_{d}$. Then the following assertions hold true: (i) One can find parameters $\kappa \in[2 \alpha, 1 / 2]$ and $p, q \geq 2$ such that almost surely, and for every $\phi \in H^{-2 \alpha}\left(\mathbb{R}^{d}\right)$, there exists a time $T_{0}>0$ for which equation (1.3) admits a unique Wick solution $u$ (in the sense of Definition 1.8) in the set

$$
\mathcal{S}_{T_{0}}:=\imath+X_{\rho}^{\alpha, \kappa,(p, q)}\left(T_{0}\right),
$$

where

$$
X_{\rho}^{\alpha, \kappa,(p, q)}(T):=\mathcal{C}\left([0, T] ; H^{-2 \alpha}\left(\mathbb{R}^{d}\right)\right) \cap L^{p}\left([0, T] ; \mathcal{W}^{-2 \alpha, q}\left(\mathbb{R}^{d}\right)\right) \cap L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}
$$

(ii) For every $n \geq 1$, let $\widetilde{u}_{n}$ denote the smooth Wick solution of (1.3), that is $\widetilde{u}_{n}$ is the solution (in the sense of Definition 1.8) associated with the pair $\left(\rho_{\Omega_{n}}, \rho^{2} \rho_{\rho_{n}}\right)$. Then, for every $\alpha$ satisfying (1.11) and every test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}$, the sequence $\left(\chi \widetilde{u}_{n}\right)_{n \geq 1}$ converges almost surely in $\mathcal{C}\left(\left[0, T_{0}\right] ; H^{-2 \alpha}\left(\mathbb{R}^{d}\right)\right)$ to $\chi u$, where $u$ is the Wick solution exhibited in item (i).

The above wellposedness result can be considered as the main novelty of our work. We are indeed not aware of any previous study of a nonlinear Schrödinger equation involving such an irregular noise, and thus forcing us to rely on both renormalization arguments and local regularization properties.

Remark 1.12. It is interesting to note that the conditions in (H2), which stem from the deterministic part of the study (as emphasized by Theorem 4.5 below), are more restrictive than those in (H2) ensuring the existence of the stochastic element q. This observation contrasts with the situation described for instance in [20, 21, 27, 35], where the application of a similar strategy is, on the contrary, limited by the scope of validity of the stochastic objects.

Remark 1.13. Due to the limitations of the local regularization effect of the Schrödinger group (as reported in Lemma 4.3), further expansions of the strategy, in the spirit of the "second-order analysis" developed e.g. in [21, 35], seem difficult to set up in this Schrödinger setting.

Remark 1.14. As can be seen from the above description of our methodology, the cut-off function $\rho$ in (1.3) is to play two fundamental roles in the study:

- First, it will allow us to bring the analysis of the equation back to compact domains, where the regularity of the driving processes is well controlled (by Propositions 1.2 and 1.6). The situation, in this regard, is somehow equivalent to studying equation (1.3) on a torus (although the definition of the space-time fractional noise on a torus is not as clear as in the current Euclidean setting).
- Secondly, thanks to the involvement of $\rho$, we can appeal to the specific local regularization properties of $S$, which, as we have explained it above, will be our key ingredient toward stability and fixed-point arguments. Observe that no similar regularization result would be available for a study of the equation on a torus.

This being said, in spite of the restriction $\left(\overline{\mathbf{F}_{\rho}}\right)$, the function $\rho$ can still be taken equal to 1 on any arbitrary compact domain, and so, at least locally (in time and in space), our solution $u$ of (1.3) can be regarded as a viable model for the (formal) dynamics

$$
\begin{equation*}
\imath \partial_{t} u-\Delta u=|u|^{2}+\dot{B} . \tag{1.17}
\end{equation*}
$$

A direct analysis of equation (1.17) may be possible through an extension of the subsequent methods to weighted Sobolev spaces (allowing to control the asymptotic behaviour of the processes), but such adaptations are clearly beyond our reach for the moment.

Remark 1.15. Our arguments and results could certainly be extended to the nonlinearity $\rho^{2} u^{2}$ or $\rho^{2} \bar{u}^{2}$ (instead of $\rho^{2}|u|^{2}$ ) through minor modifications of the stochastic constructions of Section 2 (the deterministic well-posedness procedure would then clearly follow the same lines). Besides, the considerations in our study may be seen as a first step towards the treatment of more general power nonlinearities of the form $\rho^{i+2 j} u^{i}|u|^{2 j}$ in (1.1) (instead of $\rho^{2}|u|^{2}$ ). This would include in particular the celebrated cubic model $u|u|^{2}$, with its many physical significances. Nevertheless, it is clear to us that any such extension would be the source of new technical difficulties at every step of the subsequent analysis (stochastic constructions, fixedpoint argument,...). Therefore, at this point, we are not able to provide any specific conjecture in this direction.

Remark 1.16. In their recent work [18], Deng, Nahmod and Yue introduced the fundamental notion of random tensors, presented as a dispersive counterpart of the new parabolic theories (regularity structures, paracontrolled approach and renormalization group techniques). The applications of this promising approach to stochastic PDEs are so far limited to the treatment of irregular random initial conditions within deterministic dynamics. However, it is reasonable to expect that this sophisticated machinery could be extended to dispersive PDEs with additive stochastic noise, such as our model (1.3), and it may offer a way to handle rougher situations regarding $\dot{B}$ (just as the paracontrolled method in [28] allow to cover a space-time white noise situation in dimension 3). We hope to get the opportunity to investigate these questions in a future work.
1.5. Notations. Fix a (space) dimension $d \geq 1$. Throughout the paper, we will call a test function (on $\mathbb{R}^{d}$ ) any function $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ that is smooth and compactly-supported. Besides, we denote by $\mathcal{S}\left(\mathbb{R}^{d}\right)$ the space of Schwartz functions on $\mathbb{R}^{d}$.

We will also refer to the scale of Sobolev spaces defined for all $s \in \mathbb{R}$ and $1 \leq p \leq \infty$ as

$$
\mathcal{W}^{s, p}\left(\mathbb{R}^{d}\right):=\left\{f \in \mathcal{S}^{\prime}\left(\mathbb{R}^{d}\right):\|f\|_{\mathcal{W}^{s, p}}=\left\|\left.\mathcal{F}^{-1}\left(\left\{1+|\cdot|^{2}\right\}^{\frac{s}{2}} \mathcal{F} f\right) \right\rvert\, L^{p}\left(\mathbb{R}^{d}\right)\right\|<\infty\right\},
$$

where the Fourier transform $\mathcal{F}$, resp. the inverse Fourier transform $\mathcal{F}^{-1}$, is defined along the convention: for all $f \in \mathcal{S}\left(\mathbb{R}^{d}\right)$ and $x \in \mathbb{R}^{d}$,

$$
\mathcal{F}(f)(x)=\hat{f}(x):=\int_{\mathbb{R}^{d}} f(y) e^{-\imath\langle x, y\rangle} d y, \quad \text { resp. } \quad \mathcal{F}^{-1}(f)(x):=\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} f(y) e^{\imath\langle x, y\rangle} d y
$$

We then set $H^{s}\left(\mathbb{R}^{d}\right):=\mathcal{W}^{s, 2}\left(\mathbb{R}^{d}\right)$.
Now, as far as space-time functions (or distributions) are concerned, and for the sake of clarity, we will occasionally use the following shortcut notation: for all $T \geq 0,1 \leq p, q \leq \infty$ and $s \in \mathbb{R}$,

$$
\begin{equation*}
L_{T}^{p} \mathcal{W}^{s, q}:=L^{p}\left([0, T] ; \mathcal{W}^{s, q}\left(\mathbb{R}^{d}\right)\right), \quad\|\cdot\|_{L_{T}^{p} \mathcal{W}^{s, q}}:=\|\cdot\|_{L^{p}\left([0, T] ; \mathcal{W}^{s, q}\left(\mathbb{R}^{d}\right)\right)} \tag{1.18}
\end{equation*}
$$

The notation $\mathcal{C}\left([0, T] ; \mathcal{W}^{s, q}\left(\mathbb{R}^{d}\right)\right)$ will refer to the set of continuous functions on $[0, T]$ with values in $\mathcal{W}^{s, q}\left(\mathbb{R}^{d}\right)$.

Let us finally introduce the aforementioned local Sobolev spaces, that will play a prominent role in the analysis of the rough situation. Namely, for all test function $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ and $s \in \mathbb{R}$, we set

$$
\begin{equation*}
H_{\rho}^{s}\left(\mathbb{R}^{d}\right):=\left\{v \in \mathcal{S}^{\prime}\left(\mathbb{R}^{d}\right) ;\left\|\rho \cdot(\operatorname{Id}-\Delta)^{\frac{s}{2}}(v)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}<\infty\right\} \tag{1.19}
\end{equation*}
$$

where the operator $(\operatorname{Id}-\Delta)^{\frac{s}{2}}$ is understood (as usual) through the Fourier transform formula

$$
\mathcal{F}\left((\operatorname{Id}-\Delta)^{\frac{s}{2}}(v)\right)(\xi):=\left\{1+|\xi|^{2}\right\}^{\frac{s}{2}} \mathcal{F} v(\xi)
$$

We endow $H_{\rho}^{s}\left(\mathbb{R}^{d}\right)$ with the natural seminorm $\|v\|_{H_{\rho}^{s}}:=\left\|\rho \cdot(\operatorname{Id}-\Delta)^{\frac{s}{2}}(v)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}$, and finally set, along the convention in 1.18),

$$
L_{T}^{p} H_{\rho}^{s}:=L^{p}\left([0, T] ; H_{\rho}^{s}\left(\mathbb{R}^{d}\right)\right), \quad\|\cdot\|_{L_{T}^{p} H_{\rho}^{s}}:=\|\cdot\|_{L^{p}\left([0, T] ; H_{\rho}^{s}\left(\mathbb{R}^{d}\right)\right)}
$$

### 1.6. Outline of the study.

The rest of the paper is organized as follows. In Section 2 we perform the stochastic constructions which allow to give a sense to the equation. In Section 3 we prove the local well-posedness result in the regular case, while Section 4 is devoted to the analysis in the irregular case. Finally, Section 5 is an appendix in which we gather the proofs of some technical results.

Throughout the paper, and for the sake of clarity, we will use the notation $A \lesssim B$ in order to signify that there exists an irrelevant constant $c$ such that $A \leq c B$.

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## 2. Stochastic constructions

As emphasized in Section 1, our analysis of equation (1.3) will be clearly splitted into a stochastic part (essentially corresponding to the construction of $?$ and $q \circ$ ) and a deterministic part (devoted to the fixed-point procedure).

In this section, we propose to deal with the stochastic step of the study. In other words, we focus here on the proofs of Proposition 1.2, Proposition 1.5 and Proposition 1.6 .

Before we go into the details, and for the sake of clarity, let us recall the specific definition of the process at the core of the model, that is the space-time fractional Brownian motion:

Definition 2.1. Let $d \geq 1$ be a space dimension, $T \geq 0$ a positive time and $(\Omega, \mathcal{F}, \mathbb{P})$ a complete filtered probability space. On this space, and for every fixed $H=\left(H_{0}, H_{1}, \ldots, H_{d}\right) \in(0,1)^{d+1}$, we call a space-time fractional Brownian motion of Hurst index $H$ any centered Gaussian process $B: \Omega \times\left([0, T] \times \mathbb{R}^{d}\right) \rightarrow \mathbb{R}$ with covariance function given by the formula:

$$
\mathbb{E}\left[B\left(s, x_{1}, \ldots, x_{d}\right) B\left(t, y_{1}, \ldots, y_{d}\right)\right]=R_{H_{0}}(s, t) \prod_{i=1}^{d} R_{H_{i}}\left(x_{i}, y_{i}\right), \quad s, t \in[0, T], x, y \in \mathbb{R}^{d}
$$

where

$$
R_{H_{i}}(x, y):=\frac{1}{2}\left(|x|^{2 H_{i}}+|y|^{2 H_{i}}-|x-y|^{2 H_{i}}\right) .
$$

Now remember that our strategy to initiate the construction of both 9 and $q$ consists in the introduction of a smooth approximation $\left(B_{n}\right)_{n \geq 0}$ of $B$ (leading immediately to a smooth approximation $\left(\dot{B}_{n}\right)_{n \geq 0}$ of the noise $\left.\dot{B}\right)$. We will rely here on a sequence derived from the socalled harmonizable representation of $B$, and which happens to be especially suited for Fourier analysis and computations in Sobolev spaces (a similar choice is made in [20], for the same reasons).

Along this idea, let us first introduce, on some complete probability space $(\Omega, \mathcal{F}, \mathbb{P})$, a spacetime white noise $W$ on $\mathbb{R}^{d+1}$. Then fix $H=\left(H_{0}, H_{1}, \ldots, H_{d}\right) \in(0,1)^{d+1}$ and set, for all $t \in[0, T]$ and $x \in \mathbb{R}^{d}$,

$$
\begin{equation*}
B\left(t, x_{1}, \ldots, x_{d}\right):=c \int_{\xi \in \mathbb{R}} \int_{\eta \in \mathbb{R}^{d}} \frac{e^{2 t \xi}-1}{|\xi|^{H_{0}+\frac{1}{2}}} \prod_{i=1}^{d} \frac{e^{2 x_{i} \eta_{i}}-1}{\left|\eta_{i}\right|^{H_{i}+\frac{1}{2}}} \widehat{W}(d \xi, d \eta), \tag{2.1}
\end{equation*}
$$

for some constant $c$, and where $\widehat{W}$ stands for the Fourier transform of $W$.
It is a well-established fact (see e.g. [40]) that for some appropriate value $c=c_{H}$ of the constant, the so-defined process $B$ is a space-time fractional Brownian motion of index $H$ (in the sense of Definition 2.1). Our approximation of $B$ (for every fixed $H=\left(H_{0}, H_{1}, \ldots, H_{d}\right) \in$ $(0,1)^{d+1}$ ) is now obtained through a basic truncation of the integration domain in (2.1): namely, we set $B_{0} \equiv 0$, and for $n \geq 1$,

$$
\begin{equation*}
B_{n}\left(t, x_{1}, \ldots, x_{d}\right)(\omega):=c_{H} \int_{|\xi| \leq 2^{2 n}} \int_{|\eta| \leq 2^{n}} \frac{e^{\imath t \xi}-1}{|\xi|^{H_{0}+\frac{1}{2}}} \prod_{i=1}^{d} \frac{e^{\imath x_{i} \eta_{i}}-1}{\left|\eta_{i}\right|^{H_{i}+\frac{1}{2}}} \widehat{W}(d \xi, d \eta) . \tag{2.2}
\end{equation*}
$$

By a quick examination of the possible singularities in $(\xi, \eta)$, it is not hard to see that, owing to the restricted integration domain, $B_{n}$ indeed defines a smooth process, for every fixed $n \geq 0$.
Remark 2.2. The choice of the scaling $\left\{|\xi| \leq 2^{2 n},|\eta| \leq 2^{n}\right\}$ in (2.2) is directly related to the structure of the Schrödinger operator, and will prove to be essential in the estimation of the renormalization constant (see Section 2.4). This choice naturally contrasts with the "hyperbolic" scaling used in [20] for the corresponding wave model (see also Remark 1.3). Note also that the approximation (2.2) is the same as the one used in [19] for the study of a (rough) parabolic model.

With approximation $\left(B_{n}\right)_{n \geq 0}$ in hand, we now would like to consider the sequence $\left(g_{n}\right)_{n \geq 0}$ of approximated linear solutions, that is the sequence of solutions to

$$
\left\{\begin{array}{l}
\partial_{t} \overbrace{n}-\Delta Q_{n}=\dot{B}_{n}, \quad t \in[0, T], \quad x \in \mathbb{R}^{d},  \tag{2.3}\\
\imath_{n}(0, .)=0,
\end{array}\right.
$$

where, for every $n \geq 0, \dot{B}_{n}$ is defined as the standard derivative $\dot{B}_{n}:=\partial_{t} \partial_{x_{1}} \cdots \partial_{x_{d}} B_{n}$.
Note however that, without further integrability assumptions, the smoothness of $B_{n}$ (for each fixed $n \geq 0$ ) is not a sufficient condition to apply classical deterministic results immediately
ensuring existence and uniqueness of $\varrho_{n}$. A possible way to circumvent the problem in this case is to rely on some stochastic interpretation of $?_{n}$, based on the Gaussianity of the processes under consideration.

In order to justify this interpretation (that is, Definition 2.3 below), let us go back to formula (2.2) for $B_{n}$. Denoting by $S$ the $d$-dimensional Schrödinger group and applying a Fubinitype theorem, the solution $?_{n}$ can (at least formally) be written as

$$
\begin{aligned}
& Q_{n}(t, x)=-\imath \int_{0}^{t} d s \int_{\mathbb{R}^{d}} d y S_{t-s}(x-y) \dot{B}_{n}(s, y) \\
& =c_{H} \imath^{d} \int_{0}^{t} d s \int_{\mathbb{R}^{d}} d y \int_{|\xi| \leq 2^{2 n}} \int_{|\eta| \leq 2^{n}} S_{t-s}(x-y) \frac{\xi}{|\xi|^{H_{0}+\frac{1}{2}}} \prod_{i=1}^{d} \frac{\eta_{i}}{\left|\eta_{i}\right|^{H_{i}+\frac{1}{2}}} e^{\imath \xi s} e^{\imath\langle\eta, y\rangle} \widehat{W}(d \xi, d \eta) \\
& =c_{H} \imath^{d} \int_{|\xi| \leq 2^{2 n}} \int_{|\eta| \leq 2^{n}} \frac{\xi}{|\xi|^{H_{0}+\frac{1}{2}}} \prod_{i=1}^{d} \frac{\eta_{i}}{\left|\eta_{i}\right|^{H_{i}+\frac{1}{2}}} e^{\imath\langle\eta, x\rangle}\left[\int_{0}^{t} d s e^{\imath \xi s}\left(\int_{\mathbb{R}^{d}} d y S_{t-s}(x-y) e^{-\imath\langle\eta, x-y\rangle}\right)\right] \widehat{W}(d \xi, d \eta) \\
& =c_{H} \imath^{d} \int_{|\xi| \leq 2^{2 n}} \int_{|\eta| \leq 2^{n}} \frac{\xi}{|\xi|^{H_{0}+\frac{1}{2}}} \prod_{i=1}^{d} \frac{\eta_{i}}{\left|\eta_{i}\right|^{H_{i}+\frac{1}{2}}} e^{\imath\langle\eta, x\rangle}\left[\int_{0}^{t} d s e^{\imath \xi(t-s)}\left(\int_{\mathbb{R}^{d}} d y S_{s}(y) e^{-\imath\langle\eta, y\rangle}\right)\right] \widehat{W}(d \xi, d \eta) .
\end{aligned}
$$

At this point, remember that the spatial Fourier transform of $S$ is explicitly given by

$$
\int_{\mathbb{R}^{d}} d x e^{-\imath\langle\xi, x\rangle} S_{t}(x)=e^{\imath t|\xi|^{2}}
$$

and so we end up with the (a priori formal) representation

$$
\begin{equation*}
\varrho_{n}(t, x)=c_{H} \imath^{d} \int_{|\xi| \leq 2^{2 n}} \int_{|\eta| \leq 2^{n}} \frac{\xi}{|\xi|^{H_{0}+\frac{1}{2}}} \prod_{i=1}^{d} \frac{\eta_{i}}{\left|\eta_{i}\right|^{H_{i}+\frac{1}{2}}} e^{\imath\langle\eta, x\rangle} \gamma_{t}(\xi,|\eta|) \widehat{W}(d \xi, d \eta) \tag{2.4}
\end{equation*}
$$

where for all $t \geq 0, \xi \in \mathbb{R}$ and $r>0$, we define the quantity $\gamma_{t}(\xi, r)$ as

$$
\begin{equation*}
\gamma_{t}(\xi, r):=e^{\imath \xi t} \int_{0}^{t} e^{\imath s\left\{r^{2}-\xi\right\}} d s \tag{2.5}
\end{equation*}
$$

Thanks to formula 2.4 , we are in a position to offer the following natural (and rigorous) stochastic definition for the solution of (2.3):

Definition 2.3. We call a solution of equation (2.3) (or linear solution associated with (1.3)) any centered complex Gaussian process

$$
\left\{१_{n}(s, x), n \geq 1, s \geq 0, x \in \mathbb{R}^{d}\right\}
$$

whose covariance is given by the relations: for all $n, m \geq 1, s, t \geq 0$ and $x, y \in \mathbb{R}^{d}$,

$$
\begin{align*}
& \mathbb{E}\left[Q_{n}(s, x) \overline{Q_{m}(t, y)}\right]=c_{H}^{2} \int_{(\xi, \eta) \in D_{n} \cap D_{m}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}} \gamma_{s}(\xi,|\eta|) \overline{\gamma_{t}(\xi,|\eta|)} e^{\langle\eta, x-y\rangle} d \xi d \eta,  \tag{2.6}\\
& \mathbb{E}\left[q_{n}(s, x) \varrho_{m}(t, y)\right]=-c_{H}^{2} \int_{(\xi, \eta) \in D_{n} \cap D_{m}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}} \gamma_{s}(\xi,|\eta|) \gamma_{t}(-\xi,|\eta|) e^{\iota\langle\eta, x-y\rangle} d \xi d \eta, \tag{2.7}
\end{align*}
$$

where $D_{n}:=B_{2 n}^{1} \times B_{n}^{d}$ with $B_{\ell}^{k}:=\left\{\lambda \in \mathbb{R}^{k}:|\lambda| \leq 2^{\ell}\right\}$.
2.1. Preliminary estimates. As a first step toward Proposition 1.2 and Proposition 1.6 , let us collect some essential estimates on the quantity $\gamma_{t}(\xi, r)$ at the core of the covariance formula (2.6) (and explicitly defined by 2.5).

To this end, we set, for all $0 \leq s \leq t, \xi \in \mathbb{R}$ and $r>0$,

$$
\gamma_{s, t}(\xi, r):=\gamma_{t}(\xi, r)-\gamma_{s}(\xi, r)
$$

Lemma 2.4. For all $T \geq 0,0 \leq s \leq t \leq T, \xi \in \mathbb{R}, r>0$ and $\kappa, \lambda \in[0,1]$, it holds that

$$
\left|\gamma_{s, t}(\xi, r)\right| \lesssim \min \left(|\xi|^{\kappa}|t-s|^{\kappa}+|t-s|, \frac{|t-s| r^{2}}{|\xi|}+\frac{|t-s|^{\kappa}\left\{1+r^{2}\right\}}{|\xi|^{1-\kappa}}, \frac{|t-s|^{\kappa}\left\{r^{2 \kappa}+|\xi|^{\kappa}\right\}}{\left||\xi|-r^{2}\right|^{1-\lambda(1-\kappa)}}\right)
$$

where the proportional constant in $\lesssim$ only depends on $T$.
Proof. To begin with, let us write

$$
\gamma_{s, t}(\xi, r)=\left\{e^{\imath \xi t}-e^{\imath \xi s}\right\} \int_{0}^{s} e^{\imath u\left\{r^{2}-\xi\right\}} d u+e^{\imath \xi t} \int_{s}^{t} e^{\imath u\left\{r^{2}-\xi\right\}} d u
$$

and so

$$
\left|\gamma_{s, t}(\xi, r)\right| \lesssim\left|e^{\imath \xi t}-e^{\imath \xi s}\right|\left|\int_{0}^{s} e^{\imath u\left\{r^{2}-\xi\right\}} d u\right|+\left|\int_{s}^{t} e^{\imath u\left\{r^{2}-\xi\right\}} d u\right| \lesssim|\xi|^{\kappa}|t-s|^{\kappa}+|t-s|
$$

Then observe that

$$
\gamma_{t}(\xi, r)=e^{\imath \xi t} \int_{0}^{t} e^{\imath s\left\{r^{2}-\xi\right\}} d s=-\frac{e^{\imath r^{2} t}-e^{\imath \xi t}}{\imath \xi}+\frac{e^{\imath \xi t} r^{2}}{\xi} \int_{0}^{t} e^{\imath s\left\{r^{2}-\xi\right\}} d s
$$

which readily entails

$$
\begin{aligned}
& \gamma_{s, t}(\xi, r) \\
& =-\frac{\left\{e^{\imath r^{2} t}-e^{\imath r^{2} s}\right\}-\left\{e^{\imath \xi t}-e^{\imath \xi s}\right\}}{\imath \xi}+\frac{r^{2}}{\xi}\left\{e^{\imath \xi t}-e^{\imath \xi s}\right\} \int_{0}^{s} e^{\imath u\left\{r^{2}-\xi\right\}} d u+\frac{e^{\imath \xi t} r^{2}}{\xi} \int_{s}^{t} e^{\imath u\left\{r^{2}-\xi\right\}} d u
\end{aligned}
$$

Thus,

$$
\begin{aligned}
\left|\gamma_{s, t}(\xi, r)\right| & \lesssim r^{2} \frac{|t-s|}{|\xi|}+\frac{|t-s|^{\kappa}}{|\xi|^{1-\kappa}}+r^{2} \frac{|t-s|^{\kappa}}{|\xi|^{1-\kappa}}+\frac{r^{2}}{|\xi|}|t-s| \\
& \lesssim r^{2} \frac{|t-s|}{|\xi|}+\left\{1+r^{2}\right\} \frac{|t-s|^{\kappa}}{|\xi|^{1-\kappa}}
\end{aligned}
$$

Finally, it can be checked that

$$
\gamma_{t}(\xi, r)=\frac{\imath}{r^{2}-\xi}\left\{e^{\imath \xi t}-e^{\imath r^{2} t}\right\}
$$

which yields

$$
\begin{aligned}
\left|\gamma_{s, t}(\xi, r)\right| & =\frac{1}{\left|\xi-r^{2}\right|}\left|\left\{e^{\imath r^{2} t}-e^{\imath r^{2} s}\right\}-\left\{e^{\imath \xi t}-e^{\imath \xi s}\right\}\right|^{\kappa}\left|\left\{e^{\imath r^{2} t}-e^{\imath \xi t}\right\}-\left\{e^{\imath r^{2} s}-e^{\imath \xi s}\right\}\right|^{1-\kappa} \\
& \lesssim \frac{|t-s|^{\kappa}}{\left|\xi-r^{2}\right|}\left\{r^{2}+|\xi|\right\}^{\kappa}\left\{\left|e^{\imath r^{2} t}-e^{\imath \xi t}\right|^{\lambda}+\left|e^{\imath r^{2} s}-e^{\imath \xi s}\right|^{\lambda}\right\}^{1-\kappa} \\
& \lesssim \frac{\mid t-s s^{\kappa}}{\left|\xi-r^{2}\right|^{1-\lambda(1-\kappa)}}\left\{r^{2 \kappa}+|\xi|^{\kappa}\right\} \lesssim \frac{|t-s|^{\kappa}}{\left||\xi|-r^{2}\right|^{1-\lambda(1-\kappa)}}\left\{r^{2 \kappa}+|\xi|^{\kappa}\right\} .
\end{aligned}
$$

Corollary 2.5. For all $T \geq 0,0 \leq s \leq t \leq T, H \in(0,1), r>0, \varepsilon \in(0,1)$ and $\kappa \in$ $\left[0, \min \left(H, \frac{1-\varepsilon}{2}\right)\right)$, it holds that

$$
\int_{\mathbb{R}} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi \lesssim \frac{|t-s|^{2 \kappa}}{1+r^{4(H-\kappa)-2-2 \varepsilon}},
$$

where the proportional constant in $\lesssim$ only depends on $T$.
Proof. We will naturally lean on the estimates exhibited in Lemma 2.4
For $0<r<1$, we have

$$
\int_{\mathbb{R}} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi \lesssim|t-s|^{2 \kappa}\left[\int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{2 H-1}}+\int_{|\xi| \geq 1} \frac{d \xi}{|\xi|^{2(H-\kappa)+1}}\right] \lesssim|t-s|^{2 \kappa} .
$$

Then, for $r>1$, let us consider the decomposition

$$
\int_{\mathbb{R}} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi=\int_{\||\xi|-r^{2} \left\lvert\, \geq \frac{\xi \xi}{2}\right.} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi+\int_{\| \xi\left|-r^{2}\right| \leq \frac{|\xi|}{2}} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi .
$$

On the one hand, it holds that

$$
\begin{aligned}
\int_{\left||\xi|-r^{2}\right| \geq \frac{|\xi|}{2}} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi & \lesssim|t-s|^{2 \kappa} \int_{\left\{|\xi| \leq \frac{2}{3} r^{2}\right\} \cup\left\{|\xi| \geq 2 r^{2}\right\}} \frac{r^{4 \kappa}+|\xi|^{2 \kappa}}{|\xi|^{2 H-1}| | \xi\left|-r^{2}\right|^{2}} d \xi \\
& \lesssim \frac{|t-s|^{2 \kappa}}{r^{4(H-\kappa)}} \int_{\left\{|\xi| \leq \frac{2}{3}\right\} \cup\{|\xi| \geq 2\}} \frac{1+|\xi|^{2 \kappa}}{|\xi|^{2 H-1}| | \xi|-1|^{2}} d \xi \lesssim \frac{|t-s|^{2 \kappa}}{r^{4(H-\kappa)}} .
\end{aligned}
$$

On the other hand, setting $\lambda=\frac{1+\varepsilon}{2(1-\kappa)}$, one has

$$
\begin{aligned}
\int_{\left||\xi|-r^{2}\right| \leq \frac{|\xi|}{2}} \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H-1}} d \xi & \lesssim|t-s|^{2 \kappa} \int_{\left\{\frac{2}{3} r^{2} \leq|\xi| \leq 2 r^{2}\right\}} \frac{r^{4 \kappa}+|\xi|^{2 \kappa}}{|\xi|^{2 H-1}| | \xi\left|-r^{2}\right|^{2-2 \lambda(1-\kappa)}} d \xi \\
& \lesssim \frac{|t-s|^{2 \kappa}}{r^{4(H-\kappa)-4 \lambda(1-\kappa)}} \int_{\left\{\frac{2}{3} \leq|\xi| \leq 2\right\}} \frac{1}{|\xi|^{2 H-1}| | \xi|-1|^{2-2 \lambda(1-\kappa)}} d \xi \\
& \lesssim \frac{|t-s|^{2 \kappa}}{r^{4(H-\kappa)-2-2 \varepsilon}} .
\end{aligned}
$$

Let us also take advantage of this preliminary section to introduce the following lemma, which accounts for the technical simplifications offered by the test function $\chi$ in the forthcoming computations.

Lemma 2.6. Let $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be a test function and fix $\sigma \in \mathbb{R}$. Then, for every $p \geq 1$ and for all $\eta_{1}, \ldots, \eta_{p} \in \mathbb{R}^{d}$, it holds that

$$
\left|\int_{\mathbb{R}^{d}} d x \prod_{i=1}^{p} \int_{\left(\mathbb{R}^{d}\right)^{2}} \frac{d \lambda_{i} d \tilde{\lambda}_{i}}{\left\{1+\left|\lambda_{i}\right|^{2}\right\}^{\frac{\sigma}{2}}\left\{1+\left|\tilde{\lambda}_{i}\right|^{2}\right\}^{\frac{\sigma}{2}}} e^{\imath\left\langle x, \lambda_{i}-\tilde{\lambda}_{i}\right\rangle} \widehat{\chi}\left(\lambda_{i}-\eta_{i}\right) \overline{\hat{\chi}\left(\tilde{\lambda}_{i}-\eta_{i}\right)}\right| \lesssim \prod_{i=1}^{p} \frac{1}{\left\{1+\left|\eta_{i}\right|^{2}\right\}^{\sigma}},
$$

where the proportional constant only depends on $\chi$ and $\sigma$.
Proof. Let us set, for all $\eta, \lambda \in \mathbb{R}^{d}$,

$$
\Gamma_{\eta}(\lambda):=\widehat{\chi}(\lambda-\eta) \frac{1}{\left\{1+|\lambda|^{2}\right\}^{\frac{\sigma}{2}}}
$$

and observe that the integral under consideration can then be written as

$$
\begin{align*}
& \int_{\mathbb{R}^{d}} d x \prod_{i=1}^{p} \int_{\left(\mathbb{R}^{d}\right)^{2}} \frac{d \lambda_{i} d \tilde{\lambda}_{i}}{\left\{1+\left|\lambda_{i}\right|^{2}\right\}^{\frac{\sigma}{2}}\left\{1+\left|\tilde{\lambda}_{i}\right|^{2}\right\}^{\frac{\sigma}{2}}} e^{\imath\left\langle x, \lambda_{i}-\tilde{\lambda}_{i}\right\rangle} \widehat{\chi}\left(\lambda_{i}-\eta_{i}\right) \overline{\widehat{\chi}\left(\tilde{\lambda}_{i}-\eta_{i}\right)} \\
& =c \int_{\mathbb{R}^{d}} d x \prod_{i=1}^{p}\left|\mathcal{F}^{-1}\left(\Gamma_{\eta_{i}}\right)(x)\right|^{2} \\
& =c \int_{\mathbb{R}^{d}} d x\left|\mathcal{F}^{-1}\left(\Gamma_{\eta_{1}} * \cdots * \Gamma_{\eta_{p}}\right)(x)\right|^{2} \\
& \lesssim\left\|\Gamma_{\eta_{1}} * \cdots * \Gamma_{\eta_{p}}\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{2} \lesssim\left\|\Gamma_{\eta_{1}}\right\|_{L^{1}\left(\mathbb{R}^{d}\right)}^{2} \cdots\left\|\Gamma_{\eta_{p-1}}\right\|_{L^{1}\left(\mathbb{R}^{d}\right)}^{2}\left\|\Gamma_{\eta_{p}}\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{2}, \tag{2.8}
\end{align*}
$$

where the last inequality is derived from Plancherel theorem.
The conclusion now comes from the fact that for all $\eta, \lambda \in \mathbb{R}^{d}$ and for every $\kappa>0$,

$$
\begin{equation*}
\left|\widehat{\chi}(\lambda)\left\{1+|\lambda+\eta|^{2}\right\}^{-\frac{\sigma}{2}}\right| \leq c_{\sigma, \chi, \kappa}\left\{1+|\lambda|^{2}\right\}^{-\kappa}\left\{1+|\eta|^{2}\right\}^{-\frac{\sigma}{2}} . \tag{2.9}
\end{equation*}
$$

Let us briefly verify (2.9) for $\sigma \geq 0$ (the proof for $\sigma<0$ being immediate). In fact, since $\chi$ is smooth and compactly-supported, one has

$$
\begin{aligned}
& \left|\hat{\chi}(\lambda)\left\{1+|\lambda+\eta|^{2}\right\}^{-\frac{\sigma}{2}}\right| \\
& =\left|\hat{\chi}(\lambda)\left\{1+|\lambda+\eta|^{2}\right\}^{-\frac{\sigma}{2}}\right| \mathbf{1}_{\left\{|\lambda| \geq \frac{1}{2}|\eta|\right\}}+\left|\hat{\chi}(\lambda)\left\{1+|\lambda+\eta|^{2}\right\}^{-\frac{\sigma}{2}}\right| \mathbf{1}_{\left\{|\lambda|<\frac{1}{2}|\eta|\right\}} \\
& \leq c_{\sigma}\left[|\hat{\chi}(\lambda)| \mathbf{1}_{\left\{|\lambda| \geq \frac{1}{2}|\eta|\right\}}+|\hat{\chi}(\lambda)|\left\{1+|\eta|^{2}\right\}^{-\frac{\sigma}{2}} \mathbf{1}_{\left\{|\lambda|<\frac{1}{2}|\eta|\right\}}\right] \\
& \leq c_{\sigma, \chi, \kappa}\left[\left\{1+|\lambda|^{2}\right\}^{-\kappa}\left\{1+|\lambda|^{2}\right\}^{-\frac{\sigma}{2}} \mathbf{1}_{\left\{|\lambda| \geq \frac{1}{2}|\eta|\right\}}+\left\{1+|\lambda|^{2}\right\}^{-\kappa}\left\{1+|\eta|^{2}\right\}^{-\frac{\sigma}{2}} \mathbf{1}_{\left\{|\lambda|<\frac{1}{2}|\eta|\right\}}\right] \\
& \leq c_{\sigma, \chi, \kappa}\left\{1+|\lambda|^{2}\right\}^{-\kappa}\left\{1+|\eta|^{2}\right\}^{-\frac{\sigma}{2}} .
\end{aligned}
$$

Going back to (2.8), and with estimate (2.9) in hand, we can check for instance that

$$
\left\|\Gamma_{\eta_{1}}\right\|_{L^{1}\left(\mathbb{R}^{d}\right)}=\int_{\mathbb{R}^{d}} d \lambda|\widehat{\chi}(\lambda)| \frac{1}{\left\{1+\left|\lambda+\eta_{1}\right|^{2}\right\}^{\frac{\sigma}{2}}} \leq \frac{c_{\sigma, \chi}}{\left\{1+\left|\eta_{1}\right|^{2}\right\}^{\frac{\sigma}{2}}},
$$

which yields the desired bound.
2.2. Proof of Proposition 1.2, For simplicity, we will assume for the whole proof that $T=1$ and we set, for all $m, n \geq 0, \varrho_{n, m}:=\varrho_{m}-\varrho_{n}$. Besides, along the statement of the proposition, we fix $\alpha$ satisfying (1.11).

Step 1: Let us show that for all $p \geq 1,1 \leq n \leq m, 0 \leq s \leq t \leq 1$ and $\varepsilon, \kappa>0$ small enough, one has

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} \mathbb{E}[|\mathcal{F}^{-1}(\left\{1+|\cdot|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}(\chi[\S_{n, m}(t, .)-\overbrace{n, m}(s, .)]))(x)|^{2 p}] d x \lesssim 2^{-4 n \varepsilon p}|t-s|^{2 \kappa p} \tag{2.10}
\end{equation*}
$$

where the proportional constant only depends on $p, \alpha$ and $\chi$.
First, we can observe that the random variable under consideration is Gaussian, and so, for every $p \geq 1$, one has

$$
\begin{align*}
& \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}\left(\chi\left[\varrho_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)\right)(x)\right|^{2 p}\right] \\
& \leq c_{p} \mathbb{E}[|\mathcal{F}^{-1}(\left\{1+|.|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}(\chi[\overbrace{n, m}(t, .)-\varrho_{n, m}(s, .)]))(x)|^{2}]^{p} \tag{2.11}
\end{align*}
$$

where the constant $c_{p}$ only depends on $p$.

Let us then expand this variable as

$$
\begin{aligned}
& \mathcal{F}^{-1}\left(\left\{1+|\cdot|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}\left(\chi\left[\varrho_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)\right)(x) \\
& =\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} d \lambda e^{\imath\langle x, \lambda\rangle}\left\{1+|\lambda|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}\left(\chi\left[\varrho_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)(\lambda) \\
& =\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} d \lambda\left\{1+|\lambda|^{2}\right\}^{-\frac{\alpha}{2}} e^{\imath\langle x, \lambda\rangle}\left(\int_{\mathbb{R}^{d}} d \beta \widehat{\chi}(\lambda-\beta) \mathcal{F}\left(\left[\Omega_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)(\beta)\right),
\end{aligned}
$$

so that

$$
\begin{align*}
& \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|\cdot|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}\left(\chi\left[\varrho_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)\right)(x)\right|^{2}\right] \\
& =\frac{1}{(2 \pi)^{2 d}} \iint_{\left(\mathbb{R}^{d}\right)^{2}} \frac{d \lambda d \tilde{\lambda}}{\left\{1+|\lambda|^{2}\right\}^{\frac{\alpha}{2}}\left\{1+|\tilde{\lambda}|^{2}\right\}^{\frac{\alpha}{2}}}{ }^{\imath\langle x, \lambda-\tilde{\lambda}\rangle} \iint_{\left(\mathbb{R}^{d}\right)^{2}} d \beta d \tilde{\beta} \widehat{\chi}(\lambda-\beta) \overline{\widehat{\chi}(\tilde{\lambda}-\tilde{\beta})} \mathcal{Q}_{n, m ; s, t}(\beta, \tilde{\beta}), \tag{2.12}
\end{align*}
$$

with

$$
\mathcal{Q}_{n, m ; s, t}(\beta, \tilde{\beta}):=\mathbb{E}\left[\mathcal{F}\left(\left[\varrho_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)(\beta) \overline{\mathcal{F}\left(\left[\varrho_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)(\tilde{\beta})}\right] .
$$

To expand the latter quantity, observe first that, using the covariance formula (2.6, we get

$$
\begin{aligned}
& \mathbb{E}\left[\left\{\left\{_{n, m}(t, y)-\varrho_{n, m}(s, y)\right\} \overline{\left\{\varrho_{n, m}(t, \tilde{y})-\varrho_{n, m}(s, \tilde{y})\right\}}\right]\right. \\
& =c \int_{\left\{(\xi, \eta) \in D_{n, m}\right\}} \frac{1}{\left.|\xi|\right|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2} e^{\imath\langle\eta, y\rangle} e^{-\iota\langle\eta, \tilde{y}\rangle} d \xi d \eta
\end{aligned}
$$

where $D_{n, m}:=\left(B_{2 m}^{1} \times B_{m}^{d}\right) \backslash\left(B_{2 n}^{1} \times B_{n}^{d}\right)$, and hence

$$
\begin{equation*}
\mathcal{Q}_{n, m ; s, t}(\beta, \tilde{\beta})=c \int_{\left\{(\xi, \eta) \in D_{n, m}\right\}} d \xi d \eta \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2} \delta_{\beta=\eta} \delta_{\tilde{\beta}=\eta} . \tag{2.13}
\end{equation*}
$$

Combining (2.11)-2.12)-2.13) and using a standard Fubini argument, we end up with the estimate

$$
\begin{aligned}
& \int_{\mathbb{R}^{d}} d x \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+\left.|\cdot|\right|^{2}\right\}^{-\frac{\alpha}{2}} \mathcal{F}\left(\chi\left[१_{n, m}(t, .)-\varrho_{n, m}(s, .)\right]\right)\right)(x)\right|^{2 p}\right] \\
& \lesssim \int_{\mathbb{R}^{d}} d x\left(\int_{\left\{(\xi, \eta) \in D_{n, m}\right\}} d \xi d \eta \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2}\right. \\
& \left.\iint_{\left(\mathbb{R}^{d}\right)^{2}} \frac{d \lambda d \tilde{\lambda}}{\left\{1+|\lambda|^{2}\right\}^{\frac{\alpha}{2}}\left\{1+|\tilde{\lambda}|^{2}\right\}^{\frac{\alpha}{2}}} e^{\langle\langle x, \lambda-\tilde{\lambda}\rangle} \widehat{\chi}(\lambda-\eta) \overline{\widehat{\widehat{\chi}}(\tilde{\lambda}-\eta)}\right)^{p} \\
& \lesssim\left(\int_{(\xi, \eta) \in D_{n, m}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left\{1+|\eta|^{2}\right\}^{-\alpha}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2} d \xi d \eta\right)^{p},
\end{aligned}
$$

where we have used Lemma 2.6 to get the last inequality.

Now we can obviously write

$$
\begin{align*}
& \left(\int_{(\xi, \eta) \in D_{n, m}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left\{1+|\eta|^{2}\right\}^{-\alpha}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2} d \xi d \eta\right)^{p} \\
& \lesssim\left(\int_{2^{2 n} \leq|\xi| \leq 2^{2 m}} \int_{|\eta| \leq 2^{m}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left\{1+|\eta|^{2}\right\}^{-\alpha}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2} d \xi d \eta\right)^{p} \\
& \quad \quad+\left(\int_{|\xi| \leq 2^{2 m}} \int_{2^{n} \leq|\eta| \leq 2^{m}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left\{1+|\eta|^{2}\right\}^{-\alpha}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2} d \xi d \eta\right)^{p} \\
& =:\left(\mathbb{I}_{n, m}(s, t)\right)^{p}+\left(\mathbb{I I}_{n, m}(s, t)\right)^{p} . \tag{2.14}
\end{align*}
$$

Let us focus on the estimation of $\mathbb{I}_{n, m}(s, t)$. To this end, we fix $0<\varepsilon<\min \left(H_{0}, \frac{1}{2}\right)$, so that

$$
\begin{equation*}
\mathbb{I}_{n, m}(s, t) \leq 2^{-4 n \varepsilon} \int_{\mathbb{R}} d \xi \int_{\mathbb{R}^{d}} d \eta \frac{1}{|\xi|^{2 H_{0}-2 \varepsilon-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left\{1+|\eta|^{2}\right\}^{-\alpha}\left|\gamma_{s, t}(\xi,|\eta|)\right|^{2}, \tag{2.15}
\end{equation*}
$$

which, through an elementary hyperspherical change of variables, leads us to

$$
\begin{equation*}
\mathbb{I}_{n, m}(s, t) \lesssim 2^{-4 n \varepsilon} \int_{0}^{\infty} d r \frac{\left\{1+r^{2}\right\}^{-\alpha}}{r^{2\left(H_{1}+\cdots+H_{d}\right)-2 d+1}}\left(\int_{\mathbb{R}} d \xi \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H_{0}-2 \varepsilon-1}}\right) . \tag{2.16}
\end{equation*}
$$

We can now apply Corollary 2.5 with $H:=H_{0}-\varepsilon$, which gives, for all $0<\kappa<\min \left(H_{0}-\varepsilon, \frac{1}{2}-\varepsilon\right)$,

$$
\begin{align*}
& \mathbb{I}_{n, m}(s, t) \\
& \lesssim 2^{-4 n \varepsilon}|t-s|^{2 \kappa} \int_{0}^{\infty} d r \frac{\left\{1+r^{2}\right\}^{-\alpha}}{r^{2\left(H_{1}+\cdots+H_{d}\right)-2 d+1}} \frac{1}{1+r^{4 H_{0}-4 \kappa-2-8 \varepsilon}} \\
& \lesssim 2^{-4 n \varepsilon}|t-s|^{2 \kappa}\left(\int_{0}^{1} \frac{1}{r^{2\left(H_{1}+\cdots+H_{d}\right)-2 d+1}} d r+\int_{1}^{\infty} \frac{1}{r^{2 \alpha+2\left(2 H_{0}+H_{1}+\cdots+H_{d}\right)-2 d-1-4 \kappa-8 \varepsilon}} d r\right) . \tag{2.17}
\end{align*}
$$

Due to our assumption (1.11), we can in fact pick $\varepsilon$ and $\kappa$ sufficiently small so that

$$
4 \varepsilon+2 \kappa<\alpha-\left[d+1-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right)\right]
$$

and for such a choice, the integrals involved in (2.17) are obviously finite, implying that

$$
\mathbb{I}_{n, m}(s, t) \lesssim 2^{-4 n \varepsilon}|t-s|^{2 \kappa} .
$$

Similar arguments can be used to bound $\mathbb{I}_{n, m}(s, t)$. Using the same hyperspherical change of variables as above, we can first write

$$
\begin{aligned}
\mathbb{I}_{n, m}(s, t) & \lesssim \int_{2^{n}}^{2^{m}} d r \frac{\left\{1+r^{2}\right\}^{-\alpha}}{r^{2\left(H_{1}+\cdots+H_{d}\right)-2 d+1}}\left(\int_{|\xi| \leq 2^{2 m}} d \xi \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H_{0}-1}}\right) \\
& \lesssim 2^{-4 n \varepsilon} \int_{0}^{\infty} d r \frac{\left\{1+r^{2}\right\}^{-\alpha}}{r^{2\left(H_{1}+\cdots+H_{d}\right)-4 \varepsilon-2 d+1}}\left(\int_{\mathbb{R}} d \xi \frac{\left|\gamma_{s, t}(\xi, r)\right|^{2}}{|\xi|^{2 H_{0}-1}}\right),
\end{aligned}
$$

for any $\varepsilon>0$. Now we are essentially in the same position as in (2.16), and so we can rely on the estimation strategy of (2.17) to derive that

$$
\mathbb{I}_{n, m}(s, t) \lesssim 2^{-4 n \varepsilon}|t-s|^{2 \kappa},
$$

for any $\varepsilon, \kappa>0$ small enough.
Going back to (2.14, we deduce the desired bound (2.10).

Step 2: The estimate obtained in the previous step can naturally be rephrased as

$$
\begin{equation*}
\mathbb{E}\left[\left\|\chi \varrho_{n, m}(t, .)-\chi \varrho_{n, m}(s, .)\right\|_{\mathcal{W}^{-\alpha, 2 p}}^{2 p}\right] \lesssim 2^{-4 n \varepsilon p}|t-s|^{2 \kappa p} \tag{2.18}
\end{equation*}
$$

for all $p \geq 1,1 \leq n \leq m, 0 \leq s \leq t \leq 1$ and $\varepsilon, \kappa>0$ small enough.
By choosing $p \geq 1$ large enough so that $2 \kappa p>1$, Kolmogorov continuity criterion allows us to assert that $\chi \varrho_{n, m} \in \mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, 2 p}\left(\mathbb{R}^{d}\right)\right)$ almost surely. In turn, this puts us in a position to use the classical Garsia-Rodemich-Rumsey estimate (see [24]) and deduce that almost surely, for all $p \geq 1,0<\kappa_{0}<\kappa, 0 \leq s \leq t \leq 1$, one has

$$
\left\|\chi \bigcap_{\varrho_{n, m}}(t, .)-\chi \varrho_{n, m}(s, .)\right\|_{\mathcal{W}^{-\alpha, 2 p}}^{2 p} \lesssim|t-s|^{2 \kappa_{0} p} \int_{[0,1]^{2}} \frac{\left\|\chi \ell_{n, m}(u, .)-\chi \varrho_{n, m}(v, .)\right\|_{\mathcal{W}^{-\alpha, 2 p}}^{2 p}}{|u-v|^{2 \kappa_{0} p+2}} d u d v
$$

for some proportional constant that only depends on $\kappa_{0}$ and $p$.
Picking $s=0$ and taking the supremum over $t \in[0,1]$, we derive

$$
\left\|\chi \varrho_{n, m}\right\|_{\mathcal{C}_{T} \mathcal{W}^{-\alpha, 2 p}}^{2 p} \lesssim \int_{[0,1]^{2}} \frac{\left\|\chi \varrho_{n, m}(u, .)-\chi \varrho_{n, m}(v, .)\right\|_{\mathcal{W}^{-\alpha, 2 p}}^{2 p}}{|u-v|^{2 \kappa_{0} p+2}} d u d v
$$

and therefore, using 2.18 again, we obtain that

$$
\mathbb{E}\left[\left\|\chi \varrho_{n, m}\right\|_{\mathcal{C}_{T} \mathcal{W}^{-\alpha, 2 p}}^{2 p}\right] \lesssim 2^{-4 n \varepsilon p} \int_{[0,1]^{2}} \frac{d u d v}{|u-v|^{-2\left(\kappa-\kappa_{0}\right) p+2}} \lesssim 2^{-4 n \varepsilon p}
$$

for any $p \geq 1$ large enough so that $-2\left(\kappa-\kappa_{0}\right) p+2<1$.
We can conclude that for any $p \geq 1$ large enough,

$$
\begin{equation*}
\left\|\chi_{\imath n, m}\right\|_{L^{2 p}\left(\Omega ; \mathcal{C}_{T} \mathcal{W}^{-\alpha, 2 p}\right)} \lesssim 2^{-2 n \varepsilon} \tag{2.19}
\end{equation*}
$$

In particular, $\left(\chi_{\ell_{n}}\right)_{n \geq 1}$ is a Cauchy sequence in $L^{2 p}\left(\Omega ; \mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, 2 p}\left(\mathbb{R}^{d}\right)\right)\right.$ ) (for any $p \geq 1$ large enough), which entails convergence in this space to a limit $\chi$ !. Going back to $(2.19)$, we also have

$$
\left\|\chi!-\chi \varrho_{n}\right\|_{L^{2 p}\left(\Omega ; \mathcal{C}_{T} \mathcal{W}^{-\alpha, 2 p}\right)} \lesssim 2^{-2 n \varepsilon}
$$

and from there, a standard Borell-Cantelli argument provides us with the desired almost sure convergence of $\left(\chi{ }_{n}\right)_{n \geq 1}$ to $\chi$ ! in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$, for every $2 \leq p<\infty$. Finally, the convergence in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, \infty}\left(\mathbb{R}^{d}\right)\right.$ ) follows from the Sobolev embedding $\mathcal{W}^{-\alpha+\frac{d}{p}+\eta, p}\left(\mathbb{R}^{d}\right) \subset$ $\mathcal{W}^{-\alpha, \infty}\left(\mathbb{R}^{d}\right)$, for any $\eta>0$.
2.3. Proof of Proposition 1.6. We will follow the same general scheme as in the proof of Proposition 1.2. Just as in Section 2.2, let us assume that $T=1$, and set, for all $m, n \geq 0$, $q_{\rho_{n, m}}:=q_{\rho_{m}}-q_{\rho_{n}}$.

Step 1: Our main objective here is to show that for all $p \geq 1,0 \leq n \leq m, 0 \leq s \leq t \leq 1$, $\varepsilon, \kappa>0$ small enough, and for every $\alpha$ satisfying

$$
\begin{equation*}
d+1-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right)<\alpha<\frac{1}{4} \tag{2.20}
\end{equation*}
$$

one has

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{n, m}(t, .)-\wp_{n, m}(s, .)\right]\right)\right)(x)\right|^{2 p}\right] d x \lesssim 2^{-4 n \varepsilon p}|t-s|^{\kappa p} \tag{2.21}
\end{equation*}
$$

where the proportional constant only depends on $p, \alpha$, and $\chi$.

For the sake of conciseness, we shall only prove estimate 2.21 for $n=0$, that is we will only focus on the estimate for the time-variation $\wp_{m}(t,)-.\wp_{m}(s,$.$) , with m \geq 1$. The extension of the result to all $m \geq n \geq 0$ could in fact be easily deduced from the combination of the subsequent estimates with the elementary bounding argument used in (2.15).

A first fundamental observation is that the contractivity argument used in (2.11 can be extended to the present setting. Indeed, the random variable under consideration clearly belongs to the first two chaoses generated by $W$ (with representation 2.2 of the noise in mind), and therefore, due to the hypercontractivity property holding in such a space (see e.g. [33]), we can assert that

$$
\begin{aligned}
& \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)\right)(x)\right|^{2 p}\right] \\
& \leq c_{p} \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)\right)(x)\right|^{2}\right]^{p}
\end{aligned}
$$

where the constant $c_{p}$ only depends on $p$.
Let us then write

$$
\begin{aligned}
& \mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)\right)(x) \\
& =\frac{1}{(2 \pi)^{d}} \int_{\mathbb{R}^{d}} d \lambda\left\{1+|\lambda|^{2}\right\}^{-\alpha} e^{\imath\langle x, \lambda\rangle}\left(\int_{\mathbb{R}^{d}} d \beta \widehat{\chi^{2}}(\lambda-\beta) \mathcal{F}\left(\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)(\beta)\right)
\end{aligned}
$$

which immediately yields

$$
\begin{align*}
& \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)\right)(x)\right|^{2}\right] \\
& =\frac{1}{(2 \pi)^{2 d}} \iint_{\left(\mathbb{R}^{d}\right)^{2}} \frac{d \lambda d \tilde{\lambda}}{\left\{1+|\lambda|^{2}\right\}^{\alpha}\left\{1+|\tilde{\lambda}|^{2}\right\}^{\alpha}} e^{\imath\langle x, \lambda-\tilde{\lambda}\rangle} \iint_{\left(\mathbb{R}^{d}\right)^{2}} d \beta d \tilde{\beta} \widehat{\chi^{2}}(\lambda-\beta) \overline{\widehat{\chi^{2}}(\tilde{\lambda}-\tilde{\beta})} \mathcal{Q}_{m ; s, t}^{(2)}(\beta, \tilde{\beta}), \tag{2.22}
\end{align*}
$$

with

$$
\mathcal{Q}_{m ; s, t}^{(2)}(\beta, \tilde{\beta}):=\mathbb{E}\left[\mathcal{F}\left(\left[\wp_{\infty}(t, .)-\wp_{m}(s, .)\right]\right)(\beta) \overline{\mathcal{F}\left(\left[\wp_{m}(t, .)-\wp_{\rho_{m}}(s, .)\right]\right)(\tilde{\beta})}\right]
$$

In order to expand $\mathcal{Q}_{m ; s, t}^{(2)}(\beta, \tilde{\beta})$, let us first recall that the expansion of the ("renormalized") quantities

$$
\mathbb{E}\left[\wp_{m}(t, y) \overline{\wp_{m}(s, \tilde{y})}\right]
$$

is governed by the following standard application of Wick formula (see e.g. [33]) :
Lemma 2.7. For all $m, n \geq 1, s, t \geq 0$ and $y, \tilde{y} \in \mathbb{R}^{d}$, it holds that

$$
\mathbb{E}\left[\wp_{m}(t, y) \overline{\wp_{n}(s, \tilde{y})}\right]=\left|\mathbb{E}\left[\wp_{m}(t, y) \overline{Q_{n}(s, \tilde{y})}\right]\right|^{2}+\left|\mathbb{E}\left[\wp_{m}(t, y) \Re_{n}(s, \tilde{y})\right]\right|^{2}
$$

Using Lemma 2.7 and the shortcut notation $\varrho_{m}(u, v ; z):=\varrho_{m}(v, z)-\varrho_{m}(u, z)$, we easily verify that

$$
\begin{aligned}
& \mathbb{E}\left[\left[\wp_{m}(t, y)-\wp_{m}(s, y)\right] \overline{\left[\wp_{m}(t, \tilde{y})-\wp_{\rho_{m}}(s, \tilde{y})\right]}\right]
\end{aligned}
$$

$$
\begin{aligned}
& +\left[\mathbb{E}\left[\varrho_{m}(s, t ; y) \varrho_{m}(t, \tilde{y})\right] \cdot \mathbb{E}\left[\overline{\varrho_{m}(t, y) \varrho_{m}(t, \tilde{y})}\right]+\mathbb{E}\left[\varrho_{m}(s, y) \varrho_{m}(t, \tilde{y})\right] \cdot \mathbb{E}\left[\overline{\varrho_{m}(s, t ; y) \varrho_{m}(t, \tilde{y})}\right]\right. \\
& \left.+\mathbb{E}\left[\varrho_{m}(t, s ; y) \varrho_{m}(s, \tilde{y})\right] \cdot \mathbb{E}\left[\overline{\varrho_{m}(t, y){ }^{\prime}(s, \tilde{y})}\right]+\mathbb{E}\left[\varrho_{m}(s, y) \varrho_{m}(s, \tilde{y})\right] \cdot \mathbb{E}\left[\overline{\varrho_{m}(t, s ; y) \varrho_{m}(s, \tilde{y})}\right]\right],
\end{aligned}
$$

which, combined with the covariance formulas (2.6)-2.7), allows us to expand 2.22 as

$$
\begin{aligned}
& \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)\right)(x)\right|^{2}\right] \\
& =c \int_{(\xi, \eta) \in D_{m}} d \xi d \eta \int_{(\tilde{\xi}, \tilde{\eta}) \in D_{m}} d \tilde{\xi} d \tilde{\eta} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}} \frac{1}{|\tilde{\xi}|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta}_{i}\right|^{2 H_{i}-1}} \\
& \quad\left(\sum_{\ell=1}^{8} \Gamma_{m ; s, t}^{\ell}(\xi,|\eta| ; \tilde{\xi},|\tilde{\eta}|)\right) \\
& \quad \iint_{\left(\mathbb{R}^{d}\right)^{2}} \frac{d \lambda d d \tilde{\lambda}}{\left\{1+|\lambda|^{2}\right\}^{\alpha}\left\{1+|\tilde{\lambda}|^{2}\right\}^{\alpha}} e^{\imath\langle x, \lambda-\tilde{\lambda}\rangle} \widehat{\chi^{2}}(\lambda-(\eta-\tilde{\eta})) \widehat{\chi^{2}}(\tilde{\lambda}-(\eta-\tilde{\eta})),
\end{aligned}
$$

with $\Gamma_{m ; s, t}^{\ell}=\Gamma_{m ; s, t}^{\ell}(\xi,|\eta| ; \tilde{\xi},|\tilde{\eta}|)$ given by

$$
\begin{array}{rll}
\Gamma_{m ; s, t}^{1}:=\gamma_{s, t}(\xi,|\eta|) \overline{\gamma_{t}(\xi,|\eta|)}\left|\gamma_{t}(\tilde{\xi},|\tilde{\eta}|)\right|^{2}, & \Gamma_{m ; s, t}^{2}:=\gamma_{s}(\xi,|\eta|) \overline{\gamma_{t}(\xi,|\eta|)} \overline{\gamma_{s, t}(\tilde{\xi},|\tilde{\eta}|)} \gamma_{t}(\tilde{\xi},|\tilde{\eta}|), \\
\Gamma_{m ; s, t}^{3}:=\gamma_{t, s}(\xi,|\eta|) \overline{\gamma_{s}(\xi,|\eta|)} \overline{\gamma_{t}(\tilde{\xi},|\tilde{\eta}|)} \gamma_{s}(\tilde{\xi},|\tilde{\eta}|), & \Gamma_{m ; s, t}^{4}:=\left|\gamma_{s}(\xi,|\eta|)\right|^{2} \overline{\gamma_{t, s}(\tilde{\xi},|\tilde{\eta}|)} \gamma_{s}(\tilde{\xi},|\tilde{\eta}|), \\
\Gamma_{m ; s, t}^{5}:=\gamma_{s, t}(\xi,|\eta|) \gamma_{t}(-\xi,|\eta|) \overline{\gamma_{t}(\tilde{\xi},|\tilde{\eta}|) \gamma_{t}(-\tilde{\xi},|\tilde{\eta}|)}, & \Gamma_{m ; s, t}^{6}:=\gamma_{s}(\xi,|\eta|) \gamma_{t}(-\xi,|\eta|) \overline{\gamma_{s, t}(\tilde{\xi},|\tilde{\eta}|) \gamma_{t}(-\tilde{\xi},|\tilde{\eta}|)} \\
\Gamma_{m ; s, t}^{7}:=\gamma_{t, s}(\xi,|\eta|) \gamma_{s}(-\xi,|\eta|) \overline{\gamma_{t}(\tilde{\xi},|\tilde{\eta}|) \gamma_{s}(-\tilde{\xi},|\tilde{\eta}|)}, & \Gamma_{m ; s, t}^{8}:=\gamma_{s}(\xi,|\eta|) \gamma_{s}(-\xi,|\eta|) \overline{\gamma_{t, s}(\tilde{\xi},|\tilde{\eta}|) \gamma_{s}(-\tilde{\xi},|\tilde{\eta}|)} .
\end{array}
$$

At this point, we can rely on the technical Lemma 2.6 to assert that

$$
\begin{equation*}
\int_{\mathbb{R}^{d}} d x \mathbb{E}\left[\left|\mathcal{F}^{-1}\left(\left\{1+|.|^{2}\right\}^{-\alpha} \mathcal{F}\left(\chi^{2}\left[\wp_{m}(t, .)-\wp_{m}(s, .)\right]\right)\right)(x)\right|^{2}\right]^{p} \lesssim\left(\sum_{\ell=1}^{8} \mathcal{J}_{m ; s, t}^{\ell}\right)^{p} \tag{2.23}
\end{equation*}
$$

where

$$
\begin{array}{r}
\mathcal{J}_{m ; s, t}^{\ell}:=\int_{(\xi, \eta) \in D_{m}} d \xi d \eta \int_{(\tilde{\xi}, \tilde{\eta}) \in D_{m}} d \tilde{\xi} d \tilde{\eta} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}} \frac{1}{|\tilde{\xi}|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta}_{i}\right|^{2 H_{i}-1}} \\
\left|\Gamma_{m ; s, t}^{\ell}(\xi,|\eta| ; \tilde{\xi},|\tilde{\eta}|)\right|\left\{1+|\eta-\tilde{\eta}|^{2}\right\}^{-2 \alpha} \tag{2.24}
\end{array}
$$

Using the trivial bound $|\eta-\tilde{\eta}| \geq||\eta|-|\tilde{\eta}||$ and going back to the expression of $\Gamma_{m ; s, t}^{\ell}(\xi,|\eta| ; \tilde{\xi},|\tilde{\eta}|)$, we can first bound these integrals along the pattern

$$
\begin{align*}
\mathcal{J}_{m ; s, t}^{\ell} \lesssim \int_{\left(\mathbb{R} \times \mathbb{R}^{d}\right)^{2}} \frac{d \xi d \eta d \tilde{\xi} d \tilde{\eta}}{\left\{1+\left||\eta|-|\tilde{\eta}|^{2}\right\}^{2 \alpha}\right.}\left(\frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{a_{1}}(\xi,|\eta|)\right|\left|\gamma_{a_{2}}( \pm \xi,|\eta|)\right|\right) \\
\quad\left(\frac{1}{|\tilde{\xi}|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta}_{i}\right|^{2 H_{i}-1}}\left|\gamma_{a_{3}}(\tilde{\xi},|\tilde{\eta}|)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\tilde{\eta}|)\right|\right) \tag{2.25}
\end{align*}
$$

for some $a_{1}, a_{2}, a_{3}, a_{4} \in\{s, t,\{s, t\}\}$ (depending on $\ell$ ) such that exactly one of the $a_{k}$ 's is $\{s, t\}$.

Now, let us decompose the latter integration domain into $\left(\mathbb{R} \times \mathbb{R}^{d}\right)^{2}:=\mathfrak{D}_{1} \cup \mathfrak{D}_{2}$, where

$$
\mathfrak{D}_{1}:=\left\{(\xi, \eta, \tilde{\xi}, \tilde{\eta}): 0 \leq|\tilde{\eta}| \leq \frac{|\eta|}{2} \text { or }|\tilde{\eta}| \geq \frac{3|\eta|}{2}\right\} \text { and } \mathfrak{D}_{2}:=\left\{(\xi, \eta, \tilde{\xi}, \tilde{\eta}): \frac{|\eta|}{2}<|\tilde{\eta}|<\frac{3|\eta|}{2}\right\}
$$

For the integral over $\mathfrak{D}_{1}$, we can rely on the inequality $\| \eta|-|\tilde{\eta}|| \geq \max \left(\frac{|\eta|}{2}, \frac{|\tilde{\eta}|}{3}\right)$ (valid for all $\left.(\xi, \eta, \tilde{\xi}, \tilde{\eta}) \in \mathfrak{D}_{1}\right)$ to write

$$
\begin{aligned}
& \mathcal{A}_{1}^{\ell}:=\int_{\mathfrak{D}_{1}} \frac{d \xi d \eta d \tilde{\xi} d \tilde{\eta}}{\left\{1+\left||\eta|-|\tilde{\eta}|^{2}\right\}^{2 \alpha}\right.}\left(\frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{a_{1}}(\xi,|\eta|)\right|\left|\gamma_{a_{2}}( \pm \xi,|\eta|)\right|\right) . \\
& \quad\left(\frac{1}{|\tilde{\xi}|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta}_{i}\right|^{2 H_{i}-1}}\left|\gamma_{a_{3}}(\tilde{\xi},|\tilde{\eta}|)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\tilde{\eta}|)\right|\right) \\
& \lesssim\left(\int_{\mathbb{R} \times \mathbb{R}^{d}} \frac{d \xi d \eta}{\left\{1+|\eta|^{2}\right\}^{\alpha}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{a_{1}}(\xi,|\eta|)\right|\left|\gamma_{a_{2}}( \pm \xi,|\eta|)\right|\right) . \\
& \quad\left(\int_{\mathbb{R} \times \mathbb{R}^{d}} \frac{d \tilde{\xi} d \tilde{\eta}}{\left\{1+|\tilde{\eta}|^{2}\right\}^{\alpha}} \frac{1}{|\tilde{\xi}|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta_{i}}\right|^{2 H_{i}-1}}\left|\gamma_{a_{3}}(\tilde{\xi},|\tilde{\eta}|)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\tilde{\eta}|)\right|\right) \\
& \lesssim \prod_{k=1}^{4}\left(\int_{\mathbb{R} \times \mathbb{R}^{d}} \frac{d \xi d \eta}{\left\{1+|\eta|^{2}\right\}^{\alpha}} \frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left.\left.\left|{\left|\eta_{i}\right|^{2 H_{i}-1}}\right| \gamma_{a_{k}}(\xi,|\eta|)\right|^{2}\right)^{\frac{1}{2}},}\right.
\end{aligned}
$$

where we have merely used Cauchy-Schwarz inequality to derive the last estimate.
Observe that we are here dealing with the same integrals as in the proof of Proposition 1.2 (see in particular (2.15) , with $\alpha$ satisfying the same assumption (1.11). Therefore we can mimic the arguments in (2.16)-2.17) to deduce the desired estimate, namely: for every $\ell=1, \ldots, 8$ and any $\kappa>0$ small enough,

$$
\mathcal{A}_{1}^{\ell} \lesssim|t-s|^{\kappa},
$$

where the $|t-s|^{\kappa}$-increment stems from the fact that one of the $a_{k}$ 's must be $\{s, t\}$.
Let us turn to the integral over $\mathfrak{D}_{2}$ and lean first on a hyperspherical change of variable (with respect to $\tilde{\eta}$ ) to write

$$
\begin{aligned}
& \int_{\frac{|\eta|}{2}<|\tilde{\eta}|<\frac{3|\eta|}{2}} \frac{d \tilde{\eta}}{\left\{1+||\eta|-|\tilde{\eta}||^{2}\right\}^{2 \alpha}}\left|\gamma_{a_{3}}(\tilde{\xi},|\tilde{\eta}|)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\tilde{\eta}|)\right| \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta}_{i}\right|^{2 H_{i}-1}} \\
& =|\eta|^{-2\left(H_{1}+\cdots+H_{d}\right)+2 d} \int_{\frac{1}{2}<|\tilde{\eta}|<\frac{3}{2}} \frac{d \tilde{\eta}}{\left\{1+|\eta|^{2}(1-|\tilde{\eta}|)^{2}\right\}^{2 \alpha}}\left|\gamma_{a_{3}}(\tilde{\xi},|\eta||\tilde{\eta}|)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\eta||\tilde{\eta}|)\right| \prod_{i=1}^{d} \frac{1}{\mid \tilde{\left.\eta_{i}\right|^{2 H_{i}-1}}} \\
& \lesssim|\eta|^{-2\left(H_{1}+\cdots+H_{d}\right)+2 d} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{\left\{1+|\eta|^{2}(1-r)^{2}\right\}^{2 \alpha}}\left|\gamma_{a_{3}}(\tilde{\xi},|\eta| r)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\eta| r)\right| .
\end{aligned}
$$

As a consequence,

$$
\left.\begin{array}{rl}
\mathcal{A}_{2}^{\ell} & :=\int_{\mathfrak{D}_{2}} \frac{d \xi d \eta d \tilde{\xi} d \tilde{\eta}}{\left\{1+\left||\eta|-|\tilde{\eta}|^{2}\right\}^{2 \alpha}\right.}\left(\frac{1}{|\xi|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{a_{1}}(\xi,|\eta|)\right|\left|\gamma_{a_{2}}( \pm \xi,|\eta|)\right|\right) . \\
& \left(\frac{1}{|\tilde{\xi}|^{2 H_{0}-1}} \prod_{i=1}^{d} \frac{1}{\left|\tilde{\eta_{i}}\right|^{2 H_{i}-1}}\left|\gamma_{a_{3}}(\tilde{\xi},|\tilde{\eta}|)\right|\left|\gamma_{a_{4}}( \pm \tilde{\xi},|\tilde{\eta}|)\right|\right) \\
\lesssim & \int_{\mathbb{R}^{d}} \frac{d \eta}{|\eta|^{2\left(H_{1}+\cdots+H_{d}\right)-2 d}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{\left\{1+|\eta|^{2}(1-r)^{2}\right\}^{2 \alpha}} \cdot \\
& \quad\left(\int_{\mathbb{R}} d \xi \frac{\left|\gamma_{a_{1}}(\xi,|\eta|)\right|\left|\gamma_{a_{2}}( \pm \xi,|\eta|)\right|}{|\xi|^{2 H_{0}-1}}\right) \cdot\left(\int_{\mathbb{R}} d \tilde{\xi} \frac{\tilde{\xi} \gamma_{a_{3}}(\tilde{\xi},|\eta| r)| | \gamma_{a_{4}}( \pm \tilde{\xi},|\eta| r) \mid}{|\tilde{\xi}|^{2 H_{0}-1}}\right) \\
& \int_{\mathbb{R}^{d}} \frac{d \eta}{|\eta|^{2\left(H_{1}+\cdots+H_{d}\right)-2 d}} \prod_{i=1}^{d} \frac{1}{\left|\eta_{i}\right|^{2 H_{i}-1}} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{\left\{1+|\eta|^{2}(1-r)^{2}\right\}^{2 \alpha}}
\end{array}\right] .
$$

Using again a hyperspherical change of variable (with respect to $\eta$ ), we obtain

$$
\begin{aligned}
& \mathcal{A}_{2}^{\ell} \lesssim \int_{0}^{\infty} \frac{d \rho}{\rho^{4\left(H_{1}+\cdots+H_{d}\right)-4 d+1}} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{\left\{1+\rho^{2}(1-r)^{2}\right\}^{2 \alpha}} \\
& \quad\left(\int_{\mathbb{R}} d \xi \frac{\left|\gamma_{a_{1}}(\xi, \rho)\right|^{2}}{|\xi|^{2 H_{0}-1}}\right)^{\frac{1}{2}} \cdot\left(\int_{\mathbb{R}} d \xi \frac{\left|\gamma_{a_{2}}(\xi, \rho)\right|^{2}}{|\xi|^{2 H_{0}-1}}\right)^{\frac{1}{2}} \cdot\left(\int_{\mathbb{R}} d \tilde{\xi} \frac{\left|\gamma_{a_{3}}(\tilde{\xi}, \rho r)\right|^{2}}{|\tilde{\xi}|^{2 H_{0}-1}}\right)^{\frac{1}{2}} \cdot\left(\int_{\mathbb{R}} d \tilde{\xi} \frac{\left.\tilde{\xi} \gamma_{a_{4}}(\tilde{\xi}, \rho r)\right|^{2}}{|\tilde{\xi}|^{2 H_{0}-1}}\right)^{\frac{1}{2}} .
\end{aligned}
$$

We can now appeal to Corollary 2.5, together with the fact that one of the $a_{k}$ 's is $\{s, t\}$, to assert that

$$
\begin{aligned}
\mathcal{A}_{2}^{\ell} \lesssim|t-s|^{\kappa}[ & \int_{0}^{1} \frac{d \rho}{\rho^{4\left(H_{1}+\cdots+H_{d}\right)-4 d+1}} \\
& \left.+\int_{1}^{\infty} \frac{d \rho}{\rho^{4\left(2 H_{0}+H_{1}+\cdots+H_{d}\right)-4 d-3-8 \varepsilon-8 \kappa}} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{\left\{1+\rho^{2}(1-r)^{2}\right\}^{2 \alpha}}\right]
\end{aligned}
$$

for all $0<\varepsilon<\frac{1}{2}$ and $0<\kappa<\min \left(H_{0}, \frac{1}{2}-\varepsilon\right)$.
At this point, observe that

$$
\begin{align*}
& \int_{1}^{\infty} \frac{d \rho}{\rho^{4\left(2 H_{0}+H_{1}+\cdots+H_{d}\right)-4 d-3-8 \varepsilon-8 \kappa}} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{\left\{1+\rho^{2}(1-r)^{2}\right\}^{2 \alpha}} \\
& \leq \int_{1}^{\infty} \frac{d \rho}{\rho^{4 \alpha+4\left(2 H_{0}+H_{1}+\cdots+H_{d}\right)-4 d-3-8 \varepsilon-8 \kappa}} \int_{\frac{1}{2}<r<\frac{3}{2}} \frac{d r}{(1-r)^{4 \alpha}} \tag{2.26}
\end{align*}
$$

Thanks to our assumption (2.20) on $\alpha$, we know on the one hand that $4 \alpha<1$, and on the other hand we can pick $\varepsilon, \kappa>0$ such that

$$
4 \alpha+4\left(2 H_{0}+H_{1}+\cdots+H_{d}\right)-4 d-3-8 \varepsilon-8 \kappa>1 .
$$

For such a choice of $\varepsilon, \kappa$, the integrals in the right-hand side of (2.26) are clearly finite, and therefore we have shown that for every $\ell=1, \ldots, 8$,

$$
\mathcal{A}_{2}^{\ell} \lesssim|t-s|^{\kappa}
$$

Going back to (2.25), we have thus proved that for every $\ell=1, \ldots, 8$, and uniformly over $m$,

$$
\mathcal{J}_{m ; s, t}^{\ell} \lesssim|t-s|^{\kappa} .
$$

Injecting the above estimates into 2.23 provides us with the desired bound 2.21 .
Step 2: Conclusion. Let $\alpha$ satisfying the condition in the statement of Proposition 1.6, that is $\alpha>d+1-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right)$.

If in addition one has $\alpha<\frac{1}{4}$, then condition 2.20 is satisfied, and so the moment estimate (2.21) holds true. Starting from this estimate, we can use the same arguments as in Step 2 of Section 2.2 to obtain that $\left(\chi^{2} \wp_{\circ_{n}}\right)_{n \geq 1}$ converges almost surely in the space $\mathcal{C}\left([0, T] ; \mathcal{W}^{-2 \alpha, p}\left(\mathbb{R}^{d}\right)\right)$, for every $2 \leq p<\infty$.

If $\alpha \geq \frac{1}{4}$, observe that due to assumption $\overline{\mathbf{H 2}}$, we can pick $\alpha^{\prime}$ satisfying $\alpha^{\prime}<\alpha$ and $d+1-\left(2 H_{0}+\sum_{i=1}^{d} H_{i}\right)<\alpha^{\prime}<\frac{1}{4}$ (that is, $\alpha^{\prime}$ satisfies 2.20). By repeating the above arguments, we deduce that the sequence $\left(\chi^{2} \wp_{\rho}\right)_{n \geq 1}$ converges almost surely in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-2 \alpha^{\prime}, p}\left(\mathbb{R}^{d}\right)\right)$, and therefore it converges almost surely in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-2 \alpha, p}\left(\mathbb{R}^{d}\right)\right)$ as well, for $2 \leq p<\infty$.

Finally, the (a.s.) convergence in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-2 \alpha, \infty}\left(\mathbb{R}^{d}\right)\right)$ can be easily derived from the Sobolev embedding $\mathcal{W}^{-2 \alpha+\frac{d}{p}+\eta, p}\left(\mathbb{R}^{d}\right) \subset \mathcal{W}^{-2 \alpha, \infty}\left(\mathbb{R}^{d}\right)$, for any $\eta>0$, which completes the proof of Proposition 1.6
2.4. Proof of Proposition 1.5. Fix $d \geq 1$ and $\left(H_{0}, \ldots, H_{d}\right) \in(0,1)^{d+1}$ such that

$$
d+\frac{3}{4}<2 H_{0}+\sum_{i=1}^{d} H_{i} \leq d+1
$$

Using (2.6), the quantity under consideration can be written as

$$
\begin{aligned}
\sigma_{n}(t, x)=\mathbb{E}\left[\left|冃_{n}(t, x)\right|^{2}\right] & =c_{H}^{2} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{2 H_{0}-1}} \int_{|\eta| \leq 2^{n}} \prod_{i=1}^{d} \frac{d \eta_{i}}{\left|\eta_{i}\right|^{2 H_{i}-1}}\left|\gamma_{t}(\xi,|\eta|)\right|^{2} \\
& =C_{H} \int_{0}^{2^{n}} \frac{d r}{r^{2\left(H_{1}+\cdots+H_{d}\right)-2 d+1}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{2 H_{0}-1}}\left|\gamma_{t}(\xi, r)\right|^{2},
\end{aligned}
$$

for some constant $C_{H}$, and where the last identity is derived from a hyperspherical change of variables. The above formula shows in particular that $\sigma_{n}$ does not depend on $x$, as stated in Proposition 1.5. Regarding the desired estimate 1.14, it is now a consequence of the following technical result (applied with $\alpha:=2 H_{0} \in(0,2)$ and $\kappa:=d+1-\left[2 H_{0}+\sum_{i=1}^{d} H_{i}\right] \geq 0$ ):

Proposition 2.8. For all $\alpha \in(0,2)$ and $\kappa \geq 0$ verifying $\alpha+\kappa>1$, one has

$$
\int_{0}^{2^{n}} \frac{d r}{r^{-2 \alpha-2 \kappa+3}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{\alpha-1}}\left|\gamma_{t}(\xi, r)\right|^{2} \underset{n \rightarrow \infty}{\sim} \begin{cases}\frac{\pi}{\kappa} 4^{n \kappa} t & \text { if } \kappa>0 \\ \pi \ln (4) \cdot n t & \text { if } \kappa=0\end{cases}
$$

Proof. It is easy to check that

$$
\left|\gamma_{t}(\xi, r)\right|^{2}=2 \frac{1-\cos \left(t\left(\xi-r^{2}\right)\right)}{\left|r^{2}-\xi\right|^{2}}
$$

which allows us to write the integral under consideration as

$$
\begin{align*}
\int_{0}^{2^{n}} \frac{d r}{r^{-2 \alpha-2 \kappa+3}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{\alpha-1}}\left|\gamma_{t}(\xi, r)\right|^{2} & =\int_{0}^{4^{n}} \frac{d r}{r^{-\alpha-\kappa+2}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{1-\cos (t(\xi-r))}{(\xi-r)^{2}} \\
& =4^{n \kappa} \int_{0}^{1} \frac{d r}{r^{-\alpha-\kappa+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{1-\cos \left(4^{n} t(\xi-r)\right)}{4^{n}(\xi-r)^{2}} \\
& =t \cdot 4^{n \kappa} \int_{0}^{1} \frac{d r}{r^{-\alpha-\kappa+2}} \int_{|\xi| \leq 1} \frac{\left.d \xi \sin ^{2} 4^{4^{n} t}(\xi-r)\right)}{|\xi|^{\alpha-1}} \frac{4^{n} t}{2}(\xi-r)^{2} \tag{2.27}
\end{align*} .
$$

First case: $\kappa>0$.
Let us set $\Phi(x):=\frac{\sin ^{2}(x)}{\pi x^{2}}$ for every $x \in \mathbb{R}$. This function is positive, even and satisfies $\int_{\mathbb{R}} \Phi(x) d x=1$, so that the sequence

$$
\Phi_{n}(x):=\frac{4^{n} t}{2} \Phi\left(\frac{4^{n} t}{2} x\right)
$$

defines an approximation of the identity. Let us now rewrite the quantity in (2.27) as

$$
\begin{aligned}
& t 4^{n \kappa} \int_{0}^{1} \frac{d r}{r^{-\alpha-\kappa+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{\sin ^{2}\left(\frac{4^{n} t}{2}(\xi-r)\right)}{\frac{4^{n} t}{2}(\xi-r)^{2}} \\
& =\pi t 4^{n \kappa} \int_{0}^{1} \frac{d r}{r^{-\alpha-\kappa+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \Phi_{n}(\xi-r)=\pi t 4^{n \kappa} \int_{\mathbb{R}} h(r)\left(\Phi_{n} * g\right)(r) d r
\end{aligned}
$$

where we have set

$$
h(r):=\mathbb{1}_{[0,1]}(r) \frac{1}{r^{-\alpha-\kappa+2}} \quad \text { and } \quad g(\xi):=\mathbb{1}_{[-1,1]}(\xi) \frac{1}{|\xi|^{\alpha-1}} .
$$

Due to our assumptions on $\kappa, \alpha$ (i.e., $\kappa>0, \alpha \in(0,2)$ and $\alpha+\kappa>1$ ), we can exhibit a pair of conjugate exponents $p, p^{\prime}>1$ such that

$$
g \in L^{p}(\mathbb{R}) \quad \text { and } \quad h \in L^{p^{\prime}}(\mathbb{R}) .
$$

Indeed, it is easy to check that any $p$ such that

$$
\max (0, \alpha-1)<\frac{1}{p}<\min (1, \alpha+\kappa-1)
$$

meets the above requirements. Let us fix such a pair $\left(p, p^{\prime}\right)$, and now write

$$
\begin{aligned}
\left|\int_{\mathbb{R}} h(r)\left(\Phi_{n} * g\right)(r) d r-\int_{\mathbb{R}} h(r) g(r) d r\right| & =\left|\int_{\mathbb{R}} h(r)\left[\left(\Phi_{n} * g\right)(r)-g(r)\right] d r\right| \\
& \leq\left(\int_{\mathbb{R}}|h(r)|^{p^{\prime}} d r\right)^{\frac{1}{p^{\prime}}}\left(\int_{\mathbb{R}}\left|\left(\Phi_{n} * g\right)(r)-g(r)\right|^{p} d r\right)^{\frac{1}{p}}
\end{aligned}
$$

which tends to 0 as $n$ tends to infinity. Since $\int_{\mathbb{R}} h(r) g(r) d r=\frac{1}{\kappa}$, we obtain the desired conclusion

$$
\int_{0}^{2^{n}} \frac{d r}{r^{-2 \alpha-2 \kappa+3}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{\alpha-1}}\left|\gamma_{t}(\xi, r)\right|^{2} \underset{n \rightarrow \infty}{\sim} \frac{\pi t}{\kappa} 4^{n \kappa} .
$$

Second case: $\kappa=0$.

Note first that, due to our assumptions, one has here $1<\alpha<2$. Besides, let us recall (see 2.27) that

$$
\begin{equation*}
\int_{0}^{2^{n}} \frac{d r}{r^{-2 \alpha+3}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{\alpha-1}}\left|\gamma_{t}(\xi, r)\right|^{2}=t \int_{0}^{1} \frac{d r}{r^{-\alpha+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{\sin ^{2}\left(\frac{4^{n} t}{2}(\xi-r)\right)}{\frac{4^{n} t}{2}(\xi-r)^{2}} \tag{2.28}
\end{equation*}
$$

Consider the function $f$ defined for all $T>0$ by

$$
f(T):=\int_{0}^{1} \frac{d r}{r^{-\alpha+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{\sin ^{2}(T(\xi-r))}{(\xi-r)^{2}}
$$

One has of course

$$
\begin{aligned}
f^{\prime}(T) & =\int_{0}^{1} \frac{d r}{r^{-\alpha+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{2 \sin (T(\xi-r)) \cos (T(\xi-r))}{(\xi-r)} \\
& =\int_{0}^{1} \frac{d r}{r^{-\alpha+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \frac{\sin (2 T(\xi-r))}{(\xi-r)}
\end{aligned}
$$

and

$$
\begin{aligned}
f^{\prime \prime}(T) & =2 \int_{0}^{1} \frac{d r}{r^{-\alpha+2}} \int_{|\xi| \leq 1} \frac{d \xi}{|\xi|^{\alpha-1}} \cos (2 T(\xi-r)) \\
& =2 \int_{0}^{1} \frac{d r}{r^{-\alpha+2}} \int_{0}^{1} \frac{d \xi}{|\xi|^{\alpha-1}}\{\cos (2 T(\xi+r))+\cos (2 T(\xi-r))\} \\
& =4\left(\int_{0}^{1} d r \frac{\cos (2 T r)}{r^{-\alpha+2}}\right)\left(\int_{0}^{1} d \xi \frac{\cos (2 T \xi)}{|\xi|^{\alpha-1}}\right)=\frac{2}{T}\left(\int_{0}^{2 T} d r \frac{\cos (r)}{r^{-\alpha+2}}\right)\left(\int_{0}^{2 T} d \xi \frac{\cos (\xi)}{|\xi|^{\alpha-1}}\right) .
\end{aligned}
$$

At this point, let us recall the following standard Fourier transform: for $\beta \in(0,1)$ and $\xi \neq 0$,

$$
\int_{\mathbb{R}} d x \frac{e^{-i x \xi}}{|x|^{\beta}}=\frac{2 \sin \left(\frac{\pi \beta}{2}\right) \Gamma(1-\beta)}{|\xi|^{1-\beta}}
$$

which immediately yields

$$
\int_{0}^{\infty} d x \frac{\cos (x)}{|x|^{\beta}}=\sin \left(\frac{\pi \beta}{2}\right) \Gamma(1-\beta)
$$

and therefore

$$
f^{\prime \prime}(T) \sim \frac{2}{T} \sin \left(\frac{\pi}{2}(2-\alpha)\right) \sin \left(\frac{\pi}{2}(\alpha-1)\right) \Gamma(\alpha-1) \Gamma(2-\alpha)
$$

Using Euler's reflection formula

$$
\Gamma(\beta) \Gamma(1-\beta)=\frac{\pi}{\sin (\pi \beta)} \quad(\beta \in(0,1))
$$

the above quantity reduces to

$$
\begin{aligned}
& \sin \left(\frac{\pi}{2}(2-\alpha)\right) \sin \left(\frac{\pi}{2}(\alpha-1)\right) \Gamma(\alpha-1) \Gamma(2-\alpha) \\
& =\pi \frac{\sin \left(\frac{\pi}{2}(2-\alpha)\right) \sin \left(\frac{\pi}{2}(\alpha-1)\right)}{\sin (\pi(\alpha-1))}=\frac{\pi}{2} \frac{\sin \left(\frac{\pi}{2}(2-\alpha)\right) \sin \left(\frac{\pi}{2}(\alpha-1)\right)}{\sin \left(\frac{\pi}{2}(\alpha-1)\right) \cos \left(\frac{\pi}{2}(\alpha-1)\right)}=\frac{\pi}{2}
\end{aligned}
$$

and so

$$
f^{\prime \prime}(T) \underset{T \rightarrow \infty}{\sim} \frac{\pi}{T}
$$

We can finally use a standard comparison argument (see Lemma 2.9 below) to successively assert that $f^{\prime}(T) \underset{T \rightarrow \infty}{\sim} \pi \ln (T)$ and

$$
f(T) \underset{T \rightarrow \infty}{\sim} \pi(T \ln (T)-T) \underset{T \rightarrow \infty}{\sim} \pi T \ln (T) .
$$

Going back to (2.28), we get the desired conclusion, namely

$$
\int_{0}^{2^{n}} \frac{d r}{r^{-2 \alpha+3}} \int_{|\xi| \leq 4^{n}} \frac{d \xi}{|\xi|^{\alpha-1}}\left|\gamma_{t}(\xi, r)\right|^{2} \underset{n \rightarrow \infty}{\sim} \pi t n \ln (4) .
$$

Lemma 2.9. Fix $a \in \mathbb{R}$ and let $g:[a,+\infty[\rightarrow \mathbb{R}, h:[a,+\infty[\rightarrow(0, \infty)$, be two continuous functions. If $g(t) \underset{t \rightarrow \infty}{\sim} h(t)$ and $\int_{a}^{+\infty} h(t) d t=\infty$, then

$$
\int_{a}^{T} g(t) d t \underset{T \rightarrow \infty}{\sim} \int_{a}^{T} h(t) d t
$$

### 2.5. Global definition of the linear solution.

At first reading, the result of Proposition 1.2 only provides us with a local definition of the process $?$. In other words, what is actually given by the statement is the set of the limit elements $\left\{\chi \ell, \chi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)\right\}$. For the sake of rigor, let us say a few words about how those elements can be glued together into a single process $?$, which we can then inject into the transformation $v:=u-$ ? of Definition 1.4 or Definition 1.8

To this end, fix $p \geq 2$ and $\alpha$ satisfying (1.11). Let us denote by $\mathcal{P}$ the set of sequences $\sigma=\left(\sigma_{k}\right)_{k \geq 1}$ such that for each $k \geq 1, \sigma_{k}: \mathbb{R}^{d} \rightarrow \mathbb{R}$ is a smooth function satisfying

$$
\sigma_{k}(x)= \begin{cases}1 & \text { if }\|x\| \leq k, \\ 0 & \text { if }\|x\| \geq k+1 .\end{cases}
$$

Given such a sequence $\sigma$, and for each fixed $k \geq 1$, let us denote by ${ }^{\circ}\left(\sigma_{k}\right)$ the limit of the sequence $\left(\sigma_{k} i_{n}\right)_{n \geq 1}$ in the space $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$, as provided by Proposition 1.2 . We know in particular that $\varrho^{\left(\sigma_{k}\right)}$ is defined on a probability space $\Omega^{\left(\sigma_{k}\right)}$ of full measure 1. Let us set $\Omega^{(\sigma)}:=\cap_{k \geq 1} \Omega^{\left(\sigma_{k}\right)}$, and note that this space is still of measure 1 .

For every fixed time $t \in[0, T]$, we now define the random distribution

$$
\varrho^{(\sigma)}(t): \Omega^{(\sigma)} \rightarrow \mathcal{D}^{\prime}\left(\mathbb{R}^{d}\right)
$$

as follows: for every test function $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ such that $\operatorname{supp}(\varphi) \subset B(0, k)$ (for some $k \geq 1$ ),

$$
\left\langle\varrho^{(\sigma)}(t), \varphi\right\rangle:=\left\langle\varrho^{\left(\sigma_{k}\right)}(t), \varphi\right\rangle .
$$

 for every test function $\varphi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ with $\operatorname{supp}(\varphi) \subset B(0, k) \subset B(0, \ell)$, one has

$$
\left\langle q^{\left(\sigma_{k}\right)}(t), \varphi\right\rangle=\left\langle q^{\left(\sigma_{\ell}\right)}(t), \varphi\right\rangle \quad \text { on } \quad \Omega^{(\sigma)} .
$$

(ii) For any test function $\chi: \mathbb{R}^{d} \rightarrow \mathbb{R}$, one has, on $\Omega^{(\sigma)}$,

$$
\chi \cdot \varrho_{n} \underset{n \rightarrow \infty}{\rightarrow} \chi \cdot \wp^{(\sigma)} \quad \text { in } \mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)
$$

(iii) If $\sigma, \gamma \in \mathcal{P}$, it holds that

$$
q^{(\sigma)}=q^{(\gamma)} \quad \text { on } \quad \Omega^{(\sigma)} \cap \Omega^{(\gamma)} \text {. }
$$

Due to the latter identification property, we simply set $\uparrow:=\varrho^{(\sigma)}$, for some fixed element $\sigma \in \mathcal{P}$.

Proof. (i) We first show that for $1 \leq k \leq \ell$

$$
\begin{equation*}
\mathrm{q}^{\left(\sigma_{k}\right)}=\sigma_{k} \vartheta^{\left(\sigma_{\ell}\right)} \quad \text { on } \Omega^{\left(\sigma_{k}\right)} \cap \Omega^{\left(\sigma_{\ell}\right)} . \tag{2.29}
\end{equation*}
$$

By Proposition 1.2 , the sequence $\sigma_{k} \sigma_{\ell} \cap_{n}=\sigma_{k} \varrho_{n}$ converges almost surely to $\varrho^{\left(\sigma_{k}\right)}$ in the space $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$, on $\Omega^{\left(\sigma_{k}\right)}$. But by continuity of the product in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$ by the test function $\sigma_{k}$, we also have that $\sigma_{k} \sigma_{\ell, n}$ converges almost surely to $\sigma_{k}{ }^{\rho}{ }^{\left(\sigma_{\ell}\right)}$ in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$ on $\Omega^{\left(\sigma_{\ell}\right)}$, and we deduce (2.29) by uniqueness of the limit. Therefore by 2.29 , for all $t \in[0, T]$ and $\varphi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\operatorname{supp}(\varphi) \subset B(0, k)$
(ii) Let $\chi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)$, then there exists $k \geq 1$ such that $\operatorname{supp}(\chi) \subset B(0, k)$. According to Proposition 1.2, $\chi \ell_{n}=\chi \sigma_{k} i_{n}$ converges almost surely to $\chi \varrho^{\mathrm{q}\left(\sigma_{k}\right)}$ in $\mathcal{C}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right)$ on $\Omega^{\left(\sigma_{k}\right)}$. But $\chi \varrho^{\left(\sigma_{k}\right)}=\chi!^{\mathrm{\rho}}{ }^{(\sigma)}$ on $\Omega^{\left(\sigma_{k}\right)}$. Indeed, for all $t \in[0, T]$, if $\varphi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)$,

$$
\left\langle\chi \varrho^{(\sigma)}(t), \varphi\right\rangle=\left\langle\varrho^{(\sigma)}(t), \chi \varphi\right\rangle=\left\langle\imath^{\left(\sigma_{k}\right)}(t), \chi \varphi\right\rangle=\left\langle\chi \varrho^{\left(\sigma_{k}\right)}(t), \varphi\right\rangle,
$$

where we have used that $\operatorname{supp}(\chi \varphi) \subset B(0, k)$ to derive the second equality.
(iii) For all $t \in[0, T]$ and $\varphi \in \mathcal{C}_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ such that $\operatorname{supp}(\varphi) \subset B(0, k)$, we have

$$
\left\langle\sigma_{k} Q_{n}(t), \varphi\right\rangle=\left\langle\sigma_{k} \gamma_{k} Q_{n}(t), \varphi\right\rangle=\left\langle\gamma_{k} i_{n}(t), \sigma_{k} \varphi\right\rangle=\left\langle\gamma_{k} i_{n}(t), \varphi\right\rangle,
$$

then by taking the limit $n \longrightarrow+\infty$, we get $\left\langle\varrho^{\left(\sigma_{k}\right)}(t), \varphi\right\rangle=\left\langle!\imath^{\left(\gamma_{k}\right)}(t), \varphi\right\rangle$ on $\Omega^{\left(\sigma_{k}\right)} \cap \Omega^{\left(\gamma_{k}\right)}$, hence the result.

Remark 2.11. The above patching procedure could of course be applied to the second-order process $\chi^{2} Q_{\rho}$ as well, leading to a well-defined distribution-valued function $q \rho$.

## 3. Deterministic analysis of the equation under condition $\mathbf{H 1}$

In this section, we propose to analyze the equation in the regular situation, that is when assumption (H1) on the Hurst index is satisfied, and the linear solution $\rho$ ! (defined by Proposition (1.2) is known to be a function of time and space. Remember that in this case, the model is interpreted through Definition 1.4. Thus, what we wish to solve in this section is the equation

$$
\begin{align*}
& v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho \rho_{\tau}\right)\right) d \tau \\
&-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot\left(\overline{\rho_{\tau}^{\circ}}\right)\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}(|\rho_{\overbrace{\tau}}|^{2}) d \tau, \quad t \in[0, T] . \tag{3.1}
\end{align*}
$$

As opposed to the stochastic arguments used in the previous section, our strategy towards a (local) solution $v$ for (3.1) will rely on deterministic estimates only. In other words, we henceforth consider $\rho$ ? as a fixed (i.e., non-random) element in the space

$$
\mathcal{E}_{\beta}:=\bigcap_{2 \leq p \leq \infty} \mathcal{C}\left([0, T] ; \mathcal{W}^{\beta, p}\left(\mathbb{R}^{d}\right)\right)=\mathcal{C}\left([0, T] ; H^{\beta}\left(\mathbb{R}^{d}\right)\right) \cap \mathcal{C}\left([0, T] ; \mathcal{W}^{\beta, \infty}\left(\mathbb{R}^{d}\right)\right)
$$

for some appropriate $0<\beta<1$ (where $\beta=-\alpha$ is given by Proposition 1.2), and try to solve the following deterministic equation: for $\boldsymbol{\Psi} \in \mathcal{E}_{\beta}$,

$$
\begin{align*}
v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath & \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot \mathbf{\Psi}_{\tau}\right) d \tau \\
& -\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot \overline{\mathbf{\Psi}}_{\tau}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left|\mathbf{\Psi}_{\tau}\right|^{2}\right) d \tau \tag{3.2}
\end{align*}
$$

Let us set the stage for this solving procedure by reporting on fundamental estimates related to the two main operations in (3.2).

### 3.1. Pointwise multiplication.

Lemma 3.1 (Fractional Leibniz rule, see [42, Proposition 1.1, p. 105]). Let $s \geq 0,1<r<\infty$ and $1<p_{1}, p_{2}, q_{1}, q_{2}<\infty$ satisfying

$$
\frac{1}{r}=\frac{1}{p_{1}}+\frac{1}{p_{2}}=\frac{1}{q_{1}}+\frac{1}{q_{2}}
$$

Then one has

$$
\|u \cdot v\|_{\mathcal{W}^{s, r}\left(\mathbb{R}^{d}\right)} \lesssim\|u\|_{\mathcal{W}^{s, p_{1}}\left(\mathbb{R}^{d}\right)}\|v\|_{L^{p_{2}\left(\mathbb{R}^{d}\right)}}+\|u\|_{L^{q_{1}}\left(\mathbb{R}^{d}\right)}\|v\|_{\mathcal{W}^{s, q_{2}}\left(\mathbb{R}^{d}\right)} .
$$

### 3.2. Convolution with the Schrödinger group.

Naturally, we also need some control on the operation $(\phi, F, t) \mapsto S_{t} \phi-\imath \int_{0}^{t} S_{t-s}\left(F_{s}\right) d s$, or otherwise stated some estimate on the solution $u$ to the general Schrödinger equation

$$
\left\{\begin{array}{l}
\imath \partial_{t} u(t, x)-\Delta u(t, x)=F(t, x), \quad t \in[0, T], x \in \mathbb{R}^{d}  \tag{3.3}\\
u(0, x)=\phi(x)
\end{array}\right.
$$

Such a control is classically provided by the so-called Strichartz inequalities, which will prove to be sufficient for our purpose in this functional setting.

Definition 3.2. A pair $(p, q)$ is said to be Schrödinger admissible if

$$
(p, q) \in[2, \infty]^{2},(p, q, d) \neq(2, \infty, 2), \frac{2}{p}+\frac{d}{q}=\frac{d}{2}
$$

Lemma 3.3 (Strichartz inequalities, see [4, Paragraph 2.3]). Fix $d \geq 1, s \in \mathbb{R}$, and let $u$ stand for the mild solution of equation (3.3).

Then for all Schrödinger admissible pairs $(p, q)$ and $(a, b)$, it holds that

$$
\begin{equation*}
\|u\|_{L^{p}\left([0, T] ; \mathcal{W}^{s, q}\left(\mathbb{R}^{d}\right)\right)} \lesssim\|\phi\|_{H^{s}\left(\mathbb{R}^{d}\right)}+\|F\|_{L^{a^{\prime}}\left([0, T] ; \mathcal{\mathcal { N }}^{s, b^{\prime}}\left(\mathbb{R}^{d}\right)\right)} \tag{3.4}
\end{equation*}
$$

where the notations $a^{\prime}, b^{\prime}$ refer to the Hölder conjugates of $a, b$.

### 3.3. Solving the equation.

Our main result regarding equation (3.2) can be stated as follows:
Theorem 3.4. Assume that $1 \leq d \leq 4$ and fix $\beta \in(0,1)$. Consider the Schrödinger admissible pair $(p, q)$ given by the formulas

$$
p=\frac{12}{d-\beta}, q=\frac{6 d}{2 d+\beta}
$$

and for every $T>0$, define the space $X^{\beta}(T)$ as

$$
X^{\beta}(T):=\mathcal{C}\left([0, T] ; H^{\beta}\left(\mathbb{R}^{d}\right)\right) \cap L^{p}\left([0, T] ; \mathcal{W}^{\beta, q}\left(\mathbb{R}^{d}\right)\right)
$$

Then for all $\phi \in H^{\beta}\left(\mathbb{R}^{d}\right)$ and $\boldsymbol{\Psi} \in \mathcal{E}_{\beta}$, there exists a time $T_{0}>0$ such that equation (3.2) admits a unique solution in $X^{\beta}\left(T_{0}\right)$.

This local well-posedness result will be derived from a standard fixed-point argument. To this end, we introduce the map $\Gamma$ defined by the right-hand side of $(3.2)$, that is: for all $\boldsymbol{\Psi} \in \mathcal{E}_{\beta}$, $\phi \in H^{\beta}\left(\mathbb{R}^{d}\right), v \in X^{\beta}(T), T \geq 0$ and $t \in[0, T]$, set

$$
\begin{aligned}
& \Gamma_{T, \boldsymbol{\Psi}}(v)_{t}:=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot \mathbf{\Psi}_{\tau}\right) d \tau \\
&-\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot \overline{\mathbf{\Psi}}_{\tau}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\left|\mathbf{\Psi}_{\tau}\right|^{2}\right) d \tau
\end{aligned}
$$

Proposition 3.5. In the setting of Theorem 3.4, the following bounds hold true: there exists $\varepsilon>0$ such that for all $0 \leq T \leq 1, \phi \in H^{\beta}\left(\mathbb{R}^{d}\right),\left(\mathbf{\Psi}_{1}, \mathbf{\Psi}_{2}\right) \in \mathcal{E}_{\beta} \times \mathcal{E}_{\beta}$ and $v, v_{1}, v_{2} \in X^{\beta}(T)$,

$$
\begin{equation*}
\left\|\Gamma_{T, \boldsymbol{\Psi}_{1}}(v)\right\|_{X^{\beta}(T)} \lesssim\|\phi\|_{H^{\beta}\left(\mathbb{R}^{d}\right)}+T^{\varepsilon}\left[\|v\|_{X^{\beta}(T)}^{2}+\left\|\boldsymbol{\Psi}_{1}\right\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}\|v\|_{X^{\beta}(T)}+\left\|\boldsymbol{\Psi}_{1}\right\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}^{2}\right] \tag{3.5}
\end{equation*}
$$

and

$$
\begin{align*}
& \left\|\Gamma_{T, \boldsymbol{\Psi}_{1}}\left(v_{1}\right)-\Gamma_{T, \boldsymbol{\Psi}_{2}}\left(v_{2}\right)\right\|_{X^{\beta}(T)} \\
& \lesssim T^{\varepsilon}\left[\left\|v_{1}-v_{2}\right\|_{X^{\beta}(T)}\left\{\left\|v_{1}\right\|_{X^{\beta}(T)}+\left\|v_{2}\right\|_{X^{\beta}(T)}\right\}+\left\|\boldsymbol{\Psi}_{1}-\boldsymbol{\Psi}_{2}\right\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}\left\|v_{1}\right\|_{X^{\beta}(T)}\right. \\
& \left.\quad+\left\|\boldsymbol{\Psi}_{2}\right\|_{L_{T}^{\infty}} \mathcal{W}^{\beta, q}\left\|v_{1}-v_{2}\right\|_{X^{\beta}(T)}+\left\|\boldsymbol{\Psi}_{1}-\boldsymbol{\Psi}_{2}\right\|_{L_{T}^{\infty}} \mathcal{W}^{\beta, q}\left\{\left\|\boldsymbol{\Psi}_{1}\right\|_{L_{T}^{\infty}} \mathcal{W}^{\beta, q}+\left\|\boldsymbol{\Psi}_{2}\right\|_{L_{T}^{\infty}} \mathcal{W}^{\beta, q}\right\}\right] \tag{3.6}
\end{align*}
$$

where the proportional constants only depend on $\beta$ and $\rho$.
Before we turn to the proof of this proposition, let us briefly recall that, once endowed (3.5)(3.6), the statement of Theorem 3.4 follows from a standard two-step procedure. Namely, using (3.5), we can first establish that for any $T=T(\phi, \boldsymbol{\Psi})>0$ small enough, there exists a ball in $X^{\beta}(T)$ that is stable through the application of $\Gamma_{T, \Psi}$. Then, thanks to (3.6) (applied with $\boldsymbol{\Psi}_{1}=\boldsymbol{\Psi}_{2}=\boldsymbol{\Psi}$ ), we can show that $\Gamma_{T, \boldsymbol{\Psi}}$ is actually a contraction on this ball (for $T>0$ possibly even smaller), which completes the proof of the assertion.

Note also that the continuity of $\Gamma_{T, \Psi}$ with respect to $\Psi$ (an immediate consequence of (3.6) will be the key ingredient toward item (ii) of Theorem 1.10 .

Proof of Proposition 3.5. Let us set, for any suitable distribution $u$ on $\mathbb{R}^{d+1}$,

$$
\mathcal{G}(u)_{t}:=-\imath \int_{0}^{t} S_{t-\tau}\left(u_{\tau}\right) d \tau
$$

which allows to recast $\Gamma_{T, \Psi}$ as

$$
\Gamma_{T, \Psi}(v)=S(\phi)+\mathcal{G}\left(\rho^{2}|v|^{2}\right)+\mathcal{G}(\rho \bar{v} \cdot \boldsymbol{\Psi})+\mathcal{G}(\rho v \cdot \overline{\mathbf{\Psi}})+\mathcal{G}\left(|\boldsymbol{\Psi}|^{2}\right) .
$$

Let us now bound each of the four above terms separately.
Bound on $S(\phi)$ : Since $(\infty, 2)$ and $(p, q)$ are both Schrödinger admissible pairs, we can apply Lemma 3.3 to assert that

$$
\begin{equation*}
\|S(\phi)\|_{X^{\beta}(T)} \lesssim\|\phi\|_{H^{\beta}} . \tag{3.7}
\end{equation*}
$$

Bound on $\mathcal{G}\left(\rho^{2}|v|^{2}\right)$ : By Lemma 3.3, we can first assert that

$$
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X^{\beta}(T)} \lesssim\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}}
$$

Let us now introduce the additional parameter $n:=\frac{3 d}{d-\beta}>1$, in such a way that

$$
\frac{1}{q^{\prime}}=\frac{1}{q}+\frac{1}{n}
$$

Using the fractional Leibniz rule given by Lemma 3.1. we get that for all $t \geq 0$,

$$
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{\beta, q^{\prime}}} \lesssim\|\rho v(t, .)\|_{\mathcal{W}^{\beta, q}}\|\rho v(t, .)\|_{L^{n}}
$$

It is easy to check that $\beta \geq d\left(\frac{1}{q}-\frac{1}{n}\right)$, and accordingly we can rely on the Sobolev embedding

$$
\begin{equation*}
\mathcal{W}^{\beta, q}\left(\mathbb{R}^{d}\right) \hookrightarrow L^{n}\left(\mathbb{R}^{d}\right) \tag{3.8}
\end{equation*}
$$

to derive that

$$
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{\beta}, q^{\prime}} \lesssim\|\rho v(t, .)\|_{\mathcal{W}^{\beta, q}}^{2}
$$

Consider the parameter $m$ defined through the relation

$$
\frac{1}{p^{\prime}}=\frac{2}{p}+\frac{1}{m}
$$

Since $1 \leq d \leq 4$ and $\beta \in(0,1)$, it can actually be verified that $\frac{1}{m}=1-\frac{d-\beta}{4}>0$. Then, by Hölder inequality, we have

$$
\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}} \lesssim T^{\frac{1}{m}}\|v\|_{L_{T}^{p} \mathcal{W}^{\beta, q}}^{2} \lesssim T^{1-\frac{d-\beta}{4}}\|v\|_{X^{\beta}(T)}^{2}
$$

and we have thus shown that

$$
\begin{equation*}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X^{\beta}(T)} \lesssim T^{1-\frac{d-\beta}{4}}\|v\|_{X^{\beta}(T)}^{2} . \tag{3.9}
\end{equation*}
$$

Bound on $\mathcal{G}(\rho \bar{v} \cdot \boldsymbol{\Psi}), \mathcal{G}(\rho v \cdot \overline{\mathbf{\Psi}})$ : Just as above, we can first apply Lemma 3.3 to get that

$$
\|\mathcal{G}(\rho \bar{v} \cdot \boldsymbol{\Psi})\|_{X^{\beta}(T)}+\|\mathcal{G}(\rho v \cdot \overline{\mathbf{\Psi}})\|_{X^{\beta}(T)} \lesssim\|\rho \bar{v} \cdot \boldsymbol{\Psi}\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}}
$$

Thanks to Lemma 3.1, it holds, for all $t \geq 0$,

$$
\begin{aligned}
\|\rho \bar{v} \cdot \boldsymbol{\Psi}(t, .)\|_{\mathcal{W}^{\beta, q^{\prime}}} & \lesssim\|\rho v(t, .)\|_{\mathcal{W}^{\beta, q}}\|\mathbf{\Psi}(t, .)\|_{L^{n}}+\|\boldsymbol{\Psi}(t, .)\|_{\mathcal{W}^{\beta, q}}\|\rho v(t, .)\|_{L^{n}} \\
& \lesssim\|\boldsymbol{\Psi}(t, .)\|_{\mathcal{W}^{\beta, q}}\|\rho v(t, .)\|_{\mathcal{W}^{\beta, q}}
\end{aligned}
$$

where we have again used the Sobolev embedding (3.8).
Then, by Hölder inequality, we deduce

$$
\begin{aligned}
\|\rho \bar{v} \cdot \boldsymbol{\Psi}\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}} & \lesssim T^{\frac{1}{p}+\frac{1}{m}}\|\boldsymbol{\Psi}\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}\|v\|_{L_{T}^{p} \mathcal{W}^{\beta, q}} \\
& \lesssim T^{\frac{1}{p}+\frac{1}{m}}\|\boldsymbol{\Psi}\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}\|v\|_{X^{\beta}(T)}
\end{aligned}
$$

and we have thus established that

$$
\begin{equation*}
\|\mathcal{G}(\rho \bar{v} \cdot \boldsymbol{\Psi})\|_{X^{\beta}(T)}+\|\mathcal{G}(\rho v \cdot \overline{\mathbf{\Psi}})\|_{X^{\beta}(T)} \lesssim T^{\frac{1}{p}+\frac{1}{m}}\|\boldsymbol{\Psi}\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}\|v\|_{X^{\beta}(T)} \tag{3.10}
\end{equation*}
$$

Bound on $\mathcal{G}\left(|\Psi|^{2}\right)$ : By Lemma 3.3.

$$
\left\|\mathcal{G}\left(|\Psi|^{2}\right)\right\|_{X^{\beta}(T)} \lesssim\left\||\Psi|^{2}\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}} .
$$

Using Lemma 3.1 and the Sobolev embedding (3.8), we get that for every $t \geq 0$,

$$
\left\||\boldsymbol{\Psi}|^{2}(t, .)\right\|_{\mathcal{W}^{\beta, q^{\prime}}} \lesssim\|\boldsymbol{\Psi}(t, .)\|_{\mathcal{W}^{\beta, q}}\|\boldsymbol{\Psi}(t, .)\|_{L^{n}} \lesssim\|\boldsymbol{\Psi}(t, .)\|_{\mathcal{W}^{\beta, q}}^{2} .
$$

Then

$$
\left\||\Psi|^{2}\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}} \lesssim T^{\frac{1}{p^{\prime}}}\|\Psi\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}^{2}
$$

and finally

$$
\begin{equation*}
\left\|\mathcal{G}\left(|\boldsymbol{\Psi}|^{2}\right)\right\|_{X^{\beta}(T)} \lesssim T^{\frac{1}{p^{\prime}}}\|\boldsymbol{\Psi}\|_{L_{T}^{\infty} \mathcal{W}^{\beta, q}}^{2} . \tag{3.11}
\end{equation*}
$$

The combination of estimates (3.7), (3.9), (3.10) and (3.11) entails the desired bound (3.5).
It is then easy to see that (3.6) can be derived from similar arguments: for instance,

$$
\left\|\mathcal{G}\left(\rho^{2}\left(\left|v_{1}\right|^{2}-\left|v_{2}\right|^{2}\right)\right)\right\|_{X^{\beta}(T)} \lesssim\left\|\rho^{2}\left(\left|v_{1}\right|^{2}-\left|v_{2}\right|^{2}\right)\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}} \lesssim\left\|\left|v_{1}\right|^{2}-\left|v_{2}\right|^{2}\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}}
$$

Combining again Lemma 3.1 and embedding (3.8), we obtain, for every $t \geq 0$,

$$
\begin{aligned}
\left\|\left(\left|v_{1}\right|^{2}-\left|v_{2}\right|^{2}\right)(t, .)\right\|_{\mathcal{W}^{\beta, q^{\prime}}} \lesssim & \left\|\left(v_{1}-v_{2}\right)(t, .)\right\|_{\mathcal{W}^{\beta, q}}\left\{\left\|v_{1}(t, .)\right\|_{L^{n}}+\left\|v_{2}(t, .)\right\|_{L^{n}}\right\} \\
& +\left\|\left(v_{1}-v_{2}\right)(t, .)\right\|_{L^{n}}\left\{\left\|v_{1}(t, .)\right\|_{\mathcal{W}^{\beta, q}}+\left\|v_{2}(t, .)\right\|_{\mathcal{W}^{\beta, q}}\right\} \\
\lesssim & \left\|\left(v_{1}-v_{2}\right)(t, .)\right\|_{\mathcal{W}^{\beta, q}}\left\{\left\|v_{1}(t, .)\right\|_{\mathcal{W}^{\beta, q}}+\left\|v_{2}(t, .)\right\|_{\mathcal{W}^{\beta}, q}\right\},
\end{aligned}
$$

and as a result

$$
\left\|\left|v_{1}\right|^{2}-\left|v_{2}\right|^{2}\right\|_{L_{T}^{p^{\prime}} \mathcal{W}^{\beta, q^{\prime}}} \lesssim T^{\frac{1}{m}}\left\|v_{1}-v_{2}\right\|_{L_{T}^{p} \mathcal{W}^{\beta, q}}\left\{\left\|v_{1}\right\|_{L_{T}^{p} \mathcal{W}^{\beta, q}}+\left\|v_{2}\right\|_{L_{T}^{p} \mathcal{W}^{\beta, q}}\right\} .
$$

### 3.4. Proof of Theorem 1.10,

At this point, the statement of Theorem 1.10 (item (i)) is of course a mere combination of the construction of $\rho$ ? as an element in $\mathcal{E}_{\beta}$ (Proposition 1.2) with the well-posedness result of Theorem 3.4. In brief, it suffices to apply Theorem 3.4 (in an almost sure way) to $\boldsymbol{\Psi}:=\rho$ ?

As for the convergence property in item (ii), it can easily be deduced from the continuity of $\Gamma_{T, \boldsymbol{\Psi}}$ with respect to $\boldsymbol{\Psi}$ (along (3.6) and the almost sure convergence of $\chi$ ? $n$ to $\chi$ ?. Additional details about this elementary procedure can be found in the proof of [20, Theorem 1.7].

## 4. Deterministic analysis of the equation under condition (H2')

It remains us to deal with the wellposedness issue in the rough case, that is to present the proof of Theorem 1.11. Therefore, we assume in this section that condition (H2') on the Hurst indexes is satisfied. We recall that in this rough situation, the equation is understood in the sense of Definition 1.8, that is as

$$
\begin{align*}
v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath & \int_{0}^{t} S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot\left(\rho \rho_{\tau}\right)\right) d \tau \\
& -\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot\left(\overline{\rho_{\tau}^{\prime}}\right)\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2} Q_{\tau}\right) d \tau \tag{4.1}
\end{align*}
$$

where the processes $\rho$ ? and $\rho^{2} Q_{\rho}$ are defined through Proposition 1.2 and Proposition 1.6
In order to handle (4.1), we intend to follow the same deterministic approach as in Section 3 . In other words, we will henceforth consider the pair ( $\rho \mathrm{Q}, \rho^{2} Q \rho$ ) as a given element in the space

$$
\begin{equation*}
\mathcal{R}_{\alpha}:=\bigcap_{2 \leq p \leq \infty} L^{\infty}\left([0, T] ; \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right)\right) \times L^{\infty}\left([0, T] ; H^{-2 \alpha}\left(\mathbb{R}^{d}\right)\right) \tag{4.2}
\end{equation*}
$$

for some $0<\alpha<1$ (provided by Propositions 1.2 and 1.6), and then try to solve the more general deterministic equation: for $\left(\boldsymbol{\Psi}^{\mathbf{1}}, \boldsymbol{\Psi}^{\mathbf{2}}\right) \in \mathcal{R}_{\alpha}$,

$$
\begin{align*}
v_{t}=S_{t}(\phi)-\imath \int_{0}^{t} S_{t-\tau}\left(\rho^{2}\left|v_{\tau}\right|^{2}\right) d \tau-\imath \int_{0}^{t} & S_{t-\tau}\left(\left(\rho \bar{v}_{\tau}\right) \cdot \mathbf{\Psi}_{\tau}^{\mathbf{1}}\right) d \tau \\
& -\imath \int_{0}^{t} S_{t-\tau}\left(\left(\rho v_{\tau}\right) \cdot \overline{\mathbf{\Psi}_{\tau}^{\mathbf{1}}}\right) d \tau-\imath \int_{0}^{t} S_{t-\tau}\left(\mathbf{\Psi}_{\tau}^{\mathbf{2}}\right) d \tau . \tag{4.3}
\end{align*}
$$

As we have already highlighted it in Section 1.4.2, the whole specificity of the situation (in comparison to Section 3 ) lies in the irregularity of $\boldsymbol{\Psi}_{\tau}^{1}$ and $\boldsymbol{\Psi}_{\tau}^{2}$, which can only be treated as negative-order distributions (note indeed that $\alpha>0$ in $(\overline{4.2})$ ). The technical ingredients towards a fixed-point argument need to be revised accordingly: this will be the purpose of the subsequent Sections 4.1.4.3, which lay the ground for our main wellposedness result, namely Theorem 4.5.

### 4.1. Pointwise multiplication and interpolation.

In view of the above considerations, our only possibility to handle the product $\left(\rho v_{\tau}\right) \cdot\left(\boldsymbol{\Psi}_{\tau}^{\mathbf{1}}\right)$ in (4.3) will be to rely on the following general multiplication property in Sobolev spaces (see e.g. [39, Section 4.4.3] for a proof of this result):

Lemma 4.1. Fix $d \geq 1$. Let $\alpha, \beta>0$ and $1 \leq p, p_{1}, p_{2}<\infty$ be such that

$$
\begin{array}{r}
\frac{1}{p}=\frac{1}{p_{1}}+\frac{1}{p_{2}} \quad \text { and } \quad 0<\alpha<\beta \\
\text { If } f \in \mathcal{W}^{-\alpha, p_{1}}\left(\mathbb{R}^{d}\right) \text { and } g \in \mathcal{W}^{\beta, p_{2}}\left(\mathbb{R}^{d}\right) \text {, then } f \cdot g \in \mathcal{W}^{-\alpha, p}\left(\mathbb{R}^{d}\right) \text { and } \\
\|f \cdot g\|_{\mathcal{W}^{-\alpha, p}} \lesssim\|f\|_{\mathcal{W}^{-\alpha, p_{1}}}\|g\|_{\mathcal{W}^{\beta, p_{2}}}
\end{array}
$$

Let us also label the following classical interpolation result for further reference:
Lemma 4.2. Fix $d \geq 1$. Let $s, s_{1}, s_{2} \in \mathbb{R}$ and $1 \leq p, p_{1}, p_{2}<\infty$ be such that, for some $\theta \in(0,1)$,

$$
s=\theta s_{1}+(1-\theta) s_{2} \quad \text { and } \quad \frac{1}{p}=\frac{\theta}{p_{1}}+\frac{1-\theta}{p_{2}} .
$$

Then for every $v \in \mathcal{W}^{s_{1}, p_{1}}\left(\mathbb{R}^{d}\right) \cap \mathcal{W}^{s_{2}, p_{2}}\left(\mathbb{R}^{d}\right)$, it holds that $v \in \mathcal{W}^{s, p}\left(\mathbb{R}^{d}\right)$ and

$$
\|v\|_{\mathcal{W}^{s, p}} \leq\|v\|_{\mathcal{W}^{s_{1}, p_{1}}}^{\theta}\|v\|_{\mathcal{W}^{s_{2}, p_{2}}}^{1-\theta} .
$$

### 4.2. A local regularization property of the Schrödinger group $S$.

It is a well-known fact that the classical Strichartz inequalities for the Schrödinger group (summed up in Lemma 3.3) do not offer any regularization effect, as can be seen from the constant derivative parameter $s$ in (3.4). This phenomenon naturally becomes a fundamental obstacle in our rough setting, where, for stability reasons, the distribution $\left(\rho v_{\tau}\right) \cdot\left(\boldsymbol{\Psi}_{\tau}^{\mathbf{1}}\right)$ in (4.3) is expected to turn into a function through the action of $S$.

A possible way to reach such a regularization property is to let local Sobolev topologies come into picture, through the consideration of the spaces $H_{\rho}^{s}\left(\mathbb{R}^{d}\right)$ defined by (1.19). Our main technical result in this direction can be stated as follows:
Lemma 4.3. Fix $d \geq 1$. Let $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be of the form $\left(\overline{\mathbf{F}_{\rho}}\right), 0 \leq \alpha, \kappa \leq \frac{1}{2}$ and $0 \leq T \leq 1$. Assume that $\phi \in H^{-\alpha}\left(\mathbb{R}^{d}\right), F \in L^{1}\left([0, T] ; H^{-\alpha}\left(\mathbb{R}^{d}\right)\right)$, and consider the solution $u$ of the following inhomogeneous Schrödinger equation on $\mathbb{R}^{d}$

$$
\left\{\begin{array}{l}
\imath \partial_{t} u(t, x)-\Delta u(t, x)=F(t, x), \quad t \in[0, T], x \in \mathbb{R}^{d} \\
u_{0}=\phi
\end{array}\right.
$$

Then it holds that

$$
\begin{equation*}
\|u\|_{L_{T}^{\frac{1}{\hbar}} H_{\rho}^{-\alpha+\kappa}} \lesssim\|\phi\|_{H^{-\alpha}\left(\mathbb{R}^{d}\right)}+\|F\|_{L_{T}^{1} H^{-\alpha}}, \tag{4.4}
\end{equation*}
$$

where the proportional constant only depends on $\rho, \alpha$ and $\kappa$.
The above property can in fact be seen as a slight extension of the result of [8, Theorem 3.1]. For the sake of clarity, we have postponed the proof of the lemma to Section 5.1.

### 4.3. A commutator estimate.

Keeping our objective in mind (that is, to settle a fixed-point argument for (4.3), the previous estimate (4.4) clearly lacks some stability: the left-hand side is indeed based on the consideration of a local Sobolev norm (in $H_{\rho}^{-\alpha+\kappa}$ ), while the right-hand side appeals to a standard Sobolev space $\left(H^{-\alpha}\right)$.

Our strategy to overcome this problem will consist in using the presence of the cut-off function $\rho$ within the model (4.3) (through $\rho v$ ), which somehow allows us to turn global Sobolev norms into local ones. To implement this idea, an additional commutator-type estimate will be required:

Lemma 4.4. For every $s>0$ and for all test functions $\rho, g: \mathbb{R}^{d} \rightarrow \mathbb{R}$, it holds that

$$
\begin{equation*}
\left\|(\operatorname{Id}-\Delta)^{\frac{s}{2}}(\rho \cdot g)-\rho \cdot(\operatorname{Id}-\Delta)^{\frac{s}{2}}(g)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)} \lesssim\|g\|_{H^{s-1}\left(\mathbb{R}^{d}\right)} \tag{4.5}
\end{equation*}
$$

where the proportional constant only depends on $\rho$ and $s$.
As a consequence, for every test function $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ and for every $g \in H_{\rho}^{s}\left(\mathbb{R}^{d}\right) \cap H^{s-1}\left(\mathbb{R}^{d}\right)$, it holds that

$$
\begin{equation*}
\|\rho \cdot g\|_{H^{s}} \lesssim\|g\|_{H_{\rho}^{s}}+\|g\|_{H^{s-1}} \tag{4.6}
\end{equation*}
$$

and

$$
\begin{equation*}
\|g\|_{H_{\rho}^{s}} \lesssim\|\rho \cdot g\|_{H^{s}}+\|g\|_{H^{s-1}} \tag{4.7}
\end{equation*}
$$

for some proportional constant depending only on $\rho$ and $s$.

Proof. See Section 5.2

### 4.4. Solving the auxiliary deterministic equation.

Let us fix (once and for all) a cut-off function $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ of the form $\left(\mathbf{\mathbf { F } _ { \rho }}\right)$, and for all $T \geq 0$, $\alpha, \kappa>0, p, q \geq 2$, define the space

$$
\begin{equation*}
X_{\rho}^{\alpha, \kappa,(p, q)}(T):=\mathcal{C}\left([0, T] ; H^{-2 \alpha}\left(\mathbb{R}^{d}\right)\right) \cap L^{p}\left([0, T] ; \mathcal{W}^{-2 \alpha, q}\left(\mathbb{R}^{d}\right)\right) \cap L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa} \tag{4.8}
\end{equation*}
$$

Besides, recall that the space $\mathcal{R}_{\alpha}$ has been introduced in 4.2.
We are finally in a position to state (and prove) the main result of this section:
Theorem 4.5. Assume that $1 \leq d \leq 3$ and that

$$
0<\alpha< \begin{cases}\frac{3}{20} & \text { if } d=1  \tag{4.9}\\ \frac{1}{10} & \text { if } d=2 \\ \frac{1}{24} & \text { if } d=3\end{cases}
$$

Then one can find parameters $\kappa \in[2 \alpha, 1 / 2]$ and $p, q \geq 2$ such that for all $\phi \in H^{-2 \alpha}\left(\mathbb{R}^{d}\right)$ and $\left(\mathbf{\Psi}^{\mathbf{1}}, \mathbf{\Psi}^{\mathbf{2}}\right) \in \mathcal{R}_{\alpha}$, there exists a time $T>0$ for which equation (4.3) admits a unique solution in the above-defined set $X_{\rho}^{\alpha, \kappa,(p, q)}(T)$.

Remark 4.6. As could be checked from a review of our arguments in the below proof, the condition (4.9) on $\alpha$ is essentially optimal with respect to the spaces and the tools that we have relied on. To be more specific, condition 4.9 is derived from an optimal choice of the four parameters $\alpha, \kappa, p, q$ in the scale of spaces (4.8), when using Lemmas 4.1 4.4 to estimate the right-hand side of 4.3).

We do not pretend that this restriction on $\alpha$ could not be alleviated by considering a different solution space, or using more sophisticated tools to control the equation.

Remark 4.7. As we mentionned it in the introduction, the well-posedness of similar (deterministic) quadratic NLS has already been studied in the literature. A recurrent ingredient consists of sharp bilinear estimates prevailing in the so-called Bourgain spaces (see [7, 2]). However, it seems to us that those techniques could not be directly applied to our problem, for two reasons.

Firstly, it is not clear how the term $\rho^{2}|u|^{2}$ could be treated through the bilinear estimates of [7], since the Bourgain spaces $X^{s, b}$ are not stable by multiplication with a $\mathcal{C}_{c}^{\infty}$ function ${ }^{1}$, Secondly, even if we replace $\rho$ with 1 in the initial problem (1.1) (thus getting access to sharp bilinear estimates for $|u|^{2}$ ), it is unlikely that the stochastic terms 9 and $q \circ$ can then be injected into Bourgain spaces, owing to the spatial asymptotic behavior of those processes.

[^1]Just as in Section 3.3, the proof of Theorem 4.5 is in fact a straightforward consequence of the following estimates for the map $\Gamma_{T, \Psi^{\mathbf{1}}, \boldsymbol{\Psi}^{\mathbf{2}}}$ defined for all $T \geq 0$ and $\left(\mathbf{\Psi}^{\mathbf{1}}, \mathbf{\Psi}^{\mathbf{2}}\right) \in \mathcal{R}_{\alpha}$ by

$$
\Gamma_{T, \boldsymbol{\Psi}^{\mathbf{1}}, \boldsymbol{\Psi}^{\mathbf{2}}}(v):=S(\phi)+\mathcal{G}\left(\rho^{2}|v|^{2}\right)+\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}^{\mathbf{1}}\right)+\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}^{\mathbf{1}}}\right)+\mathcal{G}\left(\mathbf{\Psi}^{\mathbf{2}}\right)
$$

where the shortcut notation $\mathcal{G}$ refers to the operator

$$
\mathcal{G}(u)_{t}:=-\imath \int_{0}^{t} S_{t-\tau}\left(u_{\tau}\right) d \tau
$$

Proposition 4.8. Assume that $1 \leq d \leq 3$ and that $\alpha$ satisfies condition (4.9). Then one can find parameters $\kappa>0, p, q \geq 2$ and $\varepsilon>0$ such that, setting $X(T):=X_{\rho}^{\alpha, \kappa,(p, q)}(T)$, the following bounds hold true: for all $0 \leq T \leq 1, \phi \in H^{-2 \alpha}\left(\mathbb{R}^{d}\right),\left(\mathbf{\Psi}_{1}^{\mathbf{1}}, \mathbf{\Psi}_{1}^{\mathbf{2}}\right) \in \mathcal{R}_{\alpha},\left(\mathbf{\Psi}_{2}^{\mathbf{1}}, \mathbf{\Psi}_{2}^{\mathbf{2}}\right) \in \mathcal{R}_{\alpha}$ and $v, v_{1}, v_{2} \in X(T)$,

$$
\begin{equation*}
\left\|\Gamma_{T, \mathbf{\Psi}_{1}^{1}, \boldsymbol{\Psi}_{1}^{2}}(v)\right\|_{X(T)} \lesssim\|\phi\|_{H^{-2 \alpha}}+T^{\varepsilon}\left[\|v\|_{X(T)}^{2}+\left\|\mathbf{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)}+\left\|\mathbf{\Psi}_{1}^{\mathbf{2}}\right\|_{L_{T}^{\infty} H^{-2 \alpha}}\right] \tag{4.10}
\end{equation*}
$$

and

$$
\begin{align*}
& \left\|\Gamma_{T, \mathbf{\Psi}_{1}^{1}, \boldsymbol{\Psi}_{1}^{2}}\left(v_{1}\right)-\Gamma_{T, \boldsymbol{\Psi}_{2}^{1}, \mathbf{\Psi}_{2}^{2}}\left(v_{2}\right)\right\|_{X(T)} \\
& \lesssim T^{\varepsilon}\left[\left\|v_{1}-v_{2}\right\|_{X(T)}\left\{\left\|v_{1}\right\|_{X(T)}+\left\|v_{2}\right\|_{X(T)}\right\}+\left\|\mathbf{\Psi}_{1}^{\mathbf{1}}-\mathbf{\Psi}_{2}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\left\|v_{1}\right\|_{X(T)}\right. \\
&  \tag{4.11}\\
& \left.\quad+\left\|\mathbf{\Psi}_{2}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\left\|v_{1}-v_{2}\right\|_{X(T)}+\left\|\mathbf{\Psi}_{1}^{\mathbf{2}}-\mathbf{\Psi}_{2}^{\mathbf{2}}\right\|_{L_{T}^{\infty} H^{-2 \alpha}}\right]
\end{align*}
$$

where $r$ depends on $\alpha$ and $\kappa$ and the proportional constants depend only on $\rho$ and $\alpha$.
The choice of the three parameters $\kappa, p, q$ in the above proposition highly depends on the space dimension $d \in\{1,2,3\}$. For the sake of clarity, let us consider each value of $d$ in a distinct subsection.

### 4.4.1. Proof of Proposition 4.8 when $d=1$.

In this situation, we pick $\kappa$ such that $3 \alpha<\kappa<\inf \left(\frac{1}{2}, \frac{3}{4}-2 \alpha\right)$ and $(p, q):=(\infty, 2)$, so that the space under consideration reduces to

$$
X(T):=\mathcal{C}\left([0, T] ; H^{-2 \alpha}(\mathbb{R})\right) \cap L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}
$$

Also, we set $\theta:=\frac{2 \alpha}{\kappa} \in\left(0, \frac{2}{3}\right)$.
We now bound each term in the expression of $\Gamma_{T, \Psi^{1}, \Psi^{2}}$ separately. In the sequel we assume that $0 \leq T \leq 1$.

Bound on $S(\phi)$ : As $S$ is a unitary operator on $H^{-2 \alpha}(\mathbb{R})$, one has

$$
\|S(\phi)\|_{L_{T}^{\infty} H^{-2 \alpha}}=\|\phi\|_{H^{-2 \alpha}} .
$$

Besides, since $\alpha \leq \frac{1}{4}$ and $\kappa \leq \frac{1}{2}$, we can apply Lemma 4.3 to assert that

$$
\|S(\phi)\|_{L_{T}^{\frac{1}{\varepsilon}} H_{\rho}^{-2 \alpha+\kappa}} \lesssim\|\phi\|_{H^{-2 \alpha}}
$$

and we have thus shown that

$$
\|S(\phi)\|_{X(T)} \lesssim\|\phi\|_{H^{-2 \alpha}}
$$

Bound on $\mathcal{G}\left(\rho^{2}|v|^{2}\right)$ : Since $\kappa>0$, and since $\rho$ is smooth and compactly-supported, one has

$$
\begin{aligned}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} & =\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\frac{1}{k}} H_{\rho}^{-2 \alpha+\kappa}} \\
& \lesssim\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\frac{1}{k}} H^{-2 \alpha+\kappa}} \\
& \lesssim\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha+\kappa}} .
\end{aligned}
$$

From here we can apply Strichartz inequality (Lemma 3.3) to assert that

$$
\begin{equation*}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} \lesssim\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{\frac{4}{3}} \mathcal{W}^{-2 \alpha+\kappa, 1}} \tag{4.12}
\end{equation*}
$$

By Lemma 3.1, one has, for every fixed $t \geq 0$,

$$
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, 1}} \lesssim\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}\|\rho v(t, .)\|_{L^{2}},
$$

and then, by Lemma 4.2,

$$
\|\rho v(t, .)\|_{L^{2}} \leq\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}^{\theta}\|\rho v(t, .)\|_{H^{-2 \alpha}}^{1-\theta},
$$

which entails, for every fixed $t \geq 0$,

$$
\begin{aligned}
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, 1}} & \lesssim\|\rho v(t, .)\|_{H-2 \alpha+\kappa}^{1+\theta}\|\rho v(t, .)\|_{H^{-2 \alpha}}^{1-\theta} \\
& \lesssim\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{1+\theta}\|v(t, .)\|_{H^{-2 \alpha}}^{1-\theta}+\|v(t, .)\|_{H^{-2 \alpha}}^{2},
\end{aligned}
$$

where we have used Lemma 4.4 to derive the second inequality.
As a result,

$$
\begin{aligned}
& \int_{0}^{T} d t\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, 1}}^{\frac{4}{3}} \\
& \lesssim\|v\|_{X(T)}^{\frac{4}{3}(1-\theta)} \int_{0}^{T} d t\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{\frac{4}{3}(1+\theta)}+T\|v\|_{X(T)}^{\frac{8}{3}} \\
& \lesssim T^{1-\frac{4}{3}(1+\theta) \kappa}\|v\|_{X(T)}^{\frac{4}{3}(1-\theta)}\left(\int_{0}^{T} d t\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{\frac{1}{\kappa}}\right)^{\frac{4}{3}(1+\theta) \kappa}+T\|v\|_{X(T)}^{\frac{8}{3}} \\
& \lesssim T^{1-\frac{4}{3}(1+\theta) \kappa}\|v\|_{X(T)}^{\frac{8}{3}},
\end{aligned}
$$

and thus, going back to 4.12), we have shown the desired estimate, that is

$$
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} \lesssim T^{\frac{3}{4}-(\kappa+2 \alpha)}\|v\|_{X(T)}^{2}
$$

Bound on $\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{1}\right), \mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{1}}\right)$ : Since $\alpha \leq \frac{1}{4}$ and $\kappa \leq \frac{1}{2}$, we can appeal to Lemma 4.3 to assert that

$$
\begin{align*}
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)\right\|_{L_{T}^{\frac{1}{k}} H_{\rho}^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{1}}\right)\right\|_{L_{T}^{\frac{1}{\hbar}} H_{\rho}^{-2 \alpha+\kappa}} & \lesssim\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{1} H^{-2 \alpha}} \\
& \lesssim\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{1} H^{-\alpha}} . \tag{4.13}
\end{align*}
$$

Let $s>\alpha$ and $2 \leq r, p<\infty$ be such that $1 / 2=1 / r+1 / p$, then by Lemma 4.1. for every $t \geq 0$

$$
\begin{equation*}
\left\|\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)(t, .)\right\|_{H^{-\alpha}} \lesssim\left\|\boldsymbol{\Psi}_{1}^{1}(t, .)\right\|_{\mathcal{W}^{-\alpha, r}}\|\rho v(t, .)\|_{\mathcal{W}^{s, p}} . \tag{4.14}
\end{equation*}
$$

Next, observe that if

$$
-2 \alpha+\kappa-s=d\left(\frac{1}{2}-\frac{1}{p}\right)=\frac{d}{r} \quad \Leftrightarrow \quad s=-2 \alpha+\kappa-\frac{d}{r},
$$

we have the Sobolev embedding $H^{-2 \alpha+\kappa}\left(\mathbb{R}^{d}\right) \hookrightarrow \mathcal{W}^{s, p}\left(\mathbb{R}^{d}\right)$. Therefore, since $-2 \alpha+\kappa>\alpha$, if $r \geq 2$ is large enough, $s=-2 \alpha+\kappa-\frac{d}{r}>\alpha$, and from 4.14 we deduce

$$
\left\|\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{\mathbf{1}}\right)(t, .)\right\|_{H^{-\alpha}} \lesssim\left\|\boldsymbol{\Psi}_{1}^{1}(t, .)\right\|_{\mathcal{W}^{-\alpha, r}}\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}
$$

By applying Lemma 4.4, we obtain that for every $t \geq 0$,

$$
\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}} \lesssim\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}+\|v(t, .)\|_{H^{-2 \alpha}}
$$

and so we deduce

$$
\begin{aligned}
\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{1} H^{-\alpha}}+\left\|\rho v \cdot \overline{\boldsymbol{\Psi}_{1}^{\mathbf{1}}}\right\|_{L_{T}^{1} H^{-\alpha}} & \lesssim\left\|\boldsymbol{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\left\{T^{1-\kappa}\|v\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}}+T\|v\|_{L_{T}^{\infty} H^{-2 \alpha}}\right\} \\
& \lesssim T^{1-\kappa}\left\|\boldsymbol{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)}
\end{aligned}
$$

which, going back to 4.13, leads us to

$$
\begin{equation*}
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{\mathbf{1}}\right)\right\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{\mathbf{1}}}\right)\right\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}} \lesssim T^{1-\kappa}\left\|\mathbf{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)} \tag{4.15}
\end{equation*}
$$

On the other hand, by applying Lemma 3.3, we get

$$
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{\mathbf{1}}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{\mathbf{1}}}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}} \lesssim\left\|\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{1} H^{-2 \alpha}} \lesssim\left\|\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{1} H^{-\alpha}}
$$

We are thus in the same position as in (4.13), and we can repeat the above arguments to obtain

$$
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{\mathbf{1}}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{\mathbf{1}}}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}} \lesssim T^{1-\kappa}\left\|\boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)}
$$

Combining this bound with 4.15, we can conclude that

$$
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{\mathbf{1}}\right)\right\|_{X(T)}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{\mathbf{1}}}\right)\right\|_{X(T)} \lesssim T^{1-\kappa}\left\|\boldsymbol{\Psi}_{1}^{\mathbf{1}}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)}
$$

Bound on $\mathcal{G}\left(\mathbf{\Psi}_{1}^{\mathbf{2}}\right)$ : First, according to Lemma 4.3. we know that

$$
\left\|\mathcal{G}\left(\mathbf{\Psi}_{1}^{\mathbf{2}}\right)\right\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}} \lesssim\left\|\mathbf{\Psi}_{1}^{\mathbf{2}}\right\|_{L_{T}^{1} H^{-2 \alpha}} \lesssim T\left\|\mathbf{\Psi}_{1}^{\mathbf{2}}\right\|_{L_{T}^{\infty} H^{-2 \alpha}}
$$

Then, by applying Lemma 3.3, we obtain that

$$
\left\|\mathcal{G}\left(\boldsymbol{\Psi}_{1}^{\mathbf{2}}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}} \lesssim\left\|\boldsymbol{\Psi}_{1}^{\mathbf{2}}\right\|_{L_{T}^{1} H^{-2 \alpha}} \lesssim T\left\|\boldsymbol{\Psi}_{1}^{\mathbf{2}}\right\|_{L_{T}^{\infty} H^{-2 \alpha}}
$$

and we have thus shown that

$$
\left\|\mathcal{G}\left(\mathbf{\Psi}_{1}^{\mathbf{2}}\right)\right\|_{X(T)} \lesssim T\left\|\mathbf{\Psi}_{1}^{\mathbf{2}}\right\|_{L_{T}^{\infty} H^{-2 \alpha}}
$$

Combining the above estimates provides us with 4.10. It is then easy to see that 4.11 can be derived from similar arguments.
4.4.2. Proof of Proposition 4.8 when $d=2$.

In this situation, we pick $\kappa$ such that $3 \alpha<\kappa<\frac{1}{2}-2 \alpha$ and $(p, q):=(4,4)$, so that the space under consideration becomes

$$
X(T):=\mathcal{C}\left([0, T] ; H^{-2 \alpha}\left(\mathbb{R}^{2}\right)\right) \cap L^{4}\left([0, T] ; \mathcal{W}^{-2 \alpha, 4}\left(\mathbb{R}^{2}\right)\right) \cap L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}
$$

Note in particular that the so-defined pair $(p, q)$ is Schrödinger admissible. Also, as in the previous section, we set $\theta:=\frac{2 \alpha}{\kappa} \in\left(0, \frac{2}{3}\right)$, and we only focus on the derivation of (4.10) (estimate (4.11) could be obtained along the same arguments).

Bound on $S(\phi)$ : The arguments are exactly the same as for $d=1$ (see Section4.4.1), and yield

$$
\|S(\phi)\|_{X(T)} \lesssim\|\phi\|_{H^{-2 \alpha}}
$$

Bound on $\mathcal{G}\left(\rho^{2}|v|^{2}\right)$ : Since $\kappa>0$, and since $\rho$ is smooth and compactly-supported, one has

$$
\begin{aligned}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} & =\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{4} \mathcal{W}^{-2 \alpha, 4}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}} \\
& \lesssim\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{4} \mathcal{W}-2 \alpha+\kappa, 4}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\frac{1}{k}} H^{-2 \alpha+\kappa}} \\
& \lesssim\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{4} \mathcal{W}^{-2 \alpha+\kappa, 4}},
\end{aligned}
$$

and from here we can apply Strichartz inequality (Lemma 3.3) to assert that

$$
\begin{equation*}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} \lesssim\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{r^{\prime}} \mathcal{W}^{-2 \alpha+\kappa, s^{\prime}}} \tag{4.16}
\end{equation*}
$$

where we define the (Schrödinger admissible) pair $(r, s)$ along the formula

$$
(r, s):=\left(\frac{4}{1+\theta}, \frac{4}{1-\theta}\right)
$$

By Lemma 3.1. one has, for every fixed $t \geq 0$,

$$
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, s^{\prime}}} \lesssim\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}\|\rho v(t, .)\|_{L^{\frac{4}{1+\theta}}} .
$$

Besides, by using Lemma 4.2, one can check that

$$
\|\rho v(t, .)\|_{L^{\frac{4}{1+\theta}}} \leq\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}^{\theta}\|\rho v(t, .)\|_{\mathcal{W}^{-2 \alpha, 4}}^{1-\theta} .
$$

Therefore, for every fixed $t \geq 0$,

$$
\begin{aligned}
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, s^{\prime}}} & \lesssim\|\rho v(t, .)\|_{H-2 \alpha+\kappa}^{1+\theta}\|\rho v(t, .)\|_{\mathcal{W}^{-2 \alpha, 4}}^{1-\theta} \\
& \lesssim\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{1+\theta}\|v(t, .)\|_{\mathcal{W}^{-2 \alpha, 4}}^{1+\theta}+\|v(t,)\|_{H^{-2 \alpha}}^{1+\theta}\|v(t,)\|_{\mathcal{W}^{-2 \alpha, 4}}^{1-\theta},
\end{aligned}
$$

where we have used Lemma 4.4 to derive the second inequality.
Then, taking $\lambda:=\frac{3-\theta}{2}>1$, we get

$$
\begin{aligned}
& \int_{0}^{T} d t\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, s^{\prime}}}^{r^{\prime}} \\
& \lesssim\left(\int_{0}^{T} d t\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{(1+\theta) r^{\prime} \lambda}\right)^{\frac{1}{\lambda}}\left(\int_{0}^{T} d t\|v(t,)\|_{\mathcal{W}^{-2 \alpha, 4}}^{(1-\theta) r^{\prime} \lambda^{\prime}}\right)^{\frac{1}{\lambda^{\prime}}}+\|v\|_{X(T)}^{r^{\prime}(1+\theta)} \int_{0}^{T} d t\left\|v_{t}\right\|_{\mathcal{W}^{-2 \alpha, 4}}^{r^{\prime}(1-\theta)} .
\end{aligned}
$$

With our choices of parameters (remember that $\kappa<\frac{1}{2}-2 \alpha$ and $\theta=\frac{2 \alpha}{\kappa}$ ), one has in fact

$$
(1+\theta) r^{\prime} \lambda=2(1+\theta)<\frac{1}{\kappa} \quad \text { and } \quad(1-\theta) r^{\prime} \lambda^{\prime}=4
$$

so that the above inequality yields

$$
\begin{aligned}
& \int_{0}^{T} d t\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, s^{\prime}}}^{r^{\prime}} \\
& \lesssim T^{\frac{1-2 \kappa(1+\theta)}{\lambda}}\left(\int_{0}^{T} d t\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{\frac{1}{\kappa}}\right)^{(1+\theta) r^{\prime} \kappa}\|v\|_{X(T)}^{\frac{4}{\lambda}}+T^{1-\frac{r^{\prime}(1-\theta)}{4}}\|v\|_{X(T)}^{r^{\prime}(1+\theta)}\|v\|_{X(T)}^{r^{\prime}(1-\theta)} \\
& \lesssim T^{\frac{1-2 \kappa-4 \alpha}{\lambda}}\|v\|_{X(T)}^{(1+\theta) r^{\prime}+\frac{4}{\lambda}}+T^{1-\frac{r^{\prime}(1-\theta)}{4}}\|v\|_{X(T)}^{2 r^{\prime}} .
\end{aligned}
$$

It is now easy to check that this estimate can be rephrased as

$$
\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{r^{\prime}} \mathcal{W}^{-2 \alpha+\kappa, s^{\prime}}} \lesssim\left\{T^{\varepsilon}+T^{\frac{1}{2}}\right\}\|v\|_{X(T)}^{2},
$$

with $\varepsilon:=\frac{1}{2}(1-2 \kappa-4 \alpha)$.

Going back to (4.16), we can conclude that

$$
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} \lesssim\left\{T^{\varepsilon}+T^{\frac{1}{2}}\right\}\|v\|_{X(T)}^{2} .
$$

Bound on $\mathcal{G}\left(\rho \bar{v} \cdot \Psi_{1}^{1}\right), \mathcal{G}\left(\rho v \cdot \overline{\Psi_{1}^{1}}\right)$ : Using the same arguments as for $d=1$, we can show first that

$$
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)\right\|_{L_{T}^{\frac{1}{k}} H_{\rho}^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{1}}\right)\right\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}} \lesssim\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{1} H^{-\alpha}},
$$

and then

$$
\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{1} H^{-\alpha}}+\left\|\rho v \cdot \overline{\boldsymbol{\Psi}_{1}^{1}}\right\|_{L_{T}^{1} H^{-\alpha}} \lesssim T^{1-\kappa}\left\|\boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)} .
$$

On the other hand, we can use Lemma 3.3 to assert that
$\left(\left\|\mathcal{G}\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\boldsymbol{\Psi}_{1}^{1}}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}\right)+\left(\left\|\mathcal{G}\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)\right\|_{L_{T}^{4} \mathcal{W}^{-2 \alpha, 4}}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\boldsymbol{\Psi}_{1}^{1}}\right)\right\|_{L_{T}^{4} \mathcal{W}^{-2 \alpha, 4}}\right)$ $\lesssim\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{1} H^{-2 \alpha}} \lesssim\left\|\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{1} H^{-\alpha}}$.
Combining the above estimates easily provides us with the desired bound

$$
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)\right\|_{X(T)}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\boldsymbol{\Psi}_{1}^{1}}\right)\right\|_{X(T)} \lesssim T^{1-\kappa}\left\|\boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)} .
$$

Bound on $\mathcal{G}\left(\boldsymbol{\Psi}_{1}^{\mathbf{2}}\right)$ : The arguments are exactly the same as for $d=1$ (see Section 4.4.1), and lead us to

$$
\left\|\mathcal{G}\left(\mathbf{\Psi}_{1}^{2}\right)\right\|_{X(T)} \lesssim T\left\|\boldsymbol{\Psi}_{1}^{2}\right\|_{L_{T}^{\infty} H^{-2 \alpha}}
$$

4.4.3. Proof of Proposition 4.8 when $d=3$.

In this situation, we pick $\kappa:=4 \alpha$ and $(p, q):=(2,6)$, so that one has

$$
X(T):=\mathcal{C}\left([0, T] ; H^{-2 \alpha}\left(\mathbb{R}^{3}\right)\right) \cap L^{2}\left([0, T] ; \mathcal{W}^{-2 \alpha, 6}\left(\mathbb{R}^{3}\right)\right) \cap L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-2 \alpha+\kappa}
$$

Observe here again that $(p, q)=(2,6)$ defines a Schrödinger admissible pair.
Bound on $S(\phi)$ : We can repeat the arguments used for $d=1,2$ to assert that

$$
\|S(\phi)\|_{X(T)} \lesssim\|\phi\|_{H^{-2 \alpha}\left(\mathbb{R}^{3}\right)} .
$$

Bound on $\mathcal{G}\left(\rho^{2}|v|^{2}\right)$ : Let us first write, just as for $d=2$,

$$
\begin{aligned}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} & =\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{2} \mathcal{W}^{-2 \alpha, 6}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\frac{1}{T}} H_{\rho}^{-2 \alpha+\kappa}} \\
& \lesssim\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{\infty} H^{-2 \alpha+\kappa}}+\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{L_{T}^{2} \mathcal{W}^{-2 \alpha+\kappa, 6}},
\end{aligned}
$$

and then apply Strichartz inequality (Lemma 3.3) to obtain

$$
\begin{equation*}
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} \lesssim\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{2} \mathcal{W}^{-2 \alpha+\kappa, \frac{6}{5}}} . \tag{4.17}
\end{equation*}
$$

By Lemma 3.1, one has, for every fixed $t \geq 0$,

$$
\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, 6} 5} \lesssim\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}\|\rho v(t, .)\|_{L^{3}} .
$$

Besides, by using Lemma 4.2, one can check that

$$
\|\rho v(t, .)\|_{L^{3}} \leq\|\rho v(t, .)\|_{H^{-2 \alpha+\kappa}}^{\frac{1}{2}}\|\rho v(t, .)\|_{\mathcal{W}^{-2 \alpha, 6}}^{\frac{1}{2}} .
$$

Therefore, for every fixed $t \geq 0$,

$$
\begin{aligned}
\left\|\rho^{2}|v|^{2}(t,)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, 6} 5} & \lesssim\|\rho v(t,)\|_{H^{-2 \alpha+\kappa}}^{\frac{3}{2}}\|\rho v(t,)\|_{\mathcal{W}^{-2 \alpha, 6}}^{\frac{1}{2}} \\
& \lesssim\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{\frac{3}{2}}\|v(t, .)\|_{\mathcal{W}^{-2 \alpha, 6}}^{\frac{1}{2}}+\|v(t, .)\|_{H^{-2 \alpha}}^{\frac{3}{2}}\|v(t,)\|_{\mathcal{W}^{-2 \alpha, 6}}^{\frac{1}{2}}
\end{aligned}
$$

where we have used Lemma 4.4 to derive the second inequality.
This entails

$$
\begin{aligned}
& \int_{0}^{T} d t\left\|\rho^{2}|v|^{2}(t, .)\right\|_{\mathcal{W}^{-2 \alpha+\kappa, \frac{6}{5}}}^{2} \\
& \lesssim\left(\int_{0}^{T} d t\|v(t, .)\|_{H_{\rho}^{-2 \alpha+\kappa}}^{6}\right)^{\frac{1}{2}}\left(\int_{0}^{T} d t\|v(t, .)\|_{\mathcal{W}^{-2 \alpha, 6}}^{2}\right)^{\frac{1}{2}}+\|v\|_{X(T)}^{3} \int_{0}^{T} d t\left\|v_{t}\right\|_{\mathcal{W}^{-2 \alpha, 6}} \\
& \lesssim\left\{T^{\frac{1}{2}(1-24 \alpha)}+T^{\frac{1}{2}}\right\}\|v\|_{X(T)}^{4},
\end{aligned}
$$

which can obviously be recast as

$$
\left\|\rho^{2}|v|^{2}\right\|_{L_{T}^{2} \mathcal{W}^{-2 \alpha+\kappa, \frac{6}{5}}} \lesssim T^{\frac{1}{4}(1-24 \alpha)}\|v\|_{X(T)}^{2} .
$$

Going back to (4.17), we can conclude that

$$
\left\|\mathcal{G}\left(\rho^{2}|v|^{2}\right)\right\|_{X(T)} \lesssim T^{\frac{1}{4}(1-24 \alpha)}\|v\|_{X(T)}^{2}
$$

Bound on $\mathcal{G}\left(\rho \bar{v} \cdot \mathbf{\Psi}_{1}^{1}\right), \mathcal{G}\left(\rho v \cdot \overline{\boldsymbol{\Psi}_{1}^{1}}\right)$ : Using the same arguments as for $d=2$ (note indeed that $-2 \alpha+\kappa=2 \alpha>\alpha)$, we get that

$$
\left\|\mathcal{G}\left(\rho \bar{v} \cdot \boldsymbol{\Psi}_{1}^{1}\right)\right\|_{X(T)}+\left\|\mathcal{G}\left(\rho v \cdot \overline{\mathbf{\Psi}_{1}^{1}}\right)\right\|_{X(T)} \lesssim T^{1-\kappa}\left\|\boldsymbol{\Psi}_{1}^{1}\right\|_{L_{T}^{\infty} \mathcal{W}^{-\alpha, r}}\|v\|_{X(T)} .
$$

Bound on $\mathcal{G}\left(\boldsymbol{\Psi}_{1}^{\mathbf{2}}\right)$ : Here again, the arguments are exactly the same as for $d=1,2$ and entail

$$
\left\|\mathcal{G}\left(\boldsymbol{\Psi}_{1}^{2}\right)\right\|_{X(T)} \lesssim T\left\|\boldsymbol{\Psi}_{1}^{2}\right\|_{L_{T}^{\infty} H^{-2 \alpha}} .
$$

### 4.5. Proof of Theorem 1.11.

With the statements of Proposition 1.6 and Theorem 4.5 in hand, we are of course in the same position as in Section 3.4, and so the desired properties follow again from an elementary combination of these results.

## 5. Appendix

We gather here the proofs of two technical lemmas that have been used in the analysis of the rough case, namely Lemma 4.3 and Lemma 4.4

### 5.1. Proof of Lemma 4.3.

The argument is based on an interpolation procedure, combined with the following result:
Proposition 5.1 (Constantin-Saut [8]). Fix $d \geq 1$. Let $\rho: \mathbb{R}^{d} \rightarrow \mathbb{R}$ be of the form $\left(\overline{\mathbf{F}_{\rho}}\right)$, $0 \leq \alpha \leq \frac{1}{2}$ and $0 \leq T \leq 1$. Assume that $\phi \in H^{-\alpha}\left(\mathbb{R}^{d}\right), F \in L^{1}\left([0, T] ; H^{-\alpha}\left(\mathbb{R}^{d}\right)\right)$, and consider the solution $u$ of the following inhomogeneous Schrödinger equation on $\mathbb{R}^{d}$

$$
\left\{\begin{array}{l}
\imath \partial_{t} u(t, x)-\Delta u(t, x)=F(t, x), \quad t \in[0, T], x \in \mathbb{R}^{d}, \\
u_{0}=\phi
\end{array}\right.
$$

Then it holds that

$$
\begin{equation*}
\|u\|_{L_{T}^{2} H_{\rho}^{-\alpha+\frac{1}{2}}} \lesssim\|\phi\|_{H^{-\alpha}}+\|F\|_{L_{T}^{1} H^{-\alpha}}, \tag{5.1}
\end{equation*}
$$

where the proportional constant only depends on $\rho$ and $\alpha$.

Proof of Lemma 4.3. For every fixed $t \in[0, T]$, we can use the commutator estimate (4.6) and the fact that $\kappa \leq 1$ to write

$$
\begin{equation*}
\|u(t, .)\|_{H_{\rho}^{-\alpha+\kappa}} \lesssim\|\rho u(t, .)\|_{H^{-\alpha+\kappa}}+\|u(t, .)\|_{H^{-\alpha}} \tag{5.2}
\end{equation*}
$$

Then, by a standard interpolation argument (see Lemma 4.2), we get that

$$
\begin{aligned}
\|\rho u(t, .)\|_{H^{-\alpha+\kappa}} & \lesssim\|\rho u(t, .)\|_{H^{-\alpha}}^{1-2 \kappa}\|\rho u(t, .)\|_{H^{-\alpha+\frac{1}{2}}}^{2 \kappa} \\
& \lesssim\|u(t, .)\|_{H^{-\alpha}}^{1-2 \kappa}\|u(t, .)\|_{H_{\rho}^{-\alpha+\frac{1}{2}}}^{2 \kappa}+\|u(t, .)\|_{H^{-\alpha}}
\end{aligned}
$$

where we have used (4.6) to derive the second inequality.
By injecting the latter bound into 5.2 , we deduce that

$$
\|u\|_{L_{T}^{\frac{1}{\kappa}} H_{\rho}^{-\alpha+\kappa}} \lesssim\|u\|_{L_{T}^{\infty} H^{-\alpha}}^{1-2 \kappa}\|u\|_{L_{T}^{2} H_{\rho}^{-\alpha+\frac{1}{2}}}^{2 \kappa}+\|u\|_{L_{T}^{\infty} H^{-\alpha}}
$$

Using estimates $(3.4)$ to bound $\|u\|_{L_{T}^{\infty} H^{-\alpha}}$, and then estimate (5.1) to bound $\|u\|_{L_{T}^{2} H_{\rho}^{-\alpha+\frac{1}{2}}}$, we get the desired inequality (4.4).

### 5.2. Proof of Lemma 4.4.

In order to establish this commutator estimate, we can essentially follow the arguments of Kato and Ponce in their proof of [32, Lemma X1]. However, in the specific case where $\rho$ is a test function (which is the situation we would like to handle here), the bound (4.5) is clearly sharper than the general estimate in [32, Lemma X1]. For this reason, let us briefly review the main modifications leading to 4.5. The bound 4.5 also follows from the theory of pseudodifferential operators (see e.g. [1]) but the proof below only relies on the classical estimate by Coifman and Meyer ([6]).
Proposition 5.2. Fix $d \geq 1$ and consider a function $\sigma \in \mathcal{C}^{\infty}\left(\left(\mathbb{R}^{d} \times \mathbb{R}^{d}\right) \backslash(0,0) ; \mathbb{R}\right)$ satisfying

$$
\begin{equation*}
\left|\partial_{\xi}^{\alpha} \partial_{\eta}^{\beta} \sigma(\xi, \eta)\right| \leq C_{\alpha, \beta}(|\xi|+|\eta|)^{-|\alpha|-|\beta|} \tag{5.3}
\end{equation*}
$$

for all $(\xi, \eta) \neq(0,0)$ and $\alpha, \beta \in \mathbb{N}^{d}$. Let us denote by $B(\sigma)$ the bilinear operator defined for all test functions $\varphi, \psi: \mathbb{R}^{d} \rightarrow \mathbb{R}$ as

$$
B(\sigma)(\varphi, \psi)(x)=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta e^{i\langle x, \xi+\eta\rangle} \sigma(\xi, \eta) \mathcal{F} \varphi(\xi) \mathcal{F} \psi(\eta)
$$

Then it holds that

$$
\|B(\sigma)(\varphi, \psi)\|_{L^{2}\left(\mathbb{R}^{d}\right)} \lesssim\|\varphi\|_{L^{\infty}\left(\mathbb{R}^{d}\right)}\|\psi\|_{L^{2}\left(\mathbb{R}^{d}\right)}
$$

where the proportional constant only depends on the coefficients $\left(C_{\alpha, \beta}\right)_{\alpha, \beta \in \mathbb{N}^{d}}$ in (5.3).
Proof of Lemma 4.4. The quantity under consideration can be written as

$$
\begin{aligned}
& \left\|(\operatorname{Id}-\Delta)^{\frac{s}{2}}(\rho \cdot g)-\rho \cdot(\operatorname{Id}-\Delta)^{\frac{s}{2}}(g)\right\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{2} \\
& =c \int_{\mathbb{R}^{d}} d \xi\left|\int_{\mathbb{R}^{d}} d \eta\left[\left\{1+|\xi|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right] \mathcal{F} \rho(\xi-\eta) \mathcal{F} g(\eta)\right|^{2}
\end{aligned}
$$

Let us introduce a smooth function $\Phi: \mathbb{R} \rightarrow[0,1]$ with support in $\left[-\frac{1}{3}, \frac{1}{3}\right]$ such that $\Phi=1$ on $\left[-\frac{1}{4}, \frac{1}{4}\right]$. Then bound the above integral as

$$
\int_{\mathbb{R}^{d}} d \xi\left|\int_{\mathbb{R}^{d}} d \eta\left[\left\{1+|\xi|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right] \mathcal{F} \rho(\xi-\eta) \mathcal{F} g(\eta)\right|^{2} \lesssim \mathcal{J}_{1}+\mathcal{J}_{2}
$$

where

$$
\mathcal{J}_{1}:=\int_{\mathbb{R}^{d}} d \xi\left|\int_{\mathbb{R}^{d}} d \eta(1-\Phi)\left(\frac{|\xi-\eta|}{|\eta|}\right)\left[\left\{1+|\xi|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right] \mathcal{F} \rho(\xi-\eta) \mathcal{F} g(\eta)\right|^{2}
$$

and

$$
\mathcal{J}_{2}:=\int_{\mathbb{R}^{d}} d \xi\left|\int_{\mathbb{R}^{d}} d \eta \Phi\left(\frac{|\xi-\eta|}{|\eta|}\right)\left[\left\{1+|\xi|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right] \mathcal{F} \rho(\xi-\eta) \mathcal{F} g(\eta)\right|^{2} .
$$

Bound on $\mathcal{J}_{1}$. Using Cauchy-Schwarz inequality, we get first

$$
\begin{equation*}
\mathcal{J}_{1} \leq\|g\|_{H^{s-1}\left(\mathbb{R}^{d}\right)}^{2}\left(\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta\left|(1-\Phi)\left(\frac{|\xi-\eta|}{|\eta|}\right)\right|^{2}|\mathcal{F} \rho(\xi-\eta)|^{2} \frac{\left|\left\{1+|\xi|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right|^{2}}{\left\{1+|\eta|^{2}\right\}^{s-1}}\right) . \tag{5.4}
\end{equation*}
$$

In order to show that the latter integral is indeed finite, observe that if $(1-\Phi)(|\xi-\eta| /|\eta|) \neq 0$, then $|\xi-\eta| \geq \frac{1}{4}|\eta|$, and so $|\xi-\eta| \geq \frac{1}{5}|\xi|$. Therefore, as $\rho$ is smooth and compactly-supported, one has, for all $\lambda, \beta \geq 0$,

$$
\left|(1-\Phi)\left(\frac{|\xi-\eta|}{|\eta|}\right)\right|^{2}|\mathcal{F} \rho(\xi-\eta)|^{2} \leq c_{\rho, \lambda, \beta}\left\{1+|\xi|^{2}\right\}^{-\lambda}\left\{1+|\eta|^{2}\right\}^{-\beta},
$$

and the finiteness of the integral in (5.4) immediately follows.
Bound on $\mathcal{J}_{2}$. By Fourier isometry, we can write this quantity as $\mathcal{J}_{2}=c\|F\|_{L^{2}\left(\mathbb{R}^{d}\right)}^{2}$, with

$$
F(x):=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta e^{\imath\langle x, \xi+\eta\rangle} \Phi\left(\frac{|\xi|}{|\eta|}\right)\left[\left\{1+|\xi+\eta|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right] \mathcal{F} \rho(\xi) \mathcal{F} g(\eta) .
$$

Then

$$
\begin{aligned}
& \Phi\left(\frac{|\xi|}{|\eta|}\right)\left[\left\{1+|\xi+\eta|^{2}\right\}^{\frac{s}{2}}-\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\right] \mathcal{F} \rho(\xi) \mathcal{F} g(\eta) \\
& =\Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{s}{2}}\left[\left(1+\frac{\langle\xi, \xi+2 \eta\rangle}{1+|\eta|^{2}}\right)^{\frac{s}{2}}-1\right] \mathcal{F} \rho(\xi) \mathcal{F} g(\eta) \\
& =\Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}}\left[\left(1+\frac{\langle\xi, \xi+2 \eta\rangle}{1+|\eta|^{2}}\right)^{\frac{s}{2}}-1\right] \mathcal{F} \rho(\xi) \mathcal{F}\left((\operatorname{Id}-\Delta)^{\frac{s-1}{2}} g\right)(\eta) .
\end{aligned}
$$

At this point, observe that if $\Phi\left(\frac{|\xi|}{|\eta|}\right) \neq 0$, then $|\xi| \leq \frac{1}{3}|\eta|$, and so $|\langle\xi, \xi+2 \eta\rangle| \leq \frac{7}{9}|\eta|^{2}$. Therefore, we can rely on the pointwise expansion

$$
\Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}}\left[\left(1+\frac{\langle\xi, \xi+2 \eta\rangle}{1+|\eta|^{2}}\right)^{\frac{s}{2}}-1\right]=\sum_{k \geq 1} a_{k}(s) \Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}-k}\langle\xi, \xi+2 \eta\rangle^{k},
$$

where $a_{k}(s):=\frac{s(s-1) \cdots(s-k+1)}{k!}$.
Since $s>0$ and $\rho, g$ are assumed to be test functions, one has

$$
\begin{aligned}
& \int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta \sum_{k \geq 1}\left|a_{k}(s) \Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}-k}\langle\xi, \xi+2 \eta\rangle^{k} \mathcal{F} \rho(\xi) \mathcal{F}\left((\operatorname{Id}-\Delta)^{\frac{s-1}{2}} g\right)(\eta)\right| \\
& \left.\lesssim\left(\sum_{k \geq 1}\left|a_{k}(s)\right|\right) \int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}}|\mathcal{F} \rho(\xi)| \right\rvert\, \mathcal{F}\left(\left.\left((\operatorname{Id}-\Delta)^{\frac{s-1}{2}} g\right)(\eta) \right\rvert\,<\infty,\right.
\end{aligned}
$$

and accordingly we can write

$$
\begin{aligned}
& F(x)=\sum_{k \geq 1} a_{k}(s) \int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta e^{\imath\langle x, \xi+\eta\rangle} \Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}-k}\langle\xi, \xi+2 \eta\rangle^{k} \mathcal{F} \rho(\xi) \mathcal{F}\left((\operatorname{Id}-\Delta)^{\frac{s-1}{2}} g\right)(\eta) \\
& =\sum_{k \geq 1} a_{k}(s) \sum_{i=1}^{d} \int_{\mathbb{R}^{d} \times \mathbb{R}^{d}} d \xi d \eta e^{\langle\langle x, \xi+\eta\rangle} \Phi\left(\frac{|\xi|}{|\eta|}\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}-k}\langle\xi, \xi+2 \eta\rangle^{k-1}\left(\xi_{i}+2 \eta_{i}\right)\left[\xi_{i} \mathcal{F} \rho(\xi)\right] \mathcal{F}\left((\operatorname{Id}-\Delta)^{\frac{s-1}{2}} g\right)(\eta) .\right.
\end{aligned}
$$

Using the notation of Proposition 5.2 and the fact that $\xi_{i} \mathcal{F} \rho(\xi)=\imath \mathcal{F}\left(\partial_{x_{i}} \rho\right)(\xi)$, the latter identity can be rephrased as

$$
F(x)=\imath \sum_{k \geq 1} a_{k}(s) \sum_{i=1}^{d} B\left(\sigma_{k, i}\right)\left(\partial_{x_{i}} \rho,(\mathrm{Id}-\Delta)^{\frac{s-1}{2}} g\right)
$$

with

$$
\sigma_{k, i}(\xi, \eta):=\Phi\left(\frac{|\xi|}{|\eta|}\right)\left\{1+|\eta|^{2}\right\}^{\frac{1}{2}-k}\langle\xi, \xi+2 \eta\rangle^{k-1}\left(\xi_{i}+2 \eta_{i}\right) .
$$

It is not hard to check that for all $k \geq 1$ and $1 \leq i \leq d$, the function $\sigma_{k, i}$ satisfies condition (5.3) with coefficients $\left(C_{\alpha, \beta}\right)_{\alpha, \beta \in \mathbb{N}^{d}}$ independent of $k$ and $i$. Consequently, we are in a position to apply Proposition 5.2 and conclude that

$$
\|F\|_{L^{2}\left(\mathbb{R}^{d}\right)} \lesssim \sum_{k \geq 1}\left|a_{k}(s)\right| \sum_{i=1}^{d}\left\|\partial_{x_{i}} \rho\right\|_{L^{\infty}\left(\mathbb{R}^{d}\right)}\|g\|_{H^{s-1}\left(\mathbb{R}^{d}\right)} \lesssim\|g\|_{H^{s-1}\left(\mathbb{R}^{d}\right)}
$$

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