# UNCONDITIONAL UNIQUENESS RESULTS FOR THE NONLINEAR SCHRÖDINGER EQUATION 

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#### Abstract

We study the problem of unconditional uniqueness of solutions to the cubic nonlinear Schrödinger equation. We introduce a new strategy to approach this problem on bounded domains, in particular on rectangular tori.

It is a known fact that solutions to the cubic NLS give rise to solutions of the Gross-Pitaevskii hierarchy, which is an infinite-dimensional system of linear equations. By using the uniqueness analysis of the Gross-Pitaevskii hierarchy, we obtain new unconditional uniqueness results for the cubic NLS on rectangular tori, which cover the full scaling-subcritical regime in high dimensions. In fact, we prove a more general result which is conditional on the domain.

In addition, we observe that well-posedness of the cubic NLS in FourierLebesgue spaces implies unconditional uniqueness.


## 1. Introduction

Consider the (defocusing, cubic) nonlinear Schrödinger equation (NLS)

$$
\left\{\begin{array}{l}
i \partial_{t} \phi+\Delta \phi=|\phi|^{2} \phi  \tag{1}\\
\left.\phi\right|_{t=0}=\phi_{0}
\end{array}\right.
$$

for given initial data $\phi_{0} \in H^{s}\left(\mathbb{R}^{d}\right)$. The standard strategy of proving well-posedness of the initial value problem is similar to that for the Picard-Lindelöf theorem for ODEs. One applies the contraction mapping principle in some suitable Banach space $X \hookrightarrow C\left([-T, T], H^{s}\left(\mathbb{R}^{d}\right)\right)$ to solve the corresponding integral equation. Given rough initial data, i.e. if $s<d / 2$, typical choices for $X$ are mixed Lebesgue spaces (Strichartz norms) [6] or Fourier restriction spaces [2]. As a consequence, this argument yields uniqueness of solutions in $X$, not unconditionally in $C\left([-T, T], H^{s}\left(\mathbb{R}^{d}\right)\right)$. Tosio Kato $[26,27]$ was the first to address this problem, see Subsection 1.2 for further results and references. The purpose of this paper is to introduce new strategies which apply in the setting of certain bounded spatial domains.
1.1. Setup of the problem and main results. We are considering a spatial domain $\mathbb{D}$ which is either a $d$-dimensional rectangular torus $\mathbb{T}_{g}^{d}$ or, more generally, a smooth compact Riemannian manifold of dimension $d$ without boundary which is connected and orientable. In particular, on $\mathbb{D}$, we have a well-defined notion of a positive self-adjoint Laplacian $-\Delta \equiv-\Delta_{\mathbb{D}}$ and we can consider (1) on the

[^0]spatial domain $\mathbb{D}$. We address the question of unconditional uniqueness of mild solutions of (1). All of the results that we discuss are valid for arbitrary defocusing and focusing coupling constants in front of the cubic nonlinearity, but we set the coupling constant to be 1 in (1) for simplicity of notation.

Before we give the precise definition of a mild solution of (1) we recall the definition of $L^{2}$-based Sobolev spaces on $\mathbb{D}$. The operator $-\Delta$ has a discrete spectrum $0=\lambda_{0}<\lambda_{1}<\cdots<\lambda_{k} \rightarrow \infty$. Moreover, $\mathcal{H}_{j}$ denotes the eigenspace corresponding to eigenvalue $\lambda_{j}$,

$$
L^{2}(\mathbb{D})=\bigoplus_{j=0}^{\infty} \mathcal{H}_{j}
$$

Given $k \in \mathbb{N}$, we let $\Lambda_{k}:=\left\{j \in \mathbb{N}_{0}: \sqrt{\left|\lambda_{j}\right|} \in[k-1, k)\right\}$, and $\chi_{k}$ be the orthogonal projection $\chi_{k}: L^{2}(\mathbb{D}) \rightarrow \bigoplus_{j \in \Lambda_{k}} \mathcal{H}_{j}$. Given $s \in \mathbb{R}$, using $\langle x\rangle:=\sqrt{1+|x|^{2}}$, we define the Sobolev space $H^{s} \equiv H^{s}(\mathbb{D})$ via

$$
\begin{equation*}
\|f\|_{H^{s}}:=\left(\sum_{k=0}^{\infty}\left\langle\lambda_{k}\right\rangle^{s}\left\|\chi_{k} f\right\|_{L^{2}}^{2}\right)^{1 / 2} \tag{2}
\end{equation*}
$$

Definition 1.1 (Mild solution of the NLS). Given $s \geq 0, T>0$, and $\phi_{0} \in H^{s}(\mathbb{D})$, we say that $\phi=\phi(t, x)$ is a mild solution of (1) in $H^{s}$ with initial data $\phi_{0}$ if $\phi \in L_{t \in[-T, T]}^{\infty} H_{x}^{s} \cap L^{3}([-T, T] \times \mathbb{D})$ and the following conditions hold:
(i) For $t \in[-T, T]$ we have

$$
\begin{equation*}
\phi(x, t)=e^{i t \Delta} \phi_{0}(x)-i \int_{0}^{t} e^{i(t-\tau) \Delta}|\phi|^{2} \phi(x, \tau) d \tau \tag{3}
\end{equation*}
$$

(ii) There exists $R>0$ such that

$$
\begin{equation*}
\sup _{t \in[-T, T]}\|\phi(t)\|_{H^{s}(\mathbb{D})} \leq R \tag{4}
\end{equation*}
$$

(iii) For $t \in[-T, T]$ we have

$$
\begin{equation*}
\|\phi(t)\|_{L^{2}(\mathbb{D})}=\left\|\phi_{0}\right\|_{L^{2}(\mathbb{D})} \tag{5}
\end{equation*}
$$

Remark 1.2. Due to our assumption $\phi \in L^{3}([-T, T] \times \mathbb{D})$ the nonlinearity $|\phi|^{2} \phi$ is integrable and defines a distribution on $[-T, T] \times \mathbb{D}$. If we have the Sobolev embedding $H^{\frac{d}{6}}(\mathbb{D}) \subset L^{3}(\mathbb{D})$ on the domain $\mathbb{D}$ (e.g. on rectangular tori), then $\phi \in$ $L^{3}([-T, T] \times \mathbb{D})$ is implied by $\phi \in L_{t \in[-T, T]}^{\infty} H_{x}^{s}$ if $s \geq \frac{d}{6}$. We always have

$$
\begin{equation*}
|\phi|^{2} \phi \in L_{t \in[-T, T]}^{\infty} H_{x}^{\beta} \text { for } \beta<-\frac{d}{2} \tag{6}
\end{equation*}
$$

in particular, the integral in the right hand-side of (3) is taken in $H_{x}^{\beta}$. Note that, a posteriori, (3) tells us that the function defined by the integral belongs to $H_{x}^{s}$.

For dyadic integers $N=2^{j}$ we define

$$
P_{N}:=\sum_{k \in[N, 2 N)} \chi_{k} .
$$

Moreover we define $P_{0}:=\chi_{0}$. In particular, we have

$$
\begin{equation*}
\sum_{N} P_{N}=\mathrm{Id} \tag{7}
\end{equation*}
$$

where $\sum_{N}$ means that we are summing over all dyadic integers $N$ and $N=0$. By (2), we have

$$
\begin{equation*}
\|f\|_{H^{s}} \sim\left(\sum_{n=0}^{\infty}\langle n\rangle^{2 s}\left\|\chi_{n} f\right\|_{L^{2}}^{2}\right)^{1 / 2} \sim\left(\sum_{N}\langle N\rangle^{2 s}\left\|P_{N} f\right\|_{L^{2}}^{2}\right)^{1 / 2} \tag{8}
\end{equation*}
$$

Furthermore we define the Besov space $B_{2,1}^{s} \equiv B_{2,1}^{s}(\mathbb{D})$

$$
\begin{equation*}
\|f\|_{B_{2,1}^{s}}:=\sum_{N}\langle N\rangle^{s}\left\|P_{N} f\right\|_{L^{2}} \tag{9}
\end{equation*}
$$

For all $\delta>0$ we have the inclusion

$$
\begin{equation*}
H^{s+\delta} \hookrightarrow B_{2,1}^{s} \hookrightarrow H^{s} \tag{10}
\end{equation*}
$$

Motivated by the multi-linear analysis of [5, 20, 23] and the references therein we define the following.

Definition 1.3. We say that the domain $\mathbb{D}$ is admissible if there exist constants $\zeta_{0} \equiv \zeta_{0}(\mathbb{D}), \alpha_{0} \equiv \alpha_{0}(\mathbb{D}), \epsilon \equiv \epsilon(\mathbb{D})>0$ such that for any fixed

$$
\begin{equation*}
0 \leq \eta \leq \zeta_{0}, \quad \zeta>\zeta_{0}, \quad \alpha>\alpha_{0} \tag{11}
\end{equation*}
$$

the following estimates hold: For all $T \in[0,1]$ and $u_{j}(t)=e^{i t \Delta} \phi_{j}, j=1,2,3$,

$$
\begin{align*}
\left\|u_{1} u_{2} u_{3}\right\|_{L_{t}^{1} B_{2,1}^{-\eta}([-T, T] \times \mathbb{D})} & \lesssim T^{\epsilon}\left\|\phi_{1}\right\|_{B_{2,1}^{-\eta}}\left\|\phi_{2}\right\|_{B_{2,1}^{\zeta}}\left\|\phi_{3}\right\|_{B_{2,1}^{\zeta}}  \tag{12}\\
\left\|u_{1} u_{2} u_{3}\right\|_{L_{t}^{1} B_{2,1}^{\zeta}([-T, T] \times \mathbb{D})} & \lesssim T^{\epsilon}\left\|\phi_{1}\right\|_{B_{2,1}^{\zeta}}\left\|\phi_{2}\right\|_{B_{2,1}^{\zeta}}\left\|\phi_{3}\right\|_{B_{2,1}^{\zeta}}  \tag{13}\\
\left\|f_{1} f_{2} f_{3}\right\|_{B_{2,1}^{-\zeta}} & \lesssim\left\|f_{1}\right\|_{H^{\alpha}}\left\|f_{2}\right\|_{H^{\alpha}}\left\|f_{3}\right\|_{H^{\alpha}} . \tag{14}
\end{align*}
$$

For $\mathbb{D}$ an admissible domain we let

$$
\begin{equation*}
s_{0}:=\max \left\{\zeta_{0}, \alpha_{0}, d / 4\right\} \tag{15}
\end{equation*}
$$

Remark 1.4. On many domains (13) is a direct consequence of (12), see e.g. Section 3 for details in the case of tori.

We now state our main results.
Theorem 1.5 (Unconditional uniqueness for admissible domains). Let $\mathbb{D}$ be an admissible domain with parameter $s_{0}$ given by (15) and $s>s_{0}$. For any $\phi_{0} \in H^{s}$ there is a unique mild solution of (1) in $H^{s}$.
Remark 1.6. We refer to the uniqueness result in Theorem 1.5 as unconditional in the sense that we are only assuming the control of the $H^{s}$ norm of the solution as in Definition 1.1 above. In particular, we do not assume that the solution belongs to some Strichartz space.

We can view Theorem 1.5 as a result that is conditional on the domain. The condition is that $\mathbb{D}$ is admissible in the sense of Definition 1.3 above. Once this assumption is satisfied, the uniqueness analysis follows by the general argument, which consists of comparing mild solutions of (1) to those of the Gross-Pitaevskii (GP) hierarchy on $\mathbb{D}$, which is defined in (19) below. In particular, we use the observation that tensor products of mild solutions of the NLS are mild solutions of the GP hierarchy, which is a well-known fact (for classical solutions). In light of this observation, one typically considers the GP hierarchy as a generalization of the NLS and proves analogues of results known for the NLS in the context of the hierarchy. This point of view was first taken in [9]. Our work is the first instance in
which we can use knowledge from the GP hierarchy to answer nontrivial questions about the NLS.

In Section 3, we show that rectangular tori of dimension $d \geq 2$ are admissible in the sense of Definition 1.3. We now state the explicit unconditional uniqueness result that we obtain from Proposition 3.1.

Corollary 1.7. Solutions to (1) are unconditionally unique in $H^{s}\left(\mathbb{T}_{g}^{d}\right)$ provided that

$$
\begin{cases}s>\frac{7}{12} & \text { for } d=2  \tag{16}\\ s>\frac{4}{5} & \text { for } d=3 \\ s>\frac{d}{2}-1 & \text { for } d \geq 4\end{cases}
$$

Remark 1.8. The required regularity in Corollary 1.7 is below the energy class for $d=2,3$. For $d \geq 4$, (1) is energy-supercritical and (16) covers the full scalingsubcritical range.

The above method also applies in one dimension, i.e. in the case when $\mathbb{D}=\mathbb{T}$ and yields unconditional uniqueness in $H^{s}(\mathbb{T})$ for $s>\frac{1}{3}$. However, we can obtain a stronger result by a different method.

Theorem 1.9. Let $p>3$. For any $\phi_{0} \in L^{2}(\mathbb{T})$ there is a unique mild solution $\phi$ of (1) in $L^{2}(\mathbb{T})$ with the property $\phi \in L^{p}([-T, T] \times \mathbb{T})$.

Note that the above uniqueness result is in an almost scaling-invariant class and, by the Sobolev embedding, implies the unconditional uniqueness result in $H^{s}(\mathbb{T})$ for $s>\frac{1}{6}$. We emphazise that this result is strictly weaker than the result obtained in [19]. Indeed, the above uniqueness statement with $p=3$ follows from [19, Section 5], which is based on a normal form method with infinitely many iterations. In comparison, the proof of Theorem 1.9 is a easy consequence of an estimate in [18].
1.2. Previously known results. The problem of unconditional uniqueness of solutions to the NLS has been extensively studied in the literature. Using Sobolev embedding and Gronwall's inequality, one can easily obtain an unconditional uniqueness result the NLS when $s>d / 2$. Note that, for $d \geq 2$, this is already above the energy space. We henceforth consider only the nontrivial regime $s<d / 2$. The first such result on $\mathbb{R}^{d}$ is that of Kato [26, 27], with subsequent extensions and improvements in dimension $d \geq 2$ obtained in [7, 14, 15, 21, 31, 38, 43, 44], we refer to these references for precise results. The inherent losses of derivatives in the Strichartz estimate do not allow one to adapt these arguments to compact domains. The first unconditional uniqueness result on the one-dimensional torus was obtained in [19] by the use of normal form techniques and it applies to data in $H^{\frac{1}{6}}(\mathbb{T})$, in agreement with the result from $[26,27]$ in $H^{\frac{1}{6}}(\mathbb{R})$. In fact, in the case $\mathbb{D}=\mathbb{R}$, the result in $[26,27]$ gives uniqueness of mild solutions $\phi \in L_{t}^{\infty} L_{x}^{2} \cap L_{t}^{12} L_{x}^{3}$ and in the case $\mathbb{D}=\mathbb{T}$ the result in [19] applies to $\phi \in L_{t}^{\infty} L_{x}^{2} \cap L^{3}$. An extension of the techniques of [19] to higher dimensions has been announced by [29]. The announced regime agrees with the regime from Corollary 1.7 when $d \geq 6$. We note that the unconditional uniqueness problem was also studied for other dispersive models on $\mathbb{R}^{d}[4,13,15,32,33,37]$.

In the method of deriving the NLS from many-body quantum dynamics developed by Spohn [41], a central role is played by the GP hierarchy (19). The final step in this approach of deriving the NLS typically consists of showing uniqueness
of solutions to (19). This is a non-trivial problem due to the large number of terms that one has to consider. Several methods have been introduced to deal with this issue. They are a Feynman graph expansion [12], the boardgame argument [30], or the quantum de Finetti theorem [8]. This has been an extensively studied problem. We refer the reader to the introduction of $[22,40]$ for further references and for a more detailed discussion of the methods. In a recent preprint [1], the authors study connections between solutions of an initial value problem and an associated hierarchy in order to obtain results about the uniqueness of the hierarchy. This is the opposite direction from the one that we are taking in the proof of Theorem 1.5.

Our second method to obtain Theorem 1.9 is based on multi-linear estimates in restriction norms based on Fourier-Lebesgue spaces and the Hausdorff-Young inequality, we refer to Section 4 for more details on this. Here, while the result is not new, the general observation is that the analysis of dispersive PDEs on tori in Fourier-Lebesgue spaces yields unconditional uniqueness results in $L^{2}$-based spaces.
1.3. Organization of the paper. In Section 2, we recall the definition of the Gross-Pitaevskii hierarchy and we prove the general result given by Theorem 1.5. In Section 3, we apply Theorem 1.5 to the setting of rectangular tori and show the explicit unconditional uniqueness result stated in Corollary 1.7 above. Section 4 is devoted to the proof of Theorem 1.9. In Appendix A, we prove the multi-linear estimate given by Proposition 2.8, which is used in the proof of Theorem 1.5. In doing so, we recall several general Sobolev embedding results. In Appendix B, we give a proof of the Bernstein inequality on rectangular tori, which is stated in Lemma 3.3. In Appendix C, we give the proof of a key linear estimate used in the proof of Theorem 1.9.

## 2. The Gross-Pitaevskil hierarchy

In this section, we study the Gross-Pitaevskii hierarchy and prove Theorem 1.5. The required notation and definitions are given in Subsection 2.1. In Subsection 2.2 , we study tensor products of mild solutions of (1) and show that they solve the Gross-Pitaevskii hierarchy in an appropriate sense. Finally, we give the proof of Theorem 1.5 in Subsection 2.3.
2.1. Definition of the hierarchy. Given $k \in \mathbb{N}$ we call functions $\sigma^{(k)}: \mathbb{D}^{k} \times \mathbb{D}^{k} \rightarrow$ $\mathbb{C}$ density matrices of order $k$, and we denote this class by $\mathcal{M}_{k}$. In (19), for each fixed $t, \gamma^{(k)}(t) \in \mathcal{M}_{k}$. Given $\sigma^{(k+1)} \in \mathcal{M}_{k+1}$ we formally define $B_{k+1} \sigma^{(k+1)} \in \mathcal{M}_{k+1}$ by

$$
\begin{equation*}
B_{k+1} \sigma^{(k+1)}:=\sum_{j=1}^{k} B_{j, k+1}\left(\sigma^{(k+1)}\right), \tag{17}
\end{equation*}
$$

where, for $1 \leq j \leq k\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}\right) \in \mathbb{D}^{k} \times \mathbb{D}^{k}$ we define

$$
\begin{align*}
& B_{j, k+1}\left(\sigma^{(k+1)}\right)\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}\right):= \\
& \int_{\mathbb{D}}\left(\delta\left(x_{j}-x_{k+1}\right)-\delta\left(x_{j}^{\prime}-x_{k+1}\right)\right) \sigma^{(k+1)}\left(\vec{x}_{k}, x_{k+1} ; \vec{x}_{k}^{\prime}, x_{k+1}\right) d x_{k+1} \tag{18}
\end{align*}
$$

$B_{k+1}$ is called the collision operator. At this point we view (17) and (18) as formal definitions on the level of density matrices ${ }^{1}$. For our purposes we are going to look

[^1]at $\sigma^{(k+1)}$ of the form
$$
\sigma^{(k+1)}\left(\vec{x}_{k+1} ; \vec{x}_{k+1}^{\prime}\right)=\prod_{\ell=1}^{k+1} f_{\ell}\left(x_{\ell}\right) g_{\ell}\left(x_{\ell}^{\prime}\right)
$$
for some $f_{1}, \ldots, f_{k}, g_{1}, \ldots, g_{k} \in L^{2}(\mathbb{D})$ or superpositions thereof. All of the objects involving the collision operator that we will consider will be well-defined by using this a priori formal definition, see Lemmas 2.4 and 2.9 below. In principle it is possible to impose additional regularity assumptions on $\sigma^{(k+1)}$ under which these operations can be written in terms of an approximate identity (see [11, Theorem 1(iv)]). We do not take this approach here.

The (defocusing, cubic) Gross-Pitaevskii (GP) hierarchy on $\mathbb{D}$ is given by

$$
\left\{\begin{array}{l}
i \partial_{t} \gamma^{(k)}+\left(\Delta_{\vec{x}_{k}}-\Delta_{\vec{x}_{k}^{\prime}}\right) \gamma^{(k)}=B_{k+1} \gamma^{(k+1)}  \tag{19}\\
\left.\gamma^{(k)}\right|_{t=0}=\gamma_{0}^{(k)}
\end{array}\right.
$$

Here, $\gamma^{(k)}$ are time-dependent density matrices of order $k$. We always assume that the right-hand side of the first equation in (19) is well-defined in the sense described above.

On density matrices of order $k$ we define the Sobolev differentiation operators $S^{(k, \alpha)}$ and the free Schrödinger evolution operators $\mathcal{U}^{(k)}(t)$ in the standard way.

$$
\begin{gathered}
S^{(k, \alpha)}:=\prod_{j=1}^{k}\left\langle\Delta_{x_{j}}\right\rangle^{\alpha / 2} \prod_{j=1}^{k}\left\langle\Delta_{x_{j}^{\prime}}\right\rangle^{\alpha / 2} \\
\mathcal{U}^{(k)}(t) \sigma^{(k)}:=e^{i t \sum_{j=1}^{k} \Delta_{x_{j}}} \sigma^{(k)} e^{-i t \sum_{j=1}^{k} \Delta_{x_{j}^{\prime}}} .
\end{gathered}
$$

Definition 2.1 (The class $\left.\mathfrak{H}^{s}\right)$. Given $s \geq 0, \mathfrak{H}^{s} \equiv \mathfrak{H}^{s}(\mathbb{D})$ denotes the set of all sequences $\left(\gamma^{(k)}\right)_{k}$ of density matrices such that, for each $k \in \mathbb{N}$ the following properties hold:
(i) $\gamma^{(k)} \in L_{\text {sym }}^{2}\left(\mathbb{D}^{k} \times \mathbb{D}^{k}\right)$ in the sense that $\gamma^{(k)} \in L^{2}\left(\mathbb{D}^{k} \times \mathbb{D}^{k}\right)$ and

$$
\gamma^{(k)}\left(x_{\pi(1)}, \ldots, x_{\pi(k)} ; x_{\pi(1)}^{\prime}, \ldots, x_{\pi(k)}^{\prime}\right)=\gamma^{(k)}\left(x_{1}, \ldots, x_{k} ; x_{1}^{\prime}, \ldots, x_{k}^{\prime}\right)
$$

for all permutations $\pi \in \mathcal{S}^{k}$. Moreover,

$$
\gamma\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}\right)=\overline{\gamma^{(k)}\left(\vec{x}_{k}^{\prime} ; \vec{x}_{k}\right)}
$$

for all $\left(\vec{x}_{k}, \vec{x}_{k}^{\prime}\right)$ in $\mathbb{D}^{k} \times \mathbb{D}^{k}$.
(ii) $S^{(k, s)} \gamma^{(k)} \in \mathcal{L}_{k}^{1}$, i.e. $S^{(k, s)} \gamma^{(k)}$ belongs to the trace class on $L^{2}\left(\mathbb{D}^{k} \times \mathbb{D}^{k}\right)$, whose norm is given by $\|\mathcal{T}\|_{\mathcal{L}_{k}^{1}}:=\operatorname{Tr}|\mathcal{T}|$, where $|\mathcal{T}|:=\sqrt{\mathcal{T}^{*} \mathcal{T}}$.
We abbreviate (i) and (ii) as $\gamma^{(k)} \in \mathfrak{H}_{k}^{s}=\mathfrak{H}_{k}^{s}(\mathbb{D})$ and we define the norm on this space as

$$
\left\|\gamma^{(k)}\right\|_{\mathfrak{H}_{k}^{s}}:=\operatorname{Tr}\left|S^{(k, s)} \gamma^{(k)}\right| .
$$

Definition 2.2 (Mild solution of the Gross-Pitaevskii hierarchy). Given $s>\frac{d}{4}$, $T>0$, and $\left(\gamma_{0}^{(k)}\right)_{k} \in \mathfrak{H}^{s}$, we say that $\left(\gamma^{(k)}\right)_{k}=\left(\gamma^{(k)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right)\right)_{k}$ is a mild solution of (19) in $\mathfrak{H}^{s}$ with initial data $\left(\gamma_{0}^{(k)}\right)_{k}$ if the following conditions hold:
(i) For all $k \in \mathbb{N}$ we have

$$
\begin{equation*}
\gamma^{(k)} \in L_{t \in[-T, T]}^{\infty} \mathfrak{H}_{k}^{s} . \tag{20}
\end{equation*}
$$

(ii) For all $k \in \mathbb{N}$ and $t \in[-T, T]$ we have

$$
\begin{equation*}
\gamma^{(k)}(t)=\mathcal{U}^{(k)}(t) \gamma_{0}^{(k)}-i \int_{0}^{t} \mathcal{U}^{(k)}(t-s) B_{k+1} \gamma^{(k+1)}(s) d s \tag{21}
\end{equation*}
$$

(iii) There exists a constant $\mathcal{R}>0$ such that

$$
\begin{equation*}
\operatorname{Tr}\left(\left|S^{(k, s)} \gamma^{(k)}(t)\right|\right) \leq \mathcal{R}^{2 k} \tag{22}
\end{equation*}
$$

for all $k \in \mathbb{N}$ and $t \in[-T, T]$.
Remark 2.3. We note that, for $t \in[-T, T]$ equality in (3) is assumed to hold in $H^{s}$, whereas in (21) it is assumed to hold in $\mathfrak{H}_{k}^{s}$.
2.2. Factorized solutions of the Gross-Pitaevskii hierarchy. We first note a general relationship between mild solutions of (1) and (19).

Lemma 2.4 (Mild factorized solutions of the Gross-Pitaevskii hierarchy). Let $s \geq$ $0, T>0$ and $\phi_{0} \in H^{s}(\mathbb{D})$ be given. Suppose that $\phi$ is a mild solution of (1) in $H^{s}$ with initial data $\phi_{0}$. Then

$$
\begin{equation*}
\gamma^{(k)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right):=|\phi\rangle\left\langle\left.\phi\right|^{\otimes k}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right)\right. \tag{23}
\end{equation*}
$$

is a mild solution of (19) in $\mathfrak{H}^{s}$ with initial data $\left(\left|\phi_{0}\right\rangle\left\langle\left.\phi_{0}\right|^{\otimes k}\right)_{k}\right.$.
It is a known fact that tensor products of classical solutions to the NLS give us classical solutions of the Gross-Pitaevskii hierarchy. Lemma 2.4 tells us that the analogous claim is true when we consider mild solutions. Note that we cannot deduce Lemma 2.4 by a density argument from classical solutions, because this would in particular rely on unconditional uniqueness of solutions to (1), which is what we are proving. Instead, we have to argue directly and verify the assumptions of Definition 2.2 above. Before proving Lemma 2.4 we note a general algebraic fact that we will use in the proof.

Lemma 2.5. Given $m \in \mathbb{N}$, numbers $F_{1}, \ldots, F_{m} \in \mathbb{C}$ and locally integrable functions $G_{1}, \cdots, G_{m}: \mathbb{R} \rightarrow \mathbb{C}$, we have for all $t \in \mathbb{R}$

$$
\begin{align*}
& \prod_{r=1}^{m}\left(F_{r}+\int_{0}^{t} G_{r}\left(t_{r}\right) d t_{r}\right)= \\
& \sum_{\substack{\mathcal{S}_{1} \cup \mathcal{S}_{2}=\\
\{1,2, \ldots, m\}}}\left(\prod_{r \in \mathcal{S}_{1}} F_{r}\right)\left(\sum_{r \in \mathcal{S}_{2}} \int_{0}^{t} G_{r}\left(\tau_{r}\right)\left(\prod_{\tilde{r} \in \mathcal{S}_{2} \backslash\{r\}} \int_{0}^{\tau_{r}} G_{\tilde{r}}\left(\tau_{\tilde{r}}\right) d \tau_{\tilde{r}}\right) d \tau_{r}\right) . \tag{24}
\end{align*}
$$

Proof of Lemma 2.5. Let $\mathcal{S}_{2} \subset\{1,2, \ldots, m\}$ be nonempty. We note that

$$
\left(\prod_{r \in \mathcal{S}_{2}} \int_{0}^{t} G_{r}\left(\tau_{r}\right) d \tau_{r}\right)=\left(\sum_{r \in \mathcal{S}_{2}} \int_{0}^{t} G_{r}\left(\tau_{r}\right)\left(\prod_{\tilde{r} \in \mathcal{S}_{2} \backslash\{r\}} \int_{0}^{\tau_{r}} G_{\tilde{r}}\left(\tau_{\tilde{r}}\right) d \tau_{\tilde{r}}\right) d \tau_{r}\right) .
$$

The above identity follows by using Fubini's theorem for the integral on the lefthand side and by decomposing the region of integration $[0, t]^{\left|\mathcal{S}_{2}\right|}$ in terms of the maximal component $\tau_{r}$ with $r \in \mathcal{S}_{2}$. The claim now follows by expanding the product in the expression on the left-hand side of Lemma 2.5.

Proof of Lemma 2.4. From (4), we deduce the condition (20) of being a mild solution, as well as (22) with $\mathcal{R}=R$. We hence have to verify (21) for fixed $k \in \mathbb{N}$ and $t \in[-T, T]$. By using (3), it follows that

$$
\begin{align*}
& \gamma^{(k)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right)=\prod_{j=1}^{k}\left[\left(e^{i t \Delta} \phi_{0}\left(x_{j}\right)+\int_{0}^{t}-i e^{i\left(t-\tau_{j}\right) \Delta}|\phi|^{2} \phi\left(x_{j}, \tau_{j}\right) d \tau_{j}\right)\right. \\
& \left.\times\left(\overline{e^{i t \Delta} \phi_{0}\left(x_{j}^{\prime}\right)}+\int_{0}^{t} \overline{-i e^{i\left(t-\tau_{j}^{\prime}\right) \Delta}|\phi|^{2} \phi\left(x_{j}^{\prime}, \tau_{j}^{\prime}\right)} d \tau_{j}^{\prime}\right)\right] \tag{25}
\end{align*}
$$

We note that

$$
\begin{equation*}
\mathcal{U}^{(k)}(t) \gamma_{0}^{(k)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}\right)=\prod_{j=1}^{k} e^{i t \Delta} \phi_{0}\left(x_{j}\right) \overline{e^{i t \Delta} \phi_{0}\left(x_{j}^{\prime}\right)} \tag{26}
\end{equation*}
$$

We compute

$$
\begin{aligned}
& -i \int_{0}^{t} \mathcal{U}^{(k)}(t-\tau) B_{k+1} \gamma^{(k+1)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, \tau\right) d \tau \\
= & -i \int_{0}^{t} \mathcal{U}^{(k)}(t-\tau) \sum_{j=1}^{k}\left[\left(|\phi|^{2} \phi\left(x_{j}, \tau\right) \overline{\phi\left(x_{j}^{\prime}, \tau\right)}-\phi\left(x_{j}, \tau\right)|\phi|^{2} \bar{\phi}\left(x_{j}^{\prime}, \tau\right)\right) \prod_{\substack{\ell=1 \\
\ell \neq j}}^{k} \phi\left(x_{\ell}, \tau\right) \overline{\phi\left(x_{\ell}^{\prime}, \tau\right)}\right] d \tau \\
= & -i \int_{0}^{t} \sum_{j=1}^{k}\left[\left(e^{i(t-\tau) \Delta}|\phi|^{2} \phi\left(x_{j}, \tau\right) \overline{e^{i(t-\tau) \Delta} \phi\left(x_{j}^{\prime}, \tau\right)}-e^{i(t-\tau) \Delta} \phi\left(x_{j}, \tau\right) \overline{e^{i(t-\tau) \Delta}|\phi|^{2} \phi\left(x_{j}^{\prime}, \tau\right)}\right)\right. \\
& \left.\times \prod_{\substack{\ell=1 \\
\ell \neq j}}^{k} e^{i(t-\tau) \Delta} \phi\left(x_{\ell}, \tau\right) \overline{e^{i(t-\tau) \Delta} \phi\left(x_{\ell}^{\prime}, \tau\right)}\right] d \tau .
\end{aligned}
$$

We again use (3) and the semigroup property to deduce that the above expression equals

$$
\begin{align*}
& \int_{0}^{t} \sum_{j=1}^{k}\left\{\left[\left(-i e^{i(t-\tau) \Delta}|\phi|^{2} \phi\left(x_{j}, \tau\right)\right)\left(\overline{e^{i t \Delta} \phi_{0}\left(x_{j}^{\prime}\right)}+\int_{0}^{\tau} \overline{-i e^{i\left(t-\tau_{j}^{\prime}\right) \Delta}|\phi|^{2} \phi\left(x_{j}^{\prime}, \tau_{j}^{\prime}\right)} d \tau_{j}^{\prime}\right)\right.\right. \\
& \left.+\left(e^{i t \Delta} \phi_{0}\left(x_{j}\right)+\int_{0}^{\tau}-i e^{i\left(t-\tau_{j}\right) \Delta}|\phi|^{2} \phi\left(x_{j}, \tau_{j}\right) d \tau_{j}\right) \overline{-i e^{i(t-\tau) \Delta}|\phi|^{2} \phi\left(x_{j}^{\prime}, \tau\right)}\right] \\
& \times \prod_{\substack{\ell=1 \\
\ell \neq j}}^{k}\left[\left(e^{i t \Delta} \phi_{0}\left(x_{\ell}\right)+\int_{0}^{\tau}-i e^{i\left(t-\tau_{\ell}\right) \Delta}|\phi|^{2} \phi\left(x_{\ell}, \tau_{\ell}\right) d \tau_{\ell}\right)\right. \\
& \left.\left.\left(\overline{e^{i t \Delta} \phi_{0}\left(x_{\ell}^{\prime}\right)}+\int_{0}^{\tau} \overline{-i e^{i\left(t-\tau_{\ell}^{\prime}\right) \Delta}|\phi|^{2} \phi\left(x_{\ell}, \tau_{\ell}^{\prime}\right)} d \tau_{\ell}^{\prime}\right)\right]\right\} d \tau \tag{27}
\end{align*}
$$

The claim now follows from (25), (26), (27) if we apply Lemma 2.5 with $m=2 k$ and functions $F_{r}, G_{r}$ given by

$$
\left\{\begin{array}{l}
F_{r}:=e^{i t \Delta} \phi_{0}\left(x_{r}\right) \text { for } r \in\{1,2, \ldots, k\}  \tag{28}\\
F_{r}:=\overline{e^{i t \Delta} \phi_{0}\left(x_{r-k}^{\prime}\right)} \text { for } r \in\{k+1, k+2, \ldots, 2 k\} \\
G_{r}(\tau):=-i e^{i(t-\tau) \Delta}|\phi|^{2} \phi\left(x_{r}, \tau\right) \text { for } r \in\{1,2, \ldots, k\} \\
G_{r}(\tau):=\overline{-i e^{i(t-\tau) \Delta}|\phi|^{2} \phi\left(x_{r-k}^{\prime}, \tau\right)} \text { for } r \in\{k+1, k+2, \ldots, 2 k\}
\end{array}\right.
$$

We note that we are applying Lemma 2.5 to a time integral in functions that take values in the Sobolev space $H_{x}^{\beta}$, for $\beta$ as in (6) above. In particular, applying Lemma 2.5 at this step is well-justified.

Remark 2.6. We note that, in the proof of Lemma 2.4, we did not use condition (5) of being a mild solution of the NLS. This condition will be used crucially at several later points, see (32), (35), and (36) below.
2.3. The uniqueness analysis: Proof of Theorem 1.5. Throughout this subsection we consider $s>s_{0}$, for $s_{0}$ as in (15). It suffices to consider $s<\frac{d}{2}$. Moreover, we fix $\phi_{0} \in H^{s}(\mathbb{D})$ and consider two mild solutions $\phi^{(1)}, \phi^{(2)}$ of (1) in $H^{s}$ with initial data $\phi_{0}$ as in Definition 1.1. Our goal is to show that $\phi^{(1)}=\phi^{(2)}$. Before proving Theorem 1.5 we first note a uniqueness result for the factorized solutions (23) of (19).

Proposition 2.7. Let $\phi^{(1)}$, $\phi^{(2)}$ be mild solutions of (1) in $\mathfrak{H}^{s}$ with initial data $\phi_{0} \in H^{s}$. We have

$$
\begin{equation*}
\left|\phi^{(1)}\right\rangle\left\langle\left.\phi^{(1)}\right|^{\otimes k}=\mid \phi^{(2)}\right\rangle\left\langle\left.\phi^{(2)}\right|^{\otimes k},\right. \tag{29}
\end{equation*}
$$

for all $k \in \mathbb{N}$.
We also note a multi-linear estimate which follows from Sobolev embedding.
Proposition 2.8. Suppose that $\varrho \in\left(\frac{d}{4}, \frac{d}{2}\right)$ is given. Then we have for all $\delta>0$

$$
\left\|f_{1} f_{2} f_{3}\right\|_{H^{3 \varrho-d}} \lesssim \delta\left\|f_{1}\right\|_{H^{\varrho+\delta}}\left\|f_{2}\right\|_{H^{\varrho+\delta}}\left\|f_{3}\right\|_{H^{\varrho+\delta}}
$$

Assuming Propositions 2.7 and 2.8 for now, we give the proof of Theorem 1.5.
Proof of Theorem 1.5. By condition (5) of being a mild solution of (1), we note that $\phi_{0}=0$ implies that $\phi^{(1)}(, t)=\phi^{(2)}(, t)=0$ for all $t \in[-T, T]$. Hence, we need to consider the case when $\phi_{0} \neq 0$.

We note that (29) for $k=1$ implies that for all $x, x^{\prime} \in \mathbb{D}$ we have

$$
\phi^{(1)}(x, t) \overline{\phi^{(1)}\left(x^{\prime}, t\right)}=\phi^{(2)}(x, t) \overline{\phi^{(2)}\left(x^{\prime}, t\right)}
$$

We multiply both sides of the above identity with $\phi^{(2)}\left(x^{\prime}, t\right)$ and we integrate in $x^{\prime}$ to deduce that

$$
\phi^{(2)}(x, t)\left\|\phi^{(2)}(t)\right\|_{L_{x}^{2}}^{2}=\phi^{(1)}(x, t)\left\langle\phi^{(2)}(t), \phi^{(1)}(t)\right\rangle_{L_{x}^{2}} .
$$

Hence, by (5) for $\phi^{(2)}$, we have

$$
\begin{equation*}
\phi^{(2)}(x, t)=a(t) \phi^{(1)}(x, t) \tag{30}
\end{equation*}
$$

where

$$
a(t):=\frac{\left\langle\phi^{(2)}(t), \phi^{(1)}(t)\right\rangle_{L_{x}^{2}}^{2}}{\left\|\phi_{0}\right\|_{L_{x}^{2}}^{2}}
$$

Setting $t=0$, and recalling that $\phi^{(1)}(0)=\phi^{(2)}(0)=\phi_{0}$, it follows that

$$
\begin{equation*}
a(0)=1 \tag{31}
\end{equation*}
$$

Substituting (30) into (5), it follows that, for all $t \in[-T, T]$ we have

$$
\begin{equation*}
|a(t)|=1 \tag{32}
\end{equation*}
$$

Let us note that (30) with $a$ satisfying conditions (31) and (32) is all that we can deduce on the level of solutions to (1) from Proposition 2.7. Namely, any $\phi^{(1)}$ and $\phi^{(2)}$ satisfying the above conditions will also satisfy (29).

Our goal now is to show that $a=1$. This claim holds trivially if $\phi^{(1)}$ and $\phi^{(2)}$ are classical solutions and if $a$ is differentiable in $t$. In our case, we have to work a bit more. We start from (3). Using (30) and (32), we can rewrite (3) for $\phi^{(2)}$ as

$$
a(t) \phi^{(1)}(x, t)=e^{i t \Delta} \phi_{0}-i \int_{0}^{t} e^{i(t-\tau) \Delta}\left|\phi^{(1)}\right|^{2} \phi^{(1)}(x, \tau) a(\tau) d \tau
$$

Taking differences of the above equality with (3) for $\phi^{(1)}$ we get

$$
\begin{equation*}
(a(t)-1) \phi^{(1)}(x, t)=-i \int_{0}^{t} e^{i(t-\tau) \Delta}\left|\phi^{(1)}\right|^{2} \phi^{(1)}(x, \tau)(a(\tau)-1) d \tau \tag{33}
\end{equation*}
$$

By recalling (15) and the assumption $s>s_{0}$, we can find $\varrho \in\left(\frac{d}{4}, \frac{d}{2}\right)$ such that $\varrho<s$. We fix such a parameter $\varrho$ for the remainder of the proof. Let us now estimate both sides of (33) in the $H^{3 \varrho-d}$ norm. We have

$$
\begin{aligned}
& \left\|(a(t)-1) \phi^{(1)}(t)\right\|_{H^{3 \varrho-d}}=\left\|-i \int_{0}^{t} e^{i(t-\tau) \Delta}\left|\phi^{(1)}\right|^{2} \phi^{(1)}(, \tau)(1-a(\tau)) d \tau\right\|_{H^{3 \varrho-d}} \\
& \leq \int_{0}^{t}\left\|\left|\phi^{(1)}\right|^{2} \phi^{(1)}(\tau)\right\|_{H^{3 \varrho-d}}|a(\tau)-1| d \tau \lesssim \delta \int_{0}^{t}\left\|\phi^{(1)}(t)\right\|_{H^{e+\delta}}^{3}|a(\tau)-1| d \tau
\end{aligned}
$$

In the last step we applied Proposition 2.8 , which is possible by the assumptions on $\varrho$. Hence, by (4) for $\phi^{(1)}$, we have for $\delta>0$ sufficiently small

$$
\begin{align*}
|a(t)-1|\left\|\phi^{(1)}(t)\right\|_{H^{3 e-d}} & \lesssim \int_{0}^{t}\left\|\phi^{(1)}(\tau)\right\|_{H^{s}}^{3}|a(\tau)-1| d \tau \\
& \lesssim \int_{0}^{t} R^{3}|a(\tau)-1| d \tau \tag{34}
\end{align*}
$$

The fact that $a=1$ follows from (34) by an application of Gronwall's inequality if we show that $\left\|\phi^{(1)}(x, t)\right\|_{H_{x}^{3 s_{0}(d)-d}}$ is bounded from below. In order to do this, we need to consider two cases.
Case 1: $3 \varrho-d \geq 0$.
In this case we note that

$$
\begin{equation*}
\left\|\phi^{(1)}(t)\right\|_{H^{3 \varrho-d}} \geq\left\|\phi^{(1)}(t)\right\|_{L^{2}}=\left\|\phi_{0}\right\|_{L_{x}^{2}} \tag{35}
\end{equation*}
$$

by using (5).

Case 2: $3 \varrho-d<0$.
In this case we interpolate. Namely, we know that $3 \varrho-d<0<s$, hence there exists $\theta \in(0,1)$ such that

$$
\left\|\phi^{(1)}(t)\right\|_{L^{2}} \leq\left\|\phi^{(1)}(t)\right\|_{H^{3 e-d}}^{\theta}\left\|\phi^{(1)}(t)\right\|_{H^{s}}^{1-\theta}
$$

Using (4) and (5) we deduce that

$$
\begin{equation*}
\left\|\phi^{(1)}(t)\right\|_{H^{3 \varrho-d}} \geq\left(\frac{\left\|\phi_{0}\right\|_{L_{x}^{2}}}{R^{1-\theta}}\right)^{1 / \theta} \tag{36}
\end{equation*}
$$

Combining (34), (35), and (36), we deduce that

$$
|a(t)-1| \lesssim \int_{0}^{t}|a(\tau)-1| d \tau
$$

from where we conclude that $a=1$. The result now follows.
The remainder of this section is devoted to the proof of Proposition 2.7. The proof of Proposition 2.8 is given in Appendix A below. We first note several auxiliary facts, which are based on the arguments of [8]. Given $t \in[-T, T]$ we let

$$
\begin{equation*}
d \mu_{t}^{(1)}:=\delta_{\phi^{(1)}(, t)}, \quad d \mu_{t}^{(2)}:=\delta_{\phi^{(2)}(, t)} . \tag{37}
\end{equation*}
$$

Note that these are Borel measures on $L^{2}(\mathbb{D})$. Following [8], we define for $k, r \in \mathbb{N}$ the set

$$
\begin{align*}
& \mathcal{M}_{k, r}:=\{\sigma:\{k+1, k+2, \ldots, k+r\} \rightarrow\{1,2, \ldots, k+r-1\}: \\
&\sigma(j)<j \text { for all } j \in\{k+1, k+2, \ldots, k+r\}\}, \tag{38}
\end{align*}
$$

and the quantity

$$
\begin{align*}
& J^{k}\left(f ; \sigma ; t, t_{1}, \ldots, t_{r} ; \vec{x}_{k} ; \vec{x}_{k}^{\prime}\right):=\mathcal{U}^{(k)}\left(t-t_{1}\right) B_{\sigma(k+1), k+1} \mathcal{U}^{(k+1)}\left(t_{1}-t_{2}\right) B_{\sigma(k+2), k+2} \\
& \cdots \mathcal{U}^{(k+r-1)}\left(t_{r-1}-t_{r}\right) B_{\sigma(k+r), k+r}\left(|f\rangle\left\langle\left. f\right|^{\otimes(k+r)}\right)\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}\right)\right. \tag{39}
\end{align*}
$$

Lemma 2.9. Let $s>s_{0}$. There exists a constant $C>0$ depending only on the domain $\mathbb{D}$ and $s$ such that for $\sigma \in \mathcal{M}_{k, r}, t \in[-T, T]$ and $q \in\{1,2\}$ we have

$$
\begin{align*}
& \int_{[0, t]^{r}} \int_{L^{2}(\mathbb{D})} d \mu_{t_{r}}^{(q)}(f) \operatorname{Tr}\left|S^{\left(k,-\zeta_{0}\right)} J^{k}\left(f ; \sigma ; t, t_{1}, \ldots, t_{r} ; \vec{x}_{k} ; \vec{x}_{k}^{\prime}\right)\right| d t_{r} \cdots d t_{1} \\
& \leq C^{r} t^{\epsilon(r-1)} \int_{0}^{t}\left\|\phi^{(q)}\left(, t_{r}\right)\right\|_{H^{s}}^{2(k+r)} d t_{r} \tag{40}
\end{align*}
$$

Here $\epsilon \equiv \epsilon(\mathbb{D})$ is the parameter from Definition 1.3.
Proof of Lemma 2.9. We can rewrite (39) as a product of 1-particle kernels as in [8, (4.26)]

$$
\begin{equation*}
J^{k}\left(f ; \sigma ; t, t_{1}, \ldots, t_{r} ; \vec{x}_{k} ; \vec{x}_{k}^{\prime}\right)=\prod_{j=1}^{k} J_{j}^{1}\left(f ; \sigma_{j} ; t, t_{\ell_{j, 1}}, \ldots, t_{\ell_{j, m_{j}}} ; x_{j} ; x_{j}^{\prime}\right) \tag{41}
\end{equation*}
$$

For the precise definitions and explanations of the notation, we refer the reader to the discussion in [8, Section 5]. Following the terminology of [8], one of the factors on the right-hand side of (41) is distinguished, whereas the other $k-1$ are regular. The distinguished term contains a factor of $|f|^{2} f$.

Let us first explain the modifications needed to analyze the contribution from the distinguished factor. We change the expression in [8, eq. (6.4)] to ${ }^{2}$

$$
\int_{[0, t]^{m_{j}-1}} \operatorname{Tr}\left|S^{\left(1,-\zeta_{0}\right)} J_{j}^{1}\left(f ; \sigma_{j} ; t, t_{\ell_{j, 1}}, \ldots, t_{\ell_{j, m_{j}}}\right)\right| d t_{\ell_{j, m_{j}-1}} d t_{\ell_{j, m_{j}-2}} \cdots d t_{\ell_{j, 1}}
$$

Therefore, the expression in [8, eq. (7.1)] gets changed to

$$
\sum_{\beta_{1}} \int_{[0, t]^{m_{j}-1}}\left\|\psi_{\beta_{1}}^{1}\right\|_{H^{-\zeta_{0}}}\left\|\chi_{\beta_{1}}^{1}\right\|_{H^{-\zeta_{0}}} d t_{\ell_{j, m_{j}-1}} d t_{\ell_{j, m_{j}-2}} \cdots d t_{\ell_{j, 1}}
$$

which is

$$
\leq \sum_{\beta_{1}} \int_{[0, t]^{m_{j}-1}}\left\|\psi_{\beta_{1}}^{1}\right\|_{B_{2,1}^{-\zeta_{0}}}\left\|\chi_{\beta_{1}}^{1}\right\|_{B_{2,1}^{\zeta}} d t_{\ell_{j, m_{j}-1}} d t_{\ell_{j, m_{j}-2}} \cdots d t_{\ell_{j, 1}}
$$

for some $\zeta \in\left(\zeta_{0}, s\right)$, which we take to be fixed for the remainder of the proof. Note that, in the last step we used (10).

The key point is that we can now also change what were originally the bound on the $L^{2}$ level [8, eq. (7.4)] and the bound on the $H^{1}$ level [8, eq. (7.5)] by using the estimates that hold by assumption from Definition 1.3. In particular, we can use (12) and (13) from Definition 1.3 and replace [8, eq. (7.4)] with

$$
\begin{align*}
& \int_{[0, t]^{d}}\left\|\psi_{\beta_{a}}^{a}\right\|_{B_{2,1}^{-\zeta_{0}}}\left\|\chi_{\beta_{a}}^{a}\right\|_{B_{2,1}^{\zeta}} \prod_{a^{\prime} \in \tau_{j, a}} d t_{\ell_{j, a^{\prime}}} \\
& \lesssim t^{\epsilon}\left(\int_{[0, t]^{d} \kappa_{-}(a)}\left\|\psi_{\beta_{\kappa_{-}(a)}(a)}^{\kappa_{-}(a)}\right\|_{B_{2,1}^{\zeta}}\left\|\chi_{\beta_{\kappa_{-}(a)}^{\kappa-(a)}}^{\kappa_{-}(a)}\right\|_{B_{2,1}^{\zeta}} \prod_{a^{\prime} \in \tau_{\tau_{j, \kappa_{-}-}(a)}} d t_{\ell_{j, a^{\prime}}}\right) \\
& \times\left(\int_{[0, t]^{d_{\kappa_{+}}(a)}}\left\|\psi_{\beta_{\kappa_{+}(a)}^{\kappa_{+}(a)}}\right\|_{B_{2,1}^{-\zeta_{0}}}\left\|\chi_{\beta_{\kappa_{+}(a)}^{\kappa_{+}(a)}}\right\|_{B_{2,1}^{\zeta}} \prod_{a^{\prime} \in \tau_{j, \kappa_{+}(a)}} d t_{\ell_{j, a^{\prime}}}\right) . \tag{42}
\end{align*}
$$

We always apply this estimate in the case where $\psi_{\beta_{\kappa_{+}(a)}}^{\kappa_{+}(a)}$ contains the factor $|f|^{2} f$. Furthermore, we can use (13) and replace [8, eq. (7.5)] with

$$
\begin{align*}
& \int_{[0, t]^{d} a}\left\|\psi_{\beta_{a}}^{a}\right\|_{B_{2,1}^{\zeta}}\left\|\chi_{\beta_{a}}^{a}\right\|_{B_{2,1}^{\zeta}} \prod_{a^{\prime} \in \tau_{j, a}} d t_{\ell_{j, a^{\prime}}} \\
& \lesssim t^{\epsilon}\left(\int_{[-T, T]^{d_{\kappa_{-}}(a)}}\left\|\psi_{\beta_{\kappa_{-}(a)}^{\kappa_{-}(a)}}\right\|_{B_{2,1}^{\zeta}}\left\|\chi_{\beta_{\kappa_{-}(a)}^{\kappa-(a)}}^{\kappa_{B}}\right\|_{B_{2,1}^{\zeta}} \prod_{a^{\prime} \in \tau_{j, \kappa_{-}(a)}} d t_{\ell_{j, a^{\prime}}}\right) \\
& \times\left(\int_{[0, t]^{d_{\kappa_{+}}(a)}}\left\|\psi_{\beta_{\kappa_{+}(a)}^{\kappa_{+}(a)}}^{\kappa_{2,1}}\right\|_{B_{2,1}^{\zeta}}\left\|\chi_{\beta_{\kappa_{+}(a)}^{\kappa_{+}(a)}}^{\kappa_{B_{2,1}}}\right\|_{B^{\prime} \in \tau_{j, \kappa_{+}(a)}^{\zeta}} \prod_{\ell_{j, a^{\prime}}}\right) . \tag{43}
\end{align*}
$$

Finally, we note that we replace applications of [8, eq. (7.2)]), i.e.

$$
\left\||f|^{2} f\right\|_{L^{2}} \lesssim\|f\|_{H^{1}}^{3}
$$

by (14) of Definition 1.3 , which implies that

$$
\begin{equation*}
\left\||f|^{2} f\right\|_{B_{2,1}^{-\zeta_{0}}} \lesssim\|f\|_{H^{s}}^{3} . \tag{44}
\end{equation*}
$$

[^2]Using (42), (43), and (44), it follows that the bound [8, eq. (8.1)] for the distinguished tree is replaced with

$$
\begin{align*}
& \quad \int_{[0, t]^{m_{j}-1}} \operatorname{Tr}\left|S^{\left(1,-\zeta_{0}\right)} J_{j}^{1}\left(f ; \sigma_{j} ; t, t_{\ell_{j, 1}}, \ldots, t_{\ell_{j, m_{j}}}\right)\right| d t_{\ell_{j, m_{j}-1}} \cdots d t_{\ell_{j, 1}}  \tag{45}\\
& \leq C_{1}^{m_{j}} t^{\left(m_{j}-1\right) \epsilon}\|f\|_{B_{2,1}^{\zeta}}^{2 m_{j}-1}\|f\|_{H^{s}}^{3} \leq C^{m_{j}} t^{\left(m_{j}-1\right) \epsilon}\|f\|_{H^{s}}^{2 m_{j}+2} \tag{46}
\end{align*}
$$

where $C_{1}, C>0$ depend on $\mathbb{D}$ and $s$. In the last step we used (10).
The contribution from each regular tree, i.e. the bound [8, eq. (8.6)] is replaced with

$$
\begin{array}{r}
\int_{[0, t]^{m_{j}}} \operatorname{Tr}\left|S^{\left(1,-\zeta_{0}\right)} J_{j}^{1}\left(f ; \sigma_{j} ; t, t_{\ell_{j, 1}}, \ldots, t_{\ell_{j, m_{j}}}\right)\right| d t_{\ell_{j, m_{j}}} d t_{\ell_{j, m_{j}-1}} \cdots d t_{\ell_{j, 1}} \\
\leq C_{2}^{m_{j}} t^{m_{j} \epsilon}\|f\|_{B_{2,1}^{\iota}}^{2 m_{j}+2} \leq C^{m_{j}} t^{m_{j} \epsilon}\|f\|_{H^{s}}^{2 m_{j}+2} \tag{47}
\end{array}
$$

where $C_{2}, C>0$ depend on $\mathbb{D}$ and $s$. In order to obtain (47), one applies analogous modifications to those needed to obtain (46). The difference is that one now only uses (13) and does not need to use (12) and (14).

Using (46) and (47) in (41), it follows that the left-hand side of (40) is

$$
\leq C^{r} t^{\epsilon(r-1)} \int_{0}^{t} \int_{L^{2}(\mathbb{D})} d \mu_{t}^{(q)}(f)\|f\|_{H^{s}}^{2(k+r)} d t_{r}=C^{r} t^{\epsilon(r-1)} \int_{0}^{t}\left\|\phi^{(q)}\left(, t_{r}\right)\right\|_{H^{s}}^{2(k+r)} d t_{r}
$$

The boardgame argument of [30, Section 3] gives an explicit equivalence relation $\sim$ on the set $\mathcal{M}_{k, r}$ defined in (38) above. Furthermore, there exists a complete set of representatives $\mathcal{N}_{k, r}$ with the property that

$$
\begin{equation*}
\sharp \mathcal{N}_{k, r} \leq 2^{k+r} . \tag{48}
\end{equation*}
$$

We omit the details of the construction and refer the reader to [30, Section 3].
Lemma 2.10. For $\sigma_{0} \in \mathcal{N}_{k, r}$, there exists $D\left(\sigma_{0}, t\right) \subset[0, t]^{r}$ such that for $q=1,2$, we have

$$
\begin{align*}
\sum_{\substack{\sigma \in \mathcal{M}_{k, r} \\
\sigma \sim \sigma_{0}}} \int_{0}^{t} \int_{0}^{t_{1}} & \cdots \int_{0}^{t_{r-1}} \int_{L^{2}(\mathbb{D})} d \mu_{t_{r}}^{(q)}(f) J^{k}\left(f ; \sigma ; t, t_{1}, \ldots, t_{r}\right) d t_{r} \cdots d t_{1} \\
& =\int_{D\left(\sigma_{0}, t\right)} \int_{L^{2}(\mathbb{D})} d \mu_{t_{r}}^{(q)}(f) J^{k}\left(f ; \sigma_{0} ; t, t_{1}, \ldots, t_{r}\right) d t_{r} \cdots d t_{1} . \tag{49}
\end{align*}
$$

Proof of Lemma 2.10. This claim follows from the analysis of [30, Section 3]. A subtle point that one should address is that this proof relies on the commutation relation [30, eq. (37)], which in turn relies on considering operator kernels. To fully justify the application of this commutation relation, we work with functions which are smooth in the spatial variables and conclude (49) by density. More precisely, we find a sequence $F_{m}$ of smooth functions on $[0, t] \times \mathbb{D}$ which converge to $\phi^{(q)}$ in $L_{[0, t]}^{2(k+r)} H_{x}^{s}$. Here we think of $\phi^{(q)}$ as a function of space and time. For $\tau \in[0, t]$, we let $d \nu_{\tau}^{(m)}:=\delta_{F_{m}(, \tau)}$. Then, since for all $\tau \in[0, t]$, the function $F(, \tau)$ is smooth in
$x$, the proof of [30, Lemma 3.1] carries over and allows us to deduce the analogue of (49) with $d \mu^{(q)}$ replaced by $d \nu^{(m)}$. The claim now follows by Lemma 2.9 and letting $m \rightarrow \infty$.

We can now prove Proposition 2.7.
Proof of Proposition 2.7. Given $k \in \mathbb{N}$ we define
$\gamma^{(k)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right):=\left|\phi^{(1)}\right\rangle\left\langle\left.\phi^{(1)}\right|^{\otimes k}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right), \widetilde{\gamma}^{(k)}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right):=\mid \phi^{(2)}\right\rangle\left\langle\left.\phi^{(2)}\right|^{\otimes k}\left(\vec{x}_{k} ; \vec{x}_{k}^{\prime}, t\right)\right.$.
By Lemma 2.4, we know that $\left(\gamma^{(k)}\right)_{k},\left(\widetilde{\gamma}^{(k)}\right)_{k}$ are mild solutions of (19) in $\mathfrak{H}^{s}$ with the same initial data. In particular

$$
\widehat{\gamma}^{(k)}:=\gamma^{(k)}-\widetilde{\gamma}^{(k)}
$$

is a mild solution of (19) in $\mathfrak{H}^{s}$ with zero initial data. It suffices to show that

$$
\begin{equation*}
\operatorname{Tr}\left|S^{\left(k,-\zeta_{0}\right)} \widehat{\gamma}^{(k)}\right|=0 \quad \text { on }[-T, T] \tag{50}
\end{equation*}
$$

for all $k \in \mathbb{N}$ provided that we choose $T$ to be sufficiently small depending on $R$ (from Definition 1.1). The claim for general $T$ follows by an iteration argument.

Recalling (37), we define

$$
\begin{equation*}
d \widehat{\mu}_{t}:=d \mu_{t}^{(1)}-d \mu_{t}^{(2)} \tag{51}
\end{equation*}
$$

By applying an iterated Duhamel expansion, we have that for all $k, r \in \mathbb{N}$ and $t \in[-T, T]$

$$
\begin{align*}
& \widehat{\gamma}^{(k)}(t)=(-i)^{r} \sum_{\sigma \in \mathcal{M}_{k, r}} \int_{0}^{t} \int_{0}^{t_{1}} \cdots \int_{0}^{t_{r-1}} \int_{L^{2}(\mathbb{D})} d \widehat{\mu}_{t_{r}}(f) \\
& J^{k}\left(f ; \sigma ; t, t_{1}, \ldots, t_{r} ; \vec{x}_{k} ; \vec{x}_{k}^{\prime}\right) d t_{r} \cdots d t_{1} \tag{52}
\end{align*}
$$

for $\mathcal{M}_{k, r}$ as in (38) and $J^{k}$ as in (39). Using Lemma 2.10 and (51) we can rewrite the expression on the right hand side of (52) as

$$
\begin{equation*}
(-i)^{r} \sum_{\sigma_{0} \in \mathcal{N}_{k, r}} \int_{D\left(\sigma_{0}, t\right)} \int_{L^{2}(\mathbb{D})} d \widehat{\mu}_{t_{r}}(f) J^{k}\left(f ; \sigma_{0} ; t, t_{1}, \ldots, t_{r}\right) d t_{r} \cdots d t_{1} \tag{53}
\end{equation*}
$$

Using the triangle inequality, (48), Lemma 2.9, and the assumptions on $\phi^{(1)}, \phi^{(2)}$ in the expression (53), it follows that for $t \in[-T, T]$ we have

$$
\operatorname{Tr}\left|S^{\left(k,-\zeta_{0}\right)} \widehat{\gamma}^{(k)}\right| \leq 2^{k+r+1} C^{r} t^{\epsilon(r-1)+1} R^{2(k+r)}=2^{k+1} t^{1-\epsilon} R^{2 k}\left(2 C R^{2} t^{\epsilon}\right)^{r}
$$

We deduce (50) by letting $r \rightarrow \infty$.
Remark 2.11. Note that the proof of Proposition 2.7 implies that mild solutions of (19) in $\mathfrak{H}^{s}$ of the form

$$
\gamma^{(k)}(t)=\int_{L^{2}(\mathbb{D})} d \mu_{t}(f)|f\rangle\left\langle\left. f\right|^{\otimes k}\right.
$$

for $d \mu_{t}$ a Borel measure on $L^{2}(\mathbb{D})$ supported on a ball centered at zero of radius independent of $t$, are uniquely determined by their initial data provided that $s>s_{0}$. In light of Proposition 3.1 below, this improves on our earlier uniqueness results for such solutions on rectangular tori [22].

Remark 2.12. Note that, in the proof of Theorem 1.5, we never directly use the quantum de Finetti theorem. For our purposes, it suffices to apply the conclusion of this theorem concerning the structure of density matrices. This, in turn, allows us to apply the techniques developed in [8]. In doing so we work in $L^{2}$-based Besov spaces, which allows us to avoid the issues that arise when working in $L^{p}$-based spaces. The latter approach was taken on $\mathbb{R}^{d}$ in $[24,25]$ and was based on the stronger dispersive properties available on the whole space.

## 3. Admissibility of the rectangular torus of dimension $d \geq 2$

In this section we consider the general rectangular torus in $d \geq 2$ dimensions, which we write as

$$
\mathbb{T}_{g}^{d}=\left(\mathbb{R} / \frac{2 \pi}{\theta_{1}} \mathbb{Z}\right) \times\left(\mathbb{R} / \frac{2 \pi}{\theta_{2}} \mathbb{Z}\right) \times \cdots \times\left(\mathbb{R} / \frac{2 \pi}{\theta_{d}} \mathbb{Z}\right)
$$

for some $\theta_{1}, \theta_{2}, \ldots, \theta_{d}>0$ and we show that it is admissible in the sense of Definition 1.3 above.

Proposition 3.1. $\mathbb{D}=\mathbb{T}_{g}^{d}$ is admissible with the following parameters:

$$
\begin{gather*}
\zeta_{0}= \begin{cases}\frac{d(d-1)}{2(d+2)} & \text { for } 2 \leq d \leq 4 \\
\frac{d}{2}-1 & \text { for } d \geq 5\end{cases}  \tag{54}\\
\alpha_{0}= \begin{cases}\frac{d(d+5)}{6(d+2)} & \text { for } 2 \leq d \leq 4 \\
\frac{d}{6}+\frac{1}{3} & \text { for } d \geq 5\end{cases}  \tag{55}\\
\epsilon=\left\{\begin{array}{cl}
\left(\frac{4-d}{2(d+2)}\right)+ & \text { for } 2 \leq d \leq 4 \\
0+ & \text { for } d \geq 5
\end{array}\right. \tag{56}
\end{gather*}
$$

From Theorem 1.5 and Proposition 3.1 we obtain Corollary 1.7 stated earlier. In the proof of Proposition 3.1 we use the Strichartz estimate on $\mathbb{T}_{g}^{d}[3,28]$.

Proposition 3.2 (Strichartz estimate on $\left.\mathbb{T}_{g}^{d}\right)$. For $p>\frac{2(d+2)}{d}$ we have

$$
\left\|e^{i t \Delta} P_{\leq N} f\right\|_{L_{t, x}^{p}\left([0,1] \times \mathbb{T}_{g}^{d}\right)} \lesssim_{d, p}\langle N\rangle^{\frac{d}{2}-\frac{d+2}{p}}\|f\|_{L_{x}^{2}}
$$

Here

$$
P_{\leq N}:=\sum_{N^{\prime} \leq N} P_{N^{\prime}}
$$

Moreover we recall that the Bernstein inequality holds for the full range of integrability exponents, as on $\mathbb{R}^{n}$. We note that for general domains $\mathbb{D}$, we only have the claim for a partial range of integrability exponents, see Lemma A. 3 below.

Let $\mathbb{Z}_{g}^{d}=\left(\theta_{1} \mathbb{Z} \times \cdots \times \theta_{d} \mathbb{Z}\right)$ and $\widetilde{P}_{N}$ be a smoothed out version of $P_{N}$, i.e.

$$
\begin{equation*}
\widetilde{P}_{N}:=\sum_{k \in \mathbb{Z}_{g}^{d}} \varphi_{N}(k) \chi_{k} \tag{57}
\end{equation*}
$$

where, if $N \geq 1, \varphi_{N}=\varphi(\cdot / N)$ and $\varphi \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ equal to one on $1 \leq|\xi|<2$ and supported on $1 / 2 \leq|\xi| \leq 4$, and $\varphi_{0} \in C_{c}^{\infty}\left(\mathbb{R}^{d}\right)$ equal to one on $|\xi|<1$ and supported on $|\xi| \leq 2$. In particular, we have $\widetilde{P}_{N} P_{N}=P_{N}$. The following is well-known, for the reader's convenience we provide a short proof in Appendix B.

Lemma 3.3 (Bernstein's inequality on the rectangular torus). Let $1 \leq p \leq q \leq \infty$. For all $N$ we have

$$
\begin{equation*}
\left\|\widetilde{P}_{N} f\right\|_{L^{q}} \lesssim\langle N\rangle^{d / q-d / p}\|f\|_{L^{p}} \tag{58}
\end{equation*}
$$

We now give the proof of Proposition 3.1
Proof of Proposition 3.1. We first prove (12) for the choice of parameters as in (54), (55), and (56) above. In the sequel we abbreviate $\|\cdot\|_{L^{p}([-T, T])}$ as $\|\cdot\|_{L_{T}^{p}}$. Let us note for $0 \leq \eta \leq \zeta_{0}$ we have

$$
\begin{equation*}
\left\|u_{1} u_{2} u_{3}\right\|_{L_{t}^{1} B_{2,1}^{-\eta}\left([-T, T] \times \mathbb{T}_{g}^{d}\right)} \leq \sum_{N}\langle N\rangle^{-\eta}\left\|P_{N}\left(u_{1} u_{2} u_{3}\right)\right\|_{L_{T}^{1} L_{x}^{2}}=: I \tag{59}
\end{equation*}
$$

We write

$$
P_{N}\left(u_{1} u_{2} u_{3}\right)=\sum_{N_{1}, N_{2}, N_{3}} P_{N}\left(u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right),
$$

where we abbreviate $u_{j, N_{j}}:=P_{N_{j}} u_{j}$ for $j=1,2,3$. We remark that there is no contribution unless $N \lesssim \max \left\{N_{1}, N_{2}, N_{3}\right\}$. We now consider two cases, depending on the relative size of $\max \left\{N_{1}, N_{2}, N_{3}\right\}$ with respect to $N$.

Case 1: $\max \left\{N_{1}, N_{2}, N_{3}\right\} \lesssim N$.
Case 1.A: $N_{1}=\max \left\{N_{1}, N_{2}, N_{3}\right\}$. In this case there is no contribution unless $N \sim N_{1}$ and we obtain

$$
\begin{equation*}
I \lesssim \sum_{\substack{N_{1}, N_{2}, N_{3} \\ N_{1} \geq N_{2}, N_{3}}}\left\langle N_{1}\right\rangle^{-\eta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{2}} . \tag{60}
\end{equation*}
$$

With the abbreviation $\phi_{j, N_{j}}:=P_{N_{j}} \phi_{j}$ for $j=1,2,3$, it suffices to show that

$$
\begin{align*}
& \left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{2}} \\
& \lesssim T^{\epsilon}\left\|\phi_{1, N_{1}}\right\|_{L_{x}^{2}}\left(\left\langle N_{2}\right\rangle^{\zeta_{0}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\zeta_{0}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right) . \tag{61}
\end{align*}
$$

In order to establish (61) we consider two subcases.
Subcase 1.A.i: $N_{1} \lesssim \max \left\{N_{2}, N_{3}\right\}$. By Hölder's inequality we have

$$
\begin{align*}
& \left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{2}} \\
& \leq\left\|u_{1, N_{1}}\right\|_{L_{T, x}^{\frac{2(d+2)}{x}}}\left\|u_{2, N_{2}}\right\|_{L_{T}^{\frac{4(d+2)}{d+4}} L_{x}^{2(d+2)}}\left\|u_{3, N_{3}}\right\|_{L_{T}^{\frac{4(d+2)}{d+4}} L_{x}^{2(d+2)}} . \tag{62}
\end{align*}
$$

We will apply (62) differently depending on whether $d \leq 4$ or $d \geq 5$.
(1) $2 \leq d \leq 4$. In this case, we have $\frac{4(d+2)}{d+\underset{\sim}{4}} \leq \frac{2(d+2)}{d} \leq 2(d+2)$. We apply Lemma 3.3 (which is possible since $u_{2, N_{2}}=\widetilde{P}_{N_{2}} u_{2, N_{2}}$ ), Hölder's inequality, and Proposition 3.2 with $p=\frac{2(d+2)}{d}+$ to deduce that

$$
\left\|u_{2, N_{2}}\right\|_{L_{T}^{\frac{4(d+2)}{d+4}} L_{x}^{2(d+2)}} \lesssim T^{\frac{4-d}{4(d+2)}+}\left\langle N_{2}\right\rangle^{\frac{d(d-1)}{2(d+2)}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}
$$

An analogous bound holds for $u_{3, N_{3}}$. Using these bounds and Proposition 3.2 with $p=\frac{2(d+2)}{d}+$ for the $u_{1, N_{1}}$ factor, we get that the expression (62) is

$$
\begin{aligned}
& \lesssim T^{\frac{4-d}{2(d+2)}+}\left(\left\langle N_{1}\right\rangle^{0+}\left\|\phi_{1, N_{1}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{2}\right\rangle^{\frac{d(d-1)}{2(d+2)}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\frac{d(d-1)}{2(d+2)}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right) \\
& \lesssim T^{\frac{4-d}{2(d+2)}+}\left\|\phi_{1, N_{1}}\right\|_{L_{x}^{2}}\left(\left\langle N_{2}\right\rangle^{\frac{d(d-1)}{2(d+2)}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\frac{d(d-1)}{2(d+2)}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right) .
\end{aligned}
$$

In the last line, we used the fact that $N_{1} \lesssim \max \left\{N_{2}, N_{3}\right\}$ to distribute the factor of $\left\langle N_{1}\right\rangle^{0+}$.
(2) $d \geq 5$. In this case, we have $\frac{2(d+2)}{d}<\frac{4(d+2)}{d+4}<2(d+2)$. We first apply Lemma 3.3 to obtain that (62) is

$$
\begin{equation*}
\leq\left\|u_{1, N_{1}}\right\|_{L_{T, x}^{\frac{2(d+2)}{d}}}\left(\left\langle N_{2}\right\rangle^{\frac{d}{4}}\left\|u_{2, N_{2}}\right\|_{L_{T, x}^{\frac{4(d+2)}{d+4}}}\right)\left(\left\langle N_{3}\right\rangle^{\frac{d}{4}}\left\|u_{3, N_{3}}\right\|_{L_{T, x}^{\frac{4(d+2)}{d+4}}}\right) \tag{63}
\end{equation*}
$$

To estimate the $u_{2, N_{2}}$ and $u_{3, N_{3}}$ factors in (63) apply Proposition 3.2 with $p=$ $\frac{4(d+2)}{d+4}$. The $u_{1, N_{1}}$ factor we estimate as we did for $d \leq 4$. Therefore, (63) is

$$
\lesssim T^{0+}\left\|\phi_{1, N_{1}}\right\|_{L_{x}^{2}}\left(\left\langle N_{2}\right\rangle^{\left(\frac{d}{2}-1\right)+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\left(\frac{d}{2}-1\right)+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right)
$$

This finishes the proof of (61) in Subcase 1.A.i.
Subcase 1.A.ii: $N_{1} \gg N_{2}, N_{3}$. By symmetry, we can assume without loss of generality that $N_{2} \geq N_{3}$. In this subcase, we can no longer distribute the factor of $N_{1}^{0+}$ into the factors involving $N_{2}$ and $N_{3}$. Instead, we need to apply an almost orthogonality argument. More precisely, let $\Gamma$ denote the collection of all nonoverlapping cubes of size $N_{2}$ with centers in the lattice $N_{2} \mathbb{Z}_{g}^{d}$. For $\mathcal{C} \in \Gamma$, we denote by $P_{\mathcal{C}}$ the projection to $\mathcal{C}$ in frequency space. Given $\mathcal{C} \in \Gamma$, we observe that

$$
\begin{equation*}
\int_{\mathbb{T}_{g}^{d}}\left(P_{\mathcal{C}} u_{1, N_{1}}\right) u_{2, N_{2}} u_{3, N_{3}} \overline{\left(P_{\mathcal{C}^{\prime}} u_{1, N_{1}}\right) u_{2, N_{2}} u_{3, N_{3}}} d x \neq 0 \tag{64}
\end{equation*}
$$

only for finitely many $\mathcal{C}^{\prime} \in \Gamma$, with a bound on this number independent of $\mathcal{C}, N_{1}, N_{2}, N_{3}$. By Galilean invariance of $e^{i t \Delta}$ and by Proposition 3.2 for $p=$ $\frac{2(d+2)}{d}+$, it follows that

$$
\begin{equation*}
\left\|P_{\mathcal{C}} u_{1, N_{1}}\right\|_{L_{T, x}^{2(d+2)}+} \lesssim\left\langle N_{2}\right\rangle^{0+}\left\|P_{\mathcal{C}} \phi_{1, N_{1}}\right\|_{L_{x}^{2}} \tag{65}
\end{equation*}
$$

By applying (65) instead of Proposition 3.2 when $p=\frac{2(d+2)}{d}+$ to estimate the factor with frequency of order $N_{1}$, the arguments from Subcase 1.A.i imply that

$$
\begin{align*}
& \left\|\left(P_{\mathcal{C}} u_{1, N_{1}}\right) u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{2}} \\
& \lesssim T^{\epsilon}\left\|P_{\mathcal{C}} \phi_{1, N_{1}}\right\|_{L_{x}^{2}}\left(\left\langle N_{2}\right\rangle^{\zeta_{0}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\zeta_{0}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right) \tag{66}
\end{align*}
$$

Taking an $\ell^{2}$ sum over $\mathcal{C} \in \Gamma$ in (66) and using (64), we deduce (61) in Subcase 1.A.ii. In summary, in Case 1.A we obtain

$$
\begin{align*}
I & \lesssim \sum_{\substack{N_{1}, N_{2}, N_{3} \\
N_{1} \geq N_{2}, N_{3}}}\left\langle N_{1}\right\rangle^{-\eta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{2}} \\
& \lesssim \sum_{\substack{N_{1}, N_{2}, N_{3} \\
N_{1} \geq N_{2}, N_{3}}} T^{\epsilon}\left(\left\langle N_{1}\right\rangle^{-\eta}\left\|\phi_{1, N_{1}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{2}\right\rangle^{\zeta_{0}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\zeta_{0}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right) \\
& \lesssim T^{\epsilon}\left\|\phi_{1}\right\|_{B_{2,1}^{-\eta}}^{-\eta}\left\|\phi_{2}\right\|_{B_{2,1}^{\zeta}}\left\|\phi_{3}\right\|_{B_{2,1}^{\zeta}} . \tag{67}
\end{align*}
$$

Case 1.B: $N_{2}=\max \left\{N_{1}, N_{2}, N_{3}\right\}$. Since $N \sim N_{2} \geq N_{1}$, it suffices to estimate

$$
\begin{equation*}
\sum_{\substack{N_{1}, N_{2}, N_{3} \\ N_{2} \geq N_{1}, N_{3}}}\left\langle N_{1}\right\rangle^{-\eta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{2}} . \tag{68}
\end{equation*}
$$

We estimate each summand in (68) by using (61) as in Case 1. The point is that, since $N_{1} \leq N_{2}$, a factor of $\left\langle N_{1}\right\rangle^{0+}$ can be absorbed into the powers of $\left\langle N_{2}\right\rangle$. Therefore, we can deduce that the contribution from Case 1.B also satisfies the bound (67).

Case 1.C: $N_{3}=\max \left\{N_{1}, N_{2}, N_{3}\right\}$. Case 1.C is analogous to Case 1.B if we reverse the roles of $N_{2}$ and $N_{3}$.

Case 2: $\max \left\{N_{1}, N_{2}, N_{3}\right\} \gg N$. We observe that $P_{N}\left(u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right)=0$ unless the largest two elements of $\left\{N_{1}, N_{2}, N_{3}\right\}$ are comparable. We consider three possible cases.

Case 2.A: $N_{1} \sim N_{2} \gtrsim N_{3}$. In this case, we know that $N_{1} \sim N_{2} \gg N$. Let $q \in[1,2]$ be fixed. We have

$$
\begin{align*}
& \langle N\rangle^{-\eta}\left\|P_{N}\left(u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right)\right\|_{L_{T}^{1} L_{x}^{2}} \leq\langle N\rangle^{-\eta}\left\|\widetilde{P}_{N}\left(u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right)\right\|_{L_{T}^{1} L_{x}^{2}} \\
& \lesssim\langle N\rangle^{-\delta}\langle N\rangle^{-\eta+\frac{d}{q}-\frac{d}{2}+\delta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}} \tag{69}
\end{align*}
$$

For the last inequality, we used Lemma 3.3. Here $\delta>0$ is a small, arbitrary constant. We now choose $q$ to be $q_{0} \equiv q_{0}(d)$ such that $\frac{d}{q_{0}}-\frac{d}{2}=\zeta_{0}$, for $\zeta_{0}$ as in (54). In particular, we have

$$
q_{0}= \begin{cases}\frac{2(d+2)}{2 d+1} & \text { for } d \leq 4  \tag{70}\\ \frac{d}{d-1} & \text { for } d \geq 5\end{cases}
$$

Then $q_{0} \in[1,2]$. Furthermore, the condition on $\eta$ in (11) can then be rewritten as

$$
0 \leq \eta \leq \frac{d}{q_{0}}-\frac{d}{2}
$$

In particular, choosing $q=q_{0}$, the expression (69) is

$$
\begin{align*}
& \lesssim\langle N\rangle^{-\delta}\left\langle N_{1}\right\rangle^{-\eta+\zeta_{0}+\delta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}} \\
& \sim\langle N\rangle^{-\delta}\left\langle N_{1}\right\rangle^{-\eta}\left\langle N_{2}\right\rangle^{\zeta_{0}+\delta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}} \tag{71}
\end{align*}
$$

Therefore, it remains to estimate

$$
\begin{equation*}
\left\langle N_{1}\right\rangle^{-\eta}\left\langle N_{2}\right\rangle^{\zeta_{0}+\delta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}} \tag{72}
\end{equation*}
$$

for fixed dyadic $N_{1}, N_{2}, N_{3}$ with $N_{1} \sim N_{2} \gtrsim N_{3}$. As in Case 1 we consider the cases $d \leq 4$ and $d \geq 5$ separately.
(1) $2 \leq d \leq 4$. By Hölder's inequality we note that, for $q$ as in (70) we have

$$
\begin{equation*}
\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}} \leq\left\|u_{1, N_{1}}\right\|_{L_{T, x}^{\frac{2(d+2)}{d}}}\left\|u_{2, N_{2}}\right\|_{L_{T, x}^{\frac{2(d+2)}{d}}}\left\|u_{3, N_{3}}\right\|_{L_{T}^{\frac{d+2}{2}} L_{x}^{2(d+2)}} \tag{73}
\end{equation*}
$$

By Lemma 3.3, Hölder's inequality, and Proposition 3.2, we can estimate the third factor in (73) as

$$
\begin{equation*}
\left\|u_{3, N_{3}}\right\|_{L_{T}^{\frac{d+2}{2}}}^{L_{x}^{2(d+2)}} \ll T^{\frac{4-d}{2(d+2)}+}\left\langle N_{3}\right\rangle^{\frac{d(d-1)}{2(d+2)}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}} \tag{74}
\end{equation*}
$$

Estimating the $u_{1, N_{1}}, u_{2, N_{2}}$ factors in (73) as in Case 1, using (74) as well as (70) we obtain that the right-hand side of (72) is

$$
\begin{equation*}
\lesssim T^{\epsilon}\left(\left\langle N_{1}\right\rangle^{-\eta}\left\|\phi_{1, N_{1}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{2}\right\rangle^{\zeta_{0}+}\left\|\phi_{2, N_{2}}\right\|_{L_{x}^{2}}\right)\left(\left\langle N_{3}\right\rangle^{\zeta_{0}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}\right) \tag{75}
\end{equation*}
$$

We now conclude the argument as in Case 1.
(2) $d \geq 5$. In this case we use

$$
\begin{align*}
& \left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{\frac{d}{d-1}}} \\
& \leq\left\|u_{1, N_{1}}\right\|_{L_{T, x}^{\frac{2(d+2)}{d}}}\left\|u_{2, N_{2}}\right\|_{L_{T, x}^{\frac{2(d+2)}{d}}}\left\|u_{3, N_{3}}\right\|_{L_{T}^{\frac{d+2}{2}} L_{x}^{\frac{d(d+2)}{d-2}}} . \tag{76}
\end{align*}
$$

By Lemma 3.3, Hölder's inequality, and Proposition 3.2 we obtain that

$$
\left\|u_{3, N_{3}}\right\|_{L_{T}^{\frac{d+2}{2}} L_{x}^{\frac{d(d+2)}{d-2}}} \lesssim T^{0+}\left\langle N_{3}\right\rangle^{\zeta_{0}+}\left\|\phi_{3, N_{3}}\right\|_{L_{x}^{2}}
$$

and we conclude the claim as in the case $d \leq 4$.
Case 2.B: $N_{1} \sim N_{3} \gtrsim N_{2}$. This is analogous to Case 2.A (interchange the roles $N_{2}$ and $\left.N_{3}\right)$.

Case 2.C: $N_{2} \sim N_{3} \gtrsim N_{1}$. In this case, we replace the power of $N_{1}$ in (71) by the corresponding power of $N_{2}$ and obtain

$$
\begin{aligned}
& \langle N\rangle^{-\eta}\left\|P_{N}\left(u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right)\right\|_{L_{T}^{1} L_{x}^{2}} \\
\lesssim & \langle N\rangle^{-\delta}\left\langle N_{2}\right\rangle^{-\eta+\zeta_{0}+\delta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}} \\
\lesssim & \langle N\rangle^{-\delta}\left\langle N_{1}\right\rangle^{-\eta}\left\langle N_{2}\right\rangle^{\zeta_{0}+\delta}\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} L_{x}^{q}}
\end{aligned}
$$

We can now argue as in Case 2.A. It is important to note that the factor of $N_{1}^{0+}$ coming from the application of Proposition 3.2 can be absorbed into the powers of $N_{2}$ and $N_{3}$ in this case. This finishes the proof of (12) for the given choice of parameters.

As mentioned in Remark 1.4, on $\mathbb{D}=\mathbb{T}_{g}^{d}$ the estimate (13) is a immediate consequence of (12). Indeed, if $S_{3}$ denotes the symmetric group, we have

$$
\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} B_{2,1}^{\zeta}} \lesssim \sum_{\sigma \in S_{3}}\left\|\left\langle N_{\sigma(1)}\right\rangle^{\zeta} u_{\sigma(1), N_{\sigma(1)}} u_{\sigma(2), N_{\sigma(2)}} u_{\sigma(3), N_{\sigma(3)}}\right\|_{L_{T}^{1} B_{2,1}^{0}}
$$

Now, (12) with $\eta=0$ yields

$$
\left\|u_{1, N_{1}} u_{2, N_{2}} u_{3, N_{3}}\right\|_{L_{T}^{1} B_{2,1}^{\zeta}} \lesssim\left\|\phi_{1, N_{1}}\right\|_{B_{2,1}^{\zeta}}\left\|\phi_{2, N_{2}}\right\|_{B_{2,1}^{\zeta}}\left\|\phi_{3, N_{3}}\right\|_{B_{2,1}^{\zeta}}
$$

and by summing up with respect to $N_{1}, N_{2}, N_{3}$ we obtain (13).
We now turn to the proof of (14). We start by noting that

$$
\begin{equation*}
\|f\|_{B_{2,1}^{-\zeta_{0}}} \lesssim\|f\|_{L_{x}^{q}} \tag{77}
\end{equation*}
$$

whenever $q \in\left(q_{0}, \infty\right]$ where $q_{0}$ is given by (70). By duality and compactness of the domain, (77) follows if we prove

$$
\begin{equation*}
\|f\|_{L_{x}^{q_{0}^{\prime}-}} \lesssim\|f\|_{B_{2, \infty}^{\zeta_{0}}} \tag{78}
\end{equation*}
$$

Let $p:=q_{0}^{\prime}-$. We then have $p>2$. In order to obtain (78) we use $P_{N}=\widetilde{P}_{N} P_{N}$ and Lemma 3.3 to deduce that

$$
\|f\|_{L_{x}^{p}} \lesssim \delta \sup _{N}\langle N\rangle^{d / 2-d / p+\delta}\left\|P_{N} f\right\|_{L_{x}^{2}}=\|f\|_{B_{2, \infty}^{d / 2-d / p+\delta}} \leq\|f\|_{B_{2, \infty}^{\zeta_{0}}},
$$

for $\delta$ sufficiently small. This proves (78) and hence (77). To conclude the claim, we consider two cases, depending on the dimension.
(1) $2 \leq d \leq 4$. In this case, we use (77) and Hölder's inequality to obtain

$$
\left\|f_{1} f_{2} f_{3}\right\|_{B_{2,1}^{-\zeta_{0}}} \lesssim\left\|f_{1}\right\|_{L_{x}^{\frac{6(d+2)}{2 d+1}}+}\left\|f_{2}\right\|_{L_{x}^{\frac{6(d+2)}{2 d+1}}+}\left\|f_{3}\right\|_{L_{x} \frac{6(d+2)}{2 d+1}+}
$$

and we deduce the claim from the Sobolev embedding

$$
H_{x}^{\frac{d(d+5)}{6(d+2)}+} \hookrightarrow L_{x}^{\frac{6(d+2)}{2 d+1}+}
$$

see Proposition A. 2 below.
(1) $d \geq 5$. In this case, we use

$$
\left\|f_{1} f_{2} f_{3}\right\|_{B_{2,1}^{-\zeta_{0}}} \lesssim\left\|f_{1}\right\|_{L_{x}^{\frac{3 d}{d-1}}+}\left\|f_{2}\right\|_{L_{x}^{\frac{3 d}{d-1}}+}\left\|f_{3}\right\|_{L_{x}^{\frac{3 d}{d-1}}+}
$$

and the Sobolev embedding

$$
H_{x}^{\frac{d}{6}+\frac{1}{3}+} \hookrightarrow L_{x}^{\frac{3 d}{d-1}+}
$$

implies the claimed estimate.
Remark 3.4. One can also adapt the methods developed in [45] to prove an unconditional uniqueness result for (1), but in this way, one obtains a worse range for $s$ than in (16). In particular, this method does not give unconditional uniqueness in the energy class when $d=3$.

## 4. Uniqueness results in the 1d case and the Proof of Theorem 1.9

Let $I=[-T, T]$ for $T>0, \mathbb{D}=\mathbb{T}$ or $\mathbb{D}=\mathbb{R}$. Consider (1) in $1 d$ with initial condition $\phi(0)=\phi_{0} \in L^{2}(\mathbb{D})$. It is well-posed, with uniqueness in some auxiliary space such as $X^{0, \frac{3}{8}}$ [2] or merely $L^{4}(I \times \mathbb{D})$. The latter is obvious in the case $\mathbb{D}=\mathbb{R}$, while in the case $\mathbb{D}=\mathbb{T}$ this follows from [35] and the Christ-Kiselev Lemma. Further, the problem is ill-posed below $L^{2}(\mathbb{T})$, see [34, 36].

If $\mathbb{D}=\mathbb{R}$, it is proved in [17] that (1) is well-posed for initial data in FourierLebesgue spaces. To be precise, let $\mu$ denote the Lebesgue measure on $\mathbb{D}^{*}=\mathbb{R}$ if $\mathbb{D}=\mathbb{R}$ and the counting measure on $\mathbb{D}^{*}=\mathbb{Z}$ if $\mathbb{D}=\mathbb{T}$. Suppose that $\psi(0)=\psi_{0}$ is a ( $2 \pi$-periodic if $\mathbb{D}=\mathbb{T}$ ) tempered distribution satisfying

$$
\left\|\phi_{0}\right\|_{\widehat{L^{r}}}:=\left(\int\left|\widehat{\phi_{0}}(\xi)\right|^{r^{\prime}} d \mu(\xi)\right)^{\frac{1}{r^{\prime}}}<+\infty
$$

for $1<r \leq 2$, and $\frac{1}{r}+\frac{1}{r^{\prime}}=1$. Following $[16,17]$ we define the space $X_{r}^{s, b}$ via

$$
\left.\|\phi\|_{X_{r}^{s, b}}=\left(\left.\int\langle\xi\rangle^{s r^{\prime}} \int_{\mathbb{R}}\langle\tau+| \xi\right|^{2}\right\rangle^{b r^{\prime}}|\widehat{\phi}(\tau, \xi)|^{r^{\prime}} d \tau d \mu(\xi)\right)^{\frac{1}{r^{\prime}}}
$$

and $X_{r}^{s, b}(I)$ is the restriction space. We refer the reader to [16, Section 2] for more details on these spaces. Now, if $\mathbb{D}=\mathbb{R}$, due to [17, Theorem 1] there exists a unique solution $\psi \in X_{r}^{0, b}(I)$ of (1), provided that $T>0$ is small enough, $1<r \leq 2$ and $b>\frac{1}{r}$.

In the case $\mathbb{D}=\mathbb{T}$, the key nonlinear estimate for the corresponding result in the case $\mathbb{D}=\mathbb{T}$ is proved in the introduction of [18]. This case is slightly more delicate and one needs to renormalize the equation first. The well-known gauge transform

$$
\psi(t, x)=\exp \left(i t \frac{1}{\pi} \int_{0}^{2 \pi}|\phi(t, y)|^{2} d y\right) \phi(t, x)
$$

leads to the renormalized NLS

$$
\begin{equation*}
i \partial_{t} \psi+\partial_{x}^{2} \psi=\left(|\psi|^{2}-\frac{1}{\pi} \int_{0}^{2 \pi}|\psi(\cdot, y)|^{2} d y\right) \psi \text { on }(0, T) \times \mathbb{T} \tag{79}
\end{equation*}
$$

The discussion in the introduction of [18] implies that there is a trilinear form $N\left(\psi_{1}, \psi_{2}, \psi_{3}\right)$ with the property that

$$
N(\psi, \psi, \bar{\psi})=\left(|\psi|^{2}-\frac{1}{\pi} \int_{0}^{2 \pi}|\psi(\cdot, y)|^{2} d y\right) \psi
$$

[18, Prop. 1.3] further implies that if $1<r \leq 2, b>\frac{1}{r}, \varepsilon>0$, then for all $\psi_{j} \in X_{r}^{0, b}$ we have $N\left(\psi_{1}, \psi_{2}, \overline{\psi_{3}}\right) \in X_{r}^{0,-\varepsilon}$ and

$$
\left\|N\left(\psi_{1}, \psi_{2}, \overline{\psi_{3}}\right)\right\|_{X_{r}^{0,-\varepsilon}} \lesssim \prod_{j=1}^{3}\left\|\psi_{j}\right\|_{X_{r}^{0, b}}
$$

Then, by using the linear estimates in Lemma C. 1 and the contraction mapping principle, there exists a unique solution $\psi \in X_{r}^{0, b}(I)$, provided that $T>0$ is small enough and $b>\frac{1}{r}$, see [17] or [16, Section 2] for details.
Proof of Theorem 1.9. By invariance under time translations we may assume that $|I|$ is sufficiently small. By assumption $\phi$ satisfies

$$
\begin{equation*}
e^{-i t \Delta} \phi(t)=\phi_{0}-i \int_{0}^{t} e^{-i t^{\prime} \Delta}|\phi|^{2} \phi\left(t^{\prime}\right) d t^{\prime}, t \in I \tag{80}
\end{equation*}
$$

in some Sobolev space of negative order and $\phi$ conserves the $L^{2}$-norm. Let

$$
\psi(t, x)=e^{i t m} \phi(t, x), \quad m=\frac{1}{\pi} \int_{0}^{2 \pi}|\phi(t, y)|^{2} d y
$$

Then,

$$
\begin{aligned}
i \int_{0}^{t} e^{-i t^{\prime} \Delta}\left(|\psi|^{2} \psi\left(t^{\prime}\right)-m \psi\left(t^{\prime}\right)\right) d t^{\prime}= & i \int_{0}^{t} e^{-i t^{\prime} \Delta}|\phi|^{2} \phi\left(t^{\prime}\right) e^{i t^{\prime} m} d t^{\prime} \\
& -\int_{0}^{t} e^{-i t^{\prime} \Delta} \phi\left(t^{\prime}\right) \partial_{t^{\prime}} e^{i t^{\prime} m} d t^{\prime}
\end{aligned}
$$

Now, (80) and integration by parts yield

$$
-\int_{0}^{t} e^{-i t^{\prime} \Delta} \phi\left(t^{\prime}\right) \partial_{t^{\prime}} e^{i t^{\prime} m} d t^{\prime}=\phi_{0}-e^{-i t \Delta} \phi(t) e^{i t m}-i \int_{0}^{t} e^{-i t^{\prime} \Delta}|\phi|^{2} \phi\left(t^{\prime}\right) e^{i t^{\prime} m} d t^{\prime}
$$

which implies

$$
i \int_{0}^{t} e^{-i t^{\prime} \Delta}\left(|\psi|^{2} \psi\left(t^{\prime}\right)-m \psi\left(t^{\prime}\right)\right) d t^{\prime}=\psi_{0}-e^{-i t \Delta} \psi(t)
$$

which is equivalent to

$$
\begin{equation*}
\psi(t)=e^{i t \Delta} \psi_{0}-i \int_{0}^{t} e^{i\left(t-t^{\prime}\right) \Delta}\left(|\psi|^{2} \psi\left(t^{\prime}\right)-m \psi\left(t^{\prime}\right)\right) d t^{\prime}, t \in I \tag{81}
\end{equation*}
$$

hence $\psi$ is a mild solution of $(79)$ in $L^{2}(\mathbb{T})$ with $\psi \in L^{p}(I \times \mathbb{T})$.
Let $f=|\psi|^{2} \psi-m \psi$, extended by zero outside of $I$. It satisfies $f \in L^{r}(\mathbb{R} \times \mathbb{T})$ with $r=p / 3>1$. The statement of the Theorem is relevant only if $p>3$ is small, so without loss we may assume $p \leq 6$. By Hausdorff-Young we have $\widehat{f} \in L^{r^{\prime}}(\mathbb{R} \times \mathbb{Z})$, in other words $f \in X_{r}^{0,0}$. Let $\chi \in C^{\infty}(\mathbb{R})$ with $\chi(t)=1$ for $t \in[-1,1], \operatorname{supp}(\chi) \subset(-2,2)$ and $\chi_{T}(t)=\chi(t / T)$ and let

$$
F(t):=\chi(t) e^{i t \Delta} \psi_{0}-i \chi_{T}(t) \int_{0}^{t} e^{i\left(t-t^{\prime}\right) \Delta} f\left(t^{\prime}\right) d t^{\prime}
$$

For $t \in I$ we have $F(t)=\psi(t)$. Observing $L^{2}(\mathbb{T}) \subset \widehat{L^{r}}(\mathbb{T})$ for $r \leq 2$, Lemma C. 1 implies $F \in X_{r}^{0,1}$ and we obtain $\psi \in X_{r}^{0,1}(I)$. Due to the uniqueness result in $X_{r}^{0,1}(I)$ explained above the proof is complete.

## Appendix A. Proof of Proposition 2.8

For completeness, we present the details of the proof of Proposition 2.8. The results here are well-known, some of them can be slightly improved on certain domains (such as tori). We first recall a general fact concerning the boundedness of the spectral projectors $\chi_{k}$ from the work of Sogge [39]. Here $\mathbb{D}$ is a general smooth compact Riemannian manifold of dimension $d$ without boundary, which is connected and orientable.

Proposition A.1. Let $k \in \mathbb{N}$ be given. We have

$$
\begin{equation*}
\left\|\chi_{k} f\right\|_{L^{\infty}} \lesssim k^{(d-1) / 2}\|f\|_{L^{2}} \tag{82}
\end{equation*}
$$

This result is shown in [39, Proposition 2.1]. Furthermore, we note that the nonendpoint Sobolev embedding holds on $\mathbb{D}$. This estimate suffices for our purposes.
Proposition A. 2 (Non-endpoint Sobolev embedding). The following results hold:
(a) Let $p \geq 2$ be given. Then we have

$$
\|f\|_{L^{p}} \lesssim_{p, s}\|f\|_{H^{s}}
$$

for all $s>\frac{d}{2}-\frac{d}{p}$.
(b) Let $1 \leq p \leq 2$ be given. Then we have

$$
\|f\|_{H^{-s}} \lesssim_{q, s}\|f\|_{L^{p}}
$$

for all $s>\frac{d}{q}-\frac{d}{2}$.
Before proceeding to the proof of Proposition A.2, we recall the Bernstein inequality on general domains $\mathbb{D}$.

Lemma A. 3 (Bernstein's inequality on $\mathbb{D}$ ). For all $N$ and $2 \leq p \leq \infty$, we have

$$
\begin{equation*}
\left\|P_{N} f\right\|_{L^{p}} \lesssim\langle N\rangle^{d / 2-d / p}\left\|P_{N} f\right\|_{L^{2}} \tag{83}
\end{equation*}
$$

Proof of Lemma A.3. We first prove (83) for $p=\infty$. Note that, for fixed $N$

$$
\left\|P_{N} f\right\|_{L^{\infty}} \leq \sum_{k \sim N}\left\|\chi_{k} f\right\|_{L^{\infty}}
$$

which by Proposition A. 1 is

$$
\lesssim \sum_{k \sim N} k^{(d-1) / 2}\left\|\chi_{k} f\right\|_{L^{2}} \lesssim\langle N\rangle^{d / 2}\left\|P_{N} f\right\|_{L^{2}}
$$

In the second step above, we used the Cauchy-Schwarz inequality in $k$. The claim for general $2 \leq p \leq \infty$ follows from (83) with $p=\infty$ by using interpolation.

Proof of Proposition A.2. We observe that part (b) follows from part (a) by duality. Therefore, it suffices to prove part (a). By Lemma A. 3 we have

$$
\|f\|_{L^{p}} \lesssim \sum_{N}\langle N\rangle^{d / 2-d / p-s}\langle N\rangle^{s}\left\|P_{N} f\right\|_{L^{2}} \lesssim\|f\|_{H^{s}}
$$

whenever $s>\frac{d}{2}-\frac{d}{p}$.
Before proving Proposition 2.8, we prove an auxiliary estimate for the product of two functions.

Lemma A.4. Suppose that $\varrho_{1}, \varrho_{2} \in\left(0, \frac{d}{2}\right)$ are given. Then, for all $\delta>0$ we have

$$
\left\|f_{1} f_{2}\right\|_{H^{e_{1}+\varrho_{2}-\frac{d}{2}}} \lesssim \delta\left\|f_{1}\right\|_{H^{\varrho_{1}+\delta}}\left\|f_{2}\right\|_{H^{\varrho_{2}+\delta}}
$$

Proof of Lemma A.4. Let us write $\varrho:=\varrho_{1}+\varrho_{2}-\frac{d}{2}$. We consider two cases.
Case 1: $\varrho>0$. We note that by (10), it suffices to estimate $\left\|f_{1} f_{2}\right\|_{B_{2,1}^{o}}$, which by (9) and the triangle inequality is

$$
\begin{array}{r}
\leq \sum_{\substack{N, N_{1}, N_{2} \\
\max \left\{N_{1}, N_{2}\right\} \gtrsim N}}\langle N\rangle^{\varrho}\left\|P_{N}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right)\right\|_{L^{2}}+\sum_{\substack{N, N_{1}, N_{2} \\
N \gg N_{1}, N_{2}}}\langle N\rangle^{\varrho}\left\|P_{N}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right)\right\|_{L^{2}} \\
=: I+I I \tag{84}
\end{array}
$$

We first estimate the term $I$. By symmetry it suffices to estimate the expression

$$
I^{\prime}:=\sum_{\substack{N, N_{1}, N_{2} \\ N_{1} \gtrsim N}}\langle N\rangle^{\varrho}\left\|P_{N}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right)\right\|_{L^{2}}
$$

Since $\varrho>0$, we can sum in $N$ and obtain that

$$
\begin{equation*}
I^{\prime} \lesssim \varrho \sum_{N_{1}, N_{2}}\left\langle N_{1}\right\rangle^{\varrho}\left\|P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right\|_{L^{2}} \tag{85}
\end{equation*}
$$

We choose $p_{1}, p_{2}$ such that $\frac{1}{p_{1}}=\frac{\varrho_{2}}{d}$ and $\frac{1}{p_{2}}=\frac{1}{2}-\frac{\varrho_{2}}{d}$. By assumption on $\varrho_{2}$, it follows that $p_{1}, p_{2} \in(2, \infty)$. Hence, applying Hölder's inequality, Lemma A.3, and the identities $\varrho+\frac{d}{2}-\frac{d}{p_{1}}=\varrho_{1}, \frac{d}{2}-\frac{d}{p_{2}}=\varrho_{2}$ in (85) it follows that

$$
I^{\prime} \lesssim \sum_{N_{1}, N_{2}}\left\langle N_{1}\right\rangle^{\varrho_{1}}\left\|P_{N_{1}} f_{1}\right\|_{L^{2}}\left\langle N_{2}\right\rangle^{\varrho_{2}}\left\|P_{N_{2}} f_{2}\right\|_{L^{2}}=\left\|f_{1}\right\|_{B_{2,1}^{\varrho_{1}}}\left\|f_{2}\right\|_{B_{2,1}^{\varrho_{2}}}
$$

which by (10) is an acceptable bound.
In order to estimate $I I$ in (84) we first consider fixed $N$ and $N_{1}, N_{2} \ll N$. Let us assume that $N_{1} \leq N_{2}$. We estimate the quantity $\left\|P_{N}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right)\right\|_{L^{2}}$ by duality. Namely, for $g \in L^{2}$ and $k \in \mathbb{N}$ we note that, by the self-adjointness of $\Delta$,

$$
\begin{align*}
& \left|\int P_{N_{1}} f_{1} P_{N_{2}} f_{2} P_{N} g d x\right|=\left|\int \Delta^{k}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right) \Delta^{-k} P_{N} g d x\right| \\
& \leq\left\|\Delta^{k}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right)\right\|_{L^{2}}\left\|\Delta^{-k} P_{N} g\right\|_{L^{2}} \\
& \lesssim_{k}\langle N\rangle^{-2 k}\left\langle N_{1}\right\rangle^{d / 2}\left\langle N_{2}\right\rangle^{2 k}\left\|P_{N_{1} f_{1}}\right\|_{L^{2}}\left\|P_{N_{2}} f_{2}\right\|_{L^{2}}\|g\|_{L^{2}} \tag{86}
\end{align*}
$$

Note that, in order to obtain (86), we used a Leibniz rule for $\Delta^{k}$, Hölder's inequality, Lemma A. 3 with $p=\infty$, together with the assumption that $N_{1} \leq N_{2}$. Moreover, we used the observation that

$$
\left\|\Delta^{-k} P_{N}\right\|_{L^{2} \rightarrow L^{2}} \lesssim_{k}\langle N\rangle^{-2 k}
$$

We hence obtain from (86) and duality that for all $k \in \mathbb{N}$

$$
\begin{equation*}
\left\|P_{N}\left(P_{N_{1}} f_{1} P_{N_{2}} f_{2}\right)\right\|_{L^{2}} \lesssim_{k}\langle N\rangle^{-2 k}\left\langle N_{1}\right\rangle^{d / 2}\left\langle N_{2}\right\rangle^{2 k}\left\|P_{N_{1} f_{1}}\right\|_{L^{2}}\left\|P_{N_{2}} f_{2}\right\|_{L^{2}} \tag{87}
\end{equation*}
$$

Let us observe that if $k>\frac{\varrho_{2}}{2}$, we have for $\delta>0$

$$
\begin{aligned}
& \sum_{\substack{N, N_{1}, N_{2} \\
N_{1} \leq N_{2} \ll N}}\langle N\rangle^{\varrho_{1}+\varrho_{2}-d / 2-2 k}\left\langle N_{1}\right\rangle^{d / 2-\varrho_{1}-\delta}\left\langle N_{2}\right\rangle^{2 k-\varrho_{2}-\delta} \\
\lesssim & \sum_{\substack{N, N_{1}, N_{2} \\
N_{1} \leq N_{2} \ll N}}\left(\frac{\left\langle N_{1}\right\rangle}{\langle N\rangle}\right)^{d / 2-\varrho_{1}}\left\langle N_{1}\right\rangle^{-\delta}\left\langle N_{2}\right\rangle^{-\delta} \lesssim 1 .
\end{aligned}
$$

Substituting this bound into (87) and using (10), we obtain the desired estimate on $I I$ in (84).

Case 2: $\varrho \leq 0$. We let $q_{1}, q_{2}$ be such that $\frac{1}{q_{1}}=\frac{1}{2}-\frac{\varrho_{1}}{d}$ and $\frac{1}{q_{2}}=\frac{1}{2}-\frac{\varrho_{2}}{d}$. By the assumptions on $\varrho_{1}, \varrho_{2}$, it follows that $q_{1}, q_{2} \in(2, \infty)$. Hence $\frac{1}{q_{1}}+\frac{1}{q_{2}}=\frac{1}{q}$ for some $q \in(1, \infty)$. For such $q$ we have $\frac{1}{2}=\frac{1}{q}+\frac{\varrho}{d}$. Since $\varrho \leq 0$ we can use Proposition A. 2 (b) with $s=-\varrho$ and $\tilde{q}=q+\delta_{0}$, with $\delta_{0}>0$, to deduce that

$$
\begin{equation*}
\left\|f_{1} f_{2}\right\|_{H^{e}} \lesssim\left\|f_{1} f_{2}\right\|_{L^{\tilde{q}}} . \tag{88}
\end{equation*}
$$

Using Hölder's inequality, it follows that the right-hand side of (88) is

$$
\leq\left\|f_{1}\right\|_{L^{q_{1}+\delta_{1}}}\left\|f_{2}\right\|_{L^{q_{2}+\delta_{2}}}
$$

for appropriate $\delta_{1}, \delta_{2}>0$. By Proposition A. 2 (a), this expression is

$$
\lesssim\left\|f_{1}\right\|_{H^{e_{1}+\delta}}\left\|f_{2}\right\|_{H^{e_{2}+\delta}},
$$

provided that we choose the parameters $\delta_{0}, \delta_{1}, \delta_{2}$ small in terms of $\delta$. This is an admissible upper bound.

We can now prove Proposition 2.8.
Proof of Proposition 2.8. We note that, by the assumptions of the proposition, we have $2 \varrho-\frac{d}{2} \in\left(0, \frac{d}{2}\right)$. By applying Lemma A. 4 to $f_{1} f_{2}$ and $f_{3}$, we obtain

$$
\left\|f_{1} f_{2} f_{3}\right\|_{H^{\left(2 \varrho-\frac{d}{2}\right)+\varrho-\frac{d}{2}}} \lesssim \delta\left\|f_{1} f_{2}\right\|_{H^{\left(2 \varrho-\frac{d}{2}\right)+\delta / 2}}\left\|f_{3}\right\|_{H^{\varrho+\delta / 2}},
$$

We can hence apply Lemma A. 4 again to estimate the first factor above and obtain the claim.

## Appendix B. Proof of Lemma 3.3

After an anisotropic rescaling of the torus we may assume that $\theta_{1}=\cdots=\theta_{d}=1$ and $\tilde{P}_{N}$ is of the form

$$
\tilde{P}_{N} f=\sum_{k \in \mathbb{Z}^{d}} \rho_{N}(k) \widehat{f}(k) e^{i x \cdot k}
$$

for $\rho_{N}(k)=\rho(k / N)$ if $N \geq 1$, with some functions $\rho, \rho_{0} \in \mathcal{S}\left(\mathbb{R}^{d}\right)$, and the standard $k$-th Fourier coefficient $\widehat{f}(k)$ of $f$. By the Poisson summation formula we have

$$
\tilde{P}_{N} f=\psi_{N} *_{\mathbb{T}^{d}} f, \text { where } \psi_{N}(x)=(2 \pi)^{d / 2} \sum_{k \in \mathbb{Z}^{d}} \mathcal{F}^{-1}\left(\rho_{N}\right)(x+k)
$$

where $\mathcal{F}$ denotes the Fourier transform on $\mathbb{R}^{d}$. Now, for any $1 \leq r<\infty$ and $N \geq 1$,

$$
\begin{aligned}
\left\|\psi_{N}\right\|_{L^{r}\left([0,2 \pi]^{d}\right)}^{r} & \lesssim \sum_{k \in \mathbb{Z}^{d}}\left\|\mathcal{F}^{-1}\left(\rho_{N}\right)\right\|_{L^{r}\left(k+[0,2 \pi]^{d}\right)}^{r} \leq\left\|\mathcal{F}^{-1}\left(\rho_{N}\right)\right\|_{L^{r}\left(\mathbb{R}^{d}\right)}^{r} \\
& =N^{d(r-1)}\left\|\mathcal{F}^{-1}(\rho)\right\|_{L^{r}\left(\mathbb{R}^{d}\right)}^{r},
\end{aligned}
$$

and similarly for $r=\infty$ or $N=0$. If $1 \leq p \leq q \leq \infty$, Young's inequality implies

$$
\left\|\tilde{P}_{N} f\right\|_{L^{q}} \lesssim\langle N\rangle^{d(1-1 / r)}\|f\|_{L^{p}}
$$

where $1-\frac{1}{r}=\frac{1}{p}-\frac{1}{q}$.

## Appendix C. Linear estimates in Fourier-Lebesgue spaces

Here, we include the linear estimates used in the proof of Theorem 1.9, which are also needed to perform the contraction argument. These are due to Grünrock, see [16, Section 2]. Here, we adapt the argument given in [18, Proof of Lemma 7.1], which follows the strategy of [10, Lemma 3.1].

Lemma C.1. Let $1<r<\infty, b \geq 0$. Let $\chi \in C^{\infty}(\mathbb{R})$ with $\chi(t)=1$ for $t \in[-1,1]$, $\operatorname{supp}(\chi) \subset(-2,2)$ and $\chi_{T}(t)=\chi(t / T)$. For all $0<T \leq 1$ and $\phi_{0} \in \widehat{L^{r}}$ we have $\chi_{T} e^{i t \Delta} \phi_{0} \in X_{r}^{0, b}$ and

$$
\left\|\chi_{T} e^{i t \Delta} \psi_{0}\right\|_{X_{r}^{0, b}} \lesssim T^{\frac{1}{r}-b}\left\|\psi_{0}\right\|_{\widehat{L^{r}}}
$$

Further, if $\beta \geq b-1$ and $\frac{1}{r}>\beta>-\frac{1}{r^{\prime}}$, for all $f \in X_{r}^{0, \beta}$ we have

$$
\chi_{T} \int_{0}^{t} e^{i\left(t-t^{\prime}\right) \Delta} f\left(t^{\prime}\right) d t^{\prime} \in X_{r}^{0, b}
$$

and

$$
\left\|\chi_{T} \int_{0}^{t} e^{i\left(t-t^{\prime}\right) \Delta} f\left(t^{\prime}\right) d t^{\prime}\right\|_{X_{r}^{0, b}} \lesssim T^{1+\beta-b}\|f\|_{X_{r}^{0, \beta}}
$$

Proof. If $\mathcal{F}$ denotes the space-time Fourier transform, we obtain

$$
\mathcal{F}\left(\chi_{T}(t) e^{i t \Delta} \psi_{0}\right)(\tau, \xi)=T \widehat{\chi}\left(T\left(\tau+|\xi|^{2}\right) \widehat{\psi_{0}}(\xi)\right.
$$

and, using $\widehat{\chi} \in \mathcal{S}(\mathbb{R})$,

$$
\begin{align*}
\left\|\chi_{T} e^{i t \Delta} \psi_{0}\right\|_{X_{r}^{0, b}} & \left.=\left(\left.\int\left|\widehat{\psi_{0}}(\xi)\right|^{r^{\prime}} \int_{\mathbb{R}}\langle\tau+| \xi\right|^{2}\right\rangle^{b r^{\prime}}\left|T \widehat{\chi}\left(T\left(\tau+|\xi|^{2}\right)\right)\right|^{r^{\prime}} d \tau d \mu(\xi)\right)^{\frac{1}{r^{\prime}}} \\
& \lesssim T^{\frac{1}{r}-b}\left\|\psi_{0}\right\|_{\widehat{L^{r}}}, \tag{89}
\end{align*}
$$

for any $b \in \mathbb{R}$. For the integrable function $\sigma_{T}\left(t^{\prime}\right)=\frac{1}{2} \chi_{4 T}\left(t^{\prime}\right) \operatorname{sign}\left(t^{\prime}\right)$ and $\left|t^{\prime}\right|,|t| \leq 2 T$ we have $\sigma_{T}\left(t^{\prime}\right)+\sigma_{T}\left(t-t^{\prime}\right)=1$ if $0 \leq t^{\prime} \leq t$ and zero otherwise. Hence, we can decompose

$$
\chi_{T}(t) \int_{0}^{t} e^{i\left(t-t^{\prime}\right) \Delta} f\left(t^{\prime}\right) d t^{\prime}=\chi_{T}(t) e^{i t \Delta} I_{1}+\chi_{T}(t) I_{2}(t)
$$

where

$$
I_{1}=\int_{\mathbb{R}} \sigma_{T}\left(t^{\prime}\right) e^{-i t^{\prime} \Delta} f\left(t^{\prime}\right) d t^{\prime}, \quad I_{2}(t)=\int_{\mathbb{R}} e^{i\left(t-t^{\prime}\right) \Delta} \sigma_{T}\left(t-t^{\prime}\right) f\left(t^{\prime}\right) d t^{\prime}
$$

An application of (89) gives

$$
\left\|\chi_{T} e^{i t \Delta} I_{1}\right\|_{X_{r}^{0, b}} \lesssim T^{\frac{1}{r}-b}\left\|I_{1}\right\|_{\widehat{L^{r}}}
$$

By Parseval's identity,

$$
\mathcal{F}_{x} \int_{\mathbb{R}} \sigma_{T}\left(t^{\prime}\right) e^{-i t^{\prime} \Delta} f\left(t^{\prime}\right) d t^{\prime}(\xi)=\int_{\mathbb{R}} \overline{\mathcal{F}_{t} \sigma_{T}\left(\tau+|\xi|^{2}\right)} \mathcal{F} f(\tau, \xi) d \tau
$$

and due to $\left|\mathcal{F}_{t} \sigma_{T}(\tau)\right| \lesssim T\langle T \tau\rangle^{-1}$ we obtain

$$
\left\|\chi_{T} e^{i t \Delta} I_{1}\right\|_{X_{r}^{0, b}} \lesssim T^{1+\frac{1}{r}-b}\left(\int\left(\int_{\mathbb{R}}\left\langle T\left(\tau+|\xi|^{2}\right)\right\rangle^{-1}|\mathcal{F} f(\tau, \xi)| d \tau\right)^{r^{\prime}} d \mu(\xi)\right)^{\frac{1}{r^{\prime}}}
$$

Hölder's inequality implies

$$
\left.\int_{\mathbb{R}}\left\langle T\left(\tau+|\xi|^{2}\right)\right\rangle^{-1}|\mathcal{F} f(\tau, \xi)| d \tau \lesssim T^{\beta-\frac{1}{r}}\left(\left.\int_{\mathbb{R}}\langle\tau+| \xi\right|^{2}\right\rangle^{\beta r^{\prime}}|\mathcal{F} f(\tau, \xi)|^{r^{\prime}} d \tau\right)^{\frac{1}{r^{\prime}}}
$$

for any $\frac{1}{r}>\beta>-\frac{1}{r^{\prime}}$, which implies

$$
\left\|\chi_{T} e^{i t \Delta} I_{1}\right\|_{X_{r}^{0, b}} \lesssim T^{1+\beta-b}\|f\|_{X_{r}^{0, \beta}}
$$

Concerning the second contribution we first observe that

$$
\left.\left.\left.\langle\tau+| \xi\right|^{2}\right\rangle\left.^{b} \lesssim\left\langle\tau^{\prime}+\right| \xi\right|^{2}\right\rangle^{b}+\left|\tau-\tau^{\prime}\right|^{b}
$$

which implies

$$
\left\|\chi_{T} I_{2}\right\|_{X_{r}^{0, b}} \lesssim\left\|\widehat{\chi_{T}}\right\|_{L^{1}}\left\|I_{2}\right\|_{X_{r}^{0, b}}+\left\||\cdot|^{b} \widehat{\chi_{T}}\right\|_{L^{1}}\left\|I_{2}\right\|_{X_{r}^{0,0}}
$$

We have $\left\|\left\|\cdot \mid{ }^{a} \widehat{\chi_{T}}\right\|_{L^{1}} \lesssim T^{-a}\right.$ for any $a \geq 0$, and

$$
\mathcal{F} I_{2}(\tau, \xi)=\mathcal{F}_{t} \sigma_{T}\left(\tau+|\xi|^{2}\right) \mathcal{F} f(\tau, \xi)
$$

We obtain

$$
\left\|I_{2}\right\|_{X_{r}^{0, b}} \lesssim T\left(\iint_{\mathbb{R}} \frac{\left.\left.\langle\tau+| \xi\right|^{2}\right\rangle^{r^{\prime} b}}{\left\langle T\left(\tau+|\xi|^{2}\right)\right\rangle^{r^{\prime}}}|\mathcal{F} f(\tau, \xi)|^{r^{\prime}} d \tau d \mu(\xi)\right)^{\frac{1}{r^{\prime}}} \lesssim T^{1+\beta-b}\|f\|_{X_{r}^{0, \beta}}
$$

which implies the claim.

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[^1]:    ${ }^{1}$ In (18) we are just formally restricting the function $\sigma^{(k+1)}$ to the set where $x_{k+1}=x_{k+1}^{\prime}=x_{j}$ and $x_{j}^{\prime}=x_{k+1}=x_{j}^{\prime}$ respectively.

[^2]:    ${ }^{2}$ We change the convention from [8] and denote elements of $L^{2}(\mathbb{D})$ by $f$ instead of $\phi$.

