# Invariant Gibbs measures for dispersive PDEs 

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Chapter 3 : Almost sure global wellposedness of the LLL equation below $L^{2}(\mathbb{C})$

In this chapter, we show how we can use a Gibbs measure to prove almost sure global existence results. We will present the method on the Landau Lowest Level (LLL) equation. The results are taken from Germain-Hani-Thomann and some analysis from Germain-Hani-Thomann and Gérard-Germain-Thomann.

For $1 \leq p \leq+\infty$ we define the Bargmann-Fock spaces

$$
F^{p}(\mathbb{C})=\left\{u(z)=\mathrm{e}^{-\frac{|z|^{2}}{2}} f(z), f \text { entire holomorphic }\right\} \cap L^{p}(\mathbb{C})
$$

In the sequel we consider the Lowest Landau Level equation which reads

$$
\left\{\begin{array}{l}
i \partial_{t} u=\Pi\left(|u|^{2} u\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C}  \tag{LLL}\\
u(0, z)=u_{0}(z)
\end{array}\right.
$$

where $\Pi$ is the orthogonal projector on $F^{2}(\mathbb{C})$.

This equation is used in the description of fast rotating Bose-Einstein condensates, see e.g. the book of Aftalion and references therein.

The equation (LLL) can be obtained as the restriction of the continuous resonant equation (CR) which was introduced by Faou-Germain-Hani and further studied by Germain-Hani-Thomann.

Let $z=x+i y$. Denote by $H$ the harmonic oscillator $H=-\partial_{x}^{2}-\partial_{y}^{2}+x^{2}+y^{2}$. A Hilbertian basis of normalized eigenfunctions of $H$ for $F^{2}(\mathbb{C})$ is given by the so-called special Hermite functions defined for $n \geq 0$ by

$$
\varphi_{n}(z)=\frac{1}{\sqrt{\pi n!}} z^{n} \mathrm{e}^{-|z|^{2} / 2},
$$

and which satisfy

$$
H \varphi_{n}=2(n+1) \varphi_{n}
$$

Therefore, every $u \in F^{2}(\mathbb{C})$ can be decomposed in a series

$$
\begin{equation*}
u=\sum_{n=0}^{+\infty} c_{n} \varphi_{n} \tag{1}
\end{equation*}
$$

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and which satisfy

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$$

Therefore, every $u \in F^{2}(\mathbb{C})$ can be decomposed in a series

$$
\begin{equation*}
u=\sum_{n=0}^{+\infty} c_{n} \varphi_{n} \tag{2}
\end{equation*}
$$

We are able to explicitly compute the kernel of $\Pi$

$$
\sum_{n=0}^{+\infty} \varphi_{n}(z) \overline{\varphi_{n}(w)}=\frac{1}{\pi}\left(\sum_{n=0}^{+\infty} \frac{1}{n!}(z \bar{w})^{n}\right) \mathrm{e}^{-|z|^{2} / 2-|w|^{2} / 2}=\frac{1}{\pi} \mathrm{e}^{z \bar{w}-|z|^{2} / 2-|w|^{2} / 2}
$$

As a consequence,

$$
[\Pi u](z)=\frac{1}{\pi} \mathrm{e}^{-\frac{|z|^{2}}{2}} \int_{\mathbb{C}} \mathrm{e}^{z \bar{w}-\frac{|w|^{2}}{2}} u(w) d L(w)
$$

where $d L$ stands for the Lebesgue measure on $\mathbb{C}$.

We define the trilinear operator $\mathcal{T}$ by

$$
\begin{equation*}
\mathcal{T}\left(u_{1}, u_{2}, u_{3}\right)=\Pi\left(u_{1} u_{2} \overline{u_{3}}\right) \tag{3}
\end{equation*}
$$

The equation (LLL) is Hamiltonian : indeed, introducing the functional

$$
\begin{aligned}
\mathcal{E}\left(u_{1}, u_{2}, u_{3}, u_{4}\right) & \stackrel{\text { def }}{=}\left\langle\mathcal{T}\left(u_{1}, u_{2}, u_{3}\right), u_{4}\right\rangle_{L^{2}(\mathbb{C})} \\
& =\int_{\mathbb{C}}\left(u_{1} u_{2} \overline{u_{3} u_{4}}\right)(z) d L(z)
\end{aligned}
$$

and setting

$$
\mathcal{E}(u):=\mathcal{E}(u, u, u, u)=\int_{\mathbb{C}}|u(z)|^{4} d L(z)=\|u\|_{L^{4}(\mathbb{C})}^{4}
$$

then (LLL) derives from the Hamiltonian $\mathcal{E}$ given the symplectic form

$$
\omega(f, g)=\mathfrak{I m} \int_{\mathbb{C}} f \bar{g} d L
$$

so that (LLL) is equivalent to

$$
i \partial_{t} u=\frac{1}{2} \frac{\partial \mathcal{E}(u)}{\partial \bar{u}}
$$

The family $\left(\varphi_{n}\right)_{n \geq 0}$ is particularly well adapted in the study of the operator $\mathcal{T}$ since on has

$$
\begin{equation*}
\mathcal{T}\left(\varphi_{n_{1}}, \varphi_{n_{2}}, \varphi_{n_{3}}\right)=\alpha_{n_{1}, n_{2}, n_{3}, n_{4}} \varphi_{n_{\mathbf{4}}}, \quad n_{4}=n_{1}+n_{2}-n_{3} \tag{4}
\end{equation*}
$$

with

$$
\alpha_{n_{1}, n_{2}, n_{3}, n_{4}}=\mathcal{E}\left(\varphi_{n_{1}}, \varphi_{n_{2}}, \varphi_{n_{3}}, \varphi_{n_{4}}\right)=\frac{\pi}{2} \frac{\left(n_{1}+n_{2}\right)!}{2^{n_{1}+n_{2}} \sqrt{n_{1}!n_{2}!n_{3}!n_{4}!}} \mathbf{1}_{n_{1}+n_{2}=n_{3}+n_{4}}
$$

We can prove that $\mathrm{e}^{i t H} \mathcal{T}\left(u_{1}, u_{2}, u_{3}\right)=\mathcal{T}\left(\mathrm{e}^{i t H} u_{1}, \mathrm{e}^{i t H} u_{2}, \mathrm{e}^{i t H} u_{3}\right)$, and therefore with the change of unknowns $v=\mathrm{e}^{i t h} u$ we see that (LLL) is equivalent to the equation

$$
\begin{equation*}
i \partial_{t} v+H v=\Pi\left(|v|^{2} v\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C} \tag{5}
\end{equation*}
$$

## Some deterministic results

Well-posedness of the LLL equation
Define the harmonic Sobolev spaces for $s \in \mathbb{R}$, by

$$
\mathcal{H}^{s}=\mathcal{H}^{s}(\mathbb{C})=\left\{u \in \mathcal{S}^{\prime}(\mathbb{C}), H^{s / 2} u \in L^{2}(\mathbb{C})\right\}
$$

This is a weighted Sobolev norm. In the Bargmann-Fock space, it simply corresponds to a weighted $L^{2}$-norm. Set $\langle z\rangle=\left(1+|z|^{2}\right)^{1 / 2}$, then we have

## Lemma

Let $s \in \mathbb{R}$. There exists $C>0$ such that for all $u \in F^{2}(\mathbb{C}) \cap \mathcal{H}^{s}(\mathbb{C})$

$$
\frac{1}{C}\left\|\langle z\rangle^{s} u\right\|_{L^{2}(\mathbb{C})} \leq\|u\|_{\mathcal{H}^{s}(\mathbb{C})} \leq C\left\|\langle z\rangle^{s} u\right\|_{L^{2}(\mathbb{C})}
$$

## Exercise

Prove the previous lemma in the particular case where $s \in 2 \mathbb{N}$. Hint : use the decomposition (2), and the relations $z \varphi_{n}=\sqrt{n+1} \varphi_{n+1}$ and $H \varphi_{n}=2(n+1) \varphi_{n}$.

## Proposition

The following quantities are conservation laws for (LLL) :

$$
\begin{array}{ll}
\mathcal{E}(u)=\int_{\mathbb{C}}|u(z)|^{4} d L(z) & \text { (Hamiltonian) } \\
M(u)=\int_{\mathbb{C}}|u(z)|^{2} d L(z) & \text { (Angular momentum) } \\
P(u)=\int_{\mathbb{C}}\left(|z|^{2}-1\right)|u(z)|^{2} d L(z) & \text { (Mass) } \\
Q(u)=\int_{\mathbb{C}} z|u|^{2}(z) d L(z) & \text { (Magnetic momentum). }
\end{array}
$$

Notice that the $\mathcal{H}^{1}$ norm is also preserved, since in coordinates we can check that

$$
\int_{\mathbb{C}}\left|H^{1 / 2} u(z)\right|^{2} d L(z)=2 \int_{\mathbb{C}}|z|^{2}|u(z)|^{2} d L(z)=2(P(u)+M(u))
$$

An important tool in the study of the (LLL) equation are the hypercontractivity inequalities of Carlen.

## Proposition

Assume that $1 \leq p \leq q \leq \infty$. Then $F^{p}(\mathbb{C}) \subset F^{q}(\mathbb{C})$ and

$$
\begin{equation*}
\left(\frac{q}{2 \pi}\right)^{1 / q}\|u\|_{L^{q}(\mathbb{C})} \leq\left(\frac{p}{2 \pi}\right)^{1 / p}\|u\|_{L^{p}(\mathbb{C})} \tag{6}
\end{equation*}
$$

with optimal constants.
This result can be understood as smoothing estimate in the $L^{p}$ scale. Compare with the Khintchine Lemma.

Proof: We prove prove the result for $p=1$ and $q=+\infty$. Write $u(z)=f(z) \mathrm{e}^{-|z|^{2} / 2}$ where $f$ is entire. By the Cauchy formula, for all $r>0$,

$$
|f(0)| \leq \frac{1}{2 \pi} \int_{0}^{2 \pi}\left|f\left(r e^{i \theta}\right)\right| d \theta
$$

Thus by integration in $r>0$

$$
|f(0)| \int_{0}^{+\infty} r \mathrm{e}^{-r^{2} / 2} d r \leq \frac{1}{2 \pi} \int_{0}^{2 \pi} \int_{0}^{+\infty}\left|f\left(r \mathrm{e}^{i \theta}\right)\right| r \mathrm{e}^{-r^{2} / 2} d r d \theta
$$

in other words

$$
\begin{aligned}
&|u(0)|=|f(0)| \leq \frac{1}{2 \pi} \int_{0}^{2 \pi} \int_{0}^{+\infty}\left|f\left(r \mathrm{e}^{i \theta}\right)\right| r \mathrm{e}^{-r^{2} / 2} d r d \theta= \\
&=\frac{1}{2 \pi} \int_{\mathbb{C}}|f(z)| \mathrm{e}^{-|z|^{2} / 2} d L(z)=\frac{1}{2 \pi}\|u\|_{L^{1}(\mathbb{C})}
\end{aligned}
$$

More generally, for any $z \in \mathbb{C}$ and $f$ we apply the previous inequality to the entire function

$$
w \longmapsto f(z-w) \mathrm{e}^{w \bar{z}-|z|^{2} / 2}
$$

and deduce the announced bound $\|u\|_{L^{\infty}(\mathbb{C})} \leq \frac{1}{2 \pi}\|u\|_{L^{1}(\mathbb{C})} . \square$

As a consequence, we observe that for all $u \in F^{2}(\mathbb{C})$

$$
\mathcal{E}(u)=\|u\|_{L^{4}(\mathbb{C})}^{4} \leq \frac{1}{2 \pi}\|u\|_{L^{2}(\mathbb{C})}^{4}
$$

We refer to the book [Zhu] for more analysis on Bargmann-Fock spaces.

## Exercise

1. Show that with a slight modification in the previous proof one can also obtain the case $q=\infty$ and any $p \geq 1$.
2. Prove directly the inequality (6) for $(q, p)=(\infty, 2)$. Hint : use the identity

$$
\int_{\mathbb{C}} e^{-|w|^{2}+a w+c \bar{w}} d L(w)=\pi e^{a c}
$$

We are now able to show that (LLL) is globally well-posed in $F^{p}(\mathbb{C})$ with $2 \leq p \leq 4$.

## Proposition (Gérard-Germain-LT)

Assume that $2 \leq p \leq 4$. The equation (LLL) is globally well-posed for data in $F^{p}(\mathbb{C})$ and such data lead to solutions in $\mathcal{C}^{\infty}\left(\mathbb{R}, F^{p}(\mathbb{C})\right)$. Moreover, there exists $C=C\left(\left\|u_{0}\right\|_{L^{p}(\mathbb{C})}\right)>0$ such that

$$
\begin{equation*}
\left\|u(t)-u_{0}\right\|_{L^{p}(\mathbb{C})} \leq C|t|^{4 / p-1}, \quad\left\|u(t)-u_{0}\right\|_{L^{2}(\mathbb{C})} \leq C|t|, \quad \forall t \in \mathbb{R} \tag{7}
\end{equation*}
$$

Proof: Local well-posedness is obtained by a fixed point argument from the following a priori estimate : using successively the boundedness of $\Pi$, Hölder's inequality, and (6),

$$
\left\|\Pi\left(|u|^{2} u\right)\right\|_{L^{p}} \leq C_{1}\left\||u|^{2} u\right\|_{L^{p}}=C_{1}\|u\|_{L^{3 p}}^{3} \leq C_{2}\|u\|_{L^{4}}^{2}\|u\|_{L^{p}}
$$

The previous inequality shows that the lifespan of the solution only depend on the $L^{4}$ norm which is preserved, hence we get global well-posedness.

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$$

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Let us now prove the bound (7). We write $u=u_{0}+v$, then for $t \geq 0$ we have

$$
v(t)=-i \int_{0}^{t} \mathcal{T}\left(u_{0}+v\right)(s) d s
$$

We take the $L^{2}$-norm and get with the help of (6)

$$
\|v(t)\|_{L^{2}(\mathbb{C})} \leq C_{1} t\left\|u_{0}+v\right\|_{L^{6}(\mathbb{C})}^{3} \leq C_{2} t\left(\left\|u_{0}\right\|_{L^{6}(\mathbb{C})}^{3}+\|v\|_{L^{6}(\mathbb{C})}^{3}\right) \leq C_{3} t\left(\left\|u_{0}\right\|_{L^{p}(\mathbb{C})}^{3}+\|v\|_{L^{4}(\mathbb{C})}^{3}\right)
$$

Therefore, by the conservation of the energy, we obtain $\|v(t)\|_{L^{2}(\mathbb{C})} \leq C t$ which is the second bound. The first bound follows from interpolation with the energy.

## KAM results for a perturbed equation

In the sequel, we consider the (non-local) perturbation of the (LLL) equation

$$
\begin{equation*}
i \partial_{t} u+\nu M u=\varepsilon \Pi\left(|u|^{2} u\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C}, \tag{8}
\end{equation*}
$$

where $\nu, \varepsilon>0$ are small and where $M$ is the (Hermite) multiplier, defined by $M \varphi_{j}=\xi_{j} \varphi_{j}$ with $-1 \leq \xi_{j} \leq 1$.
Notice that $M$ and $H$ commute and that we have the following conservation laws:
$\int_{\mathbb{C}}|u(z)|^{2} d L(z), \quad \int_{\mathbb{C}} \bar{u} H u(z) d L(z), \quad \nu \int_{\mathbb{C}} \bar{u} M u(z) d L(z)+\varepsilon \int_{\mathbb{C}}|u(z)|^{4} d L(z)$,
which are the $L^{2}$ and $\mathcal{H}^{1}$ norms as well the Hamiltonian (there are other conservation laws).

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\end{equation*}
$$

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which are the $L^{2}$ and $\mathcal{H}^{1}$ norms as well the Hamiltonian (there are other conservation laws).

Using the commutation of $M$ and $H$, as well as the relation

$$
\mathrm{e}^{i t H} \mathcal{T}\left(u_{1}, u_{2}, u_{3}\right)=\mathcal{T}\left(\mathrm{e}^{i t H} u_{1}, \mathrm{e}^{i t H} u_{2}, \mathrm{e}^{i t H} u_{3}\right)
$$

we see that (10) is equivalent to the equation $\left(v=e^{i t h} u\right)$

$$
\begin{equation*}
i \partial_{t} v+H v+\nu M v=\Pi\left(|v|^{2} v\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C} \tag{11}
\end{equation*}
$$

The abstract KAM result of Grébert-Thomann can directly be applied to the equation (11) and hence (10)

## Theorem

Let $n \geq 1$ be an integer and set $\mathcal{A}=[-1,1]^{n+1}$. There exist $\varepsilon_{0}>0, \nu_{0}>0$, $C_{0}>0$ and, for each $\varepsilon<\varepsilon_{0}$, a Cantor set $\mathcal{A}_{\varepsilon} \subset \mathcal{A}$ of asymptotic full measure when $\varepsilon \rightarrow 0$, such that for each $\xi \in \mathcal{A}_{\varepsilon}$ and for each $C_{0} \varepsilon \leq \nu<\nu_{0}$, the solution of

$$
\begin{equation*}
i \partial_{t} u+\nu M u=\varepsilon \Pi\left(|u|^{2} u\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C} \tag{12}
\end{equation*}
$$

with initial datum

$$
\begin{equation*}
u_{0}(z)=\sum_{j=0}^{n} I_{j}^{1 / 2} e^{i \theta_{j}} \varphi_{j}(z) \tag{13}
\end{equation*}
$$

with $\left(I_{0}, \cdots, I_{n}\right) \subset(0,1]^{n+1}$ and $\theta \in \mathbb{T}^{n+1}$, is quasi periodic with a quasi period $\omega^{\star}$ close to $\omega_{0}=(2 j+2)_{j=0}^{n}:\left|\omega^{\star}-\omega_{0}\right|<C \nu$.

## Control of Sobolev norms for a perturbed equation

We define the Hermite multiplier $M$ by $M \varphi_{j}=m_{j} \varphi_{j}$, where $\left(m_{j}\right)_{j \in \mathbb{N}}$ is a bounded sequence of real numbers chosen in the following classes : for any $k \geq 1$, we define the class

$$
\mathcal{W}_{k}=\left\{\left(m_{j}\right)_{j \in \mathbb{N}}: \quad m_{j}=\frac{\tilde{m}_{j}}{(j+1)^{k}} \text { with } \quad \tilde{m}_{j} \in[-1 / 2,1 / 2]\right\}
$$

which is endowed with the product Lebesgue (probability) measure. Consider the problem

$$
\begin{equation*}
i \partial_{t} u+M u=\Pi\left(|u|^{2} u\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C} \tag{14}
\end{equation*}
$$

The following almost global existence result is proved in [Grébert-Imekraz-Paturel].

## Theorem

Let $k, r \in \mathbb{N}$. There exists a set $\mathcal{B}_{k} \subset \mathcal{W}_{k}$ of measure 1 such that if $\left(m_{j}\right)_{j \in \mathbb{N}} \in \mathcal{B}_{k}$ there exists $s_{0} \in \mathbb{N}$ such that for any $s \geq s_{0}$, there are $\varepsilon_{0}>0$, $c>0$, such that for any $\varepsilon \in\left(0, \varepsilon_{0}\right)$, for any $u_{0} \in \mathcal{H}^{5}(\mathbb{C})$ with

$$
\left\|u_{0}\right\|_{\mathcal{H}^{s}(\mathbb{C})} \leq \varepsilon,
$$

the equation (14) with initial datum $u_{0}$ has a unique global solution $u \in \mathcal{C}^{\infty}\left(\mathbb{R}, \mathcal{H}^{s}(\mathbb{C})\right)$ and it satisfies

$$
\|u(t)\|_{\mathcal{H}^{s}(\mathbb{C})} \leq 2 \varepsilon, \quad|t| \leq c \varepsilon^{-r} .
$$

To prove this result, we apply the result of Grébert-Imekraz-Paturel to the equation $i \partial_{t} v+H v+M v=\Pi\left(|v|^{2} v\right)$, obtained with the change unknown $v=\mathrm{e}^{i t H} u$.
This result shows that if the initial condition is strongly localised in space, then the corresponding solution remains localised for large times.

## Statement of the probabilistic results

Set

$$
X_{h o l}^{0}(\mathbb{C}):=\left(\cap_{\sigma>0} \mathcal{H}^{-\sigma}(\mathbb{C})\right) \cap\left(\mathcal{O}(\mathbb{C}) \mathrm{e}^{-|z|^{2} / 2}\right)
$$

Define $\gamma \in L^{2}\left(\Omega ; X^{0}(\mathbb{C})\right)$ by

$$
\gamma(\omega, z)=\sum_{n=0}^{+\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)
$$

and for $\beta>0$ we define $\gamma_{\beta}=\gamma / \sqrt{\beta}$. Consider the Gaussian probability measure

$$
\mu_{\beta}=\left(\gamma_{\beta}\right)_{\#} \mathbf{p}:=\mathbf{p} \circ \gamma_{\beta}^{-1}
$$

We will check later that $\mu_{\beta}$ is a probability measure on $X_{h o l}^{0}(\mathbb{C})$. Let $2<p \leq+\infty$, then for almost all $\omega \in \Omega$,

$$
\gamma(\omega, .) \in F^{p}(\mathbb{C}) \quad \text { but } \quad \gamma(\omega, .) \notin F^{2}(\mathbb{C})
$$

As a consequence $\mu_{\beta}\left(L^{2}(\mathbb{C})\right)=0$.

Notice that since (LLL) conserves the $\mathcal{H}^{1}(\mathbb{C})$ norm, $\mu_{\beta}$ is formally invariant by its flow. More generally, we can define a family $\left(\rho_{\beta}\right)_{\beta>0}$ of probability measures on $X_{h o l}^{0}(\mathbb{C})$ which are formally invariant by (LLL) in the following way: define for $\beta>0$ the measure $\rho_{\beta}$ by

$$
\begin{equation*}
d \rho_{\beta}(u)=C_{\beta} \mathrm{e}^{-\beta \mathcal{E}(u)} d \mu_{\beta}(u) \tag{15}
\end{equation*}
$$

where $C_{\beta}>0$ is a normalising constant. By the Kakutani theorem and its corollary, the measures $\rho_{\beta}$ are mutually singular. Actually, the $\left(\rho_{\beta}\right)_{\beta>0}$ are the Gibbs measures of the equation (5).

We are now able to state the following global existence result, which also gives some qualitative information on the long time dynamics.

## Theorem (Germain-Hani-LT)

Let $\beta>0$. There exists a set $\Sigma \subset X_{h o l}^{0}(\mathbb{C})$ of full $\rho_{\beta}$ measure so that for every $u_{0} \in \Sigma$ the equation (LLL) with initial condition $u(0)=u_{0}$ has a unique global solution $u(t)=\Phi\left(t, u_{0}\right)$ such that for any $0<s<1 / 2$

$$
u(t)-u_{0} \in \mathcal{C}\left(\mathbb{R} ; \mathcal{H}^{s}(\mathbb{C})\right)
$$

Moreover, for all $\sigma>0$ and $t \in \mathbb{R}$

$$
\begin{equation*}
\|u(t)\|_{L^{3}(\mathbb{C})}+\|u(t)\|_{\mathcal{H}^{-\sigma}(\mathbb{C})} \leq C\left(\Lambda\left(u_{0}, \sigma\right)+\ln ^{1 / 2}(1+|t|)\right) \tag{16}
\end{equation*}
$$

where the constant $\Lambda\left(u_{0}, \sigma\right)$ satisfies the bound
$\mu_{\beta}\left(u_{0}: \Lambda\left(u_{0}, \sigma\right)>\lambda\right) \leq C e^{-c \lambda^{2}}$. Furthermore, the measure $\rho_{\beta}$ is invariant by $\Phi$ : for any $\rho_{\beta}$ measurable set $A \subset \Sigma$ and for any $t \in \mathbb{R}$,

$$
\rho_{\beta}(A)=\rho_{\beta}(\Phi(t, A))
$$

Finally, for all $t \in \mathbb{R}$

$$
\|u(t)\|_{L^{\mathbf{4}}(\mathbb{C})}=\left\|u_{0}\right\|_{L^{4}(\mathbb{C})}
$$

The same result (with the ad hoc measures $\mu$ and $\rho$ ) holds for the perturbed equations (11) and (14).

## Remark

By the Birkhoff-Kintchine Theorem we have for all $k \geq 1$

$$
\begin{equation*}
\frac{1}{T} \int_{0}^{T}\|u(t)\|_{\mathcal{H}^{-\sigma}(\mathbb{C})}^{k} d t \longrightarrow G_{k}\left(u_{0}\right), \quad \text { when } \quad T \longrightarrow+\infty \tag{17}
\end{equation*}
$$

and the fonction $G_{k}$ is a conservation law : for all $t \in \mathbb{R}, G_{k}(u(t))=G_{k}\left(u_{0}\right)$. Moreover

$$
\int_{\mathcal{H}^{-\sigma}} G_{k}(u) d \mu(u)=\int_{\mathcal{H}^{-\sigma}}\|u\|_{\mathcal{H}^{-\sigma}(\mathbb{C})}^{k} d \mu(u)
$$

One even has

$$
\frac{1}{T} \int_{0}^{T} e^{\frac{1}{2}\|u(t)\|_{\mathcal{H}}^{2}-\sigma(\mathbb{C})} d t \longrightarrow G_{\infty}\left(u_{0}\right), \quad \text { when } \quad T \longrightarrow+\infty
$$

By the theorem, there may be initial conditions such that $\|u(t)\|_{\mathcal{H}^{-\sigma}(\mathbb{C})}$ may grow like $\ln ^{1 / 2}(t)$, but not many since in mean it stays bounded, by (17). Compare with the bound obtained from the deterministic result.

## Remark

Formally, the (LLL) equation looks like the Szegö equation introduced and studied by Gérard and Grellier, but their properties are different. For instance, unlike (16) there is no nonlinear smoothing for the Szegö equation, as was shown by Oh, therefore it is not clear if an analogous result holds for the Szegö equation.

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## Remark

Let us compare the types of results given by the KAM method, the Birkhoff normal form method and the probabilistic methods :

- Smooth vs rough solutions
- Randomness
- Resonances?

Let us conclude this section with a few reference concerning the use of Gibbs measure in the construction of global strong solutions to PDEs. In a compact setting : Lebowitz-Rose-Speer, Bourgain, Zhidkov, Tzvetkov, Burq-Tzvetkov, Oh, Burq-Thomann-Tzvetkov, Deng, Nahmod et al, Suzzoni, Deng-Tzvetkov-Visciglia, Bourgain-Bulut, Richards and others.

There are also other types of a.s. global wellposedness results, without the use of invariant measures, mainly for the wave equation, but we do not comment on them.

## Sketch of the proof of the global wellposedness result

In the sequel we fix $\beta=1$ (say) and write $\mu=\mu_{\beta}, \rho=\rho_{\beta}$. We only prove the result for $s=0$.

## Lemma

The measure $\mu$ is a probability measure on $X_{h o l}^{0}(\mathbb{C})$.
Proof: It is enough to show that $\gamma \in X_{h o l}^{0}(\mathbb{C})$, p-a.s. First, for all $\sigma>0$ we have $\int_{\Omega}\|\gamma\|_{\mathcal{H}^{-\sigma}(\mathbb{C})}^{2} d \mathbf{p}(\omega)=\int_{\Omega} \sum_{n=0}^{+\infty} \frac{\left|g_{n}\right|^{2}}{(2(n+1))^{\sigma+1}} d \mathbf{p}(\omega)=C \sum_{n=0}^{+\infty} \frac{1}{(n+1)^{\sigma+1}}<+\infty$, therefore $\gamma \in \bigcap_{\sigma>0} L^{2}\left(\Omega ; \mathcal{H}^{-\sigma}(\mathbb{C})\right)$. Next, for all $A \geq 1$ there exists a set $\Omega_{A} \subset \Omega$ such that $\mathbf{p}\left(\Omega_{A}^{c}\right) \leq \exp \left(-A^{\delta}\right)$ and for all $\omega \in \Omega_{A}, \varepsilon>0, n \geq 0$

$$
\left|g_{n}(\omega)\right| \leq C A(n+1)^{\varepsilon}
$$

Then for $\omega \in \bigcup_{A \geq 1} \Omega_{A}, \sum_{n=0}^{+\infty} \frac{z^{n} g_{n}(\omega)}{\sqrt{(n+1)!}} \in \mathcal{O}(\mathbb{C})$.

We first define a smooth version of the usual spectral projector. Let $\chi \in \mathcal{C}_{0}^{\infty}(-1,1)$, so that $0 \leq \chi \leq 1$, with $\chi=1$ on $\left[-\frac{1}{2}, \frac{1}{2}\right]$. We define the operators $S_{N}=\chi\left(\frac{H}{N+1}\right)$ as

$$
S_{N}\left(\sum_{n=0}^{\infty} c_{n} \varphi_{n}\right)=\sum_{n=0}^{\infty} \chi\left(\frac{n+1}{N+1}\right) c_{n} \varphi_{n}
$$

Then for all $1<p<+\infty$, the operator $S_{N}$ is bounded in $L^{p}(\mathbb{C})$. This result does not hold true if one replaces $S_{N}$ with a crude frequency truncation.

## Local existence

Recall the definition of $\mathcal{T}$ in (3). It will be useful to work with an approximation of (LLL). We consider the dynamical system given by the Hamiltonian $\mathcal{E}_{N}(u):=\mathcal{E}\left(S_{N} u\right)$. This system reads

$$
\left\{\begin{array}{l}
i \partial_{t} u_{N}=\mathcal{T}_{N}\left(u_{N}\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C}  \tag{18}\\
u_{N}(0, z)=u_{0}(z)
\end{array}\right.
$$

and $\mathcal{T}_{N}\left(u_{N}\right):=S_{N} \mathcal{T}\left(S_{N} u, S_{N} u, S_{N} u\right)$.
Denote by $E_{k}$ the space on $\mathbb{C}$ spanned by $\varphi_{k}$. Observe that (18) is a finite dimensional dynamical system on $\bigoplus_{k=0}^{N} E_{k}$ and that the projection of $u_{N}(t)$ on its complement is constant. For $N \geq 0$ we define the measures $\rho_{N}$ by

$$
d \rho_{N}(u)=C_{N} \mathrm{e}^{-\mathcal{E}_{N}(u)} d \mu(u)
$$

where $C_{N}>0$ is a normalising constant. We have the following result

$$
\left\{\begin{array}{l}
i \partial_{t} u_{N}=\mathcal{T}_{N}\left(u_{N}\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C}  \tag{15}\\
u_{N}(0, z)=u_{0}(z)
\end{array}\right.
$$

## Lemma

The system (18) is globally well-posed in $L^{2}(\mathbb{C})$. Moreover, the measures $\rho_{N}$ are invariant by its flow denoted by $\Phi_{N}$.

Proof: The global existence follows from the conservation of $\left\|u_{N}\right\|_{L^{2}(\mathbb{C})}$. The invariance of the measures is a consequence of the Liouville theorem and the conservation of $\sum_{k=0}^{\infty} \lambda_{k}\left|c_{k}\right|^{2}$ by the flow of (LLL).

We now state a result concerning dispersive bounds of Hermite functions

## Lemma

For all $2 \leq p \leq+\infty$,

$$
\begin{equation*}
\left\|\varphi_{n}\right\|_{L^{p}(\mathbb{C})} \leq C n^{\frac{1}{2 p}-\frac{1}{4}} \tag{19}
\end{equation*}
$$

Proof: By Stirling, we easily get that $\left\|\varphi_{n}\right\|_{L^{\infty}(\mathbb{C})} \leq C n^{-\frac{1}{4}}$, which is (19) for $p=+\infty$; the estimate for $2 \leq p \leq \infty$ follows by interpolation. $\square$

## Lemma

(i) For all $2<p<+\infty$

$$
\begin{align*}
\exists C>0, & \exists c>0, \forall \lambda \geq 1, \forall N \geq 1, \\
& \mu\left(u \in X_{h o l}^{0}(\mathbb{C}):\left\|S_{N} u\right\|_{L^{P}(\mathbb{C})}>\lambda\right) \leq C e^{-c \lambda^{2}}, \\
& \mu\left(u \in X_{h o l}^{0}(\mathbb{C}):\|u\|_{L^{P}(\mathbb{C})}>\lambda\right) \leq C e^{-c \lambda^{2}} . \tag{20}
\end{align*}
$$

(ii) For all $2<p<+\infty$, there exists $\delta>0$ such that

$$
\begin{align*}
& \exists C>0, \exists c>0, \forall \lambda \geq 1, \forall N \geq N_{0} \geq 1, \\
& \quad \mu\left(u \in X_{h o l}^{0}(\mathbb{C}):\left\|\left(S_{N}-S_{N_{0}}\right) u\right\|_{L^{p}(\mathbb{C})}>\lambda\right) \leq C e^{-c N_{0}^{\delta} \lambda^{2}} . \tag{21}
\end{align*}
$$

Proof: We have that

$$
\mu\left(u \in X_{h o l}^{0}(\mathbb{C}):\|u\|_{L^{p}(\mathbb{C})}>\lambda\right)=\mathbf{p}\left(\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L^{p}(\mathbb{C})}>\lambda\right) .
$$

Let $q \geq p \geq 2$. Recall here the Khintchine inequality: there exists $C>0$ such that for all real $k \geq 2$ and $\left(a_{n}\right) \in \ell^{2}(\mathbb{N})$

$$
\begin{equation*}
\left\|\sum_{n \geq 0} g_{n}(\omega) a_{n}\right\|_{L_{\mathrm{p}}^{k}} \leq C \sqrt{k}\left(\sum_{n \geq 0}\left|a_{n}\right|^{2}\right)^{\frac{1}{2}} \tag{22}
\end{equation*}
$$

if the $g_{n}$ are iid normalized Gaussians.

Applying it to (22) we get

$$
\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L_{\omega}^{q}} \leq C \sqrt{q}\left(\sum_{n=0}^{\infty} \frac{\left|\varphi_{n}(z)\right|^{2}}{2(n+1)}\right)^{1 / 2}
$$

and using twice the Minkowski inequality for $q \geq p$ gives

$$
\begin{align*}
\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L_{\omega}^{q} L_{z}^{p}} & \leq\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L_{z}^{p} L_{\omega}^{q}} \\
& \leq C \sqrt{q}\left(\sum_{n=0}^{\infty} \frac{\left\|\varphi_{n}(z)\right\|_{L^{p}(\mathbb{C})}^{2}}{\langle n\rangle}\right)^{1 / 2} . \tag{23}
\end{align*}
$$

We are now ready to prove (20). Since we have $\left\|\varphi_{n}\right\|_{L^{p}(\mathbb{C})} \leq C n^{\frac{1}{2 p}-\frac{1}{4}}$, we get from (23)

$$
\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L_{\omega}^{q} L_{z}^{p}} \leq C \sqrt{q} .
$$

The Bienaymé-Tchebichev inequality gives then

$$
\begin{array}{r}
\mathbf{p}\left(\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L^{p}(\mathbb{C})}>\lambda\right) \leq\left(\lambda^{-1}\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{\left.L_{\omega}^{q} L_{z}^{p}\right)^{q}}\right. \\
\leq\left(C \lambda^{-1} \sqrt{q}\right)^{q} .
\end{array}
$$

Thus by choosing $q=\delta \lambda^{2} \geq 4$, for $\delta$ small enough, we get the bound

$$
\mathbf{p}\left(\left\|\sum_{n=0}^{\infty} \frac{g_{n}(\omega)}{\sqrt{2(n+1)}} \varphi_{n}(z)\right\|_{L^{p}(\mathbb{C})}>\lambda\right) \leq C e^{-c \lambda^{2}}
$$

which was the claim.

## Remark

From the previous result we deduce that on the support of $\mu$ (resp. $\rho$ ) we have $u \in L^{4}(\mathbb{C})$, thus we get a global existence result. However the invariance of the measures is not directly implied.

## Lemma

Let $p \in[1, \infty[$, then when $N \longrightarrow+\infty$.

$$
C_{N} e^{-\mathcal{E}_{N}(u)} \longrightarrow C e^{-\mathcal{E}(u)} \quad \text { in } \quad L^{p}(d \mu(u)) .
$$

In particular, for all measurable sets $A \subset X_{\text {hol }}^{0}(\mathbb{C})$,

$$
\rho_{N}(A) \longrightarrow \rho(A) .
$$

Proof: Denote by $G_{N}(u)=\mathrm{e}^{-\mathcal{E}_{N}(u)}$ and $G(u)=\mathrm{e}^{-\mathcal{E}(u)}$. By (21), we deduce that $\mathcal{E}_{N}(u) \longrightarrow \mathcal{E}(u)$ in measure, w.r.t. $\mu$. In other words, for $\varepsilon>0$ and $N \geq 1$ we denote by

$$
A_{N, \varepsilon}=\left\{u \in X_{h o l}^{0}(\mathbb{C}):\left|G_{N}(u)-G(u)\right| \leq \varepsilon\right\}
$$

then $\mu\left(A_{N, \varepsilon}^{c}\right) \longrightarrow 0$, when $N \longrightarrow+\infty$. Since $0 \leq G, G_{N} \leq 1$,

$$
\begin{aligned}
\left\|G-G_{N}\right\|_{L_{\mu}^{p}} & \leq\left\|\left(G-G_{N}\right) \mathbf{1}_{A_{N, \varepsilon}}\right\|_{L_{\mu}^{p}}+\left\|\left(G-G_{N}\right) \mathbf{1}_{A_{N, \varepsilon}^{c}}\right\|_{L_{\mu}^{p}} \\
& \leq \varepsilon\left(\mu\left(A_{N, \varepsilon}\right)\right)^{1 / p}+2\left(\mu\left(A_{N, \varepsilon}^{c}\right)\right)^{1 / p} \leq C \varepsilon,
\end{aligned}
$$

for $N$ large enough. Finally, we have when $N \longrightarrow+\infty$

$$
C_{N}=\left(\int \mathrm{e}^{-\mathcal{E}_{N}(u)} d \mu(u)\right)^{-1} \longrightarrow\left(\int \mathrm{e}^{-\mathcal{E}(u)} d \mu(u)\right)^{-1}=C
$$

and this ends the proof.

We look for a solution to (LLL) of the form $u=u_{0}+v$, thus $v$ has to satisfy

$$
\left\{\begin{array}{l}
i \partial_{t} v=\mathcal{T}\left(u_{0}+v\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C}  \tag{24}\\
v(0, z)=0
\end{array}\right.
$$

with $\mathcal{T}(u)=\mathcal{T}(u, u, u)$. Similarly, we introduce

$$
\left\{\begin{array}{l}
i \partial_{t} v_{N}=\mathcal{T}_{N}\left(u_{0}+v_{N}\right), \quad(t, z) \in \mathbb{R} \times \mathbb{C}  \tag{25}\\
v(0, z)=0
\end{array}\right.
$$

Recall that equation (25) is globally well posed in $L^{2}(\mathbb{C})$, and its flowmap is denoted by $\Phi_{N}$.
Let $\sigma>0$ and let us define

$$
A(R)=\left\{u_{0} \in X_{h o l}^{0}(\mathbb{C}):\left\|u_{0}\right\|_{\mathcal{H}^{-\sigma}(\mathbb{C})}+\left\|u_{0}\right\|_{L^{3}(\mathbb{C})} \leq R^{1 / 2}\right\}
$$

Then we have the following result

## Lemma

There exist $c, C>0$ so that for all $N \geq 0$

$$
\begin{equation*}
\rho_{N}\left(A(R)^{c}\right) \leq C e^{-c R}, \quad \rho\left(A(R)^{c}\right) \leq C e^{-c R}, \quad \mu\left(A(R)^{c}\right) \leq C e^{-c R} . \tag{26}
\end{equation*}
$$

Proof: Observe that we have $\rho_{N}\left(A(R)^{c}\right), \rho\left(A(R)^{c}\right) \leq C \mu\left(A(R)^{c}\right)$. The result is therefore given by (20).

## Proposition

There exists $c>0$ such that, for any $R>1 c_{0}>0$, setting $\tau(R)=c R^{-2}$, for any $u_{0} \in A(R)$ there exists a unique solution $v \in L^{\infty}\left([-\tau, \tau] ; L^{2}(\mathbb{C})\right)$ to the equation (24) and a unique solution $v_{N} \in L^{\infty}\left([-\tau, \tau] ; L^{2}(\mathbb{C})\right)$ to the equation (25) which furthermore satisfy

$$
\|v\|_{L^{\infty}\left([-\tau, \tau] ; L^{2}(\mathbb{C})\right)} \leq c_{0} R^{-1 / 2}, \quad\left\|v_{N}\right\|_{L^{\infty}\left([-\tau, \tau] ; L^{2}(\mathbb{C})\right)} \leq c_{0} R^{-1 / 2}
$$

As a consequence, for all $|t| \leq c R^{-2}$, if $c_{0} \ll 1$

$$
\begin{equation*}
\Phi\left(t, u_{0}\right) \in A(R+1), \quad \Phi_{N}\left(t, u_{0}\right) \in A(R+1) \tag{27}
\end{equation*}
$$

Proof: We only consider the equation (24), the other case being similar by the boundedness of $S_{N}$ on $L^{p}(\mathbb{C})$. We define the space

$$
Z(\tau)=\left\{v \in \mathcal{C}\left([-\tau, \tau] ; L^{2}(\mathbb{C})\right) \text { s.t. } v(0)=0 \text { and }\|v\|_{Z(\tau)} \leq c_{0} R^{-1 / 2}\right\}
$$

with $\|v\|_{Z(\tau)}=\|v\|_{L_{[-\tau, \tau]}^{\infty}}^{L^{2}(\mathbb{C})}$, and for $u_{0} \in A(R)$ we define the operator

$$
K(v)=-i \int_{0}^{t} \mathcal{T}\left(u_{0}+v\right) d s
$$

We will show that $K$ has a unique fixed point $v \in Z(\tau)$.
We have

$$
\begin{aligned}
&\|K(v)\|_{Z(\tau)} \leq \tau\left\|\mathcal{T}\left(u_{0}+v\right)\right\|_{Z(\tau)} \\
& \leq C \tau\left(\left\|\mathcal{T}\left(u_{0}, u_{0}, u_{0}\right)\right\|_{z}+\left\|\mathcal{T}\left(u_{0}, u_{0}, v\right)\right\|_{z}+\left\|\mathcal{T}\left(u_{0}, v, v\right)\right\|_{z}\right. \\
&\left.+\|\mathcal{T}(v, v, v)\|_{z}\right)
\end{aligned}
$$

We estimate each term. The conjugation plays no role, so we forget it. We only detail the first and the last term.

- Estimate of the trilinear term in $v$ : by the hypercontractivity estimates

$$
\|\mathcal{T}(v, v, v)\|_{L^{2}(\mathbb{C})} \leq C\|v\|_{L^{6}(\mathbb{C})}^{3} \leq C\|v\|_{L^{2}(\mathbb{C})}^{3}
$$

- Estimate of the constant term in $v:$ for $u_{0}$ in $A(R)$

$$
\left\|\mathcal{T}\left(u_{0}, u_{0}, u_{0}\right)\right\|_{L^{2}(\mathbb{C})} \leq C\left\|u_{0}\right\|_{L^{6}(\mathbb{C})}^{3} \leq C\left\|u_{0}\right\|_{L^{3}(\mathbb{C})}^{3} \leq C R^{3 / 2}
$$

(recall here that the bound $\left\|u_{0}\right\|_{L^{2}(\mathbb{C})}$ is forbidden since $\left\|u_{0}\right\|_{L^{2}(\mathbb{C})}=+\infty$ on the support of $\mu$.)

With these estimates at hand, the result follows by the Picard fixed point theorem.

## Approximation and invariance of the measure

## Lemma

Fix $R \geq 0$. Then for all $\varepsilon>0$, there exists $N_{0} \geq 0$ such that for all $u_{0} \in A(R)$ and $N \geq N_{0}$

$$
\left\|\Phi\left(t, u_{0}\right)-\Phi_{N}\left(t, u_{0}\right)\right\|_{L^{\infty}\left(\left[-\tau_{\mathbf{1}}, \tau_{1}\right] ; L^{\mathbf{2}}(\mathbb{C})\right)} \leq \varepsilon
$$

where $\tau_{1}=c R^{-2}$ for some $c>0$.

Proof: We have

$$
v-v_{N}=-i \int_{0}^{t}\left[S_{N}\left(\mathcal{T}\left(u_{0}+v\right)-\mathcal{T}\left(u_{0}+v_{N}\right)\right)+\left(1-S_{N}\right) \mathcal{T}\left(u_{0}+v\right)\right] d s
$$

Then we get

$$
\left\|v-v_{N}\right\|_{Z(\tau)} \leq C \tau R^{2}\left\|v-v_{N}\right\|_{Z(\tau)}+\int_{-\tau}^{\tau}\left\|\left(1-S_{N}\right) \mathcal{T}\left(u_{0}+v\right)\right\|_{L^{2}(\mathbb{C})} d s
$$

which in turn implies when $C \tau R^{2} \leq 1 / 2$

$$
\left\|v-v_{N}\right\|_{Z(\tau)} \leq 2 \int_{-\tau}^{\tau}\left\|\left(1-S_{N}\right) \mathcal{T}\left(u_{0}+v\right)\right\|_{L^{2}(\mathbb{C})} d s
$$

Here we need a bit a compactness to conclude. We refer to [Germain-Hani-Thomann] for the details.

Let $D_{i, j}=\left(i+j^{1 / 2}\right)^{1 / 2}$, with $i, j \in \mathbb{N}$ and set $T_{i, j}=\sum_{\ell=1}^{j} \tau_{1}\left(D_{i, \ell}\right)$. Let

$$
\Sigma_{N, i}:=\left\{u_{0}: \forall j \in \mathbb{N}, \quad \Phi_{N}\left( \pm T_{i, j}, u_{0}\right) \in A\left(D_{i, j+1}\right)\right\},
$$

and

$$
\Sigma_{i}:=\limsup _{N \rightarrow+\infty} \Sigma_{N, i}, \quad \Sigma:=\bigcup_{i \in \mathbb{N}} \Sigma_{i} .
$$

## Proposition

The following holds true:
(i) The set $\Sigma$ is of full $\rho$ measure.
(ii) For all $u_{0} \in \Sigma$, there exists a unique global solution $u=u_{0}+v$ to (LLL). This defines a global flow $\Phi$ on $\Sigma$.
(iii) For all measurable set $A \subset \Sigma$, and all $t \in \mathbb{R}$,

$$
\rho(A)=\rho(\Phi(t, A))
$$

The proof of (ii) relies on the invariance of the measure $\rho_{N}$ under the flow $\Phi_{N}$.
A repeated use of the approximation result will be crucial to prove (iii).

Let us show how one uses the Gibbs measure to define a global flow and to get the quantitative bound in $\ln ^{1 / 2}(t)$ in the main theorem.
Let $c>0$ be given by (26). For $T \leq \mathrm{e}^{c R / 2}$ we define the set of the good data

$$
\begin{equation*}
\Sigma_{R}=\bigcap_{k=-[T / \tau]}^{[T / \tau]} \Phi_{N}\left(-k \tau, B_{R}\right) \tag{28}
\end{equation*}
$$

Now we crucially use the invariance of the measure and get

$$
\begin{aligned}
\rho_{N}\left(X_{h o l}^{0}(\mathbb{R}) \backslash \Sigma_{R}\right) & \leq(2[T / \tau]+1) \rho_{N}\left(X^{0}(\mathbb{R}) \backslash B_{R}\right) \\
& \leq C R^{2} \mathrm{e}^{c R / 2} \mathrm{e}^{-c R} \leq C \mathrm{e}^{-c R / 4}
\end{aligned}
$$

which shows that $\Sigma_{R}$ is a big subset of $X_{h o l}^{0}(\mathbb{R})$ when $R \longrightarrow+\infty$.

Let us show how one uses the Gibbs measure to define a global flow and to get the quantitative bound in $\ln ^{1 / 2}(t)$ in the main theorem.
Let $c>0$ be given by (26). For $T \leq \mathrm{e}^{c R / 2}$ we define the set of the good data

$$
\begin{equation*}
\Sigma_{R}=\bigcap_{k=-[T / \tau]}^{[T / \tau]} \Phi_{N}\left(-k \tau, B_{R}\right) . \tag{29}
\end{equation*}
$$

Now we crucially use the invariance of the measure and get

$$
\begin{aligned}
\rho_{N}\left(X_{h o l}^{0}(\mathbb{R}) \backslash \Sigma_{R}\right) & \leq(2[T / \tau]+1) \rho_{N}\left(X^{0}(\mathbb{R}) \backslash B_{R}\right) \\
& \leq C R^{2} \mathrm{e}^{c R / 2} \mathrm{e}^{-c R} \leq C \mathrm{e}^{-c R / 4}
\end{aligned}
$$

which shows that $\Sigma_{R}$ is a big subset of $X_{h o l}^{0}(\mathbb{R})$ when $R \longrightarrow+\infty$. Now, by the definition (29) of $\Sigma_{R}$ and (27), we deduce that for all $|t| \leq T$ and $u_{0} \in \Sigma_{R}$

$$
\left\|\Phi_{N}\left(t, u_{0}\right)\right\|_{L^{\mathbf{3}}(\mathbb{C})}+\left\|\Phi_{N}\left(t, u_{0}\right)\right\|_{\mathcal{H}^{-\sigma}(\mathbb{C})} \leq(R+1)^{1 / 2}
$$

In particular, for $|t|=T \sim e^{c R / 2}$

$$
\left\|\Phi_{N}\left(t, u_{0}\right)\right\|_{L^{3}(\mathbb{C})}+\left\|\Phi_{N}\left(t, u_{0}\right)\right\|_{\mathcal{H}^{-\sigma}(\mathbb{C})} \leq C(\ln |t|+1)^{1 / 2}
$$

and this bound is uniform in $N \geq 1$. The term $\ln ^{1 / 2}(t)$ is reminiscent from the large deviation estimates involving Gaussian random variables.

