

New invariant surface measures for the Cubic Schrödinger Equation

Jean-Baptiste Casteras, Ana Bela Cruzeiro, Annie Millet

▶ To cite this version:

Jean-Baptiste Casteras, Ana Bela Cruzeiro, Annie Millet. New invariant surface measures for the Cubic Schrödinger Equation. 2025. hal-04934726

HAL Id: hal-04934726 https://hal.science/hal-04934726v1

Preprint submitted on 7 Feb 2025

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers. L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

NEW INVARIANT SURFACE MEASURES FOR THE CUBIC SCHRÖDINGER EQUATION

JEAN-BAPTISTE CASTERAS, ANA BELA CRUZEIRO, AND ANNIE MILLET

ABSTRACT. We construct new invariant measures supported on mass level sets for the cubic defocusing nonlinear Schrödinger equation in dimensions 1 and 2.

1. INTRODUCTION

Consider the following equation for $\lambda \in \{-1, +1\}$

(1)
$$i\frac{\partial v}{\partial t} = -\Delta v + \lambda |v|^2 v,$$

in the torus \mathbb{T}^d of dimension d = 1, 2, with initial conditions $v_0 \in H^s$, for some $s \in \mathbb{R}$. When $\lambda = 1$ (resp. $\lambda = -1$), the equation is defocusing (resp. focusing). Recall that (1) can be seen as a Hamiltonian PDE, i.e.,

$$\frac{du}{dt} = i \frac{\partial H}{\partial \overline{u}},$$

where

$$H(u) = \frac{1}{2} \int_{\mathbb{T}^d} |\nabla u|^2 dx + \frac{\lambda}{4} \int_{\mathbb{T}^d} |u|^4 dx.$$

From this reformulation, it is easily seen that H(u) is an invariant quantity. One can also check that the mass

$$M(u) = \int_{\mathbb{T}^d} |u|^2 dx,$$

and the momentum

$$P(u) = i \int_{\mathbb{T}^d} u \overline{u}_x dx$$

are conserved. In general, these three conservation laws are the only known ones. But in some very special cases, others are available. For instance, if d = 1, the cubic Schrödinger equation possesses a Lax pair structure which allows to construct an infinity of conservation laws see [10, 29]. In this case, the equation is said to be completely integrable. Let us point out that very recently people were able to construct a continuum of conservation laws (whereas the earlier results only provided infinite but countable ones). We

J.-B.C. is supported by FCT - Fundação para a Ciência e a Tecnologia, under the project: UID/04561/2025. A. B. C. is supported by FCT - Fundação para a Ciência e a Tecnologia, under the project: UIDB/00208/2020. The work of A.M. is conducted within the FP2M federation (CNRS FR 2036).

refer to [28] for results regarding the Benjamin-Ono equation, [15] for the cubic NLS in the whole line, [7] for the KdV and modified KdV equations (the non-linear term $v\partial_x v$ is replaced by $v^2\partial_x v$).

Our goal in this paper is to construct new invariant measures for the cubic NLS equation in dimension 1 and 2. First, let us describe in a few words the situation in the *d*-dimensional case. Let $H(p,q) = H(p_1, \ldots, p_d, q_1, \ldots, q_d)$ be a Hamiltonian and let us consider its associated flow in \mathbb{R}^{2d} i.e.

(2)
$$\dot{p}_i = \frac{\partial H}{\partial q_j}, \ \dot{q}_i = -\frac{\partial H}{\partial p_j}, \quad i, j = 1, ..., d.$$

Then, using the conservation of H, Liouville's theorem states that the Gibbs measure $e^{-H(p,q)}\Pi_{j=1}^d dp_j dq_j$ is invariant under the flow generated by (2). Notice that, for any (reasonable) function F(p,q) which is conserved by the flow (2), the measure $F(p,q)e^{-H(p,q)}\Pi_{j=1}^d dp_j dq_j$ is also invariant. The situation is much more delicate in the infinite dimensional case.

In analogy with the Euclidean case \mathbb{R}^d , by a scaling argument, one expects (1) to be (at least locally) well-posed for initial data in H^s , for $s > s_c = \frac{d}{2} - 1$. The space H^{s_c} is called critical. In [12], it is proved that (1) for d = 1 in the focusing and defocusing case is globally well-posed for initial data in H^s , $s > -\frac{1}{2}$. The same result holds for the 1d modified KdV equation. When d = 2, Bourgain [2] obtained the local well-posedness of (1) for initial data in H^s , s > 0. The global well-posedness in this exact situation was obtained in [13] for small initial data in H^s , for any s > 0. In dimensions 3 and 4, for energy-critical non-linearity, local well posedness was proved in H^1 , while global well-posedness was shown for "small" initial data; see [14] and the references therein. For higher dimensions, we refer to the very recent reference [18] concerning local well-posedness results in critical spaces.

Starting from the work of Lebowitz, Rose and Speer [19], and later on from the paper of Bourgain [3], the statistical approach to Schrödinger equations has been considerably developed. It consists in constructing invariant probability (Gibbs) measures for the solutions of these equations, based on their invariant quantities. The solutions are proved to exist starting almost everywhere from the support of the measures and can be rather singular, typically living in Sobolev spaces of low regularity.

This approach uses techniques from probability theory, where transport properties of Gaussian measures under linear and nonlinear transformations have an old tradition. In the framework of Hamiltonian partial differential equations, of which Schrödinger systems are an example, invariant (or quasiinvariant) measures are associated to conservation laws and are weighted Gaussian measures. We refer to the works [23] and [25] where the authors prove the quasi-invariance of Gaussian measures under the flow generated by the quintic NLS in the defocusing case on \mathbb{T}^3 (which is energy critical). For the cubic non-linear Schrödinger equation (NLS), in the defocusing case and on the two-dimensional torus, a renormalisation must be introduced. This equation was studied in [4], with the construction of a Gibbs measure associated to the Wick ordered Hamiltonian and the existence of the corresponding flow defined on the support of the measure was proved. We also refer to [24] for related results.

In the recent paper [9] the authors generalize Bourgain's work and consider Wick ordered power nonlinearities in the Schrödinger equation which are of higher degree, proving the existence of strong solutions with respect to the associated Gibbs measure.

The Gibbs measures are absolutely continuous with respect to Gaussian ones. We consider the corresponding abstract Wiener spaces in the sense of Gross [11]. Then we use techniques from quasi-sure analysis on these spaces, as established by Malliavin [20]. For a non degenerate functional on an abstract Wiener space we can consider its level sets, which are (finite codimensional) submanifolds of that space. In [1] Airault and Malliavin defined a disintegration of the Gaussian measure into a family of conditional laws or surface measures and proved a geometric Federer-type coarea formula providing this disintegration. Such measures do not charge sets of capacity zero.

In this work we study the defocusing cubic renormalised NLS on the two dimensional torus (as in [4]) and construct surface measures for the mass, which is a conserved quantity of the equation. Then we prove existence of weak solutions living on the level sets of the mass and defined almost everywhere with respect to the surface measures. We therefore provide a more precise result than those obtained with respect to the Gibbs measures constructed so far.

We refer also to the work [21], where a different kind of conditioning and corresponding measures for NLS have been considered.

Let us describe more precisely the result of [21]. The authors construct a measure of the form

$$d\mu_0 = Z_0^{-1} e^{\frac{1}{4} \int_{\mathbb{T}} |u|^4} dP_0,$$

where Z_0 is a normalisation constant and P_0 is defined by

$$P_0(E) = \lim_{\varepsilon \to 0} P\Big(E\Big|\int_{\mathbb{T}} |u|^2 \in A_{\varepsilon}(a), \ i \int_{\mathbb{T}} u \overline{u_x} \in B_{\varepsilon}(b)\Big),$$

where $A_{\varepsilon}(a)$, $B_{\varepsilon}(b)$ are neighborhoods shrinking nicely to a and b as $\varepsilon \to 0$. In fact, this measure is shown to exist for polynomial nonlinearities $\pm |u|^{p-2}u$ with 2 (the mass needing to be small if <math>p = 6). We also refer to [5] for related results for the derivative NLS.

Let us also mention that an active line of work in the construction of Gibbs measure is to find optimal integrability thresholds for such a measure to exist; see for instance [22].

Our methods are similar to those in [8], where surface measures supported by level sets of the enstrophy were constructed for the Euler equation on the torus.

We are interested in the cubic defocusing Schrödinger equation in the one and two dimensional torus, namely

(3)
$$i\frac{\partial v}{\partial t} = -\Delta v + |v|^2 v$$

with initial conditions v_0 of low regularity. If the solution is assumed to be at least L^2 , we can use the following change of variables $v = e^{it(1-2||u||_{L^2}^2)}u$. In this case, (3) is equivalent to

(4)
$$i\frac{\partial u}{\partial t} = -\Delta u + \left(|u|^2 - 2\|u\|_{L^2}^2\right)u + u.$$

Notice that it is not anymore the case if the solution does not have a finite mass. Our aim is to take initial data in the support of a surface measure supported on mass level sets. Formally, let

$$d\mu_{2s} = Z_s^{-1} e^{-\frac{1}{2} \|u\|_{H^s}^2} du,$$

where Z_s is a normalisation constant and du is formally the Lebesgue measure, which is well-known to not exist on infinite dimensional vector spaces. The Gaussian measure μ_{2s} can also be seen as the law of the random variable

$$\omega \to \sum_{n \in \mathbb{Z}^d} \frac{g_n(\omega)}{(1+|n|^2)^{s/2}} e^{inx},$$

where $\{g_n\}$ are independent standard complex-valued Gaussian variables on a probability space $(\Omega, \mathcal{A}, \mathbb{P})$. For $\sigma \in \mathbb{R}$, it is easy to check that

$$\mathbb{E}\Big[\sum_{n \in \mathbb{Z}^d} (1+|n|^2)^{\sigma} \Big| \frac{g_n}{(1+|n|^2)^{s/2}} \Big|^2 \Big] < \infty \text{ iff } \sigma < s - \frac{d}{2}.$$

So this random series converges in $\mathbb{L}^2(\Omega, H^{\sigma}(\mathbb{T}^d))$ if and only if $\sigma < s - \frac{d}{2}$. In what follows, we will always take s = 1. We see that, when d = 2, solutions distributed according to μ_2 do not have a finite mass, which implies that equation (3) and (4) are not equivalent. The previous remarks also show that the support of this Gibbs measure is necessarily rough. We point out that an alternative method to construct invariant measures is the so-called Inviscid-Infinite-dimensional limits method due to Kuksin and Shirikyan [16, 17]. We refer to [6, 26, 27] for applications related to NLS.

We set $V_r = \{ \varphi \in L^2 | E(\varphi) = r \}$, for r > 0 where

(5)
$$E(\varphi) = \int_{\mathbb{T}^d} |\varphi|^2 dx - \mathbb{1}_{\{d=2\}} \sum_{k \neq 0} \frac{2}{|k|^2}$$

In section 4, we prove that this quantity is well-defined in dimension 2. Our main result reads as follows :

INVARIANT MEASURE NLS

Theorem 1.1. Let d = 1 or d = 2. Then, for all r > 0, there exists a probability measure σ_r defined on L^2 with support on V_r , a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and a $C(\mathbb{R}, V_r)$ -valued stochastic process u such that

- $u = u^{\omega}$ is a global weak solution to (1) if d = 1 or (4) if d = 2 for \mathbb{P} a.e. $\omega \in \Omega$.
- σ_r is invariant in the sense that the law of u(t) is σ_r for all $t \in \mathbb{R}$.

Let us also point out that our result can be easily generalised to more general dispersions namely in (1), we could have taken any power of the Laplace operator $(-\Delta)^s$, s > 0 as well as any odd power nonlinearities. In the case of a more general dispersion of the form $(-\Delta)^s$, the difference would have been the value of the parameter a = 2 in the above renormalisation to define $E(\varphi)$ (see (5)). Instead of taking it equal to 2, it would have been chosen as a = 2s.

In section 2, we describe the Gaussian measures we will use, define the Sobolev and Wiener spaces. Finally, we state a result by Airault and Malliavin [1] which constructs a surface measure supported by a level set of some functional defined on a Wiener space; this will be the main ingredient in the construction of our invariant measure. In section 3, we prove Sobolev estimates on the non linear term projected on low frequencies; in dimension 2 this will require a renormalisation of this cubic term. A similar renormalisation was also considered in [4]. In Section 4, we prove several estimates for a renormalisation of the mass. Finally, we prove Theorem 1.1 in Section 5. Some technical lemma is proved in the appendix.

2. Preliminaries

Let us first state several basic facts about Gaussian measures. Next, we rewrite the cubic Schrödinger equation into frequency.

2.1. The Gaussian measures. Given a function $\varphi \in L^2(\mathbb{T}^d)$, set $\varphi = \sum_{k \in \mathbb{Z}^d} \varphi_k e_k$, where $e_k(x) = \frac{1}{2} e^{ikx}$, $k \in \mathbb{Z}^d$, with $kx = k_1 x_1 + \ldots + k_d x_d$. For a > 0 let us consider the Gaussian measures on $L^2(\mathbb{T}^d)$

$$d\mu_a = \prod_{k \in \mathbb{Z}^d} d\mu_a^k, \quad d\mu_a^k(z) = \frac{|k|^a}{2\pi} e^{-\frac{1}{2}|k|^a|z|^2} dx dy,$$

with $z = x + iy \in \mathbb{C}, k = (k_1, \dots, k_d) \in \mathbb{Z}^d$ and $|k|^2 = k_1^2 + \dots + k_d^2$.

Suppose that $\{\varphi_k, k \in \mathbb{Z}^d\}$ are independent complex-valued Gaussian random variables with distribution μ_a^k . For later purpose, let us notice that

$$E_{\mu_a}(\varphi_k) = 0, \ E_{\mu_a}(\varphi_k \varphi_{k'}) = 0, \ E_{\mu_a}(\varphi_k \overline{\varphi_{k'}}) = \delta_{k,k'} \ \frac{2}{|k|^a},$$

and, for any $p \ge 1$,

(6)
$$\mathbb{E}_{\mu_a}(|\varphi_k|^{2p}) = \frac{c_p}{|k|^{ap}}, \quad \text{where} \quad c_p = 2^p p!.$$

Indeed, letting $x = Re(\varphi_k), y = Im(\varphi_k)$, we have

$$\mathbb{E}_{\mu_a}(|\varphi_k|^{2p}) = \frac{|k|^a}{2\pi} \int_{\mathbb{R}\times\mathbb{R}} (x^2 + y^2)^p e^{-|k|^a \frac{x^2 + y^2}{2}} dx dy$$
$$= |k|^a \int_{\mathbb{R}^+} r^{2p} e^{-|k|^a \frac{r^2}{2}} r dr$$
$$= 2^p p! |k|^{-ap}.$$

2.2. The Sobolev spaces. We next recall some definitions of functional spaces such as Sobolev spaces and Wiener spaces. The Sobolev spaces of order β on the torus are defined by

$$H^{\beta} = \left\{ \varphi \in L^2 : \sum_{k \neq 0} |k|^{2\beta} |\varphi_k|^2 < +\infty \right\}$$

and are endowed with the corresponding Hilbert scalar product.

We will work with the Gaussian measure μ_2 . This measure is supported in the space H^{β} with $\beta < 0$. This is a consequence of the fact that $\int \sum_{k\neq 0} |k|^{2\beta} |\varphi_k|^2 d\mu_2 = 2 \sum_{k\neq 0} \frac{|k|^{2\beta}}{|k|^2}$ converges for $\beta < 0$; it is also a consequence of μ_2 being the measure induced by the process $\sum_{|n|>0} \frac{1}{n} G_n e_n(x)$, where $\{G_n\}_n$ are independent standard complex Gaussian random variables.

For $\beta < 0$, the triple (H^{β}, H^1, μ_2) is an abstract Wiener space in the sense of Gross (see [11]).

Given a Hilbert space G, we consider the Malliavin derivative of a functional $F: H^{\beta} \to G$, in the Cameron-Martin direction $v \in H^1$, namely the (a.e.) limit

$$\lim_{\varepsilon \to 0} \frac{1}{\varepsilon} (F(\varphi + \varepsilon v) - F(\varphi)), \quad x \in H^1.$$

By the Riesz theorem this derivative gives rise to a linear gradient operator ∇ such that the limit above coincides with $\langle \nabla F(\varphi), v \rangle$. One can iterate this procedure and obtain higher order gradients ∇^r . Let us for instance consider the second derivative of F. For $\varphi \in H^{\beta}$, and $v, w \in H^1$, iterating the first derivative, we obtain

$$\nabla^2 F(\varphi)(v,w) = D_v D_w F(\varphi).$$

We define the Hilbert-Schmidt norm

$$\|\nabla^2 F(\varphi)\|_{H.S.(H^1 \otimes H^1, H^\beta)}^2 = \sum_{j,k} \|D_{\hat{e}_j} D_{\hat{e}_k} F(\varphi)\|_{H^\beta}^2,$$

where $\{\hat{e}_j\}_j$ is an orthonormal basis of H^1 . The corresponding Sobolev spaces are $W^p_r(H^\beta; G)$, that is the space of maps $F : H^\beta \to G$ such that $F \in L^p_{\mu_2}(H^\beta; G)$, the gradients $\nabla^s F : H^\beta \to H.S.(\otimes_s H^1, G)$ are defined for every $1 \leq s \leq r$ and belong to $L^p_{\mu_2}(H^\beta; H.S.(\otimes_s H^1, G))$. We equip this space with the norm

$$\|f\|_{W^p_r(H^\beta;G)} = \|f\|_{L^p_{\mu_2}(H^\beta;G)} + \sum_{1 \le s \le r} \|f\|_{L^p_{\mu_2}(H^\beta;H.S.(\otimes_s H^1,G))}.$$

We also denote $W_{\infty} = \bigcap_{r,p} W_r^p(H^{\beta}; \mathbb{R})$ and $W_{\infty}(H^{\beta}) = \bigcap_{r,p} W_r^p(H^{\beta}; H^{\beta}).$

2.3. The nonlinear Schrödinger equation. Using the decomposition of the solution $\varphi(x,t) = \sum_k \varphi_k(t) e_k(x)$ into frequencies, the cubic Schrödinger equation (3) for φ can be rewritten as the family of equations

$$\frac{d}{dt}\varphi_k = -i|k|^2\varphi_k - iB_k(\varphi), \quad |k| > 0,$$

where

$$B_k(\varphi) = \sum_{l,m} \varphi_{k-m} \varphi_{l+m} \overline{\varphi}_l$$

The functions $\{e_k\}_k$ are eigenfunctions for the operator $-\Delta$ with eigenvalues $|k|^2 = k_1^2 + \ldots + k_d^2$; we set $A_k(\varphi) = -|k|^2 \varphi_k$, corresponding to the Laplacian $A = \Delta$.

For any $n \in \mathbb{N}$, we define

$$B^{n}(\varphi) = \Pi_{n}(B(\Pi_{n}(\varphi))),$$

where

$$\Pi_n \left(\sum_k \varphi_k e_k \right) = \sum_{k: 1 \le |k|^2 \le n} \varphi_k e_k.$$

2.4. Surface measures. When $g \in W_{\infty}$, it is possible to consider a redefinition of g, namely a function g^* that coincides with g almost surely, and is defined and continuous outside sets of capacity zero (c.f [20]). Using such redifinitions a construction of surface measures on Wiener spaces was made in [1]. These measures are defined on level sets of W_{∞} functionals, and do not charge zero capacity sets.

Let *E* be a real-valued functional defined on H^{β} and consider the surfaces $V_r = \{\varphi | E(\varphi) = r\}$ for r > 0. Assume that *E* and $\|\nabla E\|_{H^1}^{-1}$ belong to the corresponding spaces W_{∞} .

Denote the C^{∞} densities of $d(E * \mu_2)$ (resp. $d(E * g\mu_2)$) with respect to the Lebesgue measure by $\rho(r) = \frac{d(E * \mu_2)}{dr}$ (resp. $\rho_g(r) = \frac{d(E * g\mu_2)}{dr}$); they are well defined and smooth, cf. [1]. We have

Theorem 2.1. [1] Let r > 0 be such that $\rho(r) > 0$. Then there exists a Borel probability measure defined on H^{β} , denoted by ν^{r} , with support on V_{r} and such that

$$\int_{V_r} g^*(\varphi) d\nu^r = \frac{\rho_g(r)}{\rho(r)},$$

for any g^* redifinition of g.

This will be a crucial ingredient to construct our invariant measures.

3. Norm estimates for B

Our goal in this section is to derive Sobolev estimates for the vector field B. We will see that we need to renormalize B when the dimension d = 2. We proceed as follows: for every $N \in \mathbb{N}$ and $k \in \mathbb{Z}^2$, we let

(7)
$$(:B:)_{k}^{N} = B_{k} - 4 \left(\sum_{m:1 \leq |m| \leq |N|} \frac{1}{|m|^{a}} \right) \varphi_{k},$$

with a = 2. Notice that this sum diverges when $|N| \to \infty$ if d = 2. In dimension 1, a renormalisation is not needed (since the sum is converging). More precisely our aim in this section is to show that in dimension 2 (: B :)^N belongs to $L^{2p}_{\mu_a}(H^{\beta}, H^{\beta})$ and $\nabla^s(: B :)^N$ belongs to $L^{2p}_{\mu_a}(H^{\beta}, H.S.(\otimes^s H^1, H^{\beta}))$, for any $p \ge 1$, $s \in \mathbb{N}^*$ and $\beta < 0$, uniformly in N. Thanks to our estimate, we will be able to prove that (: B :)^N converges in $L^{2p}_{\mu_a}(H^{\beta}, H^{\beta})$. We will denote its limit by : B :. We begin by estimating (: B :)^N.

Proposition 3.1. Let a = 2 and $\beta < 0$ and $p \ge 1$ an integer.

In dimension 1, B is L^p -integrable with respect to the measure μ_a and to the norm H^{β} .

In dimension 2, $(:B:)^N$ is L^p -integrable with respect to the measure μ_a and to the norm H^β , uniformly in N. We denote by :B: its limit as $N \to \infty$.

Notice that our choice of parameters is related to the fact that we want to work on the abstract Wiener space (H^{β}, H^1, μ_2) , with $\beta < 0$.

Proof. By abuse of notation, in this proof, we will drop the subcript N in $: B :^N$ in dimension 2. We will work at fixed N and will notice at the end that, thanks to the normalisation, our estimates do not depend on N.

Let p be odd; we are going to show that $: B :\in L^{2p}_{\mu_a}(H^\beta, H^\beta).$

In a first step, we compute $E_{\mu_a}(|B|^{2p})$.

For every $k \in \mathbb{Z}$ with $k \neq 0$ we have

$$|B_{k}(\varphi)|^{2} = \sum_{l,l'\neq 0,m,m',|l|\vee|m|\vee|l'|\vee|m'|\leqslant N} T_{k}(l,m,l',m'),$$

where, for $m \notin \{k, -l\}$ and $m' \notin \{k, -l'\}$,

$$T_k(l,m,l',m') := \varphi_{k-m}\varphi_{l+m}\varphi_{l'}\overline{\varphi_l\varphi_{k-m'}\varphi_{l'+m'}}.$$

Note that $E_{\mu a}[T_k(l, m, l', m')^p] = 0$ except in three cases: k - m = l, k - m = k - m' and k - m = l' + m'.

Case 1 Let k - m = l; then we have to impose that either l + m = k - m'or l + m = l' + m'.

Case 1.1 If l + m = k - m', we deduce that, since k - l = m = m + m', we have m' = 0. Therefore, $T_k(l, k - l, l', 0) = |\varphi_k|^2 |\varphi_l|^2 |\varphi_{l'}|^2$. Thus,

- if $m \neq 0$, $k \neq l$ and $k \neq l'$, the indices k, l, l' are pairwise different;
- if $m \neq 0$ and l = l', then $T_k(l, k l, l, 0) = |\varphi_k|^2 |\varphi_l|^4$;

INVARIANT MEASURE NLS

- if $m \neq 0$ and l' = k, then $T_k(l, k l, k, 0) = |\varphi_k|^4 |\varphi_l|^2$;
- if m = m' = 0 and $l' \neq k$, then $T_k(k, 0, l', 0) = |\varphi_k|^4 |\varphi_{l'}|^2$; if $m = m' \neq 0$ and $l = l' \neq k$, then $T_k(l, k l, l, k l) = |\varphi_k|^2 |\varphi_l|^4$;
- finally, if m = m' = 0 and l' = k, then $T_k(k, 0, k, 0) = |\varphi_k|^6$.

Case 1.2 If l + m = l' + m', we deduce that k = l + m = l' + m' and hence k-m'=l'. Therefore, $T_k(l,k-l,l',k-l')=|\varphi_k|^2|\varphi_l|^2|\varphi_{l'}|^2$. Thus,

- if $m \neq 0$, $m' \neq 0$ and $l \neq l'$, the indices k, l, l' are pairwise different;
- if m = 0 and $m' \neq 0$, then $T_k(k, 0, l', k l') = |\varphi_k|^4 |\varphi_l'|^2$.

In the other cases, that is $m \neq 0$ and m' = 0 (resp. m = m' = 0), we have again $T_k(l, k - l, k, 0)$ (resp. $T_k(k, 0, k, 0)$).

Case 2 Let k-m = k-m', which implies m = m'. Then we have to impose either l + m = l or l + m = l' + m'.

Case 2.1 If l + m = l, we deduce m = m' = 0, so that l' = l' + m'. Therefore, $T_k(l, 0, l', 0) = |\varphi_k|^2 |\varphi_l|^2 |\varphi_{l'}|^2$. Thus

- the indices k, l, l' can be pairwise distinct:
- if $k = l' \neq l$, then $T_k(l, 0, k, 0) = |\varphi_k|^4 |\varphi_l|^2$;

In the other cases $k = l \neq l'$ (resp. k = l = l'), we have again $T_k(k, 0, l', 0)$ (resp. $T_k(k, 0, l, 0)$).

Case 2.2 If l + m = l' + m', since m = m' we conclude that l = l'. Therefore, $T_k(l,m,l,m) = |\varphi_{k-m}|^2 |\varphi_{l+m}|^2 |\varphi_l|^2$. Thus,

- if $m \neq 0$, $l \neq k m$ and $l \neq l 2m$, the indices k m, l, l + m are pairwise distinct;
- if $m \neq 0$ and l + m = k m, then $l \neq k m$ and $T_k(k 2m, m, k m)$ $(2m,m) = |\varphi_{k-m}|^4 |\varphi_{k-2m}|^2;$
- if m = 0 and $l \neq k$, then $T_k(l, 0, l, 0) = |\varphi_k|^2 |\varphi_l|^4$;

Finally, when m = 0 and k = l (resp. $m \neq 0$ and l = k - m) we have again $T_k(k, 0, k, 0)$ (resp. $T_k(l, k - l, l, k - l)$).

Case 3 Let k - m = l' + m'; then we have to impose either l + m = l or l+m=k-m'.

Case 3.1 If l + m = l, that is m = 0, we have k = l' + m'. Therefore, $T_k(l, 0, l', k - l') = |\varphi_k|^2 |\varphi_l|^2 |\varphi_{l'}|^2$. Thus,

• the indices k, l, l' can be pairwise distinct;

• if $k \neq l = l'$, then $T_k(l, 0, l, k - l) = |\varphi_k|^2 |\varphi_l|^4$.

When $k = l \neq l'$ (resp. $k = l' \neq l$), we have again $T_k(k, 0, l', k - l')$ (resp. $T_k(l, 0, k, 0)$).

Finally, if k = l = l', we have again $T_k(k, 0, k, 0)$.

Case 3.2 If l + m = k - m', we deduce that m + m' = k - l = k - l', so that l' = l. Therefore, $T_k(l, m, l, k - m - l) = |\varphi_{k-m}|^2 |\varphi_{l+m}|^2 |\varphi_l|^2$. Thus

- if $m \neq 0$, $l \neq k m$ and $l \neq k 2m$, the indices k m, l, l + m are pairwise distinct;
- if $m \neq 0$ and l = k m, then $T_k(l, k l, l, 0) = |\varphi_l|^4 |\varphi_k|^2$.

The other cases $m \neq 0$ and l+m = k-m (resp. m = 0 and $l \neq k$) give again $T_k(k-2m,m,k-2m,m)$ (resp. $T_k(l,0,l,k-l)$). Finally, the case m = 0 and l = k gives again $T_k(k,0,k,0)$.

As a summary, we have proven that

$$E_{\mu_{a}}(|B_{k}|^{2p}) = E[|\varphi_{k}|^{6p}] + 4 \sum_{|l| \leq N, |l'| \leq N, l \neq k, l' \neq k, l \neq l'} E[|\varphi_{k}|^{2p}]E[|\varphi_{l}|^{2p}]E[|\varphi_{l'}|^{2p}]$$

$$+ 2 \sum_{|l| \neq N, l \neq k, 1 \leq |m| \leq N, l+m \neq k-m} E[|\varphi_{k-m}|^{2p}]E[|\varphi_{l+m}|^{2p}]E[|\varphi_{l}|^{2p}]$$

$$+ 4 \sum_{|l| < N, l \neq k} E[|\varphi_{k}|^{4p}]E[|\varphi_{l}|^{2p}] + 4 \sum_{|l| \leq N, l \neq k} E[|\varphi_{k}|^{2p}]E[|\varphi_{l}|^{4p}]$$
(8)
$$+ \sum_{1 \leq |m| \leq N, m \neq k, 2m \neq k} E[|\varphi_{k-m}|^{2p}]E[|\varphi_{k-2m}|^{4p}].$$
We first deal with $p = 1$. Set

$$D_N := \sum_{l:0 < |l| \le N} \int |\varphi_l|^2 d\mu_a.$$

The above discussion implies that

(9)
$$E_{\mu_a}[|B_k|^2] = 4D_N^2 E_{\mu_a}[|\varphi_k|^2] + 4D_N E_{\mu_a}[|\varphi_k|^4] + \mathcal{I}_{N,1}(k) + \mathcal{I}_{N,2}(k),$$

where

$$\begin{split} \mathcal{I}_{N,1}(k) &= E_{\mu_a}[|\varphi_k|^6] + 4E_{\mu_a}[|\varphi_k|^2] \sum_{0 < |l| \le N, l \neq k} E_{\mu_a}[|\varphi_l|^4] \\ &+ \sum_{0 < |m| \le N} E_{\mu_a}[|\varphi_{k-m}|^4] E[|\varphi_{k-2m}|^2], \\ \mathcal{I}_{N,2}(k) &= 2 \sum_{0 < |m| \le N, 0 < |l| \le N, l+m \neq k, l+2m \neq k} E_{\mu_a}[|\varphi_{k-m}|^2] E_{\mu_a}[|\varphi_l|^2] E_{\mu_a}[|\varphi_{l+m}|^2] \end{split}$$

In dimension 1, it is easy to see that $\sup_N [D_N + \mathcal{I}_{N,1}(k) + \mathcal{I}_{N,2}(k)] < \infty$.

In dimension 2, the identity (6) and the upper estimate (12) in Lemma 6.1 imply that for $\beta < 0$ we have

$$\sup_{N}\sum_{0<|k|\leqslant N}|k|^{2\beta}\mathcal{I}_{N,1}(k)<\infty,$$

and the upper estimate (11) implies

$$\sup_{N} \sum_{0 < |k| \le N} |k|^{2\beta} \mathcal{I}_{N,2}(k) < \infty.$$

A similar (easier) computation, based on m = 0 and $l = k - m \neq k$, implies that for $\tilde{T}_k(l,m) = \varphi_{k-m}\varphi_{l+m}\overline{\varphi_l}\varphi_k$, we have

(10)
$$\sum_{|l|,|m| \leq N} E_{\mu_a} [\operatorname{Re} (\tilde{T}_k(l,m)] = E[|\varphi_k|^4] + 2 \sum_{1 \leq |l| \leq N, l \neq k} E[|\varphi_k|^2] E[|\varphi_l|^2].$$

Using (9) and (10) we deduce that

$$E_{\mu_a}[|B_k(\varphi) - C_N\varphi_k|^2] = \mathcal{I}_{N,1}(k) + 4D_N^2 E[|\varphi_k|^2] + 4D_N E[|\varphi_k|^4] - 2C_N E[|\varphi_k|^4] - 2C_N E[|\varphi_k|^2] 2D_N + \mathcal{I}_{N,2}(k) + C_N^2 E[|\varphi_k|^2].$$

Therefore, if $C_N = 2D_N$, we have

$$E_{\mu a}[|B_k(\varphi) - C_N \varphi_k|^2] = \mathcal{I}_{N,1}(k) + \mathcal{I}_{N,2}(k)$$

and, for $\beta < 0$,

$$\sup_{N} \sum_{0 < |k| \leq N} |k|^{-2\beta} E_{\mu_a}[|B_k(\varphi) - C_N \varphi_k|^2] < \infty.$$

For more general integers p, we proceed recursively. Indeed, we have formally (to simplify computations, we use erroneously the independence of $|B_k(\varphi) - C_N \varphi_k|^{2p}$ and $|B_k|^2 - C_N^2 |\varphi_k|^2$, but the argument can be made rigorous)

$$\begin{split} E_{\mu_a} \left(|B_k(\varphi) - C_N \varphi_k|^{2p+2} \right) \\ &= E_{\mu_a} \left(|B_k(\varphi) - C_N \varphi_k|^{2p} [|B_k(\varphi)|^2 - 2C_N Re(\varphi_k \overline{B_k}) + C_N^2 |\varphi_k|^2] \right) \\ &= E_{\mu_a} \left(|B_k(\varphi) - C_N \varphi_k|^{2p} [|B_k(\varphi)|^2 - C_N^2 |\varphi_k|^2] \right) \\ &= E_{\mu_a} \left(|B_k(\varphi) - C_N \varphi_k|^{2p} [\mathcal{I}_{N,1}(k) + \mathcal{I}_{N,2}(k)] \right) \\ &\leq C E_{\mu_a} \left(|B_k(\varphi) - C_N \varphi_k|^{2p} \right). \end{split}$$

This concludes the proof.

Next, we turn to estimate the first derivative of : B :. We have

Proposition 3.2. Let d = 1, 2. For a = 2 and $\beta < 0$, the gradient of the vector field B for d = 1 (resp. : B : for d = 2) belongs to $L^{2p}_{\mu_a}(H^{\beta}; H.S.(H^1, H^{\beta}))$ for any integer $p \ge 1$.

Proof. As in the previous proof, we work at fixed N but we drop the subscript N in $(:B:)^N$. A straight-forward computation gives

$$D_{e_j}B_k(\varphi) = \sum_{l,m} (\varphi_{l+m}\overline{\varphi_l}\delta_{j,k-m} + \varphi_{k-m}\overline{\varphi_l}\delta_{j,l+m} + \varphi_{l+m}\varphi_{k-m}\delta_{-j,l})$$
$$= \sum_{l\notin\{0,j-k\}} \varphi_{l+k-j}(2\overline{\varphi_l} + \varphi_{-l}).$$

Thus, we have

$$|D_{e_j}B_k(\varphi)|^2 = \sum_{l,l'\notin\{0,j-k\}} \varphi_{l+k-j} (2\overline{\varphi_l} + \varphi_l) \overline{\varphi_{l'+k-j}} (2\varphi_{l'} + \overline{\varphi_{l'}}).$$

Let us first consider the 1-dimensional case. Since

$$\|\nabla B(\varphi)(\hat{e}_j)\|_{H^{\beta}}^2 \leqslant c \sum_k \frac{|k|^{2\beta}}{|j|^2} \Big| \sum_{l \notin \{0,j-k\}} \varphi_{l+k-j} (2\overline{\varphi_l} + \varphi_{-l}) \Big|^2,$$

we obtain that

$$E_{\mu_a} \Big[\|\nabla B(\varphi)\|_{H.S.(H^b, H^\beta)}^2 \Big] \leqslant c \sum_{k,j} \frac{|k|^{2\beta}}{|j|^2} E_{\mu_a} \Big[\sum_{l,l' \notin \{0, l+k-j\}} T_{j,k}(l,l') \Big],$$

where $T_{j,k}(l,l') = \varphi_{l+k-j}(2\overline{\varphi_l} + \varphi_{-l})\overline{\varphi_{l'+k-j}}(2\varphi_{l'} + \overline{\varphi_{-l'}}).$

Case 1:
$$l = l', T_{j,k}(l,l) = |\varphi_{l+k-j}|^2 (4|\varphi_l|^2 + 2\overline{\varphi_l}\varphi_{-l} + 2\varphi_{-l}\varphi_l + |\varphi_{-l}|^2).$$

• If $k = j, E_{\mu_a}[T_{j,j}(l,l)] = 4E_{\mu_a}[|\varphi_l|^4] + E_{\mu_a}[|\varphi_l|^2]E_{\mu_a}[|\varphi_{-l}|^2].$
• If $k \neq j$ and $2l + k - j \neq 0$. $E_{\mu_a}[T_{j,k}(l,l)] = 4E_{\mu_a}[|\varphi_{-l}|^2]E_{\mu_a}[|\varphi_{-l}|^2].$

• If $k \neq j$ and $2l + k - j \neq 0$, $E_{\mu_a}[T_{j,k}(l,l)] = 4E_{\mu_a}[|\varphi_{l+k-j}|^2]E_{\mu_a}[|\varphi_l|^2] + E_{\mu_a}[|\varphi_{l+k-j}|^2]E_{\mu_a}[|\varphi_{-l}|^2].$ • If $k \neq j$ and 2l + k - j = 0, $E_{\mu_a}[T_{j,k}(l,l)] = 4E_{\mu_a}[|\varphi_{l+k-j}|^2]E_{\mu_a}[|\varphi_l|^2] + E_{\mu_a}[|\varphi_{l+k-j}|^4].$

Case 2: $l \neq l'$. • If k = j, $T_{j,j}(l, l') = 4|\varphi_l|^2 |\varphi_{l'}|^2 + 2|\varphi_{l'}|^2 \varphi_l \varphi_{-l} + 2|\varphi_l|^2 \overline{\varphi_{l'} \varphi_{-l'}} + \overline{\varphi_{l'} \varphi_{-l'}}$. Therefore, $E_{\mu_a}[T_{j,j}(l, l')] = 4E_{\mu_a}[|\varphi_l|^2]E_{\mu_a}[|\varphi_{l'}|^2]$. • If $k \neq j$, $E_{\mu_a}[T_{j,k}(l, l')] = 0$, except if l + k - j = -l'. In that case

• If $k \neq j$, $E_{\mu_a}[T_{j,k}(l,l')] = 0$, except if l + k - j = -l'. In that case we have $E_{\mu_a}[T_{j,k}(l,l')] = E_{\mu_a}[|\varphi_{l+k-j}|^2]E_{\mu_a}[|\varphi_{-l}|^2]$ if $l + k - j \neq -l$, and $E_{\mu_a}[|\varphi_{l+k-j}|^4]$ if l + k - j = -l.

For a = 2 and $\beta < 0$ we have

$$\sum_{k,j\neq 0} \frac{|k|^{2\beta}}{|j|^2} \Big[\sum_{l\neq 0} |l|^{-2a} + \sum_{l\notin\{0,j-k\}} |l|^{-a} |l+k-j|^{-a} + \sum_{l,l'\neq 0,l\neq l'} |l|^{-a} |l'|^{-a} \Big] < \infty.$$

Next, we deal with the 2-dimensional case. We need to be more careful in this case and we begin by computing $\|D_{e_j}B_k(\varphi)\|_{L^2_{\mu_a}}$. As in the one-dimensional case, we have

$$\begin{split} \int |D_{e_j} B_k(\varphi)|^2 d\mu_a &= \delta_{j,k} \Big(4 \sum_{l \neq 0} E_{\mu_a} [|\varphi_l|^4] + \sum_{l \neq 0} E_{\mu_a} [|\varphi_l|^2] E_{\mu_a} [|\varphi_{-l}|^2] \\ &+ 4 \sum_{l,l' \neq 0, l \neq l'} E_{\mu_a} [|\varphi_l|^2] E_{\mu_a} [|\varphi_{l'}|^2] \Big) \\ &+ 1_{\{j \neq k\}} \Big(4 \sum_{l \notin \{0, j - k\}} E_{\mu_a} [|\varphi_{l+k-j}|^2] E_{\mu_a} [|\varphi_l|^2] + 21_{\{j-k \text{ even}\}} E_{\mu_a} [|\varphi_{(k-j)/2}|^4] \\ &+ 2 \sum_{2l \notin \{0, j-k\}} E_{\mu_a} [|\varphi_{l+k-j}|^2] E_{\mu_a} [|\varphi_{-l}|^2] \Big). \end{split}$$

$$\begin{aligned} \mathcal{J}_{N}(j,k) &= \delta_{j,k} \Big(4 \sum_{l \neq 0} E_{\mu_{a}}[|\varphi_{l}|^{4}] + \sum_{l \neq 0} E_{\mu_{a}}[|\varphi_{l}|^{2}] E_{\mu_{a}}[|\varphi_{-l}|^{2}] \Big) \\ &+ 1_{\{j \neq k\}} \Big(4 \sum_{l \notin \{0,j-k\}} E_{\mu_{a}}[|\varphi_{l+k-j}|^{2}] E_{\mu_{a}}[|\varphi_{l}|^{2}] + 21_{\{j-k \text{ even}\}} E_{\mu_{a}}[|\varphi_{(k-j)/2}|^{4}] \\ &+ 2 \sum_{2l \notin \{0,j-k\}} E_{\mu_{a}}[|\varphi_{l+k-j}|^{2}] E_{\mu_{a}}[|\varphi_{-l}|^{2}] \Big). \end{aligned}$$

Using (7), we deduce that

$$D_{e_j}(B)_k(\varphi) = D_{e_j}B_k(\varphi) - C_N\delta_{j,k}.$$

Therefore, since $C_N = 2D_N$,

$$\int |D_{e_j}B_k(\varphi)|^2 d\mu_a = \mathcal{J}_N(j,k) + \delta_{j,k} \left(4D_N^2 - 2C_N 2D_N + C_N^2\right) = \mathcal{J}_N(j,k).$$

Using (6) and (11) in Lemma 6.1, we deduce that

$$\sup_{N} \sum_{j,k} \frac{|k|^{2\beta}}{|j|^2} \mathcal{J}_N(j,k) < \infty.$$

For more general integers p, we proceed recursively. Indeed, we have formally (as previously, we use erroneously the independence of $|D_{e_j}B_k(\varphi)|^2$ and $|D_{e_j}B_k(\varphi)|^{2p}$, but the argument can be made rigorous)

$$E_{\mu_a}(|D_{e_j}B_k(\varphi)|^{2p+2}) = E_{\mu_a}(|D_{e_j}B_k(\varphi)|^{2p}|D_{e_j}B_k(\varphi)|^2)$$
$$= E_{\mu_a}(|D_{e_j}B_k(\varphi)|^{2p}\mathcal{J}_N(j,k))$$
$$\leqslant CE_{\mu_a}(|D_{e_j}B_k(\varphi)|^{2p}).$$

This concludes the proof.

Using the definition of the second derivatives and corresponding norms presented in section 2, we have the following result:

Proposition 3.3. Let $a = 2, \beta < 0$; the second derivative of B belongs to $L^{2p}_{\mu_a}(H^{\beta}, H.S.(H^1 \otimes H^1, H^{\beta}))$, for any integer $p \ge 1$.

Proof. A straight-forward computation gives

$$\|D_{\hat{e}_l} D_{\hat{e}_m} B(\varphi)\|_{H^{\beta}}^2 \lesssim \frac{1}{(|l||m|)^2} \sum_k |k|^{2\beta} (|\varphi_{l+m-k}|^2 + |\varphi_{l-m+k}|^2).$$

By definition of the Hilbert-Schmidt norm, this leads to

$$\|\nabla^2 B(\varphi)\|_{H.S.(H^1 \times H^1, H^\beta)}^2 \lesssim \sum_{l,m} \frac{1}{(|l||m|)^2} \sum_k |k|^{2\beta} (|\varphi_{l+m-k}|^2 + |\varphi_{l-m+k}|^2).$$

Using once more (6), we deduce that

$$E_{\mu_a} \|\nabla^2 B(\varphi)\|_{H.S.(H^1 \times H^1, H^\beta)}^2 \lesssim \sum_{l,m} \frac{1}{(|l||m|)^2} \sum_{k \neq l+m} |k|^{2\beta} \frac{1}{|l+m-k|^a}.$$

Let

This last term converges when $\beta < 0$ (see (11) in Lemma 6.1). As previously, the case for more general integers p can be made by induction.

More generally, the following holds:

Proposition 3.4. For any $s \in \mathbb{N}$ and any integer $p \ge 1$, $\nabla^s : B :$ belong to $L^p_{\mu_a}(H^\beta; H.S.(\otimes^s H^1, H^\beta))$ when a = 2 and $\beta < 0$.

4. The renormalised mass

We will consider level surfaces for the renormalised mass E, where

$$E(\varphi) = \sum_{k} (|\varphi_k|^2 - Z_k).$$

Let us recall that the mass is an invariant of the NLS flow. We choose Z_k such that $E \in L^1_{\mu_a}$, i.e., $Z_k = c_1 |k|^{-a}$ where c_p was defined in (6). As for B, this renormalisation is not needed in dimension 1 since the serie is convergent. We recall that $W_{\infty} = \bigcap_{r,p} W^p_r(H^{\beta}, \mathbb{R})$.

Our goal in this section is to show that E and $\|\nabla E\|_{H.S.(H^l,\mathbb{R})}^{-1}$ belong to the space W_{∞} , as it was defined before, namely for a = 2 and l = 1 (let us recall that the norm of this space depends on μ_a). As previously, our choice of parameters are linked with the fact that we want to work on the abstract Wiener space $(H^{\beta}, H^1, \mu_2), \beta < 0$. We begin by estimating E.

Proposition 4.1. The mass E belongs to W_{∞} .

Proof. A straight-forward computation gives

$$\begin{split} &\int |E(\varphi)|^2 d\mu_a \\ &= \int \sum_k (|\varphi_k|^2 - Z_k)^2 d\mu_a + \sum_{k \neq l} \int |\varphi_k|^2 |\varphi_l|^2 d\mu_a - 2 \sum_{k \neq l} \int Z_l |\varphi_k|^2 d\mu_a + \sum_{k \neq l} Z_k Z_l \\ &= \int \sum_k |\varphi_k|^4 d\mu_a - 2 \int \sum_k Z_k |\varphi_k|^2 d\mu_a + \sum_k Z_k^2 \\ &\quad + \sum_{k \neq l} \int |\varphi_k|^2 |\varphi_l|^2 d\mu_a - 2 \sum_{k \neq l} \int Z_l |\varphi_k|^2 d\mu_a + \sum_{k \neq l} Z_k Z_l \\ &= \sum_k c_2 |k|^{-2a} - 2 \sum_k Z_k c_1 |k|^{-a} + \sum_k Z_k^2 \\ &\quad + \sum_{k \neq l} c_1 |k|^{-a} c_1 |l|^{-a} - 2 \sum_{k \neq l} Z_l c_1 |k|^{-a} + \sum_{k \neq l} Z_k Z_l \\ &= \sum (c_2 |k|^{-2a} - Z_k^2) \lesssim |k|^{-2a}. \end{split}$$

To deal with a general integer $p \ge 2$, using the independence of the $\varphi_{k_j}, j = 1, ..., p$ and an induction argument, we have

$$\begin{split} \int |E(\varphi)|^{p} d\mu_{a} &= \int \sum_{k} (|\varphi_{k}|^{2} - Z_{k})^{p} d\mu_{a} + \sum_{k_{1} \neq k_{j}} \int (|\varphi_{k_{1}}|^{2} - Z_{k_{1}}) \Pi_{k_{j}} (|\varphi_{k_{j}}|^{2} - Z_{k_{j}}) d\mu_{a} \\ &= \int \sum_{k} (|\varphi_{k}|^{2} - Z_{k})^{p} d\mu_{a} + \sum_{k_{1} \neq k_{j}} \int |\varphi_{k_{1}}|^{2} \Pi_{k_{j}} (|\varphi_{k_{j}}|^{2} - Z_{k_{j}}) d\mu_{a} \\ &- \sum_{k_{1} \neq k_{j}} \int Z_{k_{1}} \Pi_{k_{j}} (|\varphi_{k_{j}}|^{2} - Z_{k_{j}}) d\mu_{a} \\ &\lesssim \sum_{k} |k|^{-ap} + \sum_{k_{1} \neq k_{j}} \int Z_{k_{1}} \Pi_{k_{j}} (|\varphi_{k_{j}}|^{2} - Z_{k_{j}}) d\mu_{a} \\ &- \sum_{k_{1} \neq k_{j}} \Pi_{k_{j}} \int Z_{k_{1}} (|\varphi_{k_{j}}|^{2} - Z_{k_{j}}) d\mu_{a} \\ &\lesssim \sum_{k} |k|^{-ap}. \end{split}$$

Next, we have

$$\nabla E(\varphi)(e_k) = 2Re(\varphi_k).$$

Therefore, we deduce

$$\|\nabla E(\varphi)\|_{L^{2m}_{\mu_a}(H.S.(H^1,\mathbb{R}))}^{1/m} \leqslant c \sum_k \left(E_{\mu_a} \left[|k|^{-2m} |\varphi_k|^{2m} \right] \right)^{1/m} \leqslant c(p) \sum_k \frac{1}{|k|^{a+2}} < \infty.$$

Finally, we also have

$$\|\nabla^2 E(\varphi)\|_{L^{2m}_{\mu_2}(H.S.(H^1 \otimes H^1, \mathbb{R})}^{1/m} \leq C \sum_k |k|^{-4} < \infty.$$

We now turn to the estimate of $\|\nabla E\|_{H.S.(H^1,\mathbb{R})}^{-1}$.

Proposition 4.2. The random variable $\|\nabla E\|_{H.S.(H^1,\mathbb{R})}^{-1}$ belongs to the space W_{∞} .

Proof. By Chebycheff's inequality, we have, for all t > 0 and $\varepsilon > 0$,

$$\begin{split} \mu_a \{ \| \nabla E(\varphi) \|_{H.S.(H^1,\mathbb{R})}^2 \leqslant \varepsilon \} \leqslant e^{t\varepsilon} E_{\mu_a} \Big[e^{-t \| \nabla E(\varphi) \|_{H.S.(H^1,\mathbb{R})}^2} \Big] \\ \leqslant C e^{t\varepsilon} \prod_{k \neq 0} |k|^a \int_0^\infty e^{-4t|k|^{-2}r^2 - \frac{|k|^a}{2}r^2} r dr \\ \leqslant C e^{t\varepsilon} \prod_{k \neq 0} \frac{1}{1 + 8t/|k|^{a+2}}. \end{split}$$

On the other hand, we have, for some c > 0,

$$\prod_{k \neq 0} \left(1 + \frac{8t}{|k|^{a+2}} \right) \ge \prod_{k \neq 0, |k| \le (8t)^{1/(a+2)}} \left(1 + \frac{8t}{|k|^{a+2}} \right)$$
$$\ge 2^{(8t)^{\frac{1}{a+2}}} = e^{ct^{\frac{1}{a+2}}}.$$

Therefore, we deduce that for all $\varepsilon > 0$,

$$\mu_a\{\|\nabla E(\varphi)\|_{H.S.(H^1,\mathbb{R})}^2 \leqslant \varepsilon\} \leqslant \inf_t e^{t\varepsilon - ct^{\frac{1}{a+2}}} = e^{-c\varepsilon^{\frac{1}{1-(a+2)}}}.$$

Thus, letting $\|\nabla E(\varphi)\|_{H.S.(H^1,\mathbb{R})}^2 = X$ and using the equality $\mu_a(1/|X| \ge \varepsilon) = \mu_a(|X| \le 1/\varepsilon)$, we have

$$E_{\mu_a}\left(\frac{1}{|X|}^p\right) = \int_0^\infty p\varepsilon^{p-1}\mu_a\left(\frac{1}{|X|} \ge \varepsilon\right)d\varepsilon$$
$$\leqslant \int_0^\infty p\varepsilon^{p-1}e^{-c\varepsilon^{\frac{1}{(a+2)-1}}}d\varepsilon < \infty$$

when 1 < a + 2 and, in particular, for a = 2.

5. The NLS equation and invariant measures

Following [4], if we consider the renormalised NLS equation and its corresponding renormalised Hamiltonian, which is the limit in N of

$$H_N(u) = \int |\nabla u|^2 + \frac{1}{2} \int |u|^4 - 2a_N \int |u|^2 + a_N^2,$$

with $a_N = 2 \sum_{k:|k| < N, k \neq 0} \frac{1}{|k|^2}$. The corresponding Gibbs measure

$$d\mu(\varphi) = \lim_{N} e^{-\frac{1}{2}\int |u|^4 - 2a_N \int |u|^2 + a_N^2} d\mu_2(\varphi)$$

is invariant for the renormalised flow (4), and the measures $e^{-H_N(\varphi)} \Pi d\varphi$ are invariant for the corresponding truncated equations. Moreover μ has a Radon-Nikodym derivative w.r.t. μ_2 which belongs to all L^p spaces (all these statements were shown in [4]).

Therefore, by Theorem 2.1 and thanks to Propositions 4.1 and 4.2, we can repeat the construction of the surface measures on level sets (V_r) of the mass E, starting from this weighted Wiener measure μ instead of μ_2 . Let us denote by σ^r this surface measure on the level set of renormalised mass equal to r.

In this section we prove the existence of a global solution of the renormalised cubic nonlinear Schrödinger equation in the two-dimensional torus in a weak sense and living in the level sets of the renormalised mass. More precisely, we have

Theorem 5.1. For each r such that $\rho(r) = \frac{d(E*\mu)}{dr} > 0$ there exists a probability space (Ω, \mathbb{F}, P) and a random process $u_{\omega} \in C(\mathbb{R}, V_r)$, $\omega \in \Omega$, such that for every t

(i)
$$u_t(\omega) = e^{itA}u_0(\omega) + \int_0^t e^{-i(t-s)A} : B : (u_s(\omega))ds \qquad P-a.s$$

where : B :* is a redifinition of : B : as defined in Section 2.4. Moreover the measure σ^r is invariant for the process, namely

(*ii*)
$$\int_{\Omega} f(u_t(\omega))dP(\omega) = \int_{V_r} f d\sigma^r$$

for every cylindrical f functional and every t.

Proof. The proof is similar to the one of [8, Theorem 4.1] for the Euler equation.

First notice that, given the estimates of all Sobolev norms of the field : B: which were previously derived, one can consider a redefinition : B:* on each level set V_r . By the definition of the surface measure, : B:* belongs also to all Sobolev spaces with respect to this measure.

From the regularity of the finite dimensional approximations : B:ⁿ for every n, we know that there exists a global solution $v_t^n(\varphi)$ of the equation, written here in its integral form,

$$v_t^n(\varphi) = e^{itA}\varphi + \int_0^t e^{-i(t-s)A} : B :^n (v_s^n(\varphi))ds \ , \varphi \in H^\beta.$$

For every cylindrical functional f and every $h \in C^{\infty}(\mathbb{R})$ with compact support we have,

$$\begin{split} \int_{\mathbb{R}} h(r)\rho(r) \int_{V_r} &< (A - :B:^n)(\varphi), \nabla f(\varphi) >_{H^1} d\sigma^r(\varphi) dr \\ &= \int_{H^\beta} < h((E(\varphi))(A - :B:^n)(\varphi), \nabla f(\varphi) >_{H^1} d\mu(\varphi) \\ &= -\int_{H^\beta} [h'(E(\varphi))(A - :B:^n), \nabla E >_{H^1} (\varphi)] f(\varphi) d\mu(\varphi) \\ &= 0, \end{split}$$

where the last equality is a consequence of the invariance of the mass and, in particular,

$$\int_{H^{\beta}} \langle (A - :B:^n), \nabla(h(E)f) \rangle \rangle_{H^1}(\varphi) d\mu(\varphi) = 0.$$

This shows the invariance of the surface measure with respect to the flows v_t^n .

Denote by Q_n^k the law of $[v_t^n]^k$ in the space $C(\mathbb{R}^+;\mathbb{R})$ endowed with the supremum norm, namely

$$Q_n^k(\Gamma) = \sigma^r(\{[v^n(\varphi)]^k \in \Gamma, \quad \Gamma \subset C(\mathbb{R}^+; \mathbb{R})\})$$

The index k stands here for component in the H^{β} basis $\{\frac{1}{k^{\beta}}e_k\}$.

On the space of probability measures over the real valued continuous path space we consider the weak topology. We show that, for $w \in C(\mathbb{R}^+; \mathbb{R})$,

$$\lim_{R \to \infty} \sup_{n} Q_n^k \Big(|w(0)| > R \Big) = 0,$$

and that for every R > 0 and T > 0

$$\lim_{\delta \to 0} \sup_{n} Q_n^k \Big(\sup_{0 \le t \le t' \le T, t' - t \le \delta} |w(t) - w(t')| \ge R \Big) = 0.$$

For negative times the proof is analogous.

The first statement is simply due to the estimate

$$\begin{split} Q_n^k \Big(|w(0)| > R \Big) &\leqslant \frac{1}{R^2} \int |w(0)|^2 dQ_n^k \\ &= \frac{1}{R^2} \int |\varphi^k|^2 d\sigma^r(\varphi), \end{split}$$

where in the first inequality we used Chebycheff's inequality. Concerning the second one, we have

$$\begin{aligned} Q_n^k \Big(\sup_{0 \leqslant t \leqslant t' \leqslant T, t' - t \leqslant \delta} |w(t) - w(t')| &\geq R \Big) \\ &\leqslant \frac{1}{R} \int_{V_r} \sup |[[v_t^n(\varphi) - v_{t'}^n(\varphi)]^k| d\sigma^r(\varphi) \\ &\leqslant \frac{1}{R} \sup \Big(\int_{V_r} \Big(|[e^{it'A}\varphi - e^{itA}\varphi]^k| + \Big| \Big[\int_0^t |(e^{-i(t-s)A}) : B :^n (v_s^n(\varphi))] ds \\ &\quad - \int_0^{t'} (e^{-i(t'-s)A}) : B :^n (v_s^n(\varphi)) \Big]^k \Big| \Big) d\sigma^r(\varphi) \Big) \end{aligned}$$

Using equalities $\frac{d}{dt}e^{itA}\varphi = itAe^{itA}\varphi$ and

$$\begin{aligned} \frac{d}{dt} \int_0^t (e^{-i(t-s)A}) &: B :^n (v_s^n(\varphi)) ds \\ &= -iAe^{-itA} \int_0^t e^{isA} : B :^n (v_s^n(\varphi)) ds + : B :^n (v_s^n(\varphi)), \end{aligned}$$

we derive

$$\begin{split} Q_n^k \Big(\sup_{0 \leq t \leq t' \leq T, t' - t \leq \delta} |w(t) - w(t')| \geq R \Big) \\ &\leqslant \frac{\delta^2}{R^2} c(k, T) \Big(\int_0^T \int_0^T [|e^{isA}\varphi|^k]^2 ds d\sigma^r(\varphi) \Big)^{\frac{1}{2}} \\ &\quad + \frac{\delta^2}{R^2} c(k, T) \Big(\int_0^T \int_0^T [|e^{-isA}e^{irA}\varphi|^k]^2 ds dr d\sigma^r(\varphi) \Big)^{\frac{1}{2}} \\ &\quad + \frac{\delta^2}{R^2} c(k, T) \Big(\int_0^T \int_0^T [|:B^n:(\varphi)]^k|^2 dt d\sigma^r(\varphi) \Big)^{\frac{1}{2}} \\ &\leqslant \frac{\delta c_r}{R^2} \int_{V_r} \int_0^T \Big(\|\varphi\|_{H^\beta} + \|:B:^n(\varphi)\|_{H^\beta} \Big) ds d\sigma^r(\varphi), \end{split}$$

where we have denoted c_r a generic constant depending on r, k and T, and the last step is a consequence of the invariance of the process with respect to the surface measure.

The limit follows from the fact that, since the sequence $\{: B : n\}_n$, as well as that of their gradients converges to : B : in all L^p_{μ} spaces, it also converges (to $: B :^*$) in $L^p_{\sigma^r}$. Hence

$$Q_n^k \Big(\sup_{0 \le t \le t' \le T, t' - t \le \delta} |w(t) - w(t')| \ge R \Big) \le \frac{\delta}{R^2} T c$$

for some constant c, independent of n. We conclude that the sequence of probability laws Q_n^k is tight on $C(R^+; R)$, and therefore that we can extract a subsequence converging weakly to a probability measure Q. By Skorohod's theorem there exists a probability space (Ω, \mathbb{F}, P) and processes $u_t^n(\omega) = \sum_k [u_t^n(\omega)]^k e_k, u_t(\omega)$ with laws Q_n^k and Q^k resp., such that $u_i^n \to u$. P-almost surely.

We now prove (ii).

For every cylindrical functional f, we have

$$\int f(u_t(\omega))dP(\omega) = \lim_n \int f(u_t^n(\omega))dP(\omega)$$
$$= \lim_n \int_{V^r} f(v_t^n(\varphi))d\sigma^r(\varphi)$$
$$= \int fd\sigma^r.$$

For the last equality we have used the invariance of the measure σ^r for the flow v_t^n .

Recalling that $\rho(r) = \frac{d(E * \mu_2)}{dr} > 0$, we have

$$\begin{split} \int_{\mathbb{R}^{+}} \psi(r)\rho(r) \int_{V_{r}} f(v_{t}^{n}(\varphi)) d\sigma^{r}(\varphi) dr &= \int \psi(E(\varphi)) f(v_{t}^{n}(\varphi)) d\mu(\varphi) \\ &= \int_{V_{r}} \psi(E(v_{-t}^{n}(\varphi)) f(\varphi) d\mu(\varphi) \\ &= \int_{V_{r}} \psi(E(\varphi)) f(\varphi) d\mu(\varphi) \\ &= \int_{\mathbb{R}^{+}} \psi(r)\rho(r) \int_{V_{r}} f(\varphi) d\sigma^{r}(\varphi) dr \end{split}$$

for every smooth function ψ . In particular the process u_t takes values in V_r .

It remains to show (i). For every k, we get for $\beta < 0$

$$\begin{split} \int \int_{0}^{t} & \Big| \Big[e^{-i(t-s)A} : B :^{n} (u_{t}^{n}(\omega)) - \int_{0}^{t} e^{-i(t-s)A} : B :^{*} (u_{t}(\omega)) \Big]^{k} \Big| dsdP(\omega) \\ & \leq \int \int_{0}^{T} \Big| \Big[e^{-i(t-s)A} \Big(: B :^{n} (u_{t}^{n}) - : B :^{*} (u_{t}^{n}) \Big) \Big]^{k} \Big| dsdP \\ & + \int \int_{0}^{T} \Big| \Big[e^{-i(t-s)A} \Big(: B :^{*} (u_{t}^{n}) - : B :^{*} (u_{t}) \Big) \Big]^{k} \Big| dsdP \\ & \leq c(k) \Big(\int \int_{0}^{T} \Big\| \Big(: B :^{n} (u_{t}^{n}) - : B :^{*} (u_{t}^{n}) \Big) \Big\|_{H^{\beta}} dsdP \\ & + \int \int_{0}^{T} \Big\| \Big(: B :^{*} (u_{t}^{n}) - : B :^{*} (u_{t}) \Big) \Big\|_{H^{\beta}} dsdP \end{split}$$

The first term converges to zero due to the invariance of the measure P and the (L^1) convergence of : B^n : to : B :*. Concerning the second term, let $\Lambda \subset H_\beta$ be a set of σ^r measure equal to zero such that : B :* is continuous in Λ^c . Define

$$\Lambda_n = \{ (t, \omega) \in [0, T] \times \Omega : u_t^n(t, \omega) \in C \},\$$

and

$$\Lambda_{\infty} = \{ (t, \omega) \in [0, T] \times \Omega : u_t(t, \omega) \in C \}.$$

We have,

$$\int_0^T \int \chi_{\Lambda}(u_t^n(\omega)) dP(\omega) dt = \int_0^T \int \chi_{\Lambda}(\varphi) d\sigma^r(\varphi) dt = 0$$

i.e, the measure of Λ_n (and, analogously, of Λ), considering the Lebesgue measure in the time coordinate, is equal to zero.

On the other hand, since u^n converges to u almost surely, consider a set $D \subset \Omega$ such that $P(D^c) = 0$ where this convergence holds, in the space H^{β} and for every $t \in [0, T]$. Then the measure of $[0, T] \times D^c$ is zero.

Using Egorov's theorem, we conclude that, for (t, ω) in $\Lambda_{\infty}^{c} \bigcap_{n} \Lambda_{n}^{c} \cap ([0, t] \times D^{c})$, and $\beta < 0$

$$\iint_{0}^{T} \left\| \left(:B:^{*}(u_{t}^{n}) - :B:^{*}(u_{t}) \right) \right\|_{H^{\beta}} ds dP \to 0,$$

which concludes the proof.

6. Appendix

In this appendix, we show a convergence result for series that were used quite often in section 3.

Lemma 6.1. Let d = 1, 2. For any $\beta < 0$, it holds that

(11)
$$\sum_{0 \neq l, m \in \mathbb{Z}^2} \frac{1}{|l|^2} \frac{1}{|m|^2} \sum_{k \notin \{0, l+m\}} \frac{|k|^{\beta}}{|l+m-k|^2} < \infty,$$

and

(12)
$$\sum_{k \neq 0, m \notin \{k, k/2\}} \frac{|k|^{\beta}}{|k - m|^4|k - 2m|^2} < \infty.$$

Proof. It is immediate to see that the above series converge in dimension 1.

Let d = 2. We first prove (11). Using a change of variable, the convergence of the serie is equivalent to the one of

$$\sum_{0 \neq l, m \in \mathbb{Z}^2} \frac{1}{|l|^2} \frac{1}{|m|^2} \sum_{k \notin \{0, -(l+m)\}} \frac{|k+l+m|^\beta}{|k|^2}.$$

Let $\lambda \in (0, 1)$. We are going to split the sum in k into three parts.

• First, we take $|k| \leq (1 - \lambda)|l + m|$; this yields

$$\sum_{\substack{0\neq k\in\mathbb{Z}^2, \ |k|\leqslant (1-\lambda)|l+m|}} \frac{|k+l+m|^{\beta}}{|k|^2} \leqslant \lambda^{\beta}|l+m|^{\beta} \sum_{1\leqslant n\leqslant (1-\lambda)|l+m|} \frac{2}{n} \leqslant C_{\lambda}|l+m|^{\beta/2}.$$

• Then, we consider $|k| \ge (1 + \lambda)|l + m|$. In this case, we get

$$\sum_{\substack{0 \neq k \in \mathbb{Z}^2, \ |k| \ge (1+\lambda)|l+m| \\ \leqslant C_{\lambda} |k| \ge 1}} \frac{|k+l+m|^{\beta}}{|k|^2} \leqslant C_{\lambda} \sum_{\substack{0 \neq k \in \mathbb{Z}^2, \ |k| \ge (1+\lambda)|l+m| \\ \leqslant C_{\lambda}|l+m|^{\beta}}}$$

• Finally, when $(1 - \lambda)|l + m| \le |k| \le (1 + \lambda)|l + m|$, i.e., $|k| \approx |l + m|$, we have

$$\sum_{\substack{0\neq k\in\mathbb{Z}^2,\ |k|\approx|l+m|\\ \leqslant C_{\lambda}|l+m|^{\beta/2}} \leq C_{\lambda}|l+m|^{\beta/2}} \sum_{\substack{0\neq k\in\mathbb{Z}^2,\ |k|\approx|l+m|\\ \leqslant C_{\lambda}|l+m|^{\beta/2}}.} |k|^{\beta/2-2}$$

Therefore, we have to upper estimate

$$\sum_{\substack{0 \neq l, m \in \mathbb{Z}^2, l+m \neq 0}} \frac{1}{|l|^2} \frac{|l+m|^{\beta/2}}{|m|^2}.$$

A similar argument implies that

$$\sum_{\substack{0 \neq l, m \in \mathbb{Z}^2, l+m \neq 0}} \frac{1}{|l|^2} \frac{|l+m|^{\beta/2}}{|m|^2} \leqslant C \sum_{\substack{0 \neq l \in \mathbb{Z}^2}} \frac{|l|^{\beta/4}}{|l|^2} < \infty,$$

which completes the proof of (11).

k

A change of variables implies that

$$\sum_{\notin \{0,m,2m\}} \frac{|k|^{\beta}}{|k-m|^4|k-2m|^2} = \sum_{k \notin \{0,-m,m\}} \frac{|k+m|^{\beta}}{|k|^4|k-m|^2}.$$

A decomposition similar to the above one implies (12).

References

- Hélène Airault and Paul Malliavin. Intégration géométrique sur l'espace de Wiener. Bull. Sci. Math. (2), 112(1):3–52, 1988.
- [2] Jean Bourgain. Fourier transform restriction phenomena for certain lattice subsets and applications to nonlinear evolution equations. I. Schrödinger equations. *Geom. Funct. Anal.*, 3(2):107–156, 1993.
- [3] Jean Bourgain. Periodic nonlinear Schrödinger equation and invariant measures. Comm. Math. Phys., 166(1):1–26, 1994.
- Jean Bourgain. Invariant measures for the 2d-defocusing nonlinear Schrödinger equation. Comm. Math. Phys., 176(2):421–445, 1996.
- [5] Justin T. Brereton. Invariant measure construction at a fixed mass. Nonlinearity, 32(2):496–558, 2019.
- [6] Jean-Baptiste Casteras and Léonard Monsaingeon. Invariant measures and global well-posedness for a fractional Schrödinger equation with Moser-Trudinger type nonlinearity. Stoch. Partial Differ. Equ. Anal. Comput., 12(1):416–465, 2024.
- [7] Andreia Chapouto and Justin Forlano. Invariant measures for the periodic kdv and mkdv equations using complete integrability. arXiv preprint arXiv:2305.14565, 2023.
- [8] Fernanda Cipriano. The two-dimensional Euler equation: a statistical study. Comm. Math. Phys., 201(1):139–154, 1999.
- [9] Yu Deng, Andrea R. Nahmod, and Haitian Yue. Invariant Gibbs measures and global strong solutions for nonlinear Schrödinger equations in dimension two. Ann. of Math. (2), 200(2):399–486, 2024.
- [10] Ludwig D. Faddeev and Leon A. Takhtajan. Hamiltonian methods in the theory of solitons. Springer Series in Soviet Mathematics. Springer-Verlag, Berlin, 1987. Translated from the Russian by A. G. Reyman [A. G. Reĭman].
- [11] Leonard Gross. Measurable functions on Hilbert space. Trans. Amer. Math. Soc., 105:372–390, 1962.
- [12] Benjamin Harrop-Griffiths, Rowan Killip, and Monica Vi,san. Sharp well-posedness for the cubic NLS and mKdV in $H^{s}(\mathbb{R})$. Forum Math. Pi, 12:Paper No. e6, 86, 2024.
- [13] Sebastian Herr and Beomjong Kwak. Strichartz estimates and global well-posedness of the cubic nls on. In *Forum of Mathematics, Pi*, volume 12, page e14. Cambridge University Press, 2024.
- [14] Rowan Killip and Monica Vi, san. Scale invariant Strichartz estimates on tori and applications. *Math. Res. Lett.*, 23(2):445–472, 2016.

INVARIANT MEASURE NLS

- [15] Herbert Koch and Daniel Tataru. Conserved energies for the cubic nonlinear Schrödinger equation in one dimension. Duke Math. J., 167(17):3207–3313, 2018.
- [16] Sergei Kuksin and Armen Shirikyan. Randomly forced CGL equation: stationary measures and the inviscid limit. J. Phys. A, 37(12):3805–3822, 2004.
- [17] Sergei Kuksin and Armen Shirikyan. Mathematics of two-dimensional turbulence, volume 194 of Cambridge Tracts in Mathematics. Cambridge University Press, Cambridge, 2012.
- [18] Beomjong Kwak and Soonsik Kwon. Critical local well-posedness of the nonlinear schrödinger equation on the torus. Annales de l'Institut Henri Poincaré C, 2024.
- [19] Joel L. Lebowitz, Harvey A. Rose, and Eugene R. Speer. Statistical mechanics of the nonlinear Schrödinger equation. J. Statist. Phys., 50(3-4):657–687, 1988.
- [20] Paul Malliavin. Stochastic analysis, volume 313 of Grundlehren der mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]. Springer-Verlag, Berlin, 1997.
- [21] Tadahiro Oh and Jeremy Quastel. On invariant Gibbs measures conditioned on mass and momentum. J. Math. Soc. Japan, 65(1):13–35, 2013.
- [22] Tadahiro Oh, Philippe Sosoe, and Leonardo Tolomeo. Optimal integrability threshold for Gibbs measures associated with focusing NLS on the torus. *Invent. Math.*, 227(3):1323–1429, 2022.
- [23] Fabrice Planchon, Nikolay Tzvetkov, and Nicola Visciglia. Transport of Gaussian measures by the flow of the nonlinear Schrödinger equation. Math. Ann., 378(1-2):389– 423, 2020.
- [24] Gigliola Staffilani. Periodic Schrödinger equations in Hamiltonian form. In HCDTE lecture notes. Part II. Nonlinear hyperbolic PDEs, dispersive and transport equations, volume 7 of AIMS Ser. Appl. Math., page 66. Am. Inst. Math. Sci. (AIMS), Springfield, MO, 2013.
- [25] Chenmin Sun and Nikolay Tzvetkov. Quasi-invariance of gaussian measures for the 3d energy critical nonlinear Schrödinger equation. arXiv preprint arXiv:2308.12758, 2023.
- [26] Mouhamadou Sy. Almost sure global well-posedness for the energy supercritical Schrödinger equations. J. Math. Pures Appl. (9), 154:108–145, 2021.
- [27] Mouhamadou Sy and Xueying Yu. Global well-posedness and long-time behavior of the fractional NLS. Stoch. Partial Differ. Equ. Anal. Comput., 10(4):1261–1317, 2022.
- [28] Nikolay Tzvetkov and Nicola Visciglia. Gaussian measures associated to the higher order conservation laws of the Benjamin-Ono equation. Ann. Sci. Éc. Norm. Supér. (4), 46(2):249–299, 2013.
- [29] Vladimir E. Zakharov and Aleksei B. Shabat. Exact theory of two-dimensional selffocusing and one-dimensional self-modulation of waves in nonlinear media. *Soviet Physics JETP* 34 (1972), no. 1, 62-69.

CEMSUL, FACULDADE DE CIÊNCIAS DA UNIVERSIDADE DE LISBOA, EDIFICIO C6, PISO 1, CAMPO GRANDE 1749-016 LISBOA, PORTUGAL

Email address: jeanbaptiste.casteras@gmail.com

GFM, DEP. DE MATEMÁTICA, INSTITUTO SUPERIOR TÉCNICO, LISBOA, PORTUGAL *Email address*: ana.cruzeiro@tecnico.ulisboa.pt

SAMM, EA 4543, UNIVERSITÉ PARIS 1 PANTHÉON SORBONNE, 90 RUE DE TOL-BIAC, 75634 PARIS CEDEX FRANCE AND LABORATOIRE DE PROBABILITÉS, STATISTIQUE ET MODÉLISATION, UMR 8001, SORBONNE UNIVERSITÉ - UNIVERSITÉ PARIS CITÉ. Email address: annie.millet@univ-paris1.fr