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NONLINEAR SCHRÖDINGER EQUATION WITH TIME DEPENDENT POTENTIAL

RÉMI CARLES

ABSTRACT. We prove a global well-posedness result for defocusing nonlinear Schrödinger equations with time dependent potential. We then focus on time dependent harmonic potentials. This aspect is motivated by Physics (Bose–Einstein condensation), and appears also as a preparation for the analysis of the propagation of wave packets in a nonlinear context. The main aspect considered here is the growth of high Sobolev norms of the solution.

1. INTRODUCTION

Let $d \geq 1$, and for $x \in \mathbf{R}^d$, consider the nonlinear Schrödinger equation

$$(1.1) \quad i\partial_t u + \frac{1}{2}\Delta u = V(t, x)u + \lambda|u|^{2\sigma}u \quad ; \quad u|_{t=0} = u_0,$$

where $V \in \mathbf{R}$ is locally bounded in time and subquadratic in space, $\lambda \in \mathbf{R}$, and the nonlinearity is energy-subcritical ($\sigma < 2/(d-2)$ if $d \geq 3$). We prove that the solution exists and is global in

$$\Sigma = \{f \in H^1(\mathbf{R}^d) \quad ; \quad x \mapsto |x|f(x) \in L^2(\mathbf{R}^d)\},$$

provided that $u_0 \in \Sigma$, $\sigma < 2/d$ (mass-subcritical nonlinearity), or $\sigma \geq 2/d$ and $\lambda \geq 0$ (defocusing nonlinearity). We then focus on the case where V is exactly quadratic in x :

$$(1.2) \quad i\partial_t u + \frac{1}{2}\Delta u = \frac{1}{2} \sum_{j=1}^d \Omega_j(t)x_j^2 u + \lambda|u|^{2\sigma}u \quad ; \quad u|_{t=0} = u_0,$$

where $\Omega_j \in \mathbf{R}$, with $\Omega_j \in C^1(\mathbf{R})$. In the isotropic case ($\Omega_j = \Omega$ for all j), we show that the above result is optimal in the sense that for all $\Omega \in C(\mathbf{R}; \mathbf{R})$, if $\lambda < 0$ and $\sigma = 2/d$, there exist blow-up solutions. We also investigate the growth of high Sobolev norms for large time.

There are at least two motivations to study (1.2). In Physics, this external potential may model a time dependent confining magnetic potential: (1.2) appears in Bose–Einstein condensation, typically for $\sigma = 1$ (or $\sigma = 2$ sometimes in the one-dimensional case $d = 1$), see e.g. [12, 22, 28]. Equation (1.2) also appears as an envelope equation in the nonlinear propagation of wave packets. In the linear case, consider

$$i\varepsilon\partial_t \psi^\varepsilon + \frac{\varepsilon^2}{2}\Delta \psi^\varepsilon = V(x)\psi^\varepsilon \quad ; \quad \psi^\varepsilon(0, x) = \frac{1}{\varepsilon^{d/4}} \varphi\left(\frac{x-x_0}{\sqrt{\varepsilon}}\right) e^{i(x-x_0)\cdot\xi_0/\varepsilon}.$$

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In the limit $\varepsilon \rightarrow 0$, ψ^ε can be approximated as follows:

$$\psi^\varepsilon(t, x) \underset{\varepsilon \rightarrow 0}{\sim} \frac{1}{\varepsilon^{d/4}} u \left(t, \frac{x - x(t)}{\sqrt{\varepsilon}} \right) e^{i\phi(t, x)/\varepsilon},$$

where $(x(t), \xi(t))$ is given by the Hamiltonian flow associated to $H = \frac{|\xi|^2}{2} + V(x)$, with initial data (x_0, ξ_0) ,

$$\phi(t, x) = (x - x(t)) \cdot \xi(t) + \int_0^t \left(\frac{1}{2} |\xi(\tau)|^2 - V(x(\tau)) \right) d\tau,$$

and u is given by the equation

$$i\partial_t u + \frac{1}{2} \Delta u = \frac{1}{2} \langle Q(t)x, x \rangle u \quad ; \quad u|_{t=0} = \varphi.$$

Here, Q is defined by $Q(t) = \text{Hess}V(x(t))$, the Hessian of V at point $x(t)$; see e.g. [15]. We note that the external potential in this case has the form presented in (1.2). To study the nonlinear propagation of wave packets, another parameter must be taken into account: the size of the initial data. In [11], it is shown that there exists a critical size (depending on the nonlinearity and the space dimension), corresponding to a certain power of ε : for initial data which are smaller (as $\varepsilon \rightarrow 0$) than this critical size, then the nonlinearity is negligible and we retrieve the same description as above; for initial data which have the critical size, we have a similar description, up to the fact that the envelope equation is now nonlinear, of the form (1.2). To study the propagation of wave packets over large times (typically, up to – an analogue of – Ehrenfest time), one has to understand the large time behavior of solutions to (1.2).

Remark 1.1 (Time dependent nonlinearity). With little modification, we could also consider the more general equation

$$(1.3) \quad i\partial_t u + \frac{1}{2} \Delta u = \frac{1}{2} \sum_{j=1}^d \Omega_j(t) x_j^2 u + h(t) |u|^{2\sigma} u \quad ; \quad u|_{t=0} = u_0,$$

where $h \in C^\infty(\mathbf{R}; \mathbf{R})$. Following [14] (see also [13]), the regularity assumption on h could be weakened. We choose to consider an autonomous nonlinearity in most of this paper though.

Remark 1.2 (Complete integrability). The cubic one-dimensional case $d = \sigma = 1$ is special: if $\dot{\Omega} = 0$, then the equation is completely integrable ([1]). However, if $\dot{\Omega} \neq 0$, there exists no Lax pair when the nonlinearity is autonomous as in (1.2); see [31, 27]. Indeed, (1.3) has a Lax pair if and only if (recall that $d = \sigma = 1$) Ω and h are related through the identities:

$$\Omega(t) = \ddot{f}(t) - \dot{f}(t)^2 \quad ; \quad h(t) = ae^{f(t)} \quad ; \quad a \in \mathbf{R}, \quad f \in C^\infty(\mathbf{R}; \mathbf{R}).$$

The case where the above relation is not satisfied is included in Proposition 1.4.

The assumption we make on the external potential V is inspired by [20]:

Assumption 1.3. $V \in L_{\text{loc}}^\infty(\mathbf{R}_t \times \mathbf{R}_x^d)$ is smooth with respect to the space variable: for (almost) all $t \in \mathbf{R}$, $x \mapsto V(t, x)$ is a C^∞ map. Moreover, it is subquadratic in space:

$$\forall T > 0, \forall \alpha \in \mathbf{N}^d, |\alpha| \geq 2, \quad \partial_x^\alpha V \in L^\infty([-T, T] \times \mathbf{R}^d).$$

The time dependent harmonic potential that we consider in (1.2) is of course a peculiar case of such potentials V .

1.1. **L^2 -subcritical case.** When the energy is L^2 -subcritical ($\sigma < 2/d$), we have:

Proposition 1.4. *Let $\lambda \in \mathbf{R}$, $0 < \sigma < 2/d$ and V satisfying Assumption 1.3. For all $u_0 \in L^2(\mathbf{R}^d)$, (1.2) has a unique solution*

$$u \in C(\mathbf{R}; L^2(\mathbf{R}^d)) \cap L_{\text{loc}}^{\frac{4\sigma+4}{d\sigma}}(\mathbf{R}; L^{2\sigma+2}(\mathbf{R}^d)).$$

Moreover, its L^2 -norm is independent of time:

$$\|u(t)\|_{L^2(\mathbf{R}^d)} = \|u_0\|_{L^2(\mathbf{R}^d)}, \quad \forall t \in \mathbf{R}.$$

Sketch of the proof. In view of [20, 21], local in time Strichartz estimates are available. Therefore, one can reproduce the original proof of [37] (see also [13, 35]), in order to infer the result. \square

1.2. **Energy subcritical case.** In order to encompass the physical case $\sigma = 1$ when $d = 2$ or 3 , we need to consider the case $\sigma \geq 2/d$. We shall restrict our attention to H^1 -subcritical nonlinearities: $\sigma < 2/(d-2)$ when $d \geq 3$. To solve (1.2), even locally in time, one needs to work in Σ , and not only H^1 : symmetry is needed on physical and frequency sides, unless V is sublinear [9]. Local existence in Σ then follows from the dispersive estimates in [20, 21]: one can work as in the case $V \equiv 0$ (where it is possible to work in $H^1(\mathbf{R}^d)$ only). Instead of considering only $(u, \nabla u)$ as the unknown function, one can consider $(u, \nabla u, xu)$: the three functions are related through a closed family of estimates, and we get:

Proposition 1.5. *Let $\lambda \in \mathbf{R}$, V satisfying Assumption 1.3, and $\sigma > 0$ with $\sigma < 2/(d-2)$ if $d \geq 3$. For $u_0 \in \Sigma$, there exists T and a unique solution u solution to (1.1), such that*

$$u, \nabla u, xu \in C([-T, T]; L^2(\mathbf{R}^d)) \cap L_{\text{loc}}^{\frac{4\sigma+4}{d\sigma}}([-T, T]; L^{2\sigma+2}(\mathbf{R}^d)).$$

Moreover, its L^2 -norm is independent of time:

$$\|u(t)\|_{L^2(\mathbf{R}^d)} = \|u_0\|_{L^2(\mathbf{R}^d)}, \quad \forall t \in [-T, T].$$

Since in [20, 21], only bounded time intervals are considered, we give more precisions about this result in §2 in order to consider global in time solutions: we have been careful in the statement of Proposition 1.5 not to write that T depends only on $\|u_0\|_{\Sigma}$. To infer global existence results, we wish to replace the initial time $t = 0$ in (1.1) with $t = t_0 \geq 0$: under the assumptions of Proposition 1.5, it is not guaranteed that the corresponding T is independent of t_0 . However, it is proved in [20, 21] local dispersive estimates are available, uniformly on finite time intervals.

The natural candidate for an energy in the case of (1.1) is

$$(1.4) \quad E(t) = \frac{1}{2} \|\nabla u(t)\|_{L^2}^2 + \frac{1}{\sigma+1} \|u(t)\|_{L^{2\sigma+2}}^{2\sigma+2} + \int_{\mathbf{R}^d} V(t, x) |u(t, x)|^2 dx.$$

Proposition 1.6. *Under the assumptions of Proposition 1.5, if in addition V is C^1 with respect to t , and $\partial_t V$ satisfies Assumption 1.3, then $E \in C^1([-T, T]; \mathbf{R})$, and its evolution is given by*

$$(1.5) \quad \frac{dE}{dt} = \int_{\mathbf{R}^d} \partial_t V(t, x) |u(t, x)|^2 dx.$$

The proof of the above result is straightforward, by following the same lines as in the justification of similar evolution laws in, e.g., [13].

Theorem 1.7 (Global existence in Σ). *Let $\lambda \in \mathbf{R}$, $\sigma > 0$ with $\sigma < 2/(d-2)$ if $d \geq 3$, and V satisfying Assumption 1.3. For $u_0 \in \Sigma$, we can take $T = +\infty$ in Proposition 1.5 in the following cases:*

- $\sigma < 2/d$.
- $\sigma \geq 2/d$ and $\lambda \geq 0$, provided V is C^1 in t and $\partial_t V$ satisfies Assumption 1.3.

Remark 1.8. This result extends the main one in [7], where typically the (time independent) potential $-\omega_1^2 x_1^2 + \omega_2^2 x_2^2$ is considered. It is established in [7] that if $\lambda > 0$ and $\omega_1 \gg 1 + \omega_2$, then the solution to (1.2) is global, and there is scattering. The present theorem extends the existence part, but of course the assumptions of Theorem 1.7 are too general to expect a scattering result: in the case of (1.2) with $\Omega_j = 1 \forall j$, for instance, one can construct periodic solutions to (1.2), of the form $u(t, x) = e^{-i\omega t} \psi(x)$. Indeed, this amounts to finding a non-trivial solution to the elliptic problem

$$(1.6) \quad \omega \psi = H\psi + \lambda |\psi|^{2\sigma} \psi, \text{ where } H = -\frac{1}{2}\Delta + \frac{|x|^2}{2}.$$

Introduce the quantities

$$I(\psi) = \frac{1}{2} \langle H\psi, \psi \rangle - \frac{\omega}{2} \langle \psi, \psi \rangle,$$

$$M = \left\{ \psi \in \Sigma ; \frac{1}{\sigma+1} \int_{\mathbf{R}^d} |\psi(x)|^{2\sigma+2} dx = 1 \right\},$$

and consider

$$\delta = \inf_{\psi \in M} I(\psi).$$

If $\omega > d/2$ (the lowest eigenvalue of the harmonic oscillator), then (1.6) has a non-trivial solution for $\lambda > 0$. If $\omega < d/2$, then (1.6) has a non-trivial solution for $\lambda < 0$. See e.g. [10] for more details.

Theorem 1.7 shows that in the usual cases where global existence is known without a potential, the introduction of a smooth subquadratic potential V does not change this property, regardless of the time dependence of V with respect to time. In the case $\sigma = 2/d$ and $\lambda < 0$, we prove that finite time blow-up does occur for time dependent potentials, like in the case with no potential:

Proposition 1.9 (Finite time blow-up). *Let $\sigma = 2/d$ and $\lambda < 0$. Consider (1.2) with an isotropic potential: $\Omega_j = \Omega \in C(\mathbf{R}; \mathbf{R})$ is independent of j . There exist blow-up solutions for (1.2): we can find $T > 0$, and $u \in C([0, T]; \Sigma)$ solving (1.2) such that*

$$\|\nabla u(t)\|_{L^2} \xrightarrow[t \rightarrow 0]{} \infty.$$

1.3. Growth of higher order Sobolev norms. We now stick to the case of time dependent harmonic potentials, (1.2). In view of the analysis of nonlinear wave packets in a semi-classical regime ([11]), the evolution of weighted Sobolev norms of u over large time intervals is needed.

Consider first the autonomous isotropic case: $\Omega_j = \Omega$ is a constant.

If $\Omega = 0$, then at least when $\lambda \geq 0$ and $2/d \leq \sigma < 2/(d-2)$, the Sobolev norms of u are bounded, $u \in L^\infty(\mathbf{R}; H^k(\mathbf{R}^d))$ (provided the nonlinearity is sufficiently smooth), since we know that there is scattering in H^1 (because we know that there is scattering in Σ , since scattering in H^1 only is not known so far in the case

$\sigma = 2/d$); see e.g. [38]. The momenta of u grow algebraically in time (see [2]). We give a short alternative proof of these properties in an appendix.

If $\Omega > 0$, then $u \in L^\infty(\mathbf{R}; \Sigma)$, as proved by (1.5). The existence of periodic solutions to the nonlinear problem (see Remark 1.8) shows that we may also have $u \in L^\infty(\mathbf{R}; H^k)$ and $|x|^k u \in L^\infty(\mathbf{R}; L^2)$ for all $k \in \mathbf{N}$.

If $\Omega < 0$ (repulsive harmonic potential), then it is proved in [6] that every defocusing H^1 -subcritical nonlinearity is short range as far as scattering theory is concerned: if $\lambda > 0$ in (1.2), then $u(t) \sim U(t)u_+$ as $t \rightarrow +\infty$, for some $u_+ \in \Sigma$, where $U(t) = \exp(-it(-\frac{1}{2}\Delta + \frac{\Omega}{2}|x|^2))$. Assume $\Omega = -1$. Using Mehler's formula, and a decomposition of $U(t)$ of the form $U = M\mathcal{D}\mathcal{F}M$ as in [36] originally (for the case $\Omega = 0$), we notice

$$U(t)u_+(x) \underset{t \rightarrow +\infty}{\sim} \frac{1}{\sinh t} \widehat{u}_+ \left(\frac{x}{\sinh t} \right) e^{i \frac{\cosh t}{\sinh t} |x|^2}.$$

This shows that the L^2 norms of $\nabla U(t)u_+$ and $xU(t)u_+$ grow exponentially in time. By the results in [6], so do the L^2 norms of ∇u and xu . Note that at least in the linear case $\lambda = 0$, we see that the H^k -norms of u grow like e^{kt} as t goes to infinity.

Definition 1.10. *Let $u \in C(\mathbf{R}; \Sigma)$ be a solution to (1.2), and $k \in \mathbf{N}$.*

- $(Alg)_k$ is satisfied if there exists A such that for all admissible pair (p, q) ,

$$\forall \alpha, \beta \in \mathbf{N}^d, |\alpha| + |\beta| \leq k, \quad \|x^\alpha \partial_x^\beta u\|_{L^p([0, t]; L^q)} \lesssim t^A.$$

- $(Exp)_k$ is satisfied if there exists C such that for all admissible pair (p, q) ,

$$\forall \alpha, \beta \in \mathbf{N}^d, |\alpha| + |\beta| \leq k, \quad \|x^\alpha \partial_x^\beta u\|_{L^p([0, t]; L^q)} \lesssim e^{Ct}.$$

We wish to consider smooth energy-subcritical nonlinearities. Since we study homogeneous nonlinearities, we have to assume: $d \leq 3$, $\sigma \in \mathbf{N}$, with $\sigma = 1$ if $d = 3$.

Corollary 1.11. *Let $d \leq 3$, $\lambda \geq 0$ and $\sigma \in \mathbf{N}$, with $\sigma = 1$ if $d = 3$. Let $k \in \mathbf{N}$,*

$$u_0 \in H^k(\mathbf{R}^d), \text{ with } |x|^k u_0 \in L^2(\mathbf{R}^d).$$

If $\Omega_j \in C^1(\mathbf{R}; \mathbf{R})$ is compactly supported for all j , then u has the property $(Alg)_k$.

Proof. We may assume that $\text{supp } \Omega \subset [-M, M]$. From Theorem 1.7, $u|_{t=M} \in \Sigma$. It is easy to check that the higher regularity is conserved as well: this is rather straightforward, since we consider an energy-subcritical nonlinearity. The corollary then follows from the case $\Omega = 0$, where $(Alg)_k$ is satisfied, as recalled in the appendix. \square

In the case where the dependence of Ω_j with respect to time is not specified, the evolution of the energy (1.4) yields no exploitable information. Even in the L^2 -subcritical case, we will see that proving exponential control requires some work.

Proposition 1.12 (Exponential growth). *Let $d \leq 3$, $\sigma \in \mathbf{N}$, with $\sigma = 1$ if $d = 3$, $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ be locally Lipschitzean, $k \in \mathbf{N}$, $k \geq 1$, and*

$$u_0 \in H^k(\mathbf{R}^d), \text{ with } |x|^k u_0 \in L^2(\mathbf{R}^d).$$

$(Exp)_k$ is satisfied (at least) in the following cases:

- $\sigma = d = 1$ (L^2 -subcritical nonlinearity), and Ω is bounded.
- $\sigma \geq 2/d$, $\lambda \geq 0$, and $\Omega_j = \Omega \leq 0$ is independent of j (isotropic repulsive potential).

Remark 1.13 (Optimality). When the potential is repulsive and time-independent ($\Omega = -1$ typically), the exponential growth is sharp, and C does depend on k ($C = k$ when $\Omega = -1$), as discussed above.

Remark 1.14. We prove that in the case of an isotropic repulsive potential, there is scattering provided $\sigma \geq 2/d$ (Proposition 6.5): morally, $(Exp)_k$ is satisfied because it is satisfied in the linear setting (case $\lambda = 0$). However, this property on the linear solution demands a justification; the key is Lemma 6.3.

1.4. Outline of the paper. In Section 2, we prove Theorem 1.7. We then focus on the study of Equation (1.2). In Section 3, we derive a generalized Mehler formula to express the fundamental solution associated to the linear equation, (1.2) with $\lambda = 0$. In Section 4, we generalize a lens transform, known in the case of isotropic time-independent quadratic potentials, to the case of isotropic time-dependent quadratic potentials. This allows us to infer Proposition 1.9. In Section 5, we introduce some vector fields, which correspond to Heisenberg derivatives, and yield interesting evolution laws when the potential is isotropic. In Section 6, we examine the large time behavior of solutions to (1.2), and prove Proposition 1.12. Finally, we show in an appendix that when $V = 0$, for large time, the solutions to (1.1) have bounded Sobolev norms and algebraically growing momenta, provided there is scattering.

2. GLOBAL EXISTENCE IN Σ : PROOF OF THEOREM 1.7

2.1. Strichartz estimates. We first recall some results established in [20, 21]. Consider V satisfying Assumption 1.3. It is established in [20] that one can define $U(t, s)$ as $u(t, x) = U(t, s)\varphi(x)$, where

$$(2.1) \quad i\partial_t u + \frac{1}{2}\Delta u = V(t, x)u \quad ; \quad u(s, x) = \varphi(x),$$

along with the following properties:

- $U(t, t) = \text{Id}$.
- The map $(t, s) \mapsto U(t, s)$ is strongly continuous.
- $U(t, s)^* = U(t, s)^{-1}$.
- $U(t, \tau)U(\tau, s) = U(t, s)$.
- $U(t, s)$ is unitary on L^2 : $\|U(t, s)\varphi\|_{L^2} = \|\varphi\|_{L^2}$.

In addition, we know from [21] that for all $T > 0$, $t, s \in [-T, T]$,

$$(2.2) \quad \|U(t, 0)U(s, 0)^*\varphi\|_{L^\infty(\mathbf{R}^d)} = \|U(t, s)\varphi\|_{L^\infty(\mathbf{R}^d)} \leq \frac{C}{|t-s|^{d/2}} \|\varphi\|_{L^1(\mathbf{R}^d)},$$

provided that $|t-s| < \eta$. It is implicitly assumed in [20] that η may depend on T ; in Example 2.4 below, we show that it does indeed, in the case of the time dependent harmonic potential, if the functions Ω_j are not bounded.

Recall the standard definition in the context of Schrödinger equations:

Definition 2.1. A pair (p, q) is admissible if $2 \leq q < \frac{2d}{d-2}$ ($2 \leq q \leq \infty$ if $d = 1$, $2 \leq q < \infty$ if $d = 2$) and

$$\frac{2}{p} = \delta(q) := d \left(\frac{1}{2} - \frac{1}{q} \right).$$

The general results on Strichartz estimates (see e.g. [26]) then yield, as a consequence of the dispersive estimate (2.2):

Proposition 2.2. *Recall that $U(t, s)$ is defined by (2.1), where V satisfies Assumption 1.3. Let $T > 0$. There exists $\eta > 0$ such that the following holds:*

(1) *For any admissible pair (p, q) , there exists C_q such that*

$$\|U(\cdot, s)\varphi\|_{L^p([s, s+\eta]; L^q)} \leq C_q \|\varphi\|_{L^2}, \quad \forall \varphi \in L^2(\mathbf{R}^d), \quad \forall s \in]-T, T - \eta[.$$

(2) *For $s \in \mathbf{R}$, denote*

$$D_s(F)(t, x) = \int_s^t U(t, \tau) F(\tau, x) d\tau.$$

For all admissible pairs (p_1, q_1) and (p_2, q_2) , there exists $C = C_{q_1, q_2}$ independent of $s \in]-T, T - \eta[$ such that

$$(2.3) \quad \|D_s(F)\|_{L^{p_1}([s, s+\delta]; L^{q_1})} \leq C \|F\|_{L^{p_2'}([s, s+\delta]; L^{q_2'})},$$

for all $F \in L^{p_2'}(I; L^{q_2'})$ and $0 \leq \delta \leq \eta$.

Example 2.3 (Standard harmonic oscillator). Assume that $V(t, x) = \frac{|x|^2}{2}$. The above result is then standard (see e.g. [13]). The fact that one has to consider finite time intervals for the above result to be valid stems for instance from the existence of eigenvalues for the harmonic oscillator: let $g(x) = e^{-|x|^2/2}$ be the ground state associated to the harmonic potential, and denote $u(t, x) = e^{-itd/2}g(x)$. It solves

$$i\partial_t u + \frac{1}{2}\Delta u = \frac{|x|^2}{2}u \quad ; \quad u|_{t=0} = g.$$

We compute $\|u\|_{L^p(I; L^q)} = |I|^{1/p} \|g\|_{L^q}$, which shows that Proposition 2.2 becomes false with $\eta = \infty$.

Example 2.4. We show that in general, the above result is false with $T = \infty$. Let

$$V(t, x) = \frac{1}{2}\Omega(t)|x|^2.$$

If Ω is not bounded, then the above uniformity with respect to s fails: let

$$\Omega(t) = n^2 \text{ if } 4n + 1 = t_n \leq t \leq 4n + 2.$$

Since we have

$$\left(-\frac{1}{2}\Delta + \frac{n^2}{2}|x|^2\right) e^{-n|x|^2/2} = \frac{nd}{2} e^{-n|x|^2/2},$$

the function $u(t, x) = e^{-ind(t-t_n)/2 - n|x|^2/2}$ solves (2.1) with $s = t_n$. If Proposition 2.2 was true with $T = \infty$, we would have:

$$\|u\|_{L^p([4n+1, 4n+1+\eta]; L^q)} = \eta^{1/p} \left(\frac{2\pi}{nq}\right)^{d/(2q)} \leq C \|u(t_n)\|_{L^2} = C \left(\frac{\pi}{n}\right)^{d/4},$$

where C does not depend on n . For all $q > 2$, letting n go to infinity leads to a contradiction. Since (2.2) implies Proposition 2.2, this shows (2.2) is valid for $|t - s| < \eta$ where η depends on T , unless Ω is bounded.

2.2. Local existence in Σ . Since (1.1) is not autonomous, we consider the same problem with a varying initial time:

$$(2.4) \quad i\partial_t u + \frac{1}{2}\Delta u = V(t, x)u + \lambda|u|^{2\sigma}u \quad ; \quad u|_{t=s} = u_0,$$

with $s \in \mathbf{R}$.

Proposition 2.5. *Let $\lambda \in \mathbf{R}$, $\sigma > 0$ with $\sigma < 2/(d-2)$ if $d \geq 3$, and V satisfying Assumption 1.3. Let $M > 0$, and $s \in]-M, M[$. For all $u_0 \in \Sigma$, there exists $T = T(\|u_0\|_\Sigma, M)$ and a unique solution u solution to (2.4), such that*

$$u, \nabla u, xu \in C(]s-T, s+T[; L^2(\mathbf{R}^d)) \cap L^{\frac{4\sigma+4}{d\sigma}}_{\text{loc}}(]s-T, s+T[; L^{2\sigma+2}(\mathbf{R}^d)).$$

Moreover, its L^2 -norm is independent of time:

$$\|u(t)\|_{L^2(\mathbf{R}^d)} = \|u_0\|_{L^2(\mathbf{R}^d)}, \quad \forall t \in]s-T, s+T[.$$

If V is C^1 in t , then the energy E (defined by (1.4)) evolves according to (1.5).

Sketch of the proof. We present here only the main steps of the classical argument. Duhamel's formulation for (2.4) reads

$$u(t) = U(t, s)u_0 - i\lambda \int_s^t U(t, \tau) (|u|^{2\sigma}u)(\tau) d\tau.$$

Denote the right hand side by $\Phi^s(u)(t)$. Proposition 2.5 follows from a fixed point argument in the space

$$X_T = \left\{ u \in C(I_T; \Sigma) ; u, xu, \nabla u \in L^{\frac{4\sigma+4}{d\sigma}}(I_T; L^{2\sigma+2}(\mathbf{R}^d)) \right\},$$

where $I_T =]s-T, s+T[$. Introduce the Lebesgue exponents

$$q = 2\sigma + 2 \quad ; \quad p = \frac{4\sigma + 4}{d\sigma} \quad ; \quad \theta = \frac{2\sigma(2\sigma + 2)}{2 - (d-2)\sigma}.$$

Then (p, q) is admissible, and

$$\frac{1}{q'} = \frac{2\sigma}{q} + \frac{1}{q} \quad ; \quad \frac{1}{p'} = \frac{2\sigma}{\theta} + \frac{1}{p}.$$

Proposition 2.2 and Hölder inequality yield

$$\begin{aligned} \|\Phi^s(u)\|_{L^p(I_T; L^q) \cap L^\infty(I_T; L^2)} &\leq C\|u_0\|_{L^2} + C\||u|^{2\sigma}u\|_{L^{p'}(I_T; L^{q'})} \\ &\leq C\|u_0\|_{L^2} + C\|u\|_{L^\theta(I_T; L^q)}^{2\sigma} \|u\|_{L^p(I_T; L^q)}, \end{aligned}$$

where C is independent of $s \in [-M, M]$ and $T \leq \eta$. Using Sobolev embedding,

$$\|\Phi^s(u)\|_{L^p(I_T; L^q) \cap L^\infty(I_T; L^2)} \leq C\|u_0\|_{L^2} + CT^{2\sigma/\theta} \|u\|_{L^\infty(I_T; H^1)}^{2\sigma} \|u\|_{L^p(I_T; L^q)}.$$

We have

$$\begin{aligned} \nabla \Phi^s(u)(t) &= U(t, s)\nabla u_0 - i\lambda \int_s^t U(t, \tau)\nabla (|u|^{2\sigma}u)(\tau) d\tau \\ &\quad - i \int_s^t U(t, \tau) (\Phi^s(u)(\tau)\nabla V(\tau)) d\tau. \end{aligned}$$

We estimate the second term of the right hand side as above, and get, since ∇V is sublinear by assumption:

$$\begin{aligned} \|\nabla\Phi^s(u)\|_{L^p(I_T;L^q)\cap L^\infty(I_T;L^2)} &\leq C\|\nabla u_0\|_{L^2} + CT^{2\sigma/\theta}\|u\|_{L^\infty(I_T;H^1)}^{2\sigma}\|\nabla u\|_{L^p(I_T;L^q)} \\ &\quad + C\|\Phi^s(u)\nabla V\|_{L^1(I_T;L^2)} \\ &\leq C\|\nabla u_0\|_{L^2} + CT^{2\sigma/\theta}\|u\|_{L^\infty(I_T;H^1)}^{2\sigma}\|\nabla u\|_{L^p(I_T;L^q)} \\ &\quad + CT\|x\Phi^s(u)\|_{L^\infty(I_T;L^2)} + CT\|\Phi^s(u)\|_{L^\infty(I_T;L^2)}, \end{aligned}$$

where, again, C does not depend on $s \in [-M, M]$. We have similarly

$$\begin{aligned} \|x\Phi^s(u)\|_{L^p(I_T;L^q)\cap L^\infty(I_T;L^2)} &\leq C\|xu_0\|_{L^2} + CT^{2\sigma/\theta}\|u\|_{L^\infty(I_T;H^1)}^{2\sigma}\|xu\|_{L^p(I_T;L^q)} \\ &\quad + CT\|\nabla\Phi^s(u)\|_{L^\infty(I_T;L^2)}. \end{aligned}$$

Choosing T sufficiently small, one can then prove that Φ^s maps a suitable ball in X_T into itself. Contraction for the norm $\|\cdot\|_{L^p(I_T;L^q)}$ is proved similarly, and one concludes by remarking that X_T equipped with this norm is complete. \square

We can now infer the analogue to the standard result (which is not straightforward since we consider a non-autonomous equation, in the presence of an external potential):

Corollary 2.6. *Let $\lambda \in \mathbf{R}$, $\sigma > 0$ with $\sigma < 2/(d-2)$ if $d \geq 3$, V satisfying Assumption 1.3, and $u_0 \in \Sigma$. Either the solution to (1.2) is global in time (in the future),*

$$u, \nabla u, xu \in C(\mathbf{R}_+; L^2(\mathbf{R}^d)) \cap L_{\text{loc}}^{\frac{4\sigma+4}{d\sigma}}(\mathbf{R}_+; L^{2\sigma+2}(\mathbf{R}^d)),$$

or there exists $T > 0$, such that

$$\|\nabla u(t)\|_{L^2} \xrightarrow[t \rightarrow T]{} +\infty.$$

Proof. Let $M > 0$. Proposition 2.5 shows that the only obstruction to well-posedness on $[0, M]$ is the existence of a time $0 < T < M$ such that

$$\|xu(t)\|_{L^2} + \|\nabla u(t)\|_{L^2} \xrightarrow[t \rightarrow T]{} +\infty.$$

So long as $u \in C([0, t]; \Sigma)$, we have (see e.g. [13] for the arguments that make the computation rigorous)

$$(2.5) \quad \frac{d}{dt} \int_{\mathbf{R}^d} x_j^2 |u(t, x)|^2 dx = 2 \operatorname{Im} \int_{\mathbf{R}^d} x_j \bar{u}(t, x) \partial_j u(t, x) dx.$$

Suppose $u \in L^\infty([0, T]; H^1)$. Then the above formula, Cauchy–Schwarz inequality and Gronwall lemma show that $xu \in L^\infty([0, T]; L^2)$, hence a contradiction. The corollary follows, since $M > 0$ is arbitrary. \square

Therefore, to prove global existence in Σ in the H^1 -subcritical case, it suffices to exhibit *a priori* bounds for ∇u in L^2 .

2.3. L^2 -subcritical case. In the case $\sigma < 2/d$, recall that the classical argument of [37] can be applied directly, to infer Proposition 1.4. The *a priori* bound for $(\nabla u, xu)$ in L^2 then follows, by resuming the computations presented in the proof of Proposition 2.5. Keeping the same notations, we have in particular

$$\begin{aligned} & \|\nabla\Phi^s(u)\|_{L^p(I_T;L^q)\cap L^\infty(I_T;L^2)} + \|x\Phi^s(u)\|_{L^p(I_T;L^q)\cap L^\infty(I_T;L^2)} \\ & \leq C\|u_0\|_\Sigma + C\|u\|_{L^\theta(I_T;L^q)}^{2\sigma} (\|\nabla u\|_{L^p(I_T;L^q)} + \|xu\|_{L^p(I_T;L^q)}) \\ & \quad + CT (\|\Phi^s(u)\|_{L^\infty(I_T;L^2)} + \|x\Phi^s(u)\|_{L^\infty(I_T;L^2)} + \|\nabla\Phi^s(u)\|_{L^\infty(I_T;L^2)}), \end{aligned}$$

where we recall that

$$q = 2\sigma + 2 \quad ; \quad p = \frac{4\sigma + 4}{d\sigma} \quad ; \quad \theta = \frac{2\sigma(2\sigma + 2)}{2 - (d - 2)\sigma},$$

and, in view of Proposition 2.5, we know that $u = \Phi^s(u)$. In the case $\sigma < 2/d$, we have $1/p < 1/\theta$, and thus

$$\|u\|_{L^\theta(I_T;L^q)} \leq (2T)^{1/\theta - 1/p} \|u\|_{L^p(I_T;L^q)} = (2T)^{\frac{(2-d\sigma)(\sigma+1)}{2\sigma(2\sigma+2)}} \|u\|_{L^p(I_T;L^q)}.$$

By Proposition 1.4, $u \in L^p_{\text{loc}}(\mathbf{R}; L^q(\mathbf{R}^d))$. Splitting any given time interval $[-M, M]$ into finitely many (tiny) pieces, we obtain an *a priori* bound for $(\nabla u, xu)$ in $L^\infty([-M, M]; L^2)$. Since $M > 0$ is arbitrary, Corollary 2.6 yields the first point of Theorem 1.7.

2.4. Defocusing energy-subcritical case. We now consider the case $\lambda \geq 0$, with $\sigma < 2/(d - 2)$ if $d \geq 3$. To complete the proof of Theorem 1.7, we resume the computation initiated in the proof of Corollary 2.6, in order to infer a virial identity:

Lemma 2.7. *Let $\lambda \in \mathbf{R}$, $\sigma > 0$ with $\sigma < 2/(d - 2)$ if $d \geq 3$, and V satisfying Assumption 1.3. Let $u_0 \in \Sigma$, and $u \in C([-T, T]; \Sigma)$ be the solution to (1.1) given by Proposition 2.5 (case $s = 0$). Denote*

$$y(t) = \int_{\mathbf{R}^d} |x|^2 |u(t, x)|^2 dx.$$

Then $y \in C^2([-T, T])$, and satisfies

$$\frac{d^2 y}{dt^2} = 2\|\nabla u(t)\|_{L^2}^2 - 2 \int_{\mathbf{R}^d} x \cdot \nabla V(t, x) |u(t, x)|^2 dx + 2\lambda \frac{d\sigma}{\sigma + 1} \|u(t)\|_{L^{2\sigma+2}}^{2\sigma+2}.$$

Proof. We present the formal part of the proof, and refer to [13] for the arguments that make the proof rigorous. We first resume the computation made in the course of the proof of Corollary 2.6. Differentiating (2.5) with respect to time again, we have:

$$\begin{aligned} \frac{d^2}{dt^2} \|x_j u\|_{L^2}^2 &= 2 \operatorname{Im} \int_{\mathbf{R}^d} x_j \partial_t \bar{u} \partial_j u + 2 \operatorname{Im} \int_{\mathbf{R}^d} x_j \bar{u} \partial_j \partial_t u \\ &= -2 \operatorname{Im} \int_{\mathbf{R}^d} (\bar{u} + 2x_j \partial_j \bar{u}) \partial_t u = 2 \operatorname{Re} \int_{\mathbf{R}^d} (\bar{u} + 2x_j \partial_j \bar{u}) i \partial_t u \\ &= 2 \operatorname{Re} \int_{\mathbf{R}^d} (\bar{u} + 2x_j \partial_j \bar{u}) \left(-\frac{1}{2} \Delta u + V(t, x)u + \lambda |u|^{2\sigma} u \right) \end{aligned}$$

The terms in factor of \bar{u} simplify easily, and we infer:

$$\begin{aligned}
\frac{d^2}{dt^2} \|x_j u\|_{L^2}^2 &= \|\nabla u\|_{L^2}^2 + 2 \int_{\mathbf{R}^d} V(t, x) |u(t, x)|^2 dx + 2\lambda \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \\
&\quad - 2 \operatorname{Re} \int_{\mathbf{R}^d} x_j \partial_j \bar{u} \Delta u + 4 \operatorname{Re} \int_{\mathbf{R}^d} V(t, x) x_j u \partial_j \bar{u} \\
&\quad + 4\lambda \operatorname{Re} \int_{\mathbf{R}^d} x_j |u|^{2\sigma} u \partial_j \bar{u} \\
&= \|\nabla u\|_{L^2}^2 + 2 \int_{\mathbf{R}^d} V(t, x) |u(t, x)|^2 dx + 2\lambda \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \\
&\quad - \|\nabla u\|_{L^2}^2 + 2\|\partial_j u\|_{L^2}^2 + 2 \int_{\mathbf{R}^d} x_j V(t, x) \partial_j (|u|^2) \\
&\quad - \frac{2\lambda}{\sigma+1} \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \\
&= 2\|\partial_j u\|_{L^2}^2 - 2 \int_{\mathbf{R}^d} x_j \partial_j V(t, x) |u(t, x)|^2 dx + 2\lambda \frac{\sigma}{\sigma+1} \|u\|_{L^{2\sigma+2}}^{2\sigma+2}.
\end{aligned}$$

The result then follows by summing over j . \square

To complete the proof of Theorem 1.7, fix $M > 0$, and for $t \in [0, M]$, let

$$f(t) = y(t) + |\dot{y}(t)|.$$

We have

$$\dot{f}(t) \leq |\dot{y}(t)| + |\ddot{y}(t)| \leq |\dot{y}(t)| + 2\|\nabla u\|_{L^2}^2 + C + Cy(t) + C\|u\|_{L^{2\sigma+2}}^{2\sigma+2},$$

where we have used Lemma 2.7, the estimate

$$|x \cdot \nabla V(t, x)| \leq C(M) (1 + |x|^2), \quad \forall (t, x) \in [0, M] \times \mathbf{R}^d,$$

and the conservation of mass. Since $u_0 \in \Sigma$, (1.4)–(1.5) (this is where we have to assume that V is C^1 in t) yield

$$\|\nabla u\|_{L^2}^2 + \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \lesssim 1 + y(t) + \sup_{0 \leq s \leq t} y(s) \lesssim 1 + \sup_{0 \leq s \leq t} y(s).$$

Gronwall lemma implies $f \in L^\infty([0, M])$. We infer $y \in L^\infty_{\text{loc}}(\mathbf{R})$. With the above inequality, this implies $\nabla u \in L^\infty_{\text{loc}}(\mathbf{R}; L^2)$, and Theorem 1.7 then follows from Corollary 2.6.

3. GENERALIZED MEHLER FORMULA

In the rest of this paper, we consider the case where V is exactly quadratic in x , and study some properties associated to (1.2).

3.1. The formula. Classically, Mehler's formula refers to the explicit formula for the fundamental solution of the linear equation

$$(3.1) \quad i\partial_t u_{\text{lin}} + \frac{1}{2} \Delta u_{\text{lin}} = \frac{1}{2} \sum_{j=1}^d \Omega_j(t) x_j^2 u_{\text{lin}} \quad ; \quad u_{\text{lin}}|_{t=0} = u_0,$$

in the case $\dot{\Omega}_j = 0$, with $\Omega_j > 0$. See e.g. [19]. It was generalized (still with $\dot{\Omega}_j = 0$) in [23] to a framework where typically, $\Omega_j \in \mathbf{R}$ has no specified sign.

The case of time dependent harmonic potentials with $d = 1$ was considered in [16], along with other terms corresponding for instance to time dependent magnetic and electric fields. Since the case $d \geq 1$ for (3.1) follows by taking the tensor product of the one dimensional case, we shall simply rewrite the results of [16] (and adapt them to our conventions).

Seek formally the solution to (3.1) as

$$(3.2) \quad u_{\text{lin}}(t, x) = \left(\prod_{j=1}^d \frac{1}{2i\pi\mu_j(t)} \right)^{1/2} \int_{\mathbf{R}^d} e^{\frac{i}{2}\phi(t, x, y)} u_0(y) dy,$$

where

$$\phi(t, x, y) = \sum_{j=1}^d (\alpha_j(t)x_j^2 + 2\beta_j(t)x_jy_j + \gamma_j(t)y_j^2 + 2\delta_j(t)x_j + 2\epsilon_j(t)y_j) + \theta(t),$$

and all the functions of time involved in this formula are real-valued. For instance, when $\Omega = 0$, we have $\mu(t) = t$, $\alpha = \beta = \gamma = 1/t$ and $\delta = \epsilon = \theta = 0$: the convergence $u_{\text{lin}}(t) \rightarrow u_0$ as $t \rightarrow 0$ is recovered (at least formally) by applying the stationary phase formula. Note that in view of the results of D. Fujiwara [20, 21], we know that there exists $\eta > 0$ such that for $|t| < \eta$, the solution to (3.1) can be expressed as

$$u_{\text{lin}}(t, x) = \frac{1}{(2i\pi t)^{d/2}} \int_{\mathbf{R}^d} e^{i\varphi(t, x, y)} a(t, x, y) u_0(y) dy,$$

where $a(0, x, y) = 1$, $\partial_x^\alpha \partial_y^\beta a \in L^\infty(\cdot) - \eta, \eta[\times \mathbf{R}^d \times \mathbf{R}^d)$ for all $\alpha, \beta \in \mathbf{N}^d$, and

$$\varphi(t, x, y) = \frac{|x - y|^2}{2t} + t\xi(t, x, y),$$

with $\partial_x^\alpha \partial_y^\beta \xi \in L^\infty(\cdot) - \eta, \eta[\times \mathbf{R}^d \times \mathbf{R}^d)$ as soon as $|\alpha + \beta| \geq 2$.

Applying the differential operator $i\partial_t + \frac{1}{2}\Delta$ to (3.2), and identifying the terms (in x_j^2, x_jy_j, \dots) in (3.1), we find:

$$\begin{aligned} x_j^2 : \quad & \dot{\alpha}_j + \alpha_j^2 + \Omega_j = 0; & x_jy_j : \quad & \dot{\beta}_j + \alpha_j\beta_j = 0. \\ y_j^2 : \quad & \dot{\gamma}_j + \beta_j^2 = 0; & x_j : \quad & \dot{\delta}_j + \alpha_j\delta_j = 0. \\ y_j : \quad & \dot{\epsilon}_j + \beta_j\delta_j = 0; & \text{Im}(\mathbf{C}) : \quad & \dot{\mu}_j = \alpha_j\mu_j. \\ \text{Re}(\mathbf{C}) : \quad & \dot{\theta} + \sum_{j=1}^d \delta_j^2 = 0. \end{aligned}$$

We infer that μ_j is given by

$$(3.3) \quad \ddot{\mu}_j + \Omega_j(t)\mu_j = 0 \quad ; \quad \mu_j(0) = 0, \quad \dot{\mu}_j(0) = 1.$$

We also have

$$\alpha_j = \frac{\dot{\mu}_j}{\mu_j}.$$

Note that as in the standard cases ($\dot{\Omega}_j = 0$), $\alpha_j(t) \sim 1/t$ as $t \rightarrow 0$. For β_j , we have

$$\dot{\beta}_j + \frac{\dot{\mu}_j}{\mu_j}\beta_j = 0, \text{ hence } \beta_j(t) = \frac{C}{\mu_j(t)},$$

and the stationary phase formula (as $t \rightarrow 0$) yields $C = -1$. We also find

$$\gamma_j(t) = \frac{1}{\mu_j(t)\dot{\mu}_j(t)} - \int_0^t \frac{\Omega_j(\tau)}{(\dot{\mu}_j(\tau))^2} d\tau.$$

Since $\delta_j(0) = \epsilon_j(0) = \theta(0) = 0$, we have $\delta_j = \epsilon_j = \theta_j \equiv 0$.

Remark 3.1. The case of the usual harmonic potential ($\Omega_j = 1$) shows that singularities may be present in the fundamental solutions for positive times, corresponding to the zeroes of μ_j ; see e.g. [17, 24, 39, 40].

Remark 3.2. The dispersive properties associated to (3.1) are measured by the μ_j 's. We will see for instance that if $\Omega_j \leq 0$ for all j , then global in time Strichartz estimates are available, as in the case $\Omega_j = 0$.

To summarize, we have:

Lemma 3.3. *Let $d \geq 1$, and $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ be locally Lipschitzean. There exists $T > 0$ such that for $u_0 \in \mathcal{S}(\mathbf{R}^d)$, the solution to (3.1) is given, for $|t| < T$, by:*

$$u_{\text{lin}}(t, x) = \left(\prod_{j=1}^d \frac{1}{2i\pi\mu_j(t)} \right)^{1/2} \int_{\mathbf{R}^d} e^{i\frac{1}{2}\sum_{j=1}^d (\alpha_j(t)x_j^2 + 2\beta_j(t)x_j y_j + \gamma_j(t)y_j^2)} u_0(y) dy,$$

where

$$\begin{aligned} \ddot{\mu}_j + \Omega_j(t)\mu_j &= 0 \quad ; \quad \mu_j(0) = 0, \quad \dot{\mu}_j(0) = 1, \\ \alpha_j &= \frac{\dot{\mu}_j}{\mu_j} \quad ; \quad \beta_j = -\frac{1}{\mu_j} \quad ; \quad \gamma_j(t) = \frac{1}{\mu_j(t)\dot{\mu}_j(t)} - \int_0^t \frac{\Omega_j(\tau)}{(\dot{\mu}_j(\tau))^2} d\tau. \end{aligned}$$

Remark 3.4. The fact that the quadratic potential has no rectangle term is not necessary in order to get such a result. If we consider

$$i\partial_t u_{\text{lin}} + \frac{1}{2}\Delta u_{\text{lin}} = \frac{1}{2} \langle M(t)x, x \rangle u_{\text{lin}},$$

where $M(t)$ is a (time dependent) symmetric matrix, then a similar formula is available. Of course, the formula is more involved, and since it does not really bring new information, we do not carry out the computation here.

3.2. Some consequences. In this paragraph, we assume that the functions Ω_j are bounded. This assumption was discussed in Example 2.4.

As a consequence of the boundedness of Ω_j , we infer a uniform local bound from below for the functions μ_j . It follows from the growth of the functions μ_j 's, which is at most exponential:

Lemma 3.5. *Assume that for all $j \in \{1, \dots, d\}$, $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ is locally Lipschitzean and bounded. For $s \in \mathbf{R}$, define μ_j^s and ν_j^s as the solutions to*

$$(3.4) \quad \ddot{\mu}_j^s + \Omega_j(t)\mu_j^s = 0 \quad ; \quad \mu_j^s(s) = 0, \quad \dot{\mu}_j^s(s) = 1.$$

$$(3.5) \quad \ddot{\nu}_j^s + \Omega_j(t)\nu_j^s = 0 \quad ; \quad \nu_j^s(s) = 1, \quad \dot{\nu}_j^s(s) = 0.$$

There exists $C > 0$ independent of $s \in \mathbf{R}$ such that

$$|\mu_j^s(t)| + |\dot{\mu}_j^s(t)| + |\nu_j^s(t)| + |\dot{\nu}_j^s(t)| \leq C e^{C|t-s|}, \quad \forall t \in \mathbf{R}.$$

Proof. Introduce $f_j^s(t) = |\dot{\mu}_j^s(t)| + |\mu_j^s(t)|$. We have

$$\begin{aligned} \dot{f}_j^s(t) &\leq |\ddot{\mu}_j^s(t)| + |\dot{\mu}_j^s(t)| = |\Omega_j(t)\mu_j^s(t)| + |\dot{\mu}_j^s(t)| \\ &\leq \|\Omega_j\|_{L^\infty} |\mu_j^s(t)| + |\dot{\mu}_j^s(t)| \lesssim f_j^s(t). \end{aligned}$$

Gronwall lemma yields, since $f_j^s(s) = 1$,

$$f_j^s(t) \lesssim e^{C|t-s|},$$

for some $C > 0$, independent of j, s and t . The first part of lemma then follows. The second estimate is similar. \square

In view of the initial data for μ_j^s and ν_j^s , we infer:

Lemma 3.6. *Assume that for all $j \in \{1, \dots, d\}$, $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ is locally Lipschitzean and bounded. There exists $\eta > 0$ such that for all j , and all $s \in \mathbf{R}$,*

$$|\mu_j^s(t)| \geq \frac{|t-s|}{2}, \quad \frac{1}{2} \leq |\nu_j^s(t)| \leq \frac{3}{2}, \quad \forall t, |t-s| < \eta,$$

where μ_j^s and ν_j^s are given by (3.4) and (3.5), respectively.

This yields a uniform local dispersion in (2.2), and we infer a property which will be crucial in the study of the large time behavior of high Sobolev norms:

Proposition 3.7. *Assume that for all j , $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ is locally Lipschitzean and bounded. Then Proposition 2.2 remains valid with $T = \infty$.*

4. GENERALIZED LENS TRANSFORM

4.1. The formula. It was noticed in [25] that in the case of the L^2 -critical nonlinearity ($\sigma = 2/d$), an explicit change of unknown function makes it possible to add or remove an *isotropic* harmonic potential: if v solves

$$(4.1) \quad i\partial_t v + \frac{1}{2}\Delta v = \lambda|v|^{4/d}v \quad ; \quad v|_{t=0} = u_0,$$

where $\lambda \in \mathbf{R}$, then u , given for $|t| < \pi/(2\omega)$ by the lens transform

$$(4.2) \quad u(t, x) = \frac{1}{(\cos(\omega t))^{d/2}} v\left(\frac{\tan(\omega t)}{\omega}, \frac{x}{\cos(\omega t)}\right) e^{-i\frac{\omega}{2}|x|^2 \tan(\omega t)}$$

solves

$$i\partial_t u + \frac{1}{2}\Delta u = \frac{\omega^2}{2}|x|^2 u + \lambda|u|^{4/d}u \quad ; \quad u|_{t=0} = u_0.$$

See also [30, 5, 34]. Note that the change for the time variable is locally invertible, not globally. The case of a repulsive harmonic potential,

$$i\partial_t u + \frac{1}{2}\Delta u = -\frac{\omega^2}{2}|x|^2 u + \lambda|u|^{4/d}u \quad ; \quad u|_{t=0} = u_0,$$

is obtained by replacing ω by $i\omega$: a formula similar to (4.2) follows, where the trigonometric functions are replaced by hyperbolic functions (and the discussion on the time interval becomes different), see [6]. A heuristic way to understand why this approach works only in the case of isotropic potentials is that even though there would be a “natural” candidate to change the space variable in the anisotropic case, there is no satisfactory candidate to change the time variable.

The lens transform can be generalized to the case of (1.2) provided that the potential is isotropic in the sense that $\Omega_j(t) = \Omega(t)$ is independent of j . Seek an extension of (4.2) of the form

$$(4.3) \quad u(t, x) = \frac{1}{b(t)^{d/2}} v \left(\zeta(t), \frac{x}{b(t)} \right) e^{\frac{i}{2} a(t) |x|^2},$$

with a, b, ζ real-valued,

$$(4.4) \quad b(0) = 1 \quad ; \quad a(0) = \zeta(0) = 0.$$

Suppose also that v solves a more general non-autonomous equation

$$(4.5) \quad i\partial_t v + \frac{1}{2} \Delta v = H(t) |v|^{2\sigma} v \quad ; \quad v|_{t=0} = u_0.$$

We want u to solve

$$(4.6) \quad i\partial_t u + \frac{1}{2} \Delta u = \frac{1}{2} \Omega(t) |x|^2 u + h(t) |u|^{2\sigma} u \quad ; \quad u|_{t=0} = u_0.$$

Apply the Schrödinger differential operator to the formula (4.3), and identify the terms so that u solves (4.6). We find:

$$\dot{b} = ab \quad ; \quad \dot{a} + a^2 + \Omega = 0 \quad ; \quad \dot{\zeta} = \frac{1}{b^2} \quad ; \quad b(t)^{d\sigma-2} H(\zeta(t)) = h(t).$$

Introduce the solution to

$$(4.7) \quad \begin{cases} \ddot{\mu} + \Omega(t)\mu = 0 & ; \quad \mu(0) = 0, \quad \dot{\mu}(0) = 1. \\ \ddot{\nu} + \Omega(t)\nu = 0 & ; \quad \nu(0) = 1, \quad \dot{\nu}(0) = 0. \end{cases}$$

Note that since the Wronskian of μ and ν is constant, we have $\dot{\mu}\nu - \mu\dot{\nu} = 1$ for all time. This relation extends the identities $\cos^2 t + \sin^2 t = 1$ and $\cosh^2 t - \sinh^2 t = 1$. In view of (4.4), we infer:

$$a = \frac{\dot{\nu}}{\nu} \quad ; \quad b = \nu \quad ; \quad \zeta = \frac{\mu}{\nu}.$$

Note that ζ is locally invertible, since $\zeta(0) = 0$ and

$$\dot{\zeta} = \frac{1}{b^2} = \frac{1}{\nu^2}, \quad \text{hence } \dot{\zeta}(0) = 1.$$

Therefore, the lens transform is locally invertible. Moreover, since $b(0) = \nu(0) = 1$, we can write, locally in time,

$$H(t) = b(\zeta^{-1}(t))^{2-d\sigma} h(\zeta^{-1}(t)) = \nu \left(\left(\frac{\mu}{\nu} \right)^{-1}(t) \right)^{2-d\sigma} h \left(\left(\frac{\mu}{\nu} \right)^{-1}(t) \right).$$

Proposition 4.1. *Let v solve*

$$i\partial_t v + \frac{1}{2} \Delta v = H(t) |v|^{2\sigma} v \quad ; \quad v|_{t=0} = u_0.$$

Let $\Omega \in C(\mathbf{R}; \mathbf{R})$. There exists $T > 0$ such that the following holds. Define u by

$$u(t, x) = \frac{1}{\nu(t)^{d/2}} v \left(\frac{\mu(t)}{\nu(t)}, \frac{x}{\nu(t)} \right) e^{i \frac{\nu(t)}{\nu(t)} \frac{|x|^2}{2}}, \quad |t| \leq T,$$

where (μ, ν) is given by (4.7). Then for $|t| < \mu(T)/\nu(T)$, u solves

$$i\partial_t u + \frac{1}{2} \Delta u = \frac{1}{2} \Omega(t) |x|^2 u + h(t) |u|^{2\sigma} u \quad ; \quad u|_{t=0} = u_0,$$

where $h(t) = \nu(t)^{d\sigma-2} H(\mu(t)/\nu(t))$.

Remark 4.2. We do not require Ω to be locally Lipschitzian: all we need is the local existence of a C^2 solution to (4.7), so we can rely on Peano existence theorem.

Remark 4.3 (Generalized Avron–Herbst formula). A similar result is available, when the quadratic potential $\Omega(t)|x|^2$ is replaced by a linear (anisotropic) one $E(t) \cdot x$, where $E \in C(\mathbf{R}; \mathbf{R}^d)$. The solutions to

$$\begin{aligned} i\partial_t v + \frac{1}{2}\Delta v &= h(t)|v|^{2\sigma}v \quad ; \quad v|_{t=0} = u_0, \\ i\partial_t u + \frac{1}{2}\Delta u &= E(t) \cdot xu + h(t)|u|^{2\sigma}u \quad ; \quad u|_{t=0} = u_0, \end{aligned}$$

are related by the formula

$$u(t, x) = v \left(t, x + \int_0^t \int_0^\tau E(s) ds d\tau \right) e^{-ix \cdot \int_0^t E(\tau) d\tau - \frac{i}{2} \int_0^t |E(\tau)|^2 d\tau}.$$

4.2. Proof of Proposition 1.9. We assume in this paragraph that the nonlinearity is *focusing*: $\lambda < 0$. By homogeneity, we can assume $\lambda = -1$. It is well known that the equation

$$(4.8) \quad i\partial_t v + \frac{1}{2}\Delta v = -|v|^{4/d}v$$

possesses solutions which blow up in finite time, with different possible rates (see e.g. [3, 13, 29, 33, 35] and references therein).

By adapting Proposition 4.1 to isolate the initial time $t = 0$, we see that the lens transform maps a solution to (4.8) which blows up at time $t = 0$ to a solution to

$$(4.9) \quad i\partial_t u + \frac{1}{2}\Delta u = \frac{1}{2}\Omega(t)|x|^2u - |u|^{4/d}u$$

which blows up at time $t = 0$. Note that the blow-up rate is not altered by the lens transform, since $\nu(t) \approx 1$ and $\mu(t)/\nu(t) \approx t$ as $t \rightarrow 0$.

Typically, consider the (unstable) minimal mass blow-up solution to (4.8):

$$v(t, x) = \frac{1}{t^{d/2}}Q \left(\frac{x}{t} \right) e^{i\frac{|x|^2}{2t} - \frac{i}{t}},$$

where Q is the ground state, defined as the unique positive, radial, solution to

$$-\frac{1}{2}\Delta Q + Q = Q^{1+4/d}.$$

The lens transform yields a corresponding blow-up solution to (4.9) given by

$$u(t, x) = \frac{1}{\mu(t)^{d/2}}Q \left(\frac{x}{\mu(t)} \right) e^{i\frac{\mu(t)}{\mu(t)}\frac{|x|^2}{2} - i\frac{\nu(t)}{\mu(t)}}.$$

To our knowledge, this gives the first example of an explicit blow-up solution in the presence of a time-dependent external potential.

Note that we have considered the explicit case of minimal mass blow-up solutions for convenience. Any blow-up solution for (4.8) gives rise to a blow-up solution for (4.9), *with the same blow-up rate*.

Note also that without extra assumption on Ω , the Sobolev norms of u may have an arbitrary growth as $t \rightarrow \infty$.

Example 4.4. Consider $\mu(t) = \exp(1 - e^t) - \exp(1 - e^{2t})$ (which satisfies $\mu(0) = 0$ and $\dot{\mu}(0) = 1$). Then the growth of Sobolev norms of the function u given by the above formula is given by a double exponential in time, since

$$\|u(t)\|_{H^s} \underset{t \rightarrow +\infty}{\sim} \frac{C_s}{|\mu(t)|^s}.$$

To determine the corresponding function Ω , we compute

$$\ddot{\mu}(t) = (e^{2t} - e^t) (\exp(1 - e^t) - 4 \exp(1 - e^{2t})),$$

and therefore

$$\Omega(t) = \frac{\exp(1 - e^t) - 4 \exp(1 - e^{2t})}{\exp(1 - e^t) - \exp(1 - e^{2t})} (e^t - e^{2t}).$$

We note that $\Omega(t) \sim -e^{2t}$ as $t \rightarrow +\infty$: the harmonic potential is repulsive ($\Omega < 0$), and becomes exponentially stronger as time increases.

5. VECTOR FIELDS

The aim of this paragraph is to show that there exists vector fields which may be useful to study the nonlinear equation (1.2), in the same spirit as in [6, 8]. Consider the solutions to

$$(5.1) \quad \begin{cases} \ddot{\mu}_j + \Omega_j(t)\mu_j = 0 & ; & \mu_j(0) = 0, \quad \dot{\mu}_j(0) = 1. \\ \ddot{\nu}_j + \Omega_j(t)\nu_j = 0 & ; & \nu_j(0) = 1, \quad \dot{\nu}_j(0) = 0. \end{cases}$$

We define

$$A_j = \dot{\mu}_j x_j + i\mu_j \partial_j = i\mu_j e^{i\frac{x_j^2}{2} \frac{\dot{\mu}_j}{\mu_j}} \partial_j \left(e^{-i\frac{x_j^2}{2} \frac{\dot{\mu}_j}{\mu_j}} \right) = i\mu_j e^{i\sum_k \frac{x_k^2}{2} \frac{\dot{\mu}_k}{\mu_k}} \partial_j \left(e^{-i\sum_k \frac{x_k^2}{2} \frac{\dot{\mu}_k}{\mu_k}} \right),$$

$$B_j = \dot{\nu}_j x_j + i\nu_j \partial_j = i\nu_j e^{i\frac{x_j^2}{2} \frac{\dot{\nu}_j}{\nu_j}} \partial_j \left(e^{-i\frac{x_j^2}{2} \frac{\dot{\nu}_j}{\nu_j}} \right) = i\nu_j e^{i\sum_k \frac{x_k^2}{2} \frac{\dot{\nu}_k}{\nu_k}} \partial_j \left(e^{-i\sum_k \frac{x_k^2}{2} \frac{\dot{\nu}_k}{\nu_k}} \right).$$

Note that the last two expressions for A or B show that A and B act on gauge invariant nonlinearities like derivatives: the modulus ignores the multiplication by the exponential.

Example 5.1 ($\dot{\Omega}_j = 0$). When $\Omega_j = 0$, $B_j = i\partial_j$ and $A_j = x_j + it\partial_j$, which are Heisenberg derivatives commonly used in the theory of nonlinear Schrödinger equations (see e.g. [13]). When $\Omega_j = \omega_j^2 > 0$, $A_j = x_j \cos(\omega_j t) + i\frac{\sin(\omega_j t)}{\omega_j} \partial_j$ and $B_j = \omega_j x_j \sin(\omega_j t) + i \cos(\omega_j t) \partial_j$: we recover classical Heisenberg derivatives (see e.g. [18]). In these two cases (as well as in the case $\Omega_j = -\omega_j^2 < 0$), we have

$$A_j = U_V(t)x_j U_V(-t) \quad ; \quad B_j = U_V(t)i\partial_j U_V(-t),$$

where $U_V(t) = \exp(-it(-\frac{1}{2}\Delta + V(x)))$, $V(x) = \sum_{k=1}^d \Omega_k x_k^2$.

More generally, consider $\dot{\eta}_j x_j + i\eta_j \partial_j$: we check that this operator commute with the linear operator

$$i\partial_t + \frac{1}{2}\Delta - \frac{1}{2} \sum_{k=1}^d \Omega_k(t)x_k^2$$

if and only if η_j satisfies $\dot{\eta}_j + \Omega_j \eta_j = 0$:

$$\begin{aligned} \left[i\partial_t + \frac{1}{2}\Delta - \frac{1}{2} \sum_{k=1}^d \Omega_k(t) x_k^2, \eta_j x_j + i\eta_j \partial_j \right] &= \left[i\partial_t + \frac{1}{2}\partial_j^2 - \frac{1}{2}\Omega_j(t) x_j^2, \eta_j x_j + i\eta_j \partial_j \right] \\ &= i\dot{\eta}_j x_j - \dot{\eta}_j \partial_j + \dot{\eta}_j \partial_j + i\eta_j \Omega_j x_j. \end{aligned}$$

This is zero if (and only if) $\dot{\eta}_j + \Omega_j \eta_j = 0$.

Remark 5.2. This computation could be extended to the case where the center of the harmonic potential depends on time:

$$i\partial_t u + \frac{1}{2}\Delta u = \frac{1}{2} \sum_{k=1}^d \Omega_k(t) (x_k - c_k(t))^2 u.$$

Replacing $\dot{\eta}_j x_j + i\eta_j \partial_j$ with $\dot{\eta}_j (x_j - y_j(t)) + i\eta_j \partial_j$, we can repeat the above computation, and check that the two operators commute if and only if $\dot{\eta}_j + \Omega_j \eta_j = 0$ and $\dot{\eta}_j y_j + \dot{\eta}_j \dot{y}_j + \eta_j \Omega_j c_j = 0$. We choose not to investigate this case further into details here.

To show that the Σ -norm of u is related to the L^2 -norms of $A_j u$ and $B_j u$, write

$$\begin{pmatrix} A_j \\ B_j \end{pmatrix} = M_j \begin{pmatrix} x_j \\ i\partial_j \end{pmatrix}, \quad \text{where } M_j = \begin{pmatrix} \dot{\mu}_j & \mu_j \\ \dot{\nu}_j & \nu_j \end{pmatrix}.$$

We note that the determinant of M_j is the Wronskian of μ_j and ν_j :

$$\det M_j = \nu_j \dot{\mu}_j - \mu_j \dot{\nu}_j \equiv 1.$$

Therefore

$$(5.2) \quad \begin{pmatrix} x_j \\ i\partial_j \end{pmatrix} = \begin{pmatrix} \nu_j & -\mu_j \\ -\dot{\nu}_j & \dot{\mu}_j \end{pmatrix} \begin{pmatrix} A_j \\ B_j \end{pmatrix}.$$

We shall use these vector fields in the isotropic case, where they provide *a priori* estimates:

$$(5.3) \quad i\partial_t u + \frac{1}{2}\Delta u = \frac{1}{2}\Omega(t)|x|^2 u + \lambda|u|^{2\sigma} u.$$

Since A commutes with the linear part of (5.3) and acts on gauge invariant nonlinearities like a gradient, we have readily

$$\frac{1}{2} \frac{d}{dt} \|Au\|_{L^2}^2 = \lambda \sigma \operatorname{Im} \int_{\mathbf{R}^d} |u|^{2\sigma-2} u^2 (\overline{Au})^2.$$

Expanding $(Au)^2$, we obtain eventually:

$$(5.4) \quad \frac{d}{dt} \left(\frac{1}{2} \|Au\|_{L^2}^2 + \frac{\lambda \mu^2}{\sigma+1} \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \right) = \frac{\lambda}{\sigma+1} \mu \dot{\mu} (2-d\sigma) \|u\|_{L^{2\sigma+2}}^{2\sigma+2},$$

$$(5.5) \quad \frac{d}{dt} \left(\frac{1}{2} \|Bu\|_{L^2}^2 + \frac{\lambda \nu^2}{\sigma+1} \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \right) = \frac{\lambda}{\sigma+1} \nu \dot{\nu} (2-d\sigma) \|u\|_{L^{2\sigma+2}}^{2\sigma+2}.$$

These evolution laws are the analogue of the pseudo-conformal conservation law (see [8] for the case $\dot{\Omega} = 0$). They will allow us to infer scattering results in the case $\Omega \leq 0$, $\lambda \geq 0$ (§6.3).

6. GROWTH OF HIGHER ORDER SOBOLEV NORMS AND MOMENTA

6.1. The linear case. In this paragraph, we assume $\lambda = 0$. We recall that in general, Mehler's formula is valid only locally in time, since singularities may appear in the fundamental solution; see e.g. [17, 24, 39, 40]. To understand the long time behavior of the solution u_{lin} to (3.1), one may use Egorov Theorem (see e.g. [4]). Since we deal with a time-dependent potential, modifications would be needed in Egorov Theorem, and we rather follow another strategy, to have some estimates in the linear case (instead of an exact asymptotic behavior, as Egorov Theorem would give us). This approach is based on the vector fields introduced in §5.

We remark that since the L^2 -norm of u_{lin} does not depend on time, and since the operator A_j and B_j introduced in §5 commute with Equation (3.1), the L^2 -norm of $A_{j_1} B_{j_2} \dots A_{j_k} u_{\text{lin}}$ is constant, for whichever combination of these vector fields. In view of (5.2), we infer, for $k \in \mathbf{N}$,

$$\| |x|^k u_{\text{lin}}(t) \|_{L^2} + \| u_{\text{lin}}(t) \|_{H^k} \lesssim \sum_{j=1}^d (|\mu_j(t)|^k + |\nu_j(t)|^k).$$

Lemma 3.5 shows that if $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ is locally Lipschitzian and bounded, then the above quantity grows at most exponentially in time. By Proposition 3.7, we conclude that $(Exp)_k$ is satisfied for all k , provided u_0 is sufficiently smooth and localized. We recall that the case $\Omega_j = -1$ shows that the exponential growth may occur, and that in $(Exp)_k$, the constant C must be expected to depend on k ($C = k$ when $\Omega = -1$ is sharp).

6.2. The L^2 -subcritical case.

Lemma 6.1. *Let $\sigma, k \in \mathbf{N}$, with $\sigma \leq 2/d$, $\Omega_j \in C(\mathbf{R}; \mathbf{R})$ be locally Lipschitzian and bounded, and*

$$u_0 \in H^k(\mathbf{R}^d), \text{ with } |x|^k u_0 \in L^2(\mathbf{R}^d).$$

Suppose that there exists $f \in C(\mathbf{R}_+; \mathbf{R}_+)$ with $f(0) = 0$ such that

$$(6.1) \quad \|u\|_{L^\theta([s, s+\tau]; L^q)} \leq f(\tau), \quad \forall s, \tau \in \mathbf{R},$$

where

$$q = 2\sigma + 2 \quad ; \quad \theta = \frac{2\sigma(2\sigma + 2)}{2 - (d - 2)\sigma}.$$

Then the solution to (1.2) satisfies $(Exp)_k$.

Proof. The first step consists in resuming the computations carried out in the proof of Proposition 2.5, in the case $k = 0$. The case $k \geq 1$ will follow by induction (recall that the constant C in the exponential growth must be expected to depend on k).

Case $k = 0$. Let us pretend that the L^2 -norm of u is not conserved, to simplify the induction. Resuming the same numerology as in the proof of Proposition 2.5, Strichartz estimates yield, for all $t \in \mathbf{R}$ and $\tau > 0$,

$$\begin{aligned} \|u\|_{L^p([t, t+\tau]; L^q) \cap L^\infty([t, t+\tau]; L^2)} &\lesssim \|u(t)\|_{L^2} + \|u\|_{L^\theta([t, t+\tau]; L^q)}^{2\sigma} \|u\|_{L^p([t, t+\tau]; L^q)} \\ &\lesssim \|u(t)\|_{L^2} + f(\tau)^{2\sigma} \|u\|_{L^p([t, t+\tau]; L^q)}. \end{aligned}$$

Fix $\tau \ll 1$ once and for all so the last term of the right hand side can be absorbed by the left hand side, up to doubling the estimating constant: at every increment of time of length τ , the L^2 norm is multiplied (at most) by some fixed constant C .

This implies that it grows at most exponentially. Using Strichartz estimates again, we conclude that $(Exp)_0$ is satisfied (and actually, $(Alg)_0$ is also true).

Case $k \geq 1$. For $k \geq 1$, suppose that $(Exp)_{k-1}$ is satisfied. To avoid a lengthy presentation, we denote by w_ℓ the family of combinations of α momenta and β space derivatives of u , with $|\alpha| + |\beta| = \ell$ ($w_0 = u$). We have, rather formally,

$$(6.2) \quad i\partial_t w_k + \frac{1}{2}\Delta w_k = \frac{1}{2} \sum_{j=1}^d \Omega_j(t) x_j^2 w_\ell + \mathbf{V}(u, w_k) + F + L(w_k),$$

where \mathbf{V} is homogeneous of degree 2σ with respect to its first argument, \mathbf{R} -linear with respect to its second argument, F satisfies the pointwise estimate

$$|F| \lesssim \sum_{0 \leq \ell_j \leq k-1} |w_{\ell_1}| \cdots |w_{\ell_{2\sigma+1}}|,$$

where the sums carries over combinations such that in addition $\sum \ell_j = k$ ($F = 0$ in the case $k = 1$), and L is linear with respect to its argument. A word of explanation is needed about L : this term stems from the fact that x and ∇ do not commute with the linear part of the equation. One might argue that we could proceed as in the linear case, and use repeatedly the vector fields A_j and B_j . The problem is that even though A_j and B_j act on gauge invariant nonlinearities like derivatives, this is not so, for instance, for $A_j B_j$ (the phases do not cancel in the factored formula). We might use the operators $A_{j_1} \dots A_{j_k}$ and $B_{j_1} \dots B_{j_k}$, but this does not suffice to recover the momenta and derivatives of u , since ‘‘rectangle’’ terms (like $A_j B_j$) would be needed.

We proceed in the same spirit as in the case $k = 0$:

$$\begin{aligned} \|w_k\|_{L^p([t, t+\tau]; L^q) \cap L^\infty([t, t+\tau]; L^2)} &\lesssim \|w_k(t)\|_{L^2} + \|u\|_{L^\theta([t, t+\tau]; L^q)}^{2\sigma} \|w_k\|_{L^p([t, t+\tau]; L^q)} \\ &+ \sum_{0 \leq \ell_j \leq k-1} \|w_{\ell_1}\|_{L^\theta([t, t+\tau]; L^q)} \cdots \|w_{\ell_{2\sigma}}\|_{L^\theta([t, t+\tau]; L^q)} \|w_{\ell_{2\sigma+1}}\|_{L^p([t, t+\tau]; L^q)} \\ &+ \|L(w_k)\|_{L^1([t, t+\tau]; L^2)} \end{aligned}$$

Fixing $\tau \ll 1$ independent of $t \in \mathbf{R}$, the second term of the right hand side is absorbed by the left hand side. The sum is treated thanks to $(Exp)_{k-1}$. We notice that since $\sigma \leq 2/d$, we have $\theta \leq p$, where we recall that (p, q) is admissible: for $1 \leq j \leq 2\sigma$,

$$\|w_{\ell_j}\|_{L^\theta([t, t+\tau]; L^q)} \leq \tau^{1/\theta-1/p} \|w_{\ell_j}\|_{L^p([t, t+\tau]; L^q)} \lesssim \tau^{1/\theta-1/p} e^{C(t+\tau)},$$

where we have used $(Exp)_{k-1}$. The last term of the sum is estimated similarly.

Finally, the term $L(w_k)$ is handled in thanks to the Gronwall lemma, and $(Exp)_k$ follows. \square

The proof of Proposition 1.12 in the one-dimensional cubic case follows readily. Since this case is L^2 -subcritical, we have $\theta < p$. Using Strichartz estimate, we infer, for $s, \tau \in \mathbf{R}$,

$$\begin{aligned} \|u\|_{L^p([s, s+\tau]; L^q)} &\leq C(p) \left(\|u_0\|_{L^2} + \|u\|_{L^\theta([s, s+\tau]; L^q)}^{2\sigma} \|u\|_{L^p([s, s+\tau]; L^q)} \right) \\ &\leq C(p) \left(\|u_0\|_{L^2} + \tau^{2\sigma(1/\theta-1/p)} \|u\|_{L^p([s, s+\tau]; L^q)}^{2\sigma+1} \right), \end{aligned}$$

where we have used the conservation of mass, and $C(p)$ is independent of s and τ . Choosing τ sufficiently small, a bootstrap argument implies that there exists $C > 0$ such that

$$\|u\|_{L^p([s, s+\tau]; L^q)} \leq C, \quad \forall s \in \mathbf{R}, \quad 0 < \tau \leq \tau_0.$$

Again since $\theta < p$, we conclude that (6.1) is satisfied with $f(\tau) = C\tau^{1/\theta-1/p}$.

Remark 6.2 (Isotropic L^2 -critical case). When $\sigma = 2/d$ and $\Omega_j = \Omega$ is independent of j , the evolution laws (5.4)–(5.5) yield the *a priori* estimate

$$(\mu(t)^2 + \nu(t)^2) \|u(t)\|_{L^{2\sigma+2}}^{2\sigma+2} \leq C,$$

for some C depending on $\|u_0\|_\Sigma$. In the cases where $\dot{\Omega}_j = 0$, we know that there exists $c > 0$ such that

$$(6.3) \quad \mu(t)^2 + \nu(t)^2 \geq c > 0, \quad \forall t \in \mathbf{R}.$$

When $\dot{\Omega}_j \neq 0$, this property remains at least on bounded time intervals, since the Wronskian is conserved: $\dot{\mu}\nu - \mu\dot{\nu} \equiv 1$. It seems natural to conjecture that (6.3) remains true for *bounded* $\Omega \in C^1$. If this is true, then Lemma 6.1 allows us to add this case in Proposition 1.12.

6.3. Isotropic repulsive potential. We assume $\sigma \geq 2/d$, and $\lambda \geq 0$ (defocusing nonlinearity). We show that in the isotropic repulsive case $\Omega_j = \Omega \geq 0$ (a case where the energy E defined in (1.4) is not a positive functional), the evolution laws derived in §5 show us that the nonlinearity is negligible for large time, and there is scattering. In this paragraph, we also assume that Ω is locally Lipschitzian, without systematically recalling this assumption. We start with the straightforward result:

Lemma 6.3. *Assume $\Omega_j(t) \leq 0$ for all $t \geq 0$. Then the solutions to (5.1) satisfy:*

$$\nu_j(t) \geq 1, \quad \mu_j(t) \geq t, \quad \dot{\nu}_j(t) \geq 0, \quad \dot{\mu}_j(t) \geq 1, \quad \forall t \geq 0.$$

Remark 6.4. As a consequence of this lemma, Proposition 2.2 remains valid with $T = \infty$, even if $\Omega \leq 0$ is not bounded.

We can then prove:

Proposition 6.5. *Assume $\Omega_j = \Omega$ is independent of j , with $\Omega(t) \leq 0$ for all $t \geq 0$. Let $2/d \leq \sigma (< 2/(d-2)$ if $d \geq 3$), $\lambda \geq 0$ and $u_0 \in \Sigma$. The solution to (1.2) is global in time, and there is scattering:*

$$\exists u_+ \in \Sigma, \quad \|\Psi(t)(u(t) - U(t, 0)u_+)\|_{L^2} \xrightarrow{t \rightarrow +\infty} 0,$$

for any $\Psi \in \{\text{Id}, A_j, B_j\}$, and where $U(t, 0)$ corresponds to the free evolution (3.1).

Proof. Since $\lambda \geq 0$ and $\sigma \geq 2/d$, (5.4), (5.5) and Lemma 6.3 yield

$$\frac{d}{dt} \left(\frac{1}{2} \|Au\|_{L^2}^2 + \frac{\lambda\mu^2}{\sigma+1} \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \right) \leq 0 \quad ; \quad \frac{d}{dt} \left(\frac{1}{2} \|Bu\|_{L^2}^2 + \frac{\lambda\nu^2}{\sigma+1} \|u\|_{L^{2\sigma+2}}^{2\sigma+2} \right) \leq 0.$$

We infer *a priori* bounds for Au and Bu in L^2 . Duhamel's formula reads, for $s \in \mathbf{R}$:

$$u(t) = U(t, s)u(s) - i\lambda \int_s^t U(t, \tau) (|u|^{2\sigma} u(\tau)) d\tau.$$

For $\Psi \in \{\text{Id}, A, B\}$, apply Ψ to the above formula:

$$\Psi(t)u(t) = U(t, s)\Psi(s)u(s) - i\lambda \int_s^t U(t, \tau)\Psi(\tau) (|u|^{2\sigma} u(\tau)) d\tau,$$

where we have used the fact that Ψ commutes with the linear part of the equation. Since Ψ acts on gauge invariant nonlinearities like a derivative, we have, thanks to Strichartz estimates:

$$\|\Psi u\|_{L^p([s,t];L^q) \cap L^\infty([s,t];L^2)} \lesssim \|\Psi(s)u(s)\|_{L^2} + \|u\|_{L^\theta([0,t];L^q)}^{2\sigma} \|\Psi u\|_{L^p([0,t];L^q)},$$

with the same numerology as in the proof of Proposition 2.5. Since u , Au and Bu belong to $L^\infty(\mathbf{R}; L^2(\mathbf{R}^d))$, we have

$$(6.4) \quad \|u(t)\|_{L^q} = \|u(t)\|_{L^{2\sigma+2}} \lesssim \frac{1}{\langle t \rangle^{d\sigma/(2\sigma+2)}},$$

where we have used the factorization formula for A and B , Gagliardo–Nirenberg inequality, and Lemma 6.3. We infer $u \in L^\theta(\mathbf{R}; L^q)$:

$$\theta \frac{d\sigma}{2\sigma+2} = \frac{2d\sigma^2}{2-(d-2)\sigma} > 1,$$

since $2d\sigma^2 + (d-2)\sigma = 2\sigma(d\sigma-1) + d\sigma > 2$. Dividing \mathbf{R} into a finite number of intervals on which the $L^\theta L^q$ -norm of u is small, we infer $\Psi u \in L^p(\mathbf{R}; L^q)$. Scattering follows easily:

$$U(0,t)u(t) = u_0 - i\lambda \int_0^t U(0,s) (|u|^{2\sigma} u(s)) ds.$$

For $\tilde{\Psi} \in \{\text{Id}, \nabla, x\}$, apply $\tilde{\Psi}$ to the above formula:

$$\begin{aligned} \tilde{\Psi}U(0,t)u(t) &= \tilde{\Psi}u_0 - i\lambda \int_0^t \tilde{\Psi}U(0,s) (|u|^{2\sigma} u(s)) ds \\ &= \Psi u_0 - i\lambda \int_0^t U(0,s) \Psi (|u|^{2\sigma} u(s)) ds, \end{aligned}$$

where $\Psi = \text{Id}$ if $\tilde{\Psi} = \text{Id}$, $\Psi = -iB$ if $\tilde{\Psi} = \nabla$, and $\Psi = A$ if $\tilde{\Psi} = x$, respectively. We have

$$\begin{aligned} \|\tilde{\Psi}U(0,t_2)u(t_2) - \tilde{\Psi}U(0,t_1)u(t_1)\|_{L^2} &\lesssim \left\| \int_{t_1}^{t_2} U(0,s) \Psi (|u|^{2\sigma} u(s)) ds \right\|_{L^{p'}([t_1,t_2];L^{q'})} \\ &\lesssim \|u\|_{L^\theta([t_1,t_2];L^q)}^{2\sigma} \|\Psi u\|_{L^p([t_1,t_2];L^q)} \xrightarrow{t_1,t_2 \rightarrow +\infty} 0. \end{aligned}$$

Therefore, $U(0,t)u(t)$ converges to some $u_+ \in \Sigma$, and the proposition follows. \square

This result strongly suggests that the solution to the nonlinear solution has the same behavior as the solution to the linear solution as time goes to infinity. It should therefore not be surprising that $(Exp)_k$ is satisfied in this case. However, the delicate issue is to measure high order Sobolev norms. To do so, we modify the argument of Lemma 6.1. We will use the operators A and B once, and just once in view of the discussion in the proof of Lemma 6.1.

We have seen in the course of the proof of Proposition 6.5 that $u \in L^\theta(\mathbf{R}; L^q)$ and $\Psi u \in L^p(\mathbf{R}; L^q)$ for $\Psi \in \{\text{Id}, A, B\}$. As announced above, we modify the induction argument of Lemma 6.1: we first apply either A or B to (1.2), and then apply a combination of x^α and ∂_x^β . We still denote by w_ℓ the family of combinations of α momenta and β space derivatives, now applied to either Au or Bu , with $|\alpha| + |\beta| = \ell - 1$. Eventually, this will not alter the conclusion, in view of (5.2) and Lemma 3.5. Despite this small change in the definition of w_ℓ , we

still have (6.2) for $k \geq 2$ (the case $k \leq 1$ is of no interest, since we know that $\Psi u \in L^p(\mathbf{R}; L^q) \cap L^\infty(\mathbf{R}; L^2)$ for $\Psi \in \{\text{Id}, A, B\}$). Resume the key estimate for w_k :

$$\begin{aligned} \|w_k\|_{L^p([t, t+\tau]; L^q) \cap L^\infty([t, t+\tau]; L^2)} &\lesssim \|w_k(t)\|_{L^2} + \|u\|_{L^\theta([t, t+\tau]; L^q)}^{2\sigma} \|w_k\|_{L^p([t, t+\tau]; L^q)} \\ &+ \sum_{0 \leq \ell_j \leq k-1} \|w_{\ell_1}\|_{L^\theta([t, t+\tau]; L^q)} \cdots \|w_{\ell_{2\sigma}}\|_{L^\theta([t, t+\tau]; L^q)} \|w_{\ell_{2\sigma+1}}\|_{L^p([t, t+\tau]; L^q)} \\ &+ \|L(w_k)\|_{L^1([t, t+\tau]; L^2)} \end{aligned}$$

Fixing $\tau \ll 1$ independent of $t \in \mathbf{R}$, the second term of the right hand side is absorbed by the left hand side. The only difficulty consists in analyzing the sum. We may assume that $\ell_{2\sigma+1}$ corresponds to the largest value of indices ℓ . For $1 \leq j \leq 2\sigma$, if $\ell_j \leq k-2$, then we simply estimate

$$\begin{aligned} \|w_{\ell_j}\|_{L^\theta([t, t+\tau]; L^q)} &\leq \tau^{1/\theta} \|w_{\ell_j}\|_{L^\infty([t, t+\tau]; L^q)} \lesssim \tau^{1/\theta} \|w_{\ell_j}\|_{L^\infty([t, t+\tau]; H^1)} \\ &\lesssim \tau^{1/\theta} \|w_{\ell_{j+1}}\|_{L^\infty([t, t+\tau]; L^2)} \lesssim e^{C(t+\tau)}, \end{aligned}$$

where we have used $(Exp)_{k-1}$. So the only case we have to examine is when $\ell_{2\sigma+1} = k-1 = \ell_{j_0}$ for some $1 \leq j_0 \leq 2\sigma$. Note that since $\sum \ell_j = k$, this may happen only when $k=2$. In that case, we can assume that the term $w_{\ell_{2\sigma+1}}$ is of the form Au or Bu (a term which is $L^p(\mathbf{R}; L^q)$), and estimate as above

$$\|w_{\ell_{j_0}}\|_{L^\theta([t, t+\tau]; L^q)} \lesssim \tau^{1/\theta} \|w_2\|_{L^\infty([t, t+\tau]; L^2)}.$$

The corresponding term in the sum can therefore be absorbed by the left hand side (like **V**). In the other cases, we estimate $\|w_{\ell_{2\sigma+1}}\|_{L^p([t, t+\tau]; L^q)}$ thanks to $(Exp)_{k-1}$. Having examined all the possibilities, we conclude that $(Exp)_k$ is satisfied. Note that Proposition 6.5 suggests that the large time behavior of higher (weighted) Sobolev norms of u is the same as in the linear case, so the exponential growth is sharp in general.

APPENDIX A. THE CASE WITH NO POTENTIAL

Consider the nonlinear Schrödinger equation without potential:

$$(A.1) \quad i\partial_t v + \frac{1}{2}\Delta v = \lambda|v|^{2\sigma}v \quad ; \quad v|_{t=0} = v_0,$$

with energy-subcritical or energy-critical nonlinearity, $\sigma \leq 2/(d-2)$ if $d \geq 3$.

Lemma A.1. *Let $\sigma \in \mathbf{N}$ with $\sigma \geq 2$ if $d=1$, and $\sigma \leq 2/(d-2)$ if $d \geq 3$. Let*

$$(p_1, q_1) = \left(2\sigma + 2, \frac{2d(\sigma + 1)}{d - 2 + d\sigma} \right).$$

Assume that (A.1) possesses a global solution $v \in L^{p_1}(\mathbf{R}; W^{1, q_1}(\mathbf{R}^d))$. Let $k \in \mathbf{N}$. If $v_0 \in H^k(\mathbf{R}^d)$, then $v \in L^\infty(\mathbf{R}; H^k(\mathbf{R}^d))$, and more generally, $v \in L^{p_0}(\mathbf{R}; W^{k, q_0}(\mathbf{R}^d))$ for all admissible pair (p_0, q_0) .

Remark A.2. The assumption $\sigma \in \mathbf{N}$ is made only to simplify the presentation. The proof could be adapted to the case where the map $z \mapsto |z|^{2\sigma}z$ is C^k .

Remark A.3. The main assumption of the lemma states essentially that asymptotic completeness holds in a suitable space. We could even assume that the nonlinearity is $(2\sigma+1)$ -homogeneous, and not necessarily gauge invariant. However, scattering is known with no size assumption on v_0 in the defocusing gauge invariant case (see

below), hence our choice. Note that in the case $d = 1$, an algebraic control of the growth of Sobolev norms is known, regardless of gauge invariance [32].

Proof. We remark that the pair (p_1, q_1) is admissible, and

$$\frac{1}{p_1'} = \frac{2\sigma + 1}{p_1} \quad ; \quad \frac{1}{q_1'} = \frac{1}{q_1} + \frac{2\sigma}{d\sigma(\sigma + 1)}.$$

We prove the lemma by induction on k . We first prove

$$v \in L^{p_1}(\mathbf{R}; W^{k, q_1}) \cap L^\infty(\mathbf{R}; H^k).$$

We start with $k = 1$: applying ∇ to (A.1), Strichartz estimates on $I = [t_0, t]$ yield

$$\begin{aligned} \|\nabla v\|_{L^\infty(I; L^2) \cap L^{p_1}(I; L^{q_1})} &\leq C \left(\|\nabla v(t_0)\|_{L^2} + \| |v|^{2\sigma} \nabla v \|_{L^{p_1'}(I; L^{q_1'})} \right) \\ &\leq C \left(\|v(t_0)\|_{L^2} + \|v\|_{L^{p_1}(I; L^{d\sigma(\sigma+1)})}^{2\sigma} \|\nabla v\|_{L^{p_1}(I; L^{q_1})} \right), \end{aligned}$$

where we have used Hölder inequality. Notice the embedding

$$W^{1, q_1}(\mathbf{R}^d) \subset L^{d\sigma(\sigma+1)}(\mathbf{R}^d).$$

In view of the assumption of the lemma, this implies $v \in L^{p_1}(\mathbf{R}; L^{d\sigma(\sigma+1)})$. Therefore, we can split \mathbf{R} into finitely many intervals on which $C\|v\|_{L^{p_1}(I; L^{d\sigma(\sigma+1)})}^{2\sigma} \leq 1/2$. On each such interval I , we have

$$\|\nabla v\|_{L^\infty(I; L^2) \cap L^{p_1}(I; L^{q_1})} \leq 2C\|\nabla v(t_0)\|_{L^2}.$$

The conclusion follows in the case $k = 1$.

Assume now that the result is known for $k \geq 1$, and that the nonlinearity is C^{k+1} . Differentiating (A.1) $k + 1$ times with respect to space variable, we find, for $|\alpha| = k + 1$,

$$\left(i\partial_t + \frac{1}{2}\Delta \right) \partial^\alpha v = N_1(v) + N_2(v),$$

with the pointwise controls

$$|N_1(v)| \lesssim |v|^{2\sigma} |\partial^\alpha v| \quad ; \quad |N_2(v)| \lesssim \sum_{|\alpha_j| \leq k} |\partial^{\alpha_1} v| \dots |\partial^{\alpha_{2\sigma+1}} v|.$$

Strichartz estimates on the time interval $I = [t_0, t]$ yield

$$\begin{aligned} \|\partial^\alpha v\|_{L^\infty(I; L^2) \cap L^{p_1}(I; L^{q_1})} &\lesssim \|\partial^\alpha v(t_0)\|_{L^2} + \sum_{j=1,2} \|N_j(v)\|_{L^{p_1'}(I; L^{q_1'})} \\ &\lesssim \|\partial^\alpha v(t_0)\|_{L^2} + \|v\|_{L^{p_1}(I; L^{q_2})}^{2\sigma} \|\partial^\alpha v\|_{L^{p_1}(I; L^{q_2})} \\ &\quad + \sum_J \|\partial^{\alpha_1} v\|_{L^{p_1}(I; L^{q_2})} \dots \|\partial^{\alpha_{2\sigma}} v\|_{L^{p_1}(I; L^{q_2})} \|\partial^{\alpha_{2\sigma+1}} v\|_{L^{p_1}(I; L^{q_1})}, \end{aligned}$$

where we have denoted $q_2 = d\sigma(\sigma + 1)$, and we have used the ordering

$$J = \left\{ |\alpha_1|, \dots, |\alpha_{2\sigma-1}| \leq k - 1 \ ; \ |\alpha_{2\sigma}|, |\alpha_{2\sigma+1}| \leq k \quad ; \quad \sum \alpha_j = \alpha \right\}.$$

Proceeding as in the case $k = 1$, we consider finitely many time intervals on which

$$\begin{aligned} \|\partial^\alpha v\|_{L^\infty(I; L^2) \cap L^{p_1}(I; L^{q_1})} &\lesssim \|\partial^\alpha v(t_0)\|_{L^2} \\ &\quad + \sum_J \|\partial^{\alpha_1} v\|_{L^{p_1}(I; L^{q_2})} \dots \|\partial^{\alpha_{2\sigma}} v\|_{L^{p_1}(I; L^{q_2})} \|\partial^{\alpha_{2\sigma+1}} v\|_{L^{p_1}(I; L^{q_1})}. \end{aligned}$$

We use the embedding $W^{1,q_1} \subset L^{q_2}$ again, and the induction assumption: when $\alpha_{2\sigma} = k$, we proceed like for the term N_1 (when summing over all α 's such that $|\alpha| = k+1$), and in all the other cases, we deal with a controllable source term. This yields the lemma for the pair $(p_0, q_0) = (\infty, 2)$. The estimate for general admissible pairs follows by using Strichartz estimates again. \square

Proposition A.4. *Let $\lambda > 0$, and $\sigma \geq 2/d$ an integer, with $\sigma \leq 2/(d-2)$ if $d \geq 3$. Suppose $v_0 \in \Sigma$. Let $k \in \mathbf{N}$, $k \geq 1$.*

- (i) *If $v_0 \in H^k(\mathbf{R}^d)$, then $v \in L^{p_0}(\mathbf{R}; W^{k,q_0}(\mathbf{R}^d))$ for all admissible pair (p_0, q_0) .*
- (ii) *If in addition $x \mapsto |x|^k v_0 \in L^2(\mathbf{R}^d)$, then $|x|^k v \in C(\mathbf{R}; L^2(\mathbf{R}^d))$ and for all admissible pair (p_0, q_0) ,*

$$\forall \alpha \in \mathbf{N}^d, |\alpha| \leq k, \quad \|x^\alpha v\|_{L^{p_0}([0,t]; L^{q_0})} \lesssim \langle t \rangle^{|\alpha|}.$$

Remark A.5. In the case $\sigma > 2/d$, (i) remains valid without assuming $v_0 \in \Sigma$. The point to notice is that one can apply Lemma A.1 as in the proof below, since the assumptions of the lemma are known to be satisfied thanks to Morawetz estimates, which yield asymptotic completeness in H^1 . In the case $\sigma = 2/d$, this aspect is still an open question.

Proof. Under our assumptions on λ and σ , we know that there exists a unique, global, solution $v \in C^\infty(\mathbf{R}; \Sigma)$ to (A.1), with $v \in L^\infty(\mathbf{R}; H^1)$. The pseudo-conformal conservation law yields

$$\frac{d}{dt} \left(\frac{1}{2} \|J(t)v\|_{L^2}^2 + \frac{\lambda t^2}{\sigma+1} \|v\|_{L^{2\sigma+2}}^{2\sigma+2} \right) = \frac{t\lambda}{\sigma+1} (2-d\sigma) \|v\|_{L^{2\sigma+2}}^{2\sigma+2},$$

where $J(t) = x + it\nabla$. The right hand side is non-positive for $t \geq 0$: this yields an *a priori* bound for $J(t)v$ in $L^\infty(\mathbf{R}; L^2)$. Since

$$(A.2) \quad J(t) = ite^{i\frac{|x|^2}{2t}} \nabla \left(e^{-i\frac{|x|^2}{2t}} \cdot \right),$$

Gagliardo–Nirenberg inequality yields

$$\|v\|_{L^\rho} \lesssim \frac{1}{|t|^\delta} \|v\|_{L^2}^{1-\delta} \|J(t)v\|_{L^2}^\delta, \quad \text{where } \delta = d \left(\frac{1}{2} - \frac{1}{\rho} \right), \text{ and } 2 \leq \rho \leq \frac{2d}{d-2}.$$

We infer $v \in L^{p_1}(\mathbf{R}; L^{q_1})$. Resume the value

$$\theta = \frac{4\sigma(\sigma+1)}{(d-2)\sigma-2} \quad \left(\theta = \infty \text{ if } \sigma = \frac{2}{d-2} \right).$$

In view of the identities

$$(p, q) = \left(\frac{4\sigma+4}{d\sigma}, 2\sigma+2 \right) \quad ; \quad \frac{1}{p'} = \frac{1}{p} + \frac{2\sigma}{\theta} \quad ; \quad \frac{1}{q'} = \frac{2\sigma+1}{q},$$

Strichartz estimates yield

$$\|\nabla v\|_{L^p(I; L^q) \cap L^{p_1}(I; L^{q_1})} \lesssim 1 + \|v\|_{L^\theta(I; L^q)}^{2\sigma} \|\nabla v\|_{L^p(I; L^q)}.$$

We note that $v \in L^\theta(\mathbf{R}; L^q)$ (and $\|v(t)\|_{L^q} \rightarrow 0$ uniformly as $t \rightarrow \infty$): splitting \mathbf{R} into finitely many intervals, we infer $\nabla v \in L^{p_1}(\mathbf{R}; L^{q_1})$: the first point of the proposition then follows from Lemma A.1.

The second point of the proposition is obtained by mimicking the proof of Lemma A.1: instead of considering ∇ and its powers, consider $J = x + it\nabla$ and its powers. In view of (A.2), we can follow the same computations, since the nonlinearity we consider is gauge invariant: $|J|^{k\nu} \in L^\infty(\mathbf{R}; L^2)$. The algebraic growth of the momenta then stems from triangle inequality. \square

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