

Evolution of Gaussian measures for the one dimensional Schrödinger equation

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Travail en collaboration avec Nicolas Burq

Motivation : Long time behaviour of the nonlinear Schrödinger equation

$$\begin{cases} i\partial_s U + \partial_y^2 U = |U|^{p-1}U, & (s, y) \in \mathbb{R} \times \mathbb{R}, \\ U(0) = U_0 \in L^2(\mathbb{R}). \end{cases} \quad (NLS_p)$$

- ▶ Almost sure global existence results ($p > 1$)
- ▶ Almost sure scattering results ($p > 3$)
- ▶ Evolution of Gaussian measures

Definition 1

Consider a space X and one parameter group $(\Phi(t, \cdot))_{t \in \mathbb{R}}$ with

$$\Phi(t, \cdot) : X \longrightarrow X.$$

A measure μ defined in the space X is called invariant with respect to $(\Phi(t, \cdot))_{t \in \mathbb{R}}$ if for any μ -measurable set A

$$\mu(\Phi(t, A)) = \mu(A), \quad t \in \mathbb{R}.$$

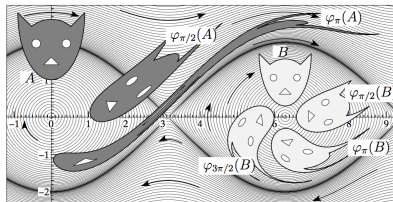


Figure: Conservation of the area for the pendulum.

Theorem 2 (Poincaré)

Let (X, \mathcal{B}, μ) be a probability space and let $\Phi(t, \cdot) : X \rightarrow X$ be a map which preserves the probability measure μ .

(i) Let $A \in \mathcal{B}$ be such that $\mu(A) > 0$, then there exists $k \geq 1$ such that

$$\mu(A \cap \Phi(k, A)) > 0.$$

(ii) Let $B \in \mathcal{B}$ be such that $\mu(B) > 0$, then for μ -almost all $x \in B$, the orbit $(\Phi(n, x))_{n \in \mathbb{N}}$ enters infinitely many times in B .

On invariant measures: the Liouville theorem

Consider the **ordinary differential equation**

$$\begin{cases} \dot{x}(t) = \frac{dx}{dt}(t) = F(x(t)), \\ x(0) = x_0. \end{cases}$$

We denote by $\Phi(t, \cdot)$ the flowmap of this system.

Theorem 3 (Liouville)

The flowmap $\Phi(t, \cdot)$ preserves the measure gdx if and only if

$$\sum_{k=1}^d \frac{\partial}{\partial x_k} (gF_k) = 0.$$

The Schrödinger equation on compact manifolds

► Let M a compact manifold. There exists a Hilbert basis $(h_n)_{n \geq 0}$ of $L^2(M)$, composed of eigenfunctions of Δ_M and we write

$$-\Delta_M h_n = \lambda_n^2 h_n \quad \text{for all } n \geq 0.$$

► Consider a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and let $(g_n)_{n \geq 0}$ be a sequence of independent **complex standard Gaussian variables** $\mathcal{N}_{\mathbb{C}}(0, 1)$.

$$g_n = \frac{1}{\sqrt{2}}(g_{1,n} + ig_{2,n}), \quad g_{1,n}, g_{2,n} \in \mathcal{N}_{\mathbb{R}}(0, 1).$$

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- Let $(\alpha_n)_{n \geq 0}$ and define the probability measure μ via the map

$$\omega \longmapsto \gamma(\omega) = \sum_{n=0}^{+\infty} \alpha_n g_n(\omega) h_n, \quad \mu = \mathbf{p} \circ \gamma^{-1},$$

defined by : for all measurable set A ,

$$\mu(A) = \mathbf{p}(\omega \in \Omega : \gamma(\omega) \in A).$$

The Schrödinger equation on compact manifolds

Proposition 4

The measure μ is invariant under the flow of the equation

$$i\partial_s U + \Delta_M U = 0, \quad (s, y) \in \mathbb{R} \times M.$$

In the nonlinear case, it is then natural to look for :

- ▶ invariant measures : **invariant Gibbs measures.**
- ▶ quasi-invariant measures.

The Schrödinger equation on compact manifolds

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Proof : For all $s \in \mathbb{R}$, the random variable

$$e^{is\Delta_M} \gamma(\omega) = \sum_{n=0}^{+\infty} \alpha_n e^{-is\lambda_n^2} g_n(\omega) h_n$$

has the same distribution as γ since for all $s \in \mathbb{R}$

$$e^{-is\lambda_n^2} g_n(\omega) \sim \mathcal{N}_{\mathbb{C}}(0, 1).$$

Proposition 5

Let $\sigma \in \mathbb{R}$ and consider a probability measure μ on $H^\sigma(\mathbb{R})$. Assume that μ is invariant under the flow Σ_{lin} of equation

$$\begin{cases} i\partial_s U + \partial_y^2 U = 0, & (s, y) \in \mathbb{R} \times \mathbb{R}, \\ U(0, \cdot) = U_0. \end{cases}$$

Then $\mu = \delta_0$.

► Same result for the equation

$$i\partial_s U + \partial_y^2 U = |U|^{p-1} U.$$

The non compact case

Proof :

► Let $\sigma \in \mathbb{R}$ and consider μ an **invariant probability measure** on $H^\sigma(\mathbb{R})$. Let $\chi \in C_0^\infty(\mathbb{R})$. By invariance of the measure,

$$\int_{H^\sigma(\mathbb{R})} \frac{\|\chi u\|_{H^\sigma}}{1 + \|u\|_{H^\sigma}} d\mu(u) = \int_{H^\sigma(\mathbb{R})} \frac{\|\chi \Sigma_{lin}(t)u\|_{H^\sigma}}{1 + \|\Sigma_{lin}(t)u\|_{H^\sigma}} d\mu(u),$$

and by **unitarity of the linear flow** in H^σ , we get

$$= \int_{H^\sigma(\mathbb{R})} \frac{\|\chi \Sigma_{lin}(t)u\|_{H^\sigma}}{1 + \|u\|_{H^\sigma}} d\mu(u). \quad (1)$$

► Assume that the r.h.s. of (1) tends to 0 when $t \rightarrow +\infty$. This implies that $\|\chi u\|_{H^\sigma} = 0$, μ -a.s., and thus $\mu = \delta_0$ since χ is arbitrary.

The non compact case

- ▶ By continuity of the product by χ in H^σ and unitarity of the linear flow in H^σ , we have

$$\frac{\|\chi \Sigma_{lin}(t)u\|_{H^\sigma}}{1 + \|u\|_{H^\sigma}} \leq C \frac{\|\Sigma_{lin}(t)u\|_{H^\sigma}}{1 + \|u\|_{H^\sigma}} = C \frac{\|u\|_{H^\sigma}}{1 + \|u\|_{H^\sigma}} \leq C.$$

- ▶ If $v \in \mathcal{C}_0^\infty(\mathbb{R})$ is smooth, by the Leibniz rule and **dispersion**

$$\|\chi \Sigma_{lin}(t)v\|_{H^\sigma} \leq \|\chi\|_{W^{\sigma,4}} \|\Sigma_{lin}(t)v\|_{W^{\sigma,4}} \leq Ct^{-1/4} \|v\|_{W^{\sigma,4/3}} \longrightarrow 0,$$

when $t \longrightarrow +\infty$.

- ▶ Conclusion with an approximation argument.

- ▶ Define the harmonic oscillator

$$H = -\partial_x^2 + x^2.$$

- ▶ There exists a Hilbert basis $(e_n)_{n \geq 0}$ of $L^2(\mathbb{R})$, composed of eigenfunctions of H and we write

$$He_n = \lambda_n^2 e_n = (2n + 1)e_n \quad \text{for all } n \geq 0.$$

- ▶ We define the harmonic Sobolev space $\mathcal{W}^{\sigma,p}$ by the norm ($\sigma > 0$)

$$\|u\|_{\mathcal{W}^{\sigma,p}} = \|H^{\sigma/2} u\|_{L^p} \equiv \|(-\Delta)^{\sigma/2} u\|_{L^p} + \|\langle x \rangle^\sigma u\|_{L^p}.$$

(Yajima-Zhang)

The Gaussian measure: definition

- Let $\varepsilon > 0$, we define the probability **Gaussian measure** μ_0 on $\mathcal{H}^{-\varepsilon}(\mathbb{R})$ as the distribution of the random variable γ

$$\begin{aligned}\Omega &\longrightarrow \mathcal{H}^{-\varepsilon}(\mathbb{R}) \\ \omega &\longmapsto \gamma(\omega) = \sum_{n=0}^{+\infty} \frac{1}{\lambda_n} g_n(\omega) e_n, \quad \mu_0 = \mathbf{p} \circ \gamma^{-1}.\end{aligned}$$

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- ▶ We denote by $X^0(\mathbb{R}) = \bigcap_{\varepsilon > 0} \mathcal{H}^{-\varepsilon}(\mathbb{R})$. Thus $L^2(\mathbb{R}) \subset X^0(\mathbb{R}) \subset \mathcal{H}^{-\varepsilon}(\mathbb{R})$.
- ▶ The measure μ_0 satisfies $\mu_0(L^2(\mathbb{R})) = 0$ and $\mu_0(X^0(\mathbb{R})) = 1$.

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- ▶ The measure μ_0 satisfies $\mu_0(L^2(\mathbb{R})) = 0$ and $\mu_0(X^0(\mathbb{R})) = 1$.
- ▶ The support of μ_0 is actually smoother

$$\mu_0(\{u_0 \in X^0(\mathbb{R}) : \|e^{-itH} u_0\|_{L^\infty((-\pi, \pi); \mathcal{W}^{1/6-, \infty})} \geq R\}) \leq C e^{-cR^2}.$$

Equivalence of Gaussian measures

- ▶ We say that $\mu \ll \nu$ if $\nu(A) = 0 \implies \mu(A) = 0$.
- ▶ Consider $\alpha_n, \beta_n > 0$ and the measures $\mu = p \circ \gamma^{-1}$ and $\nu = p \circ \psi^{-1}$ with

$$\gamma(\omega) = \sum_{n=0}^{+\infty} \frac{1}{\alpha_n} g_n(\omega) e_n, \quad \psi(\omega) = \sum_{n=0}^{+\infty} \frac{1}{\beta_n} g_n(\omega) e_n.$$

Then the measures μ and ν are **absolutely continuous with respect to each other (same zero measure sets)** if and only if

$$\sum_{n=0}^{+\infty} \left(\frac{\alpha_n}{\beta_n} - 1 \right)^2 < +\infty \quad (\text{Kakutani}).$$

- ▶ If the measures μ and ν are **not** absolutely continuous with respect to each other, they are **mutually singular** (Hajek-Feldman theorem).

The global existence result

Consider

$$\begin{cases} i\partial_s U + \partial_y^2 U = |U|^{p-1} U, & (s, y) \in \mathbb{R} \times \mathbb{R}, \\ U(0) = U_0 \in X^0(\mathbb{R}). \end{cases} \quad (NLS_p)$$

Theorem 6

Let $p > 1$.

► For μ_0 -almost every initial data $U_0 \in X^0(\mathbb{R})$, there exists a unique, global in time, solution $U = \Psi(s, 0)U_0$ to (NLS_p) .

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► The measures $\Psi(s, 0)_\# \mu_0$ and $\Psi_{lin}(s, 0)_\# \mu_0$ are equivalent:

$$\Psi_{lin}(s, 0)_\# \mu_0 \ll \Psi(s, 0)_\# \mu_0 \ll \Psi_{lin}(s, 0)_\# \mu_0.$$

► In the compact case: $\Psi_{lin}(s, 0)_\# \mu_0 = \mu_0$.

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$$\Psi_{lin}(s, 0)_\# \mu_0 \ll \Psi(s, 0)_\# \mu_0 \ll \Psi_{lin}(s, 0)_\# \mu_0.$$

► For all $s' \neq s$, the measures $\Psi(s, 0)_\# \mu_0$ and $\Psi(s', 0)_\# \mu_0$ are mutually singular.

► In the compact case: $\Psi_{lin}(s, 0)_\# \mu_0 = \mu_0$.

The scattering result

Theorem 7

► Assume that $p > 1$. Then for μ_0 -almost every initial data $U_0 \in X^0(\mathbb{R})$, there exists a constant $C > 0$ such that for all $s \in \mathbb{R}$

$$\|\Psi(s, 0)U_0\|_{L^{p+1}(\mathbb{R})} \leq \begin{cases} C \frac{(1+\log\langle s \rangle)^{1/(\rho+1)}}{\langle s \rangle^{\frac{1}{2} - \frac{1}{\rho+1}}} & \text{if } 1 < \rho < 5 \\ \frac{C}{\langle s \rangle^{\frac{1}{2} - \frac{1}{\rho+1}}} & \text{if } \rho \geq 5. \end{cases}$$

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► Assume now that $p > 3$. Then there exist $\eta > 0$ and $W_{\pm} \in L^2(\mathbb{R})$ such that for all $s \in \mathbb{R}$

$$\|\Psi(s, 0)U_0 - e^{is\partial_y^2}(U_0 + W_{\pm})\|_{L^2(\mathbb{R})} \leq C\langle s \rangle^{-\eta}.$$

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► For all $\varphi \in \mathcal{C}_0^\infty(\mathbb{R})$ we have the dispersion bound

$$\|e^{is\partial_y^2}\varphi\|_{L^{p+1}(\mathbb{R})} \leq \frac{C}{|s|^{\frac{1}{2} - \frac{1}{p+1}}} \|\varphi\|_{L^{(p+1)'(\mathbb{R})}}, \quad s \neq 0,$$

therefore, the power decay in s is optimal.

Gross-Pitaevskii equation with random initial conditions and Gibbs measures

- ▶ Thomann 2009
- ▶ Burq–Thomann–Tzvetkov 2013
- ▶ Deng 2012
- ▶ Poiret–Robert–Thomann 2014
- ▶ Burq–Thomann 2023
- ▶ Latocca 2022
- ▶ Robert–Seong–Tolomeo–Wang 2023
- ▶ Dinh–Rougerie 2023
- ▶ Dinh–Rougerie–Tolomeo–Wang 2023

Gross-Pitaevskii equation with noise

- ▶ de Bouard–Debussche–Fukuizumi 2018, 2021 & 2022

Deterministic scattering results for NLS

- ▶ Barab 1984 → never scattering when $p \leq 3$
- ▶ Tsutumi–Yajima 1984 → scattering in L^2 with H^1 data
- ▶ Nakanishi 1999 → scattering in H^σ
- ▶ Dodson 2016 → scattering in L^2 when $p = 5$

Probabilistic scattering results for NLS

- ▶ Burq–Thomann–Tzvetkov 2013 → case $d = 1$ and $p \geq 5$
- ▶ Poiret–Robert–Thomann 2014 → case $d \geq 2$ and $p \geq 3$
- ▶ Killip–Murphy–Visan 2019 → case $d = 4$ in the radial setting
- ▶ Latocca 2022 → case $d \geq 2$ in the radial setting

Quasi-invariance for NLS and dispersive PDEs

- ▶ Tzvetkov 2015
- ▶ Oh-Tsutsumi-Tzvetkov 2019
- ▶ Oh-Tzvetkov 2017 & 2020
- ▶ Planchon-Tzvetkov-Visciglia 2020
- ▶ Debussche-Tsutsumi 2021

The lens transform: compactification in time and space

$$\begin{cases} i\partial_s U + \partial_y^2 U = |U|^{p-1} U, & (s, y) \in \mathbb{R} \times \mathbb{R}, \\ U(0) = U_0. \end{cases} \quad (NLS_p)$$

If $U(s, y)$ is a solution of the problem (NLS_p) , then the function $u(t, x)$ defined for $|t| < \frac{\pi}{4}$ and $x \in \mathbb{R}$ by

$$u(t, x) = \mathcal{L}(U)(t, x) := \frac{1}{\cos^{\frac{1}{2}}(2t)} U\left(\frac{\tan(2t)}{2}, \frac{x}{\cos(2t)}\right) e^{-\frac{ix^2 \tan(2t)}{2}}$$

solves the problem

$$\begin{cases} i\partial_t u - Hu = \cos^{\frac{p-5}{2}}(2t) |u|^{p-1} u, & |t| < \frac{\pi}{4}, x \in \mathbb{R}, \\ u(0, \cdot) = U_0. \end{cases}$$

Introduction

We consider the equation

$$\begin{cases} i\partial_t u - Hu = \cos \frac{p-5}{2}(2t)|u|^{p-1}u, & (t, x) \in \left(-\frac{\pi}{4}, \frac{\pi}{4}\right) \times \mathbb{R}, \\ u(0) = u_0. \end{cases}$$

► The energy

$$\mathcal{E}(t, u(t)) = \frac{1}{2} \|\sqrt{H} u(t)\|_{L^2(\mathbb{R})}^2 + \frac{\cos \frac{p-5}{2}(2t)}{p+1} \|u(t)\|_{L^{p+1}(\mathbb{R})}^{p+1},$$

is **not conserved**.

Introduction

We consider the equation

$$\begin{cases} i\partial_t u - Hu = \cos \frac{\rho-5}{2}(2t)|u|^{p-1}u, & (t, x) \in \left(-\frac{\pi}{4}, \frac{\pi}{4}\right) \times \mathbb{R}, \\ u(0) = u_0. \end{cases}$$

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is **not conserved**.

► For $-\frac{\pi}{4} < t < \frac{\pi}{4}$, we define the measure

$$\begin{aligned} d\nu_t &= e^{-\mathcal{E}(t,u)} du \\ &= e^{-\frac{\cos \frac{\rho-5}{2}(2t)}{p+1} \|u\|_{L^{p+1}(\mathbb{R})}^{p+1}} d\mu_0 \end{aligned}$$

which is **not invariant**.

Monotonicity of the measure ν_t

We define the measure

$$d\nu_t = e^{-\frac{\cos \frac{p-5}{2}(2t)}{p+1} \|u\|_{L^{p+1}(\mathbb{R})}^{p+1}} d\mu_0, \quad -\frac{\pi}{4} < t < \frac{\pi}{4}.$$

Proposition 8

For all $0 \leq |t| < \frac{\pi}{4}$

$$\nu_0(A) \leq \begin{cases} \left[\nu_t(\Phi(t, 0)A) \right]^{(\cos(2t)) \frac{5-p}{2}} & \text{if } 1 \leq p \leq 5 \\ \nu_t(\Phi(t, 0)A) & \text{if } p \geq 5. \end{cases}$$

- ▶ Allows to extend the globalisation argument of Bourgain relying on invariant measures.
- ▶ See the link with the Radon-Nikodym theorem.

A quantitative Radon-Nikodym theorem

Proposition 8

Let μ, ν be two finite measures on a measurable space (X, \mathcal{T}) . Assume that

$$\mu \ll \nu,$$

and more precisely

$$\exists 0 < \alpha \leq 1, \quad \exists C > 0, \quad \forall A \in \mathcal{T}, \quad \mu(A) \leq C\nu(A)^\alpha.$$

By the Radon-Nikodym theorem, there exists a $f \in L^1(d\nu)$ with $f \geq 0$, such that $d\mu = fd\nu$, and we write $f = \frac{d\mu}{d\nu}$.

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By the Radon-Nikodym theorem, there exists a $f \in L^1(d\nu)$ with $f \geq 0$, such that $d\mu = fd\nu$, and we write $f = \frac{d\mu}{d\nu}$.

(i) The assertion (2) is satisfied with $0 < \alpha < 1$ iff $f \in L_w^p(d\nu) \cap L^1(d\nu)$ with $p = \frac{1}{1-\alpha}$. In other words, $f \in L^1(d\nu)$ and

$$\nu(\{x : |f(x)| \geq \lambda\}) \leq C' \langle \lambda \rangle^{-p}, \quad \forall \lambda > 0.$$

(ii) The assertion (2) is satisfied with $\alpha = 1$ iff $f \in L^\infty(d\nu) \cap L^1(d\nu)$.

Monotonicity of the measure ν_t : proof

Proof :

► Set

$$\mathcal{E}(t, u(t)) = \frac{1}{2} \|\sqrt{H} u(t)\|_{L^2(\mathbb{R})}^2 + \frac{\cos \frac{p-5}{2}(2t)}{p+1} \|u(t)\|_{L^{p+1}(\mathbb{R})}^{p+1}.$$

A direct computation shows that

$$\frac{d}{dt}(\mathcal{E}(t, u(t))) = \frac{(5-p) \sin(2t) \cos \frac{p-7}{2}(2t)}{p+1} \|u(t)\|_{L^{p+1}(\mathbb{R})}^{p+1}.$$

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► Set $F(t) = \nu_t(\Phi(t, 0)A)$. Then

$$\frac{d}{dt}F(t) = (p-5) \tan(2t) \int_A \alpha(t, u(t)) e^{-\mathcal{E}(t, u(t))} du_0,$$

where $\alpha(t, u) = \frac{\cos \frac{p-5}{2}(2t)}{p+1} \|u\|_{L^{p+1}(\mathbb{R})}^{p+1}$.

Monotonicity of the measure ν_t : proof

- ▶ By Hölder, for all $k \geq 1$

$$\frac{d}{dt} F(t) \leq (p-5) \tan(2t) \frac{k}{e} (F(t))^{1-\frac{1}{k}}.$$

- ▶ Optimisation with $k = -\log(F(t))$ yields

$$\frac{d}{dt} F(t) \leq -(p-5) \tan(2t) \log(F(t)) F(t).$$

- ▶ Integration of the differential inequality

The Bourgain argument revisited

$$\begin{cases} i\partial_t u - Hu = \cos \frac{p-5}{2}(2t)|u|^{p-1}u, & (t, x) \in (-\frac{\pi}{4}, \frac{\pi}{4}) \times \mathbb{R}, \\ u|_{t=t_0} = u_0 \in X^0(\mathbb{R}). \end{cases}$$

- ▶ There exists a flow Φ such that the time of existence τ on the ball

$$B_R = \{u \in X^0(\mathbb{R}) : \|u\| \leq R^{1/2}\},$$

is uniform and such that $\tau \sim R^{-\kappa}$ for some $\kappa > 0$.

- ▶ Moreover, for all $|t| \leq \tau$

$$\Phi(t, 0)(B_R) \subset \{u \in X^0(\mathbb{R}) : \|u\| \leq (R+1)^{1/2}\}.$$

- ▶ We have the large deviation estimate $\mu_0(X^0(\mathbb{R}) \setminus B_R) \leq Ce^{-cR}$.

The Bourgain argument revisited

- For $T \leq e^{cR/2}$ fixed, we define the **set of the good data**

$$\Sigma_R = \bigcap_{k=-\lceil T/\tau \rceil}^{\lfloor T/\tau \rfloor} \Phi(k\tau, 0)^{-1}(B_R).$$

- By the **monotonicity** of $\nu_t(\Phi(t, 0)A)$:

$$\begin{aligned} \nu_0(X^0(\mathbb{R}) \setminus \Sigma_R) &\leq \sum_{k=-\lceil T/\tau \rceil}^{\lfloor T/\tau \rfloor} \nu_0\left(\Phi(k\tau, 0)^{-1}(X^0(\mathbb{R}) \setminus B_R)\right) \\ &\leq \sum_{k=-\lceil T/\tau \rceil}^{\lfloor T/\tau \rfloor} \nu_{k\tau}(X^0(\mathbb{R}) \setminus B_R) \\ &\leq (2\lceil T/\tau \rceil + 1)\mu_0(X^0(\mathbb{R}) \setminus B_R) \\ &\leq ce^{-cR/2} \end{aligned}$$

which shows that Σ_R is a big set of $X^0(\mathbb{R})$ when $R \rightarrow +\infty$.

The Bourgain argument revisited

- ▶ We deduce that for all $|t| \leq T$ and $u \in \Sigma_R$

$$\|\Phi(t, 0)(u)\| \leq (R + 1)^{1/2}.$$

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

$$\|\Phi(t, 0)(u)\| \leq C(\ln |t| + 1)^{1/2},$$