

Ergodicity for the Randomly Forced Navier–Stokes System in a Two-Dimensional Unbounded Domain

Vahagn Nersesyan

▶ To cite this version:

Vahagn Nersesyan. Ergodicity for the Randomly Forced Navier–Stokes System in a Two-Dimensional Unbounded Domain. Annales Henri Poincaré, 2022, 23 (6), pp.2277-2294. 10.1007/s00023-022-01163-3. hal-03775067

HAL Id: hal-03775067

https://hal.science/hal-03775067

Submitted on 26 Oct 2022

HAL is a multi-disciplinary open access archive for the deposit and dissemination of scientific research documents, whether they are published or not. The documents may come from teaching and research institutions in France or abroad, or from public or private research centers.

L'archive ouverte pluridisciplinaire **HAL**, est destinée au dépôt et à la diffusion de documents scientifiques de niveau recherche, publiés ou non, émanant des établissements d'enseignement et de recherche français ou étrangers, des laboratoires publics ou privés.

Ergodicity for the randomly forced Navier–Stokes system in a two-dimensional unbounded domain

Vahagn Nersesyan*

March 4, 2019

Abstract

The ergodic properties of the randomly forced Navier–Stokes system have been extensively studied in the literature during the last two decades. The problem has always been considered in bounded domains, in order to have, for example, suitable spectral properties for the Stokes operator, to ensure some compactness properties for the resolving operator of the system and the associated functional spaces, etc. In the present paper, we consider the Navier–Stokes system in an unbounded domain satisfying the Poincaré inequality. Assuming that the system is perturbed by a bounded non-degenerate noise, we establish uniqueness of stationary measure and exponential mixing in the dual-Lipschitz metric. The proof is carried out by developing the controllability approach of the papers [Shi19, KNS18] and using the asymptotic compactness of the dynamics.

AMS subject classifications: 35Q30, 37A25, 93B05.

Keywords: Navier–Stokes system, unbounded domain, exponential mixing, controllability, asymptotic compactness.

Contents

0	Introduction	2
1	Abstract criterion	6
2	Proof of the Main Theorem	10
	2.1 Preliminaries	. 10
	2.2 Regularity condition	. 11
	2.3 Asymptotic compactness	. 12
	2.4 Approximate controllability	. 14

^{*}Laboratoire de Mathématiques, UMR CNRS 8100, UVSQ, Université Paris-Saclay, 45, av. des Etats-Unis, F-78035 Versailles, France & Centre de Recherches Mathématiques, CNRS UMI 3457, Université de Montréal, Montréal, QC, H3C 3J7, Canada; e-mail: Vahagn.Nersesyan@math.uvsq.fr

3	Appendix			
	3.1	Proof of Proposition 2.1	15	
	3.2	Proof of Proposition 2.2	15	
	3.3	Proof of Lemma 2.3	16	

0 Introduction

The ergodicity of randomly forced 2D Navier–Stokes (NS) system has been widely studied in the literature in the case of bounded domains (see the papers [FM95, KS00, EMS01, BKL02] for the first results and the book [KS12] and the reviews [Fla08, Deb13] for a detailed discussion of different methods and for further references). This paper is concerned with the ergodic behaviour of the NS system in an *unbounded* domain D in \mathbb{R}^2 with smooth boundary ∂D :

$$\partial_t u - \nu \Delta u + \langle u, \nabla \rangle u + \nabla p = \eta(t, x), \quad x \in D,$$
 (0.1)

$$\operatorname{div} u = 0, \quad u|_{\partial D} = 0, \tag{0.2}$$

$$u(0) = u_0. (0.3)$$

Here $\nu > 0$ is the kinematic viscosity of the fluid, $u = (u_1(t, x), u_2(t, x))$ is the velocity field, p = p(t, x) is the pressure, and η is an external random force. To have a suitable dissipativity property for solutions, we assume that D is a *Poincaré domain*, i.e., there is a number $\lambda_1 > 0$ such that

$$\int_{D} |v|^{2} dx \le \lambda_{1}^{-1} \int_{D} |\nabla v|^{2} dx, \quad v \in C_{0}^{\infty}(D, \mathbb{R}^{2}).$$
 (0.4)

The random force η is a process of the form

$$\eta(t,x) = \sum_{k=1}^{\infty} \mathbb{I}_{[k-1,k)}(t)\eta_k(t-k+1,x), \quad t \ge 0, \ x \in D,$$
 (0.5)

where $\mathbb{I}_{[k-1,k)}$ is the indicator function of the interval [k-1,k) and $\{\eta_k\}$ is a sequence of independent and identically distributed (i.i.d.) random variables in the space² $E := L^2([0,1],H)$. Moreover, the law of η_k is assumed to be decomposable in the following sense.

Decomposability. There is an orthonormal basis $\{e_i\}$ in E such that

$$\eta_k = \sum_{j=1}^{\infty} b_j \xi_{jk} e_j \tag{0.6}$$

for some real-valued independent random variables ξ_{jk} verifying $|\xi_{jk}| \leq 1$ and some positive numbers b_j such that $\sum_{j=1}^{\infty} b_j^2 < \infty$. The law of the random variable ξ_{jk} is absolutely continuous with respect to the Lebesgue measure and the corresponding density ρ_j is C^1 -smooth and $\rho_j(0) > 0$ for all $j \geq 1$.

¹E.g., we can assume that D is bounded in some direction $a \in \mathbb{R}^2_*$, i.e., $\sup_{x \in D} |\langle a, x \rangle| < \infty$. ²We denote by H the usual functional space for the NS system defined by (0.9).

The restriction to integer times of the velocity field u_t defines a family of Markov processes (u_k, \mathbb{P}_u) parametrised by the initial condition $u_0 = u \in H$. The associated Markov operators are denoted by \mathfrak{P}_k and \mathfrak{P}_k^* . Recall that a measure $\mu \in \mathcal{P}(H)$ is stationary for the family (u_k, \mathbb{P}_u) if $\mathfrak{P}_1^*\mu = \mu$. In this paper, we prove the following result.

Main Theorem. Under the above assumptions, the family (u_k, \mathbb{P}_u) has a unique stationary measure $\mu \in \mathcal{P}(H)$. Moreover, it is exponentially mixing in the following sense: for any compact set $\mathcal{H} \subset H$, there are numbers C > 0 and c > 0 such that

$$\|\mathfrak{P}_k^*\lambda - \mu\|_{L(H)}^* \le Ce^{-ck}, \quad k \ge 0$$
 (0.7)

for any initial measure $\lambda \in \mathcal{P}(H)$ with supp $\lambda \subset \mathcal{H}$. Here $\|\cdot\|_{L(H)}^*$ is the dual-Lipschitz metric defined by (0.11).

To the best of our knowledge, this is the first result that establishes uniqueness of stationary measure and exponential mixing for the NS system in an unbounded domain. In this unbounded setting there are at least two additional difficulties compared to the case of a bounded domain. First, the resolving operator of the equation is not compact and does not have a compact attracting set. The second and more important difficulty is related to the presence of a continuous component in the spectrum of the Stokes operator. Let us note that the existence of a stationary measure for the NS system perturbed by the random force (0.5) can be established by combining the *Bogolyubov–Krylov argument* and the *asymptotic compactness* of the dynamics. When the driving force is a white noise, existence results are obtained in [EH01], for a real Ginzburg–Landau equation, in [Rou02], for a complex Ginzburg–Landau equation, in [BL06], for the NS system, and in [EKZ17], for a damped Schrödinger equation.

The uniqueness of stationary measure and mixing properties for PDEs in unbounded domains have been studied previously only for the Burgers equation. The case of inviscid equation on the real line has been considered in the paper [BCK14], under the assumption that the driving force is a space-time homogeneous Poisson point process. The proof is based on a combination of Lagrangian methods and first/last passage percolation theory. This result is generalised to the viscous case in the recent paper [BL19], where the perturbation is a space-time homogeneous random kick force. In both papers, the stationary measure is space translation invariant. In the case of the NS system perturbed by a homogeneous noise, the uniqueness of a space-time translation invariant measure (and in some cases also the existence) remains an open problem.

The proof of the Main Theorem uses a controllability approach of the papers [Shi15, Shi19], where exponential mixing is established for the NS system with a space-time and boundary localised forcing. In [KNS18, KNS19], methods of control theory are used to study a family of parabolic PDEs with a random perturbation that is highly degenerate in the Fourier space; see also the paper [KZ18], where a non-degenerate version of these results is presented. In [JNPS19], controllability is used to derive large deviations principle for the Lagrangian trajectories of the NS system.

The main novelty of the controllability argument we use here is that the deterministic operator is not supposed to be regularising to a space that is compactly embedded into the main phase space (see Theorem 1.1). That regularisation condition is replaced by two weaker properties: we assume that the nonlinear dynamics is asymptotically compact and the linearised operator can be decomposed into a sum of two operators one of which is *dissipative* and the other is *compact* (see (1.3)). It is proved that these weaker conditions are still sufficient for the exponential mixing. We make an essential use of the form of the nonlinearity to verify these properties for the NS system. The reader is referred to Section 1 for further discussion of the abstract controllability criterion and to Section 2 for the verification of the conditions of this criterion for the NS system.

Let us close this section with two remarks. Note that the convergence rate c in (0.7) depends on the initial compact \mathcal{H} . Indeed, this is related to the fact that the NS system in the unbounded case does not have a compact invariant attracting set. Without going into the details, let us mention that the existence of such attracting set (and a uniform convergence rate) can be recovered if we take initial condition and forcing in weighted Sobolev spaces as in [Bab92].

The second remark is about the NS system with the *Ekman damping*. By literally repeating the arguments of the proof of the Main Theorem, one can establish exponential mixing in the case of the whole space $D = \mathbb{R}^2$ or arbitrary unbounded³ domain $D \subset \mathbb{R}^2$ with smooth boundary for the following system:

$$\partial_t u - \nu \Delta u + au + \langle u, \nabla \rangle u + \nabla p = \eta(t, x),$$

 $\operatorname{div} u = 0, \quad u|_{\partial D} = 0,$

where a>0 is the damping parameter. The damping ensures a dissipativity property for the solutions without any assumption on the domain. We shall not discuss the details of this generalisation in this paper.

Acknowledgments

The author thanks Armen Shirikyan for helpful discussions and suggestions. This research was supported by Agence Nationale de la Recherche grant ANR-17-CE40-0006-02 and CNRS PICS grant Fluctuation theorems in stochastic systems.

Notation

Let $D \subset \mathbb{R}^2$ be an unbounded Poincaré domain with smooth boundary. In this paper, we use the following functional spaces.

 $C_0^{\infty}(D, \mathbb{R}^2)$ is the space of compactly supported smooth functions $u: D \to \mathbb{R}^2$, and

$$\mathcal{V} = \{ u \in C_0^{\infty}(D, \mathbb{R}^2) : \text{div } u = 0 \}.$$
 (0.8)

 $^{{}^{3}}$ I.e., domain D which does not necessarily satisfy the Poincaré inequality.

 $H^s(D,\mathbb{R}^2)$ and $L^p(D,\mathbb{R}^2)$ are the Sobolev and Lebesgue spaces on D. We consider the NS system in the usual functional spaces:

$$H = \text{closure of } \mathcal{V} \text{ in } L^2(D, \mathbb{R}^2),$$
 (0.9)

$$V = \text{closure of } \mathcal{V} \text{ in } H^1(D, \mathbb{R}^2)$$
 (0.10)

endowed with the scalar products

$$\langle u, v \rangle = \int_D u \cdot v \, dx, \qquad \langle u, v \rangle_1 = \int_D (\nabla u_1 \cdot \nabla v_1 + \nabla u_2 \cdot \nabla v_2) \, dx$$

and the corresponding norms $\|\cdot\| = \sqrt{\langle \cdot, \cdot \rangle}$ and $\|\cdot\|_1 = \sqrt{\langle \cdot, \cdot \rangle_1}$. The Poincaré inequality (0.4) implies that $\|\cdot\|_1$ is equivalent to the norm inherited from $H^1(D, \mathbb{R}^2)$. The dual of V with respect to the scalar product in H is denoted by V'.

Let X be a Polish space with metric d.

 $\mathcal{B}(X)$ denotes the Borel σ -algebra on X.

 $C_b(X)$ is the space of continuous functions $g: X \to \mathbb{R}$ endowed with the supnorm $\|g\|_{\infty}$. When X is compact, we write C(X).

 $L_b(X)$ is the space of functions $g \in C_b(X)$ for which the following norm is finite:

$$||g||_{L(X)} = ||g||_{\infty} + \sup_{u \neq v} \frac{|g(u) - g(v)|}{d(u, v)}.$$

 $\mathcal{P}(X)$ is the set of Borel probability measures on X endowed with the metric

$$\|\mu_1 - \mu_2\|_{L(X)}^* = \sup_{\|g\|_{L(X)} \le 1} |\langle g, \mu_1 \rangle - \langle g, \mu_2 \rangle|, \quad \mu_1, \mu_2 \in \mathcal{P}(X),$$
 (0.11)

where $\langle g, \mu \rangle = \int_X g(u) \mu(\mathrm{d}u)$.

Let E be a Banach space endowed with a norm $\|\cdot\|_E$, and let $J_T = [0, T]$. $L^p(J_T, E), 1 \le p < \infty$ is the space of measurable functions $u: J_T \to E$ such that

$$||u||_{L^p(J_T,E)} = \left(\int_0^T ||u(s)||_E^p ds\right)^{\frac{1}{p}} < \infty.$$

 $C(J_T, E)$ is the space of continuous functions $u: J_T \to E$ with the norm

$$||u||_{C(J_T,E)} = \sup_{t \in J_T} ||u(t)||_E.$$

 $\mathcal{L}(E,Y)$ is the space of bounded linear operators from E to another Banach space Y. We write $\mathcal{L}(E)$ when E=Y.

 $B_E(a,R)$ is the closed ball in E of radius R centered at a. We write $B_E(R)$ when a=0.

 $\mathcal{D}(\eta)$ is the law of E-valued random variable η .

1 Abstract criterion

Let H be a separable Hilbert space, and E be a separable Banach space. In this section, we consider a $random\ dynamical\ system$ of the form

$$u_k = S(u_{k-1}, \eta_k), \quad k \ge 1,$$
 (1.1)

where $S: H \times E \to H$ is a continuous mapping and $\{\eta_k\}$ is a sequence of i.i.d. random variables in E. Let $\mathcal{K} \subset E$ be the support of the law $\ell := \mathcal{D}(\eta_k)$. For any sequence $\{\zeta_k\}$ in \mathcal{K} , let us denote by $S_k(u; \zeta_1, \ldots, \zeta_k)$ the trajectory of (1.1) corresponding to the initial condition

$$u_0 = u \in H \tag{1.2}$$

and the vectors $\eta_i = \zeta_i$, i = 1, ..., k. For any set $\mathcal{H} \subset \mathcal{H}$, we define the set of attainability in time $k \geq 1$:

$$\mathcal{A}_k(\mathcal{H}) := \{ S_k(u; \zeta_1, \dots, \zeta_k) : u \in \mathcal{H}, \ \zeta_1, \dots, \zeta_k \in \mathcal{K} \}$$

and the set of attainability in infinite time:

$$\mathcal{A}(\mathcal{H}) := \overline{\cup_{k=1}^{\infty} \mathcal{A}_k(\mathcal{H})}^H.$$

The following five conditions are assumed to be satisfied for the mapping S and the measure ℓ .

(i) Regularity. The mapping $S: H \times E \to H$ is twice continuously differentiable, and its derivatives are bounded on bounded subsets. Moreover, for any $(u, \eta) \in H \times K$, the derivative $(D_u S)(u, \eta)$ can be represented as

$$(D_u S)(u, \eta) = \Psi_1 + \Psi_2(u, \eta),$$
 (1.3)

where the operators⁴ $\Psi_1, \Psi_2(u, \eta) \in L(H)$ are such that

$$\|\Psi_1\|_{\mathcal{L}(H)} =: \varkappa < 1 \tag{1.4}$$

and $\Psi_2(u,\eta)$ is compact.

- (ii) Asymptotic compactness. For any bounded sequence $\{u_0^n\}$ in H, any integers $l_n \geq 1$ such that $l_n \to \infty$, and any family $\{\zeta_m^n : m, n \geq 1\} \subset \mathcal{K}$, the sequence $\{S_{l_n}(u_0^n; \zeta_1^n, \ldots, \zeta_{l_n}^n)\}$ is precompact in H.
- (iii) Approximate controllability to a point. There is a point $\hat{u} \in H$ with the following property: for any $\varepsilon > 0$ and any compact \mathcal{H} in H, there is an integer $n \geq 1$ such that, for any initial point $u \in \mathcal{H}$, there are vectors $\zeta_1, \ldots, \zeta_n \in \mathcal{K}$ satisfying

$$||S_n(u;\zeta_1,\ldots,\zeta_n) - \hat{u}|| \le \varepsilon. \tag{1.5}$$

 $^{^{4}\}Psi_{1}$ does not depend on (u, η) .

- (iv) Approximate controllability of the linearisation. For any $u \in H$ and $\eta \in \mathcal{K}$, the image of the linear mapping $(D_{\eta}S)(u,\eta): E \to H$ is dense in H.
- (v) Decomposability. The set K is compact in E. Moreover, there are sequences of closed subspaces $\{F_n\}$ and $\{G_n\}$ in E satisfying the following properties:
 - F_n are finite-dimensional, $F_n \subset F_{n+1}$ for any $n \ge 1$, and $E = \overline{\bigcup_n F_n}$.
 - E is the direct sum of the spaces F_n and G_n , and the norms of the corresponding projections P_n and Q_n are bounded in $n \ge 1$.
 - ℓ is the product of projections $\mathsf{P}_{n*}\ell$ and $\mathsf{Q}_{n*}\ell$ for any $n \geq 1$. Moreover, the measures $\mathsf{P}_{n*}\ell$ have C^1 -smooth densities with respect to the Lebesgue measure on F_n .

System (1.1) defines a family of Markov processes (u_k, \mathbb{P}_u) in H parametrised by the initial condition (1.2). Let $\mathfrak{P}_k : C_b(H) \to C_b(H)$ and $\mathfrak{P}_k^* : \mathcal{P}(H) \to \mathcal{P}(H)$ be the associated Markov operators.

Theorem 1.1. Under Conditions (i)-(v), the family (u_k, \mathbb{P}_u) has a unique stationary measure $\mu \in \mathcal{P}(H)$. Moreover, for any compact set \mathcal{H} in H, there are numbers C > 0 and c > 0 such that

$$\|\mathfrak{P}_k^*\lambda - \mu\|_{L(H)}^* \le Ce^{-ck}, \quad k \ge 1$$
 (1.6)

for any initial measure $\lambda \in \mathcal{P}(H)$ with supp $\lambda \subset \mathcal{H}$.

See the papers [Shi15, Shi19, KNS18, KNS19, KZ18, JNPS19] for related abstract criteria for uniqueness of stationary measure and mixing. The formulation of Theorem 1.1 is close to the results in [KNS19, JNPS19], but the proof is based on a theorem obtained in [Shi19]. There are two main differences in our formulation. First, in Condition (i), we do not suppose that the mapping S (or its linearisation) takes values in a space that is compactly embedded into H. Instead, we assume that $(D_uS)(u,\eta)$ is a sum of dissipative and compact operators. The second difference is Condition (ii), which allows to recover some compactness properties⁵ for the dynamics. These two new conditions make Theorem 1.1 applicable to the NS system in unbounded domains. Moreover, this theorem can be applied to dissipative PDEs without parabolic regularisation (this will be addressed in a subsequent publication).

As in this paper the random perturbation is non-degenerate, the image of the mapping $(D_{\eta}S)(u,\eta)$ is assumed to be dense for any η in \mathcal{K} ; degenerate versions of this criterion can also be envisaged (cf. [KNS18, KNS19]).

Proof of Theorem 1.1. Step 1: Existence of stationary measure. The set $X := \mathcal{A}(\mathcal{H})$ is compact in H. Indeed, it suffices to show that any sequence $\{v_n\}$ of the form

$$v_n = S_{l_n}(u_0^n; \zeta_1^n, \dots, \zeta_{l_n}^n)$$

⁵The asymptotic compactness is a well-known property in the study of the attractors for deterministic PDEs (e.g., see [Lad91, Ros98]).

with some $u_0^n \in \mathcal{H}, \zeta_1^n, \ldots, \zeta_{l_n}^n \in \mathcal{K}$, and $l_n \geq 1$, is precompact in H. If the sequence $\{l_n\}$ is bounded, then $v_n \in \mathcal{A}_m(\mathcal{H})$ for all $n \geq 1$, where $m = \max\{l_n\}$. The fact that $\mathcal{A}_m(\mathcal{H})$ is compact (as image of a compact set by a continuous mapping) implies that $\{v_n\}$ is precompact. In the case $l_n \to \infty$, the conclusion follows from Condition (ii).

The compactness of X, combined with the invariance property $S(X \times \mathcal{K}) \subset X$ and the usual Bogolyubov–Krylov argument (e.g., see Section 2.5 in [KS12]), implies the existence of a stationary measure $\mu \in \mathcal{P}(X)$.

Step 2: Limit (1.6). According to Theorem 1.1 in [Shi19], limit (1.6) will be established if we verify the following property.

Local stabilisation. Let $D_{\delta} := \{(u, u') \in X \times X : ||u - u'|| \leq \delta\}$. For any R > 0 and any compact $K \subset E$, there is a finite-dimensional subspace $\mathcal{E} \subset E$, and a continuous mapping

$$\Phi: \mathbf{D}_{\delta} \times B_E(R) \to \mathcal{E}, \quad (u, u', \eta) \mapsto \eta',$$

which is continuously differentiable in η and satisfies the inequalities

$$\sup_{\eta \in B_E(R)} (\|\Phi(u, u', \eta)\|_E + \|D_{\eta}\Phi(u, u', \eta)\|_{\mathcal{L}(E)}) \le C \|u - u'\|, \tag{1.7}$$

$$\sup_{\eta \in \mathcal{K}} \|S(u,\eta) - S(u',\eta + \Phi(u,u',\eta))\| \le q \|u - u'\|, \quad (u,u') \in \mathbf{D}_{\delta}$$
 (1.8)

for some positive constants C, δ , and q < 1.

Let us show that Conditions (i), (iv), and (v) imply this local stabilisation property. We use a construction of approximate right inverse for linear operators from Section 2.2 in [KNS18]. For any $u \in X$ and $\eta \in E$, let $A(u,\eta) : E \to H$ be given by $A(u,\eta) := (D_{\eta}S)(u,\eta)$. Then $G(u,\eta) := A(u,\eta) A(u,\eta)^* : H \to H$ is non-negative self-adjoint operator and $\mathrm{Im}(G(u,\eta))$ is dense in H by Condition (iv). Thus $(G(u,\eta) + \gamma I)^{-1}$ is well defined for any $\gamma > 0$, and we have the limit

$$G(u,\eta)(G(u,\eta) + \gamma I)^{-1}f \to f \text{ as } \gamma \to 0^+, f \in H,$$
 (1.9)

by Lemma 2.4 in [KNS18]. This shows that $A(u,\eta)^*(G(u,\eta)+\gamma I)^{-1}$ is an approximate right inverse for $A(u,\eta)$. We truncate it to obtain an operator with finite-dimensional image:

$$\mathcal{R}_{M,\gamma}(u,\eta) := \mathsf{P}_M A(u,\eta)^* (G(u,\eta) + \gamma I)^{-1},$$

where P_M is a projection as in Condition (v) and $M \geq 1$ and $\gamma > 0$ are parameters that will be chosen later. It is straightforward to see that

$$\|\mathcal{R}_{M,\gamma}(u,\eta)\|_{\mathcal{L}(H,E)} + \|(D_{\eta}\mathcal{R}_{M,\gamma})(u,\eta)\|_{\mathcal{L}(H\times E,E)} \le C_1(R,M,\gamma)$$
 (1.10)

for any $u \in X$ and $\eta \in B_E(R)$, where the constant $C_1(R, M, \gamma) > 0$ does not depend on (u, η) . By the Taylor formula, for any $u, u' \in X$ and $\eta, \eta' \in B_E(R)$, we have

$$S(u', \eta') - S(u, \eta) = (D_u S)(u, \eta)(u' - u) + (D_\eta S)(u, \eta)(\eta' - \eta) + r(u, u', \eta, \eta')$$
(1.11)

where

$$||r(u, u', \eta, \eta')|| \le C_2(R) \left(||u - u'||^2 + ||\eta - \eta'||_F^2\right). \tag{1.12}$$

The mapping Φ is defined by

$$\Phi(u, u', \eta) := -\mathcal{R}_{M, \gamma}(u, \eta)\Psi_2(u, \eta)(u' - u).$$

Then (1.7) is verified due to (1.10) and

$$C_3 := \sup_{(u,\eta) \in X \times B_E(R)} \left(\|\Psi_2(u,\eta)\|_{\mathcal{L}(H)} + \|D_{\eta}\Psi_2(u,\eta)\|_{\mathcal{L}(H \times E,H)} \right) < \infty. \quad (1.13)$$

The fact that C_3 is finite follows from the boundedness on $X \times B_E(R)$ of the norms of the derivatives of $S(u, \eta)$ and the representation (1.3). Assume that, for any $\varepsilon > 0$, we are able to find numbers M > 1 and $\gamma > 0$ such that

$$\|(D_{\eta}S)(u,\eta)\mathcal{R}_{M,\gamma}(u,\eta)\Psi_2(u,\eta)g - \Psi_2(u,\eta)g\| \le \varepsilon \|g\|$$
(1.14)

for any $u \in X$, $\eta \in \mathcal{K}$, and $g \in H$. Then combining (1.3), (1.4), and (1.11)-(1.14), we obtain

$$||S(u,\eta) - S(u', \eta + \Phi(u, u', \eta))||$$

$$\leq \varkappa ||u' - u|| + \varepsilon ||u' - u|| + ||r(u, u', \eta, \eta + \Phi(u, u', \eta))||$$

$$< (\varkappa + \varepsilon + C_4(\varepsilon, R) ||u' - u||)||u' - u|| < q ||u' - u||$$

for any $(u, u') \in \mathbf{D}_{\delta}$ and $\eta \in \mathcal{K}$, where $\varepsilon > 0$ and $\delta > 0$ are sufficiently small and q < 1. This completes the proof of the local stabilisation and limit (1.6).

Let us prove inequality (1.14). By (1.9), the continuous dependence of the operators $(D_{\eta}S)(u,\eta)$ and $\Psi_2(u,\eta)$ on (u,η) , and a simple compactness argument, we find a large integer $M \geq 1$ and a small number $\gamma > 0$ such that⁶

$$||(D_{\eta}S)(u,\eta)\mathcal{R}_{M,\gamma}(u,\eta)f - f|| \leq \varepsilon$$

for any $u \in X$, $\eta \in \mathcal{K}$, and $f \in \Psi_2(u, \eta)(B_H(1))$. Then inequality (1.14) follows by homogeneity.

Step 3: Uniqueness of stationary measure. To complete the proof, it remains to show the uniqueness of stationary measure in $\mathcal{P}(H)$. Assume that μ is the stationary measure supported⁷ in $\mathcal{A}(\{\hat{u}\})$, and let λ_1 be any stationary measure for (u_k, \mathbb{P}_u) in $\mathcal{P}(H)$. As $\lambda_1(\mathcal{A}(H)) = 1$, there is a sequence $\{\mathcal{H}_n\}$ of compacts in H such that

$$\lambda_1(\mathcal{A}(\mathcal{H}_n)) > 1 - 1/n$$
 for any $n \ge 1$.

By Condition (iii), we have $\mathcal{A}(\{\hat{u}\}) \subset \mathcal{A}(\mathcal{H}_n)$, and μ and $\lambda_1/\lambda_1(\mathcal{A}(\mathcal{H}_n))$ are stationary measures for (u_k, \mathbb{P}_u) in $\mathcal{P}(\mathcal{A}(\mathcal{H}_n))$. So $\mu = \lambda_1/\lambda_1(\mathcal{A}(\mathcal{H}_n))$ by the uniqueness of stationary measure on $X = \mathcal{A}(\mathcal{H}_n)$. Therefore,

$$\lambda_1(\Gamma) = \lim_{n \to \infty} \lambda_1(\Gamma \cap \mathcal{A}(\mathcal{H}_n)) = \lim_{n \to \infty} \lambda_1(\Gamma \cap \mathcal{A}(\mathcal{H}_n)) / \lambda_1(\mathcal{A}(\mathcal{H}_n)) = \mu(\Gamma)$$

for any
$$\Gamma \in \mathcal{B}(H)$$
. Thus $\lambda_1 = \mu$.

⁶Here we use also the boundedness of the projections P_M (see Condition (v)), which implies the limit $\mathsf{P}_M \to I$ as $M \to \infty$ in the operator topology.

⁷Using Condition (iii), it is easy to see that supp $\mu = \mathcal{A}(\{\hat{u}\})$.

2 Proof of the Main Theorem

In this section, we prove the Main Theorem by applying Theorem 1.1. We begin with a short discussion of the deterministic NS system, then turn to the verification of Conditions (i)-(v) in an appropriate functional setting.

2.1 Preliminaries

Applying the Leray projection⁸ Π to Eq. (0.1), we eliminate the pressure term and consider the evolution system

$$\dot{u} + \nu L u + B(u) = \eta, \tag{2.1}$$

where $L = -\Pi \Delta$ is the Stokes operator and $B(u) = \Pi(\langle u, \nabla \rangle u)$. Let us define a bilinear symmetric form $[\cdot, \cdot] : V \times V \to \mathbb{R}$ by

$$[u,v]:=\langle u,v\rangle_1-\frac{\lambda_1}{2}\langle u,v\rangle,\quad u,v\in V.$$

The Poincaré inequality (see (0.4))

$$||u||^2 \le \lambda_1^{-1} ||u||_1^2, \quad u \in V$$
 (2.2)

implies that

$$\frac{1}{2}||u||_1^2 \le [u]^2 := [u, u] \le ||u||_1^2. \tag{2.3}$$

Thus $[\cdot,\cdot]$ defines a scalar product on V with norm $[\cdot]$ equivalent to $\|\cdot\|_1$.

Proposition 2.1. For any T > 0, $u_0 \in H$, and $\eta \in L^2(J_T, H)$, there is a unique solution $u \in C(J_T, H) \cap L^2(J_T, V)$ of problem (2.1), (0.3). It satisfies the following inequalities

$$||u(t)||^2 \le e^{-\nu\lambda_1 t} ||u_0||^2 + \nu^{-2} \lambda_1^{-1} ||\eta||_{L^2(J_T, V')}^2, \quad t \in J_T,$$
 (2.4)

$$||u||_{L^{2}(J_{T},V)}^{2} \le \nu^{-1}||u_{0}||^{2} + \nu^{-2}||\eta||_{L^{2}(J_{T},V')}^{2}, \tag{2.5}$$

and the equality

$$||u(t)||^{2} = e^{-\nu\lambda_{1}t}||u_{0}||^{2} + 2\int_{0}^{t} e^{-\nu\lambda_{1}(t-s)} \left(\langle \eta(s), u(s) \rangle - \nu[u(s)]^{2}\right) ds, \quad t \in J_{T}.$$
(2.6)

Let us define the mapping

$$S_t: H \times L^2(J_T, H) \to H, \quad (u_0, \eta) \mapsto u(t), \quad t \in J_T.$$

The following stability result holds.

⁸That is the orthogonal projection in $L^2(D, \mathbb{R}^2)$ onto H.

Proposition 2.2. Assume that the sequence $\{u_0^n\}$ converges weakly to u_0 in H and the sequence $\{\eta^n\}$ converges strongly to η in $L^2(J_T, H)$. Then

$$S_t(u_0^n, \eta^n) \rightharpoonup S_t(u_0, \eta)$$
 weakly in H for $t \in J_T$, (2.7)

$$S_{\cdot}(u_0^n, \eta^n) \rightharpoonup S_{\cdot}(u_0, \eta)$$
 weakly in $L^2(J_T, V)$. (2.8)

The proofs of these two propositions are carried out by standard methods and are given in the Appendix.

Let us now describe the functional setting in which Theorem 1.1 is applied. The space H is defined by (0.9), $E := L^2([0,1], H)$, and the mapping

$$S := S_1 : H \times E \to H, \quad (u_0, \eta) \mapsto u(1)$$

is the time-one shift along trajectories of Eq. (2.1). Then the restriction to integer times of the solution of (2.1), (0.3), (0.5) satisfies

$$u_k = S(u_{k-1}, \eta_k), \quad k \ge 1,$$
 (2.9)

where $\{\eta_k\}$ is a sequence of i.i.d. random variables as in the decomposability condition in the Introduction. It is straightforward to see that Condition (v) satisfied.⁹ In the next three subsections, we check Conditions (i)-(iv).

2.2 Regularity condition

The smoothness of the mapping $S: H \times E \to H$ and the boundedness of its derivatives on bounded subsets of $H \times E$ are proved using well-known methods (e.g., see Chapters I and VII in [BV92] for the case of bounded domain D).

Let us consider the linearisation of Eq. (2.1) around the trajectory $\tilde{u}(t) = S_t(u, \eta)$ corresponding to an initial condition $\tilde{u}(0) = u \in H$ and control $\eta \in \mathcal{K}$:

$$\dot{w} + \nu L w + Q(\tilde{u}, w) = 0, \tag{2.10}$$

$$w(0) = w_0, (2.11)$$

where

$$Q(a,b) = \Pi(\langle a, \nabla \rangle b) + \Pi(\langle b, \nabla \rangle a). \tag{2.12}$$

Then $(D_uS)(u,\eta)w_0 = w(1)$ for any $w_0 \in H$, and we can write $w = v_1 + v_2$, where v_1 and v_2 are the solutions of the problems

$$\dot{v}_1 + \nu L v_1 = 0, \quad v_1(0) = w_0,$$
 (2.13)

$$\dot{v}_2 + \nu L v_2 + Q(\tilde{u}, w) = 0, \quad v_2(0) = 0. \tag{2.14}$$

The representation (1.3) holds with the linear operators

$$\Psi_1: H \to H, \quad w_0 \mapsto v_1(1),$$

 $\Psi_2(u, \eta): H \to H, \quad w_0 \mapsto v_2(1).$

⁹Note that $\mathcal{K} := \sup \mathcal{D}(\eta_k)$ is compact in E, since it is contained in a Hilbert cube.

Inequality

$$||v_1(t)|| \le e^{-\nu\lambda_1 t/2} ||w_0||$$

implies (1.4) with $\varkappa := e^{-\nu\lambda_1/2}$. The following lemma (whose proof is given in the Appendix) completes the verification of Condition (i).

Lemma 2.3. For any $(u, \eta) \in H \times K$, the operator $\Psi_2(u, \eta) : H \to H$ is compact.

2.3 Asymptotic compactness

For any integer $k \geq 1$ and any function $\eta \in L^2(J_k, H)$ of the form (0.5) with some sequence $\{\eta_n\} \subset E$, we write

$$u(t) = S_t(u_0, \eta) =: S_t(u_0; \eta_1, \dots, \eta_k), \quad t \in J_k.$$
 (2.15)

Proposition 2.4. For any bounded sequence $\{u_0^n\}$ in H, any sequence of integers $l_n \geq 1$ such that $l_n \to \infty$, and any family $\{\zeta_m^n : m, n \geq 1\} \subset \mathcal{K}$, the sequence

$$v_n = S_{l_n}(u_0^n; \zeta_1^n, \dots, \zeta_{l_n}^n)$$

is precompact in H.

Proof. Let us first explain the scheme of the proof. Using the dissipativity of the system, we show that $\{v_n\}$ is bounded in H. As H is a Hilbert space, there is a subsequence $\{v_{k_n}\}$ such that

$$v_{k_n} \rightharpoonup w$$
 weakly in H . (2.16)

Hence, we have

$$||w|| \le \liminf_{n \to \infty} ||v_{k_n}||. \tag{2.17}$$

On the other hand, using energy equality (2.6) and Proposition 2.2, we show that

$$\lim_{n \to \infty} \sup_{w \to \infty} \|v_{k_n}\| \le \|w\|. \tag{2.18}$$

Inequalities (2.17) and (2.18) together imply that

$$v_{k_n} \to w$$
 strongly in H ,

which gives the required result.

Step 1: Boundedness. Let us show that, for any bounded set $\mathcal{H} \subset \mathcal{H}$, the attainability set $\mathcal{A}(\mathcal{H}) \subset \mathcal{H}$ is bounded. Indeed, let

$$M_1 := \sup_{u \in \mathcal{H}} \|u\|^2 < \infty.$$

For any $u_0 \in \mathcal{H}$ and $\zeta_1, \ldots, \zeta_k \in \mathcal{K}$, let $u_k := S_k(u_0; \zeta_1, \ldots, \zeta_k)$. Then by inequality (2.4), we have

$$||u_k||^2 \le \varkappa ||u_{k-1}||^2 + M_2,$$

where $\varkappa := e^{-\nu \lambda_1} < 1$ and $M_2 := \nu^{-2} \lambda_1^{-1} \sup_{\eta \in \mathcal{K}} \|\eta\|_{L^2([0,1],V')}^2$. Iterating this, we get

$$||u_k||^2 \le \varkappa^k ||u_0||^2 + M_2(1-\varkappa)^{-1} \le M_1 + M_2(1-\varkappa)^{-1} =: M.$$

This shows that

$$\sup_{u \in \mathcal{A}(\mathcal{H})} \|u\|^2 \le M.$$

Step 2: Proof of (2.18). From the previous step it follows that $\{v_n\}$ is bounded in H. Let $\{v_{k_n}\}$ be a subsequence verifying (2.16). It is of the form

$$v_{k_n} = S_{\ell_n}(u_0^n; \eta_1^n, \dots, \eta_{\ell_n}^n)$$

for some vectors $\eta_1^n, \ldots, \eta_{\ell_n}^n \in \mathcal{K}$ and integers $\ell_n \geq 1$ such that $\ell_n \to \infty$. Using the boundedness of the set $\mathcal{A}(\{u_0^n, n \geq 1\})$ and passing to a subsequence if necessary (applying the diagonal process), we can assume that

$$w_{n,m} := S_{\ell_n - m}(u_0^n; \eta_1^n, \dots, \eta_{\ell_n - m}^n) \rightharpoonup w_m \quad \text{weakly in } H$$
 (2.19)

for any $m \geq 1$ and some $w_m \in H$. Using the compactness of \mathcal{K} and again passing to a subsequence if necessary, we can assume that

$$\eta_{\ell_n-i}^n \to \xi_{m-i}^m$$
 strongly in $E, i = 0, \dots m-1$.

Combining this with (2.19) and Proposition 2.2, we get

$$v_{k_n} = S_m(w_{n,m}; \eta_{\ell_n-m+1}^n, \dots, \eta_{\ell_n}^n) \rightharpoonup S_m(w_m; \xi_1^m, \dots, \xi_m^m)$$
 weakly in H .

From (2.16) we derive the equality

$$w = S_m(w_m; \xi_1^m, \dots, \xi_m^m)$$
 for any $m \ge 1$.

This and (2.6) imply that

$$||w||^{2} = ||u^{m}(m)||^{2}$$

$$= e^{-\nu\lambda_{1}m}||w_{m}||^{2} + 2\int_{0}^{m} e^{-\nu\lambda_{1}(m-s)} \left(\langle \xi^{m}(s), u^{m}(s) \rangle - \nu[u^{m}(s)]^{2}\right) ds,$$
(2.20)

where

$$u^{m}(t) = S_{t}(w_{m}, \xi^{m}), \qquad \xi^{m}(t) = \sum_{k=1}^{m} \mathbb{I}_{[k-1,k)}(t)\xi_{k}^{m}(t-k+1), \quad t \in J_{m}.$$

On the other hand, again by (2.6), we have

$$||v_{k_n}||^2 = ||u^{n,m}(m)||^2$$

$$= e^{-\nu\lambda_1 m} ||w_{n,m}||^2 + 2 \int_0^m e^{-\nu\lambda_1 (m-s)} (\langle \eta^{n,m}(s), u^{n,m}(s) \rangle - \nu [u^{n,m}(s)]^2) ds,$$
(2.21)

where

$$u^{n,m}(t) = S_t(w_{n,m}, \eta^{n,m}), \quad \eta^{n,m}(t) = \sum_{k=1}^m \mathbb{I}_{[k-1,k)}(t)\eta^n_{k_n-m+k}(t-k+1), \quad t \in J_m.$$

Note that

$$\int_0^m e^{-\nu\lambda_1(m-s)} \langle \eta^{n,m}(s), u^{n,m}(s) \rangle \,\mathrm{d}s \to \int_0^m e^{-\nu\lambda_1(m-s)} \langle \xi^m(s), u^m(s) \rangle \,\mathrm{d}s,$$
(2.22)

as $\eta^{n,m} \to \xi^m$ strongly in $L^2(J_m,H)$ and $u^{n,m} \rightharpoonup u^m$ weakly in $L^2(J_m,V)$. Since

$$\left(\int_0^m e^{-\nu\lambda_1(m-s)} [\cdot]^2 \,\mathrm{d}s\right)^{1/2}$$

is a norm in $L^2(J_m, V)$ which is equivalent to the original one, the following inequality holds

$$\int_0^m e^{-\nu\lambda_1(m-s)} [u^m(s)]^2 ds \le \liminf_{n \to \infty} \int_0^m e^{-\nu\lambda_1(m-s)} [u^{n,m}(s)]^2 ds.$$

Combining this with (2.20)-(2.22), we get

$$\limsup_{n \to \infty} \|v_{k_n}\|^2 \le \limsup_{n \to \infty} \left(e^{-\nu \lambda_1 m} \|w_{n,m}\|^2 \right) + \|w\|^2 - e^{-\nu \lambda_1 m} \|w_m\|^2$$

$$\le e^{-\nu \lambda_1 m} M + \|w\|^2.$$

As $m \ge 1$ is arbitrary, we arrive at (2.18).

2.4 Approximate controllability

By the decomposability assumption, we have $0 \in \mathcal{K}$. In view of (2.4), Condition (iii) is verified with $\hat{u} = 0$, $\zeta_1 = \ldots = \zeta_n = 0$, and sufficiently large $n \ge 1$.

To check Condition (iv), we consider the following linearisation of Eq. (2.1) around the same trajectory \tilde{u} as in Section 2.2:

$$\dot{w} + \nu L w + Q(\tilde{u}, w) = \zeta(t),$$

$$w(0) = 0.$$
(2.23)

where Q is given by (2.12). Then $(D_{\eta}S)(u,\eta)\zeta = w(1)$ for any $\zeta \in E$. For any smooth function $w_1 \in H$, we can find a smooth function $w: [0,1] \times D \to \mathbb{R}^2$ such that w(0) = 0, $w(1) = w_1$, and $w(t) \in H$ for all $t \in [0,1]$. Replacing w into Eq. (2.23), we find explicitly a control $\zeta \in E$ such that $(D_{\eta}S)(u,\eta)\zeta = w_1$. This shows that the image of the mapping $(D_{\eta}S)(u,\eta): E \to H$ is dense in H for any $(u,\eta) \in H \times K$. Thus Conditions (i)-(v) are verified. Applying Theorem 1.1, we complete the proof of the Main Theorem.

3 Appendix

3.1 Proof of Proposition 2.1

The existence and uniqueness of solution is proved, e.g., in Chapter III of [Tem77]. Here we give a formal derivation of inequalities (2.4) and (2.5) and equality (2.6). Taking the scalar product in H of Eq. (2.1) with 2u and using the identity

$$\langle B(u), u \rangle = 0, \quad u \in V,$$

we get

$$\frac{\mathrm{d}}{\mathrm{d}t} \|u\|^2 + 2\nu \|u\|_1^2 = 2\langle \eta, u \rangle \le 2\|\eta\|_{V'} \|u\|_1 \le \nu \|u\|_1^2 + \nu^{-1} \|\eta\|_{V'}^2. \tag{3.1}$$

Combining this with the Poincaré inequality, we obtain (2.4). From (3.1) we also derive the inequality

$$\int_0^t \|u\|_1^2 \, \mathrm{d}s \le \nu^{-1} \|u_0\|^2 + \nu^{-2} \int_0^t \|\eta\|_{V'}^2 \, \mathrm{d}s \quad t \in J_T$$
 (3.2)

which implies (2.5). To prove (2.6), we rewrite the equality in (3.1) in the form

$$\frac{\mathrm{d}}{\mathrm{d}t} \|u\|^2 + \nu \lambda_1 \|u\|^2 = 2 \left(\langle \eta, u \rangle - \nu [u]^2 \right)$$

and apply the variation of constants formula.

3.2 Proof of Proposition 2.2

By Proposition 2.1, the sequence $u^n = S(u_0^n, \eta^n)$ is bounded in the space $C(J_T, H) \cap L^2(J_T, V)$. From the inequality

$$||B(u)||_{V'} \le C||u|||u||_1, \quad u \in V$$

and Eq. (2.1) we derive that

$$\|\dot{u}^n\|_{L^2(J_T,V')} \le M_1, \quad n \ge 1.$$
 (3.3)

For any R > 0, let us set $D_R := D \cap B_{\mathbb{R}^2}(R)$. Then the space $H^1(D_R)$ is compactly embedded into $L^2(D_R)$. Applying Theorem 2.1 of Chapter III in [Tem77] with spaces $X_0 = H^1(D_R)$, $X = L^2(D_R)$, $X_1 = H^{-1}(D_R)$ and numbers $\alpha_0 = \alpha_1 = 2$, and using the diagonal process, we find a subsequence $\{u^{k_n}\}$ such that

$$u^{k_n} \rightharpoonup u$$
 weak-star in $L^{\infty}(J_T, H)$, (3.4)

$$u^{k_n} \rightharpoonup u$$
 weakly in $L^2(J_T, V)$, (3.5)

$$u^{k_n} \to u$$
 strongly in $L^2(J_T, L^2(D_R))$ (3.6)

for any R > 0. Passing to the limit in the equation for $u^{k_n}(t)$, we conclude that u(t) is the solution $S_t(u_0, \eta)$, $t \in J_T$. Moreover, by the uniqueness of the limit, it is easy to see that limits (3.4)-(3.6) hold for the full sequence $\{u^n\}$. This proves (2.8).

Let us take any $\varphi \in \mathcal{V}$ (see (0.8)). By inequality (2.4), we have

$$|\langle u^n(t), \varphi \rangle| \le M_2, \quad n \ge 1, t \in J_T,$$

and by inequality (3.3),

$$|\langle u^n(t+\tau) - u^n(t), \varphi \rangle| \le \int_t^{t+\tau} |\langle \dot{u}^n(s), \varphi \rangle| \, \mathrm{d}s$$
$$\le \tau^{\frac{1}{2}} ||\dot{u}^n||_{L^2(J_T, V')} ||\varphi||_1 \le \tau^{\frac{1}{2}} M_1 ||\varphi||_1$$

for any $t \in J_T$ and $\tau \in (0, T - t)$. Thus the Arzelà-Ascoli theorem implies that

$$\langle u^n(t), \varphi \rangle \to \langle u(t), \varphi \rangle$$

uniformly in $t \in J_T$. Using the fact that \mathcal{V} is dense in H, we get (2.7).

3.3 Proof of Lemma 2.3

Let us set J := [0,1]. By Proposition 2.1, we have $\tilde{u} = S.(u,\eta) \in L^2(J,V)$ for any $u \in H$ and $\eta \in \mathcal{K}$. Using standard methods, one can show that the mapping $\Psi : w_0 \mapsto w$ (i.e., the resolving operator of problem (2.10), (2.11)) is continuous from H to $\mathcal{X} := L^2(J,V) \cap W^{1,2}(J,H^{-1})$. Let us show that the linear mapping

$$\Phi^{\tilde{u}}: H \to L^2(J, H^{-1}), \quad w_0 \mapsto \langle \tilde{u}, \nabla \rangle w + \langle w, \nabla \rangle \tilde{u}$$

is compact. Indeed, let $\{w_0^n\}$ be a bounded sequence in H. Then the sequence $\{\Psi(w_0^n)\}$ is bounded in \mathcal{X} . As in the previous subsection, we apply Theorem 2.1 of Chapter III in [Tem77] with spaces $X_0 = H^1(D_R)$, $X = L^2(D_R)$, $X_1 = H^{-1}(D_R)$ and numbers $\alpha_0 = \alpha_1 = 2$, and use the diagonal process to find a subsequence $\{\Psi(w_0^{k_n})\}$ converging strongly in $L^2(J, L^2(D_R))$ for any R > 0.

On the other hand, for any $\varepsilon > 0$, there is a number R > 0 and a smooth function $\varphi : J \times D \to \mathbb{R}^2$ with supp $\varphi(t, \cdot) \subset D_R$ for any $t \in J$ and

$$\|\tilde{u} - \varphi\|_{L^2(J,V)} < \varepsilon. \tag{3.7}$$

Using the boundedness of the sequence $\{\Phi^{\tilde{u}}(w_0^{k_n})\}$ in $L^2(J,V)$ and the estimate

$$\|\langle a, \nabla \rangle b\|_{H^{-1}} \le C_1 \|a\|_{H^{1/2}} \|b\|_{H^{1/2}} \le C_2 \|a\|_1 (\|b\| \|b\|_1)^{\frac{1}{2}}, \quad a, b \in V,$$

we obtain

$$\begin{split} \|\Phi^{\tilde{u}}(w_0^{k_n}) - \Phi^{\tilde{u}}(w_0^{k_m})\|_{L^2(J,H^{-1})} &\leq \|\Phi^{\varphi}(w_0^{k_n}) - \Phi^{\varphi}(w_0^{k_m})\|_{L^2(J,H^{-1})} \\ &+ \|\Phi^{\tilde{u}-\varphi}(w_0^{k_n}) - \Phi^{\tilde{u}-\varphi}(w_0^{k_m})\|_{L^2(J,H^{-1})} \\ &\leq C_3 \|\Psi(w_0^{k_n}) - \Psi(w_0^{k_m})\|_{L^2(J,L^2(D_R))}^{1/2} \\ &+ C_3 \|\tilde{u} - \varphi\|_{L^2(J,V)}, \end{split}$$

where the constant $C_3 > 0$ does not depend on the numbers n, m, ε . Combining this with (3.7) and choosing n and m sufficiently large, we see that

$$\|\Phi^{\tilde{u}}(w_0^{k_n}) - \Phi^{\tilde{u}}(w_0^{k_m})\|_{L^2(J,H^{-1})} < 2C_3\varepsilon.$$

This shows that $\Phi^{\tilde{u}}$ is compact. Thus the map $\Psi_2: H \to H$ is compact as a composition of $\Phi^{\tilde{u}}$ with some linear continuous map.

References

- [Bab92] A. V. Babin. The attractor of a Navier-Stokes system in an unbounded channel-like domain. *J. Dynam. Differential Equations*, 4(4):555–584, 1992.
- [BCK14] Y. Bakhtin, E. Cator, and K. Khanin. Space-time stationary solutions for the Burgers equation. J. Amer. Math. Soc., 27(1):193–238, 2014.
- [BKL02] J. Bricmont, A. Kupiainen, and R. Lefevere. Exponential mixing of the 2D stochastic Navier–Stokes dynamics. Comm. Math. Phys., 230(1):87–132, 2002.
- [BL06] Z. Brzeźniak and Y. Li. Asymptotic compactness and absorbing sets for 2D stochastic Navier-Stokes equations on some unbounded domains. *Trans. Amer. Math. Soc.*, 358(12):5587–5629, 2006.
- [BL19] Y. Bakhtin and L. Li. Thermodynamic Limit for Directed Polymers and Stationary Solutions of the Burgers Equation. Comm. Pure Appl. Math., 72(3):536-619, 2019.
- [BV92] A. V. Babin and M. I. Vishik. Attractors of Evolution Equations. 1992.
- [Deb13] A. Debussche. Ergodicity results for the stochastic Navier-Stokes equations: an introduction. In *Topics in mathematical fluid mechanics*, volume 2073, pages 23–108. Springer, Heidelberg, 2013.
- [EH01] J.-P. Eckmann and M. Hairer. Invariant measures for stochastic partial differential equations in unbounded domains. *Nonlinearity*, 14(1):133–151, 2001.
- [EKZ17] I. Ekren, I. Kukavica, and M. Ziane. Existence of invariant measures for the stochastic damped Schrödinger equation. Stoch. Partial Differ. Equ. Anal. Comput., 5(3):343–367, 2017.
- [EMS01] W. E, J. C. Mattingly, and Ya. Sinai. Gibbsian dynamics and ergodicity for the stochastically forced Navier–Stokes equation. Comm. Math. Phys., 224(1):83–106, 2001.

- [Fla08] F. Flandoli. An introduction to 3D stochastic fluid dynamics. In SPDE in hydrodynamic: recent progress and prospects, volume 1942, pages 51–150. Springer, Berlin, 2008.
- [FM95] F. Flandoli and B. Maslowski. Ergodicity of the 2D Navier–Stokes equation under random perturbations. *Comm. Math. Phys.*, 172(1):119–141, 1995.
- [JNPS19] V. Jakšić, V. Nersesyan, C.-A. Pillet, and A. Shirikyan. Large deviations and entropy production in viscous fluid flows. *Preprint*, arXiv:1902.03278, 2019.
- [KNS18] S. Kuksin, V. Nersesyan, and A. Shirikyan. Exponential mixing for a class of dissipative PDEs with bounded degenerate noise. *Preprint*, arXiv:1802.03250, 2018.
- [KNS19] S. Kuksin, V. Nersesyan, and A. Shirikyan. Mixing via controllability for randomly forced nonlinear dissipative PDEs. Preprint, arXiv:1902.00494, 2019.
- [KS00] S. Kuksin and A. Shirikyan. Stochastic dissipative PDEs and Gibbs measures. *Comm. Math. Phys.*, 213(2):291–330, 2000.
- [KS12] S. Kuksin and A. Shirikyan. *Mathematics of Two-Dimensional Tur-bulence*. Cambridge University Press, Cambridge, 2012.
- [KZ18] S. Kuksin and H. Zhang. Exponential mixing for dissipative PDEs with bounded non-degenerate noise. *Preprint*, *arXiv:1812.11706*, 2018.
- [Lad91] O. Ladyzhenskaya. Attractors for semigroups and evolution equations. Lezioni Lincee. Cambridge University Press, Cambridge, 1991.
- [Ros98] R. Rosa. The global attractor for the 2D Navier-Stokes flow on some unbounded domains. *Nonlinear Anal.*, 32(1):71–85, 1998.
- [Rou02] J. Rougemont. Space-time invariant measures, entropy, and dimension for stochastic Ginzburg-Landau equations. Comm. Math. Phys., 225(2):423-448, 2002.
- [Shi15] A. Shirikyan. Control and mixing for 2D Navier-Stokes equations with space-time localised noise. Ann. Sci. Éc. Norm. Supér. (4), 48(2):253–280, 2015.
- [Shi19] A. Shirikyan. Controllability implies mixing II. Convergence in the dual-Lipschitz metric. *To appear in J. Eur. Math. Soc.*, 2019.
- [Tem77] R. Temam. Navier-Stokes equations. Theory and numerical analysis. North-Holland, 1977.