



Electron. J. Probab. **25** (2020), article no. 69, 1–24.
ISSN: 1083-6489 <https://doi.org/10.1214/20-EJP474>

Stationary solutions of damped stochastic 2-dimensional Euler's equation

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Abstract

Existence of stationary point vortices solution to the damped and stochastically driven Euler's equation on the two dimensional torus is proved, by taking limits of solutions with finitely many vortices. A central limit scaling is used to show in a similar manner the existence of stationary solutions with white noise marginals.

Keywords: point vortices; invariant measure; Euler equations.

AMS MSC 2010: Primary 35Q35, Secondary 35R60; 60G10; 60H15; 76B03; 76M35.

Submitted to EJP on January 22, 2019, final version accepted on June 2, 2020.

1 Introduction

The present work concerns a particular class of solutions to the 2-dimensional incompressible Euler's equation with frictional damping, on the torus $T^2 = \mathbb{R}^2/\mathbb{Z}^2$,

$$\partial_t u_t + u_t \cdot \nabla u_t + \nabla p_t = -\theta u_t + F_t, \quad \nabla \cdot u_t = 0, \quad (1.1)$$

where u_t is the velocity vector field, p_t is the (scalar) pressure, $\theta > 0$ and F_t is a stochastic forcing term. Our motivation stems from works on 2-dimensional turbulence: our model can be regarded as an inviscid version of the one considered in [6], which aimed to describe the energy cascades phenomena in stationary, energy-dissipated, 2-dimensional turbulence. Inspired by recent renewed theoretical interest for point vortices methods in the study of 2-dimensional Euler's equation stemming from [11], we will study solutions to (1.1) obtained as systems of interacting point vortices, and Gaussian limits of the latter ones. Even if our models are not able to capture turbulence phenomena such as the celebrated energy spectrum decay law of inverse cascade predicted by Kolmogorov, we believe that the mechanism of creation and damping of point vortices we describe might contribute to provide a description of experimental behaviours of models such as the ones in [6]. Moreover, the mathematical treatment of measure- or distribution-valued solution to Euler's equation is not a trivial task, due to the need of quite weak notions of solution in presence of a singular nonlinearity.

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From the mathematical viewpoint, equation (1.1) has been widely investigated especially as inviscid limit of driven and damped Navier-Stokes equation, see for instance [5], [8] and references therein. Aside from the fact that we are dealing directly with the inviscid case, a substantial difference of this work with respect to those ones is the space regularity of solutions. Indeed, existence and uniqueness for 2-dimensional Euler equations are well established facts in spaces of suitably regular function spaces, while the interesting case of solutions taking values in signed measures or distributions remains quite open, especially in the uniqueness part: we refer to [22] for a general overview of the theory. The results of [11], which we review in subsection 2.3, established an important link between the theory of point vortices models and Gaussian invariant measures to Euler's equation. We refer to [22, 23] and to [1] for reviews on, respectively, the former and latter ones. We also mention that limits of Gibbsian point vortices ensembles (originally proposed by Onsager, [26]) converging to Gaussian invariant measures were already considered for instance in [4] (the similarities between the two being already pointed out by Kraichnan [21]). However, Flandoli [11] was the first, as far as we know, to prove convergence of the system evolving in time, as opposed to the simple convergence of invariant measures of the other ones. His approach was based on the weak vorticity formulation of [30] (see also its references), which had already been considered in the point vortices model, [31], and turned out to be suitable to treat solutions with white noise marginals. Further developments include the study of limits of point vortices ensembles in the Canonical and Microcanonical ensembles, see [18] and [15] respectively, and point vortices approximation to Navier-Stokes equations, [14].

Our results generalise the ones of [11] by combining a stochastic forcing term (already considered in the vortices setting in [13], or in function spaces in [7]) and damping. Stationary solutions are regarded with particular interest in the theory, and the invariant distributions we consider are also invariants of Euler's equation with no damping of forcing (see Theorem 2.15): to our knowledge, the Poissonian invariant distributions with infinite vortices we introduce below are new, while their Gaussian counterpart (the *enstrophy measure*, more generally known as white noise) have been an object of interest since the works of Hopf [19]. For a more general discussion on invariant measures of Lévy type we refer to [2], in which most of the basic ideas we rely upon are finely presented, although their arguments then proceed along point of view of Dirichlet forms theory.

We will treat our model equation in vorticity form,

$$\partial_t \omega_t = -\theta \omega_t + u_t \cdot \nabla \omega_t + \Pi_t, \quad \omega_t = \nabla^\perp \cdot u_t, \quad (1.2)$$

where $\nabla^\perp = (\partial_2, -\partial_1)$. The idea is to exhibit solutions by adapting the point vortices model for Euler's equation, which, in absence of forcing and damping, we recall to be the measure valued solutions

$$\omega_t = \sum_{i=1}^N \xi_i \delta_{x_{i,t}}, \quad \dot{x}_{i,t} = \sum_{j \neq i} \xi_j \nabla^\perp \Delta^{-1}(x_{i,t}, x_{j,t}), \quad (1.3)$$

where $x_i \in \mathbb{T}^2$ is the position and $\xi_i \in \mathbb{R}$ the intensity of a vortex, to Euler's equation

$$\begin{cases} \partial_t \omega_t + u_t \cdot \nabla \omega_t = 0 \\ \nabla^\perp \cdot u_t = \omega_t, \end{cases} \quad (1.4)$$

(see section 2 for the appropriate notion of solution). Inclusion of the damping term in our model will amount to an exponential quenching of the vortex intensities, with rate θ .

Because of dissipation due to friction (which physically results from the 3-dimensional environment in which the 2-dimensional flow is embedded), a forcing term is necessary in order for the model to exhibit stationary behaviour. We will choose as Π_t a Poisson point process, so to add new vortices and rekindle the system. The linear part of (1.2), which is a Poissonian Ornstein-Uhlenbeck equation, suggests that stationary distributions are made of countable vortices with exponentially decreasing intensity, but in fact dealing with solutions of (1.2) having such marginals seems to be as hard as the white noise marginals case. The latter will be also addressed, taking as in [11] a “central limit” scaling of the vortices model, resulting in solutions of (1.2) with space white noise marginal, and space-time white noise as forcing term.

Our main result will be the *existence* of solutions to (1.2) in these two cases: infinite vortices marginals and Poisson point process forcing; white noise marginals and space-time white noise forcing. The latter one draws us closer to the models in [6], where the forcing term was Gaussian with delta time-correlations. We will apply a compactness method: our approximant processes will not be approximated solutions (as in Faedo-Galerkin methods), but true point vortices solutions with finitely many vortices, for which we are able to prove well-posedness thanks to the techniques of [23].

We regard the following results as a first step in the analysis of equation (1.2) by point vortices methods, the natural prosecution being the study of driving noises with more complicated space correlations, such as the ones used in numerical simulations reviewed in [6].

2 Preliminaries and main result

Consider the 2-dimensional torus $\mathbb{T}^2 = \mathbb{R}^2/\mathbb{Z}^2$, a metric space with the distance $d(x, y) = \min_{k \in \mathbb{Z}^2} |x - y - k|_{\mathbb{R}^2}$, $|\cdot|_{\mathbb{R}^2}$ being the usual Euclidean distance. We denote by $H^\alpha = H^\alpha(\mathbb{T}^2) = W^{\alpha, 2}(\mathbb{T}^2)$, for $\alpha \in \mathbb{R}$ the $L^2 = L^2(\mathbb{T}^2)$ -based Sobolev spaces, which enjoy the compact embeddings $H^\alpha \hookrightarrow H^\beta$ whenever $\beta < \alpha$, the injections being furthermore Hilbert-Schmidt if $\alpha > \beta + 1$. Sobolev spaces are conveniently represented in terms of Fourier series: let $e_k(x) = e^{2\pi i k \cdot x}$, $x \in \mathbb{T}^2$, $k \in \mathbb{Z}^2$, be the usual Fourier orthonormal basis: then

$$H^\alpha = \left\{ u(x) = \sum_{k \in \mathbb{Z}^2} \hat{u}_k e_k(x) : \|u\|_{H^\alpha}^2 = \sum_{k \in \mathbb{Z}^2} (1 + |k|^2)^\alpha |\hat{u}_k|^2 < \infty \right\}, \quad (2.1)$$

where $\hat{u}_k = \bar{\hat{u}}_{-k} \in \mathbb{C}$ (we only consider real spaces). We will often consider functions of two space variables, *i.e.* functions on $\mathbb{T}^2 \times \mathbb{T}^2 = \mathbb{T}^{2 \times 2}$, and denote by $L^2_{sym}(2 \times 2)$ the space of symmetric square integrable functions,

$$L^2_{sym}(2 \times 2) = \{f \in L^2(2 \times 2) : f(x, y) = f(y, x) \forall x, y \in \mathbb{T}^2\}.$$

Analogously, $H_{sym}^\alpha(\mathbb{T}^{2 \times 2})$, $\alpha \in \mathbb{R}$ will be the $L^2_{sym}(2 \times 2)$ -based Sobolev space of symmetric functions.

We denote by $\mathcal{M} = \mathcal{M}(\mathbb{T}^2)$ the space of finite signed measures on \mathbb{T}^2 : recall that measures have Sobolev regularity $\mathcal{M} \subset H^{-1-\delta}$ for all $\delta > 0$ (for instance, because by dominated convergence their Fourier coefficients converge to constants). The brackets $\langle \cdot, \cdot \rangle$ will stand for L^2 duality couplings, such as the one between measures and continuous functions, or between Sobolev spaces of opposite orders, unless we specify otherwise.

The capital letter C will denote (possibly different) constants, and subscripts will point out occasional dependences of C on other parameters. Lastly, we write $X \sim Y$ when the random variables X, Y have the same law.

2.1 Random variables

In order to lighten notation, in this paragraph we denote random variables (or stochastic processes) and their laws with the same symbols. Let us also fix $H := H^{-1-\delta}$, with $\delta > 0$, the Sobolev space in which we embed our random measures and distributions. We will deal with stochastic objects of Gaussian and Poissonian nature: the former are likely to be the more familiar ones, so we begin our review with them. We refer to [27, 29] for a complete discussion of the underlying classical theory.

Let W_t be the cylindrical Wiener process on $L^2(\mathbb{T}^2)$, that is $\langle W_t, f \rangle$ is a real-valued centred Gaussian process indexed by $t \in [0, \infty)$ and $f \in L^2(\mathbb{T}^2)$ with covariance

$$\mathbb{E}[\langle W_t, f \rangle, \langle W_s, g \rangle] = t \wedge s \langle f, g \rangle_{L^2(\mathbb{T}^2)} \quad (2.2)$$

for any $t, s \in [0, \infty)$ and $f, g \in L^2(\mathbb{T}^2)$. Since the embedding $L^2(\mathbb{T}^2) \hookrightarrow H^{-1-\delta}(\mathbb{T}^2)$ is Hilbert-Schmidt, W_t defines a $H^{-1-\delta}$ -valued Wiener process. The law η of W_1 is called the *white noise* on \mathbb{T}^2 , and it can thus be regarded as a Gaussian probability measure on $H^{-1-\delta}$. Analogously, the law ζ of the (distributional) time derivative of W can be identified both with a centred Gaussian process indexed by $L^2([0, \infty) \times \mathbb{T}^2)$ and identity covariance operator or with a centred Gaussian probability measure on $H^{-3/2-\delta}([0, \infty) \times \mathbb{T}^2)$; ζ is called the *space-time white noise* on \mathbb{T}^2 .

Besides those Gaussian distributions, we will be interested in a number of Poissonian variables, which we now define in the framework of [27]. For $\lambda > 0$, let π^λ be the Poisson random measure on $[0, \infty) \times H^{-1-\delta}$ with intensity measure ν given by the product of the measure λdt on $[0, \infty)$ and the image of $\sigma \delta_x$ where $\sigma = \pm 1$ and $x \in \mathbb{T}^2$ are chosen uniformly at random. In other terms, one can define the compound Poisson process on $H^{-1-\delta}$ (in fact on \mathcal{M}),

$$\Sigma_t^\lambda = \sum_{i: t_i \leq t} \sigma_i \delta_{x_i} = \int_0^t d\pi^\lambda, \quad (2.3)$$

starting from the jump times t_i of a Poisson process of parameter λ , a sequence σ_i of i.i.d. ± 1 -valued Bernoulli variable of parameter $1/2$ and a sequence x_i of i.i.d uniform variables on \mathbb{T}^2 . Notice that, since its intensity measure has 0 mean, π^λ is a compensated Poisson measure, or equivalently Σ_t^λ is a $H^{-1-\delta}$ -valued martingale. Moreover, Σ_t^λ has the same covariance of the cylindrical Wiener process W_t (up to the factor λ):

$$\mathbb{E}[\langle \Sigma_t^\lambda, f \rangle \langle \Sigma_s^\lambda, g \rangle] = \lambda(t \wedge s) \langle f, g \rangle_{L^2}^2, \quad (2.4)$$

and also the same quadratic variation,

$$[\langle \Sigma^\lambda, f \rangle]_t = \lambda t \|f\|_{L^2}^2. \quad (2.5)$$

We will need a symbol for another Poissonian integral, the $H^{-1-\delta}$ -valued (in fact \mathcal{M} -valued) variable

$$\Xi_M^{\lambda, \theta} = \sum_{i: t_i \leq M} \sigma_i e^{-\theta t_i} \delta_{x_i} = \int_0^M e^{-\theta t} d\pi^\lambda, \quad (2.6)$$

where $M, \theta > 0$. Thanks to the negative exponential, the above integrals converge also when $M = \infty$, defining a random measure: we will call it $\Xi^{\lambda, \theta} = \Xi_\infty^{\lambda, \theta}$.

Remark 2.1. By (2.6), a sample of the random measure $\Xi_M^{\lambda, \theta}$ is a finite sum of *point vortices* $\xi_i \delta_{x_i}$ with $\xi_i \in \mathbb{R}, x_i \in \mathbb{T}^2$. We will say that the random vector $(\xi_i, x_i)_{i=1\dots N} \in (\mathbb{R} \times \mathbb{T}^2)^N$ (with random length N) is sampled under $\Xi_M^{\lambda, \theta}$ if $\sum_{i=1}^N \xi_i \delta_{x_i}$ has the law of $\Xi_M^{\lambda, \theta}$. Analogously (and in a sense more generally speaking), the sequence $(t_i, \sigma_i, x_i)_{i \in \mathbb{N}}$ is sampled under π^λ if the sum of $\sigma_i \delta_{t_i} \delta_{x_i}$ has the law of the Poisson point process π^λ .

These Poissonian measures are characterised by their Laplace transforms: for any measurable and bounded $f : \mathbb{T}^2 \rightarrow \mathbb{R}$,

$$\mathbb{E} [\exp (\alpha \langle f, \Sigma_t^\lambda \rangle)] = \exp \left(\lambda t \int_{\{\pm 1\} \times \mathbb{T}^2} (e^{\alpha \sigma f(x)} - 1) d\sigma dx \right), \quad (2.7)$$

$$\mathbb{E} [\exp (\alpha \langle f, \Xi_M^{\lambda, \theta} \rangle)] = \exp \left(\lambda \int_{[0, M] \times \{\pm 1\} \times \mathbb{T}^2} (e^{\alpha \sigma e^{-\theta t} f(x)} - 1) dt d\sigma dx \right), \quad (2.8)$$

where $d\sigma$ denotes the uniform measure on ± 1 . By the isometry property of Poissonian integrals, the second moments of Σ_t^λ and $\Xi_M^{\lambda, \theta}$ are given by

$$\mathbb{E} [\|\Sigma_t^\lambda\|_{H^{-1-\delta}}^2] = C \lambda t, \quad \mathbb{E} [\|\Xi_M^{\lambda, \theta}\|_{H^{-1-\delta}}^2] = C \frac{\lambda}{\theta} (1 - e^{-\theta M}),$$

where $C = \|\delta\|_{H^{-1-\delta}}^2$ is the Sobolev norm of Dirac's delta.

An important link between the objects we have defined so far is the following:

Proposition 2.2 (Ornstein-Uhlenbeck process). *Consider the linear stochastic differential equation in $H^{-1-\delta}$*

$$du_t = -\theta u_t dt + d\Pi_t. \quad (2.9)$$

If $\Pi_t = \sqrt{\lambda} W_t$, there exists a unique stationary solution with invariant measure $\sqrt{\frac{\lambda}{2\theta}} \eta$, and if $u_0 \sim C\eta$ ($C > 0$), the invariant measure is approached exponentially fast, $u_t \sim \sqrt{\frac{\lambda}{2\theta}} (1 - e^{-2\theta t}) (1 - C^2) \eta$.

Analogously, if $\Pi_t = \Sigma_t^\lambda$, there exists a unique stationary solution with invariant measure $\Xi_\infty^{\theta, \lambda}$, and if $u_0 \sim \Xi_M^{\theta, \lambda}$, then u_t will have law $\Xi_{M+t}^{\theta, \lambda}$ for any later time $t > 0$.

The linear equation (2.9), in both the outlined cases, has a unique $H^{-1-\delta}$ -valued strong solution, with continuous trajectories in the Gaussian case, and *cadlag* trajectories in the Poissonian one. Well-posedness of the linear equation and uniqueness of the invariant measure are part of the classical theory (see [27]), and they descend from the explicit solution by stochastic convolution:

$$u_t = e^{-\theta t} u_0 + \int_0^t e^{-\theta(t-s)} d\Pi_s, \quad (2.10)$$

from which it is not difficult to derive also the last statement of the Proposition.

2.2 Stochastic double integrals

Let η be the space white noise on \mathbb{T}^2 as in the previous section. Considering η as a random distribution in $H^{-1-\delta}$, the tensor product $\eta \otimes \eta$ is defined as a distribution in $H^{-2-2\delta}(\mathbb{T}^{2 \times 2})$, so for $h \in H^{2+\delta}(\mathbb{T}^{2 \times 2})$ we can couple $\langle h, \eta \otimes \eta \rangle$.

The couplings of η against $L^2(\mathbb{T}^2)$ functions are only defined as Ito-Wiener integrals: *double* Ito-Wiener integrals play a crucial role in our discussion, so let us recall their definition (for which we refer to [20]). The double stochastic integral with respect to η is the isometry $I^2 : L^2_{sym}(\mathbb{T}^{2 \times 2}) \rightarrow L^2(\eta)$ (which is not onto, the image being the second Wiener chaos) defined as follows.

Lemma 2.3. *The following are equivalent:*

- I^2 is the extension by density in $L^2(\mathbb{T}^{2 \times 2})$ of the expression on symmetric products:

$$\forall f, g \in L^2(\mathbb{T}^2) \quad I^2(f \odot g) = : \langle \eta, f \rangle \langle \eta, g \rangle : = \langle \eta, f \rangle \langle \eta, g \rangle - \langle f, g \rangle \in L^2(\eta),$$

$$\text{with } f \odot g(x, y) = \frac{f(x)g(y) + f(y)g(x)}{2}.$$

- I^2 is the extension by density in $L^2(\mathbb{T}^{2 \times 2})$ of the map

$$\begin{aligned} L^2(\mathbb{T}^2) \ni h(x, y) &= \sum_{\substack{i_1, i_2 = 1, \dots, n \\ i_1 \neq i_2}} a_{i_1, i_2} \mathbf{1}_{A_{i_1} \times A_{i_2}}(x, y) \\ &\mapsto \sum_{\substack{i_1, i_2 = 1, \dots, n \\ i_1 \neq i_2}} a_{i_1, i_2} \eta(A_{i_1}) \eta(A_{i_2}) \in L^2(\eta), \end{aligned} \quad (2.11)$$

where $n \geq 0$, $A_1, \dots, A_n \subset \mathbb{T}^2$ are disjoint Borel sets and $a_{i,j} \in \mathbb{R}$, so that functions $h(x, y)$ of the form above vanish on the diagonal $x = y \in \mathbb{T}^2$.

- I^2 is the extension by density in $L^2(\mathbb{T}^{2 \times 2})$ of the map

$$\forall h \in C^\infty(\mathbb{T}^{2 \times 2}) : \forall x \in \mathbb{T}^2 h(x, x) = 0, \quad I^2(h) = \langle h, \eta \otimes \eta \rangle. \quad (2.12)$$

For any $h \in H_{sym}^{2+\delta}(\mathbb{T}^{2 \times 2})$ it holds as an equality between $L^2(\eta)$ variables

$$\langle h, \eta \otimes \eta \rangle = I^2(h) + \int_{\mathbb{T}^2} h(x, x) dx \quad (2.13)$$

(since it is true for the dense subset of symmetric products) where we remark that $\int_{\mathbb{T}^2} h(x, x) dx$ makes sense since h has a continuous version by Sobolev embedding. We thus see that Ito-Wiener integration corresponds to “subtract the diagonal contribution” to the tensor product. In order to make the dependence of double Ito-Wiener integrals on η , and motivated by the above discussion, we will use in the following the notation

$$: \langle h, \eta \otimes \eta \rangle := I^2(h).$$

Let $\lambda, \theta, M > 0$. In the Poissonian case, we can define double integrals against continuous functions $h \in C(\mathbb{T}^{2 \times 2})$ \mathbb{P} -almost surely as

$$\left\langle h, \Xi_M^{\lambda, \theta} \otimes \Xi_M^{\lambda, \theta} \right\rangle = \sum_{i, j: t_i, t_j \leq M} \sigma_i \sigma_j e^{-\theta(t_i + t_j)} h(x_i, x_j),$$

where x_i, σ_i, t_i are distributed as in the definition of $\Xi_M^{\lambda, \theta}$, (2.6). If we consider in the Poissonian case the third approximation procedure of Theorem 2.3, we obtain a different, “renormalised” Poissonian double integral.

Lemma 2.4. Let $\mathcal{A} \subset C(\mathbb{T}^{2 \times 2})$ be the family of continuous functions vanishing on the diagonal, $h(x, x) = 0$ for all $x \in \mathbb{T}^2$. Then for all $h \in \mathcal{A}$

$$\mathbb{E} \left[\left| \left\langle h, \Xi_M^{\lambda, \theta} \otimes \Xi_M^{\lambda, \theta} \right\rangle \right|^2 \right] = \frac{\lambda^2}{\theta} (1 - e^{-\theta M})^2 \|h\|_{L^2(\mathbb{T}^{2 \times 2})}^2. \quad (2.14)$$

As a consequence, the map $\mathcal{A} \ni h \mapsto \left\langle h, \Xi_M^{\lambda, \theta} \otimes \Xi_M^{\lambda, \theta} \right\rangle \in L^2(\Xi_M^{\lambda, \theta})$ extends by continuity to a map

$$L^2(\mathbb{T}^{2 \times 2}) \ni h \mapsto : \left\langle h, \Xi_M^{\lambda, \theta} \otimes \Xi_M^{\lambda, \theta} \right\rangle : \in L^2(\Xi_M^{\lambda, \theta})$$

which satisfies (2.14), and which is given, for functions $h \in L^2(\mathbb{T}^{2 \times 2})$ continuous outside the diagonal set $(x, x) : x \in \mathbb{T}^2 \subset \mathbb{T}^{2 \times 2}$, but possibly discontinuous or singular on it, by

$$: \left\langle h, \Xi_M^{\lambda, \theta} \otimes \Xi_M^{\lambda, \theta} \right\rangle : = \sum_{\substack{i, j: t_i, t_j \leq M \\ i \neq j}} \sigma_i \sigma_j e^{-\theta(t_i + t_j)} h(x_i, x_j). \quad (2.15)$$

The proof of the latter (as well as the one of Theorem 2.3) is a straightforward computation. In a sense, in the Poissonian case the “subtraction of diagonal contributions” is made even more evident than in the Gaussian case by (2.15), where the sum runs over distinct indices.

2.3 Weak solutions of 2D Euler equation

We now review some definitions of measure-valued and distribution-valued solutions to the 2D Euler's equation: the point is how to make sense of the multiplication appearing in the nonlinearity. The equation in terms of the vorticity ω is (1.4)

$$\begin{cases} \partial_t \omega_t + u_t \cdot \nabla \omega_t = 0 \\ \nabla^\perp \cdot u_t = \omega_t, \end{cases}$$

and it has to be complemented with boundary conditions: on the torus \mathbb{T}^2 one should impose that ω_t have zero average. However, since we are dealing with a conservation law, the space average is not involved in the dynamics (it is constant).

Remark 2.5. We will henceforth deliberately ignore the zero average condition: it will always be possible to subtract a constant number (constant in time and space, but possibly a random variable) to take care of it, but we refrain from doing so to avoid a superfluous notational burden.

Let G be the Green function of Δ on \mathbb{T}^2 with zero average, and let $K = \nabla^\perp G$ be the Biot-Savart kernel; the former has the representation

$$G(x, y) = -\frac{1}{2\pi} \log d(x, y) + g(x, y), \quad g(x, y) \in C_{sym}(\mathbb{T}^{2 \times 2}).$$

We will use the fact that $|\nabla G(x, y)|, |K(x, y)| \leq \frac{C}{d(x, y)}$ for all x, y , with C a universal constant. The second equation of (1.4) can be inverted by means of the Biot-Savart kernel: we can write $u_t = K * \omega_t$, and thus obtain an equation where only ω appears. Its integral form against a smooth test function f is

$$\langle f, \omega_t \rangle = \langle f, \omega_0 \rangle + \int_0^t \int_{\mathbb{T}^{2 \times 2}} K(x, y) \cdot \nabla f(x) \omega_s(x) \omega_s(y) dx dy ds \quad (2.16)$$

(keeping in mind that $\nabla \cdot \nabla^\perp \omega \equiv 0$ to perform integration by parts), which can be symmetrised (swapping x and y) into

$$\langle f, \omega_t \rangle = \langle f, \omega_0 \rangle + \int_0^t \int_{\mathbb{T}^{2 \times 2}} H_f(x, y) \omega_s(x) \omega_s(y) dx dy ds \quad (2.17)$$

where $H_f(x, y) = \frac{1}{2} K(x, y) (\nabla f(x) - \nabla f(y))$ is a bounded symmetric function, smooth outside the diagonal set

$$\Delta^2 := \{(x, x) \in \mathbb{T}^{2 \times 2}\}.$$

These three formulations are equivalent for smooth ω_t , but the integral forms, especially the symmetrised one, have been used to define more general solutions of Euler's equation, see [30]. One such solution is the system of (finitely many) *Euler's point vortices*: the evolution of the vorticity $\omega_t = \sum_{i=1}^N \xi_i \delta_{x_{i,t}}$ (with $\xi_i \in \mathbb{R}$ and $x_i \in \mathbb{T}^2$) is given by (1.3),

$$\dot{x}_{i,t} = \sum_{j \neq i} \xi_j K(x_{i,t}, x_{j,t}).$$

This model is thoroughly discussed for instance in [23], where it is remarked that it satisfies (2.16) if the double space integral is taken outside the diagonal Δ^2 , where K is singular:

$$\begin{aligned} & \int_{\mathbb{T}^{2 \times 2} \setminus \Delta^2} K(x, y) \cdot \nabla f(x) \omega_s(x) \omega_s(y) dx dy \\ &= \sum_{i \neq j} \xi_i \xi_j K(x_i, x_j) \cdot \nabla f(x_i) = \sum_{i \neq j} \xi_i \xi_j H_f(x_i, x_j) \\ &= \int_{\mathbb{T}^{2 \times 2} \setminus \Delta^2} H_f(x, y) \omega_s(x) \omega_s(y) dx dy. \end{aligned}$$

It is thus possible, by slightly abusing the notation we introduced in (2.15), extending it to this deterministic point distribution, to formulate Euler's equation in the point vortices case as follows: if $\omega_t = \sum_{i=1}^N \xi_i \delta_{x_i, t}$ and we denote

$$:\langle H_f, \omega_t \otimes \omega_t \rangle := \sum_{i \neq j}^N \xi_i \xi_j H_f(x_i, x_j),$$

then it holds

$$\langle f, \omega_t \rangle = \langle f, \omega_0 \rangle + \int_0^t :\langle H_f, \omega_s \otimes \omega_s \rangle : ds. \quad (2.18)$$

The need to avoid the diagonal set Δ^2 in order to give meaning to singular solutions is going to be crucial in what follows, as it is in the proof of the forthcoming important well-posedness result.

Proposition 2.6 (Dürr-Pulvirenti). *Let $\xi_1, \dots, \xi_n \in \mathbb{R}$ and $x_1, \dots, x_n \in \mathbb{T}^2$. For almost every initial data $x_{1,0}, \dots, x_{n,0} \in \mathbb{T}^2$ under the n -fold product of Lebesgue's measure, the system of differential equations (1.3) has a smooth, global in time solution $x_{1,t}, \dots, x_{n,t}$, which preserves the product measure on the initial condition. The measure-valued process $\omega_t = \sum_{i=1}^n \xi_i \delta_{x_i}$ then satisfies (2.18) in the sense above.*

(In fact the latter is a slight generalisation of the results in [10, 23], which will be a consequence of the further generalisation we will prove in section 3.)

In [11], Flandoli performed a scaling limit of the point vortices system to exhibit (stationary) solutions with space white noise marginals: the meaning of the equation for such irregular vorticity processes was understood by carrying to the limit the formulation (2.18), since, as we have seen in the last paragraph, the Wiener-Ito interpretation of the nonlinear term makes perfect sense in the case of white noise. To proceed rigorously, let us give the following:

Definition 2.7. *Let $(\omega_t)_{t \in [0, T]}$ be a $H^{-1-\delta}$ -valued continuous stochastic process defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, with fixed time marginals ω_t having the law of white noise η for all $t \in [0, T]$. We say that ω is a weak solution to Euler's equation if for any $f \in C^\infty(\mathbb{T}^2)$, \mathbb{P} -almost surely, for any $t \in [0, T]$,*

$$\langle f, \omega_t \rangle = \langle f, \omega_0 \rangle + \int_0^t :\langle H_f, \omega_s \otimes \omega_s \rangle : ds, \quad (2.19)$$

the coupling $:\langle H_f, \omega_s \otimes \omega_s \rangle :$ being defined as in Theorem 2.3 since $\omega_s \sim \eta$.

Remark 2.8. Notice that the Ito-Wiener integrals (in space) appearing in the definition are almost surely integrable in time since their $L^2(\mathbb{P})$ norms are uniformly bounded in t . The latter definition coincides with the one of [11], only, in that article, it was not observed that the approximation procedure used to define the nonlinear term in fact coincides with the classic Ito-Wiener integral.

The formulation (2.19) is in fact quite general: interpreting the colons as "subtraction of the diagonal contribution", this formulation might include all deterministic solutions, both in the classical and weak formulation (2.17) (cf. [30]), the point vortices solution of Theorem 2.6, and it is the sense in which the limit process with white noise marginals of [11] solves Euler's equation.

Proposition 2.9 (Flandoli). *There exists a stationary stochastic process ω_t with fixed-time marginals $\omega_t \sim \eta$ and trajectories of class $C([0, T], H^{-1-\delta})$ for any $\delta > 0$ which is a solution of Euler's equation in the sense of Theorem 2.7.*

Remark 2.10. In fact, [11] proves more generally existence of non-stationary solutions with fixed time marginals absolutely continuous with respect to η . The above definition

of weak solution does not encompass this case: we refer to [11, Section 2.5]. A thorough discussion of non-stationary solutions in such regimes can be found in [12].

For the sake of completeness, we recall that solutions to Euler's equation with white noise marginals were first built in [3], by means of Galerkin approximation on \mathbb{T}^2 .

Remark 2.11. We already referred to [2] mentioning possible invariant measures of Euler's equations with Lévy distributed marginals more general than the Gaussian or Poissonian ones considered here. A noticeable obstacle to such generality is the fact that the renormalised double integrals appearing in the latter cases do not find easily an analogue in a context where, for instance, distributions might not have second order moments.

2.4 Main results

Fix $\lambda, \theta > 0$. Our model is the stochastic differential equation

$$d\omega = -\theta\omega dt + (K * \omega) \cdot \nabla\omega dt + d\Pi_t, \quad (2.20)$$

where $d\Pi_t$ is either the Poisson process $d\Sigma_t^\lambda$ or the space-time white noise dW_t . We have seen in Theorem 2.2 how the linear part of the equation behaves; the intuition provided by the point vortices system suggests that, thanks to the Hamiltonian form of the nonlinearity, the latter only contributes to "shuffle" the vorticity without changes to the fixed time statistics. This intuition can be motivated as follows. Since the point vortices system preserves the product Lebesgue measure, the system must preserve the Poissonian random measures $\Xi_M^{\lambda, \theta}$ we introduced in subsection 2.1, because the positions of vortices under those measures are uniformly, independently scattered (this fact will be rigorously proved in section 3 for $M < \infty$). Building Gaussian solutions by approximation with Poissonian ones thus must produce the same phenomenon. In other words, with an eye towards stationary solutions, we expect to be able to build a Poissonian stationary solution with $\omega_t \sim \Xi_\infty^{\theta, \lambda}$ in the case $\Pi_t = \Sigma_t^\lambda$, and a stationary Gaussian solution with $\omega_t \sim \sqrt{\frac{\lambda}{2\theta}}\eta$ in the case $\Pi_t = \sqrt{\lambda}W_t$.

Remark 2.12. These claims are deeply related with the fact that 2D Euler's equation preserves enstrophy, $\int_{\mathbb{T}^2} \omega(x)^2 dx$, when smooth solutions are considered. The quadratic form associated to enstrophy, that is the $L^2(\mathbb{T}^2)$ product, is (up to multiplicative constants) the covariance of random fields $\Xi_M^{\lambda, \theta}$ and η : as already remarked in [2], one should expect all random fields with such covariance to be invariant for Euler's equation, even if the very meaning of the latter sentence has to be clarified.

First and foremost, we need to specify a suitable concept of solution: inspired by the discussion of the last paragraph, we give the following one.

Definition 2.13. Fix $T, \delta > 0$, and let $(\Omega, \mathcal{F}, \mathbb{P}, \mathcal{F}_t)$ be a probability space with a filtration \mathcal{F}_t satisfying the usual hypothesis, and let $(\omega_t)_{t \in [0, T]}$ be a $H^{-1-\delta}$ -valued \mathcal{F}_t -predictable process, with trajectories of class

$$L^2([0, T], H^{-1-\delta}) \cap \mathbb{D}([0, T], H^{-3-\delta}) \quad (2.21)$$

($\mathbb{D}([0, T], S)$ denotes the space of S -valued cadlag functions into a metric space S). On $(\Omega, \mathcal{F}, \mathbb{P}, \mathcal{F}_t)$ we also consider a $H^{-1-\delta}$ -valued \mathcal{F}_t -martingale $(\Pi_t)_{t \in [0, T]}$. We consider the following cases: for $\theta, \lambda > 0$,

- (P) $\Pi_t = \Sigma_t^\lambda$ and $\omega_t \sim \Xi_{M+t}^{\lambda, \theta}$ (defined respectively in (2.3) and (2.6)) for all $t \in [0, T]$, with $0 \leq M < \infty$;
- (Ps) $\Pi_t = \Sigma_t^\lambda$ and $\omega_t \sim \Xi_\infty^{\lambda, \theta}$ for all $t \in [0, T]$;
- (G) $\Pi_t = \sqrt{\lambda}W_t$ and $\omega_t \sim \sqrt{\frac{\lambda}{2\theta}(1 - e^{-2\theta(M+t)})}\eta$ for all $t \in [0, T]$, with $0 \leq M < \infty$;

(Gs) $\Pi_t = \sqrt{\lambda}W_t$ and $\omega_t \sim \sqrt{\frac{\lambda}{2\theta}}\eta$ for all $t \in [0, T]$.

We say that $(\Omega, \mathcal{F}, \mathbb{P}, \mathcal{F}_t, \Pi_t, \omega_0, (\omega_t)_{t \in [0, T]})$ is a weak solution of (2.20) if for any $f \in C^\infty(\mathbb{T}^2)$ it holds \mathbb{P} -almost surely for any $t \in [0, T]$:

$$\langle f, \omega_t \rangle = e^{-\theta t} \langle f, \omega_0 \rangle + \int_0^t e^{-\theta(t-s)} : \langle H_f, \omega_s \otimes \omega_s \rangle : ds + \int_0^t e^{-\theta(t-s)} \langle f, d\Pi_s \rangle, \quad (2.22)$$

where $: \langle H_f, \omega_s \otimes \omega_s \rangle :$ is defined as in Theorem 2.4 in cases (P), (Ps) and as in Theorem 2.3 in (G), (Gs). If instead, given $(\Omega, \mathcal{F}, \mathbb{P}, \mathcal{F}_t, W_t)$ there exists a process ω_t as above, we call it a strong solution.

Remark 2.14. Equation (2.22) is motivated in sight of (2.10) and (2.19). The “variation of constants” expression in the above definition is equivalent to the “integral” one

$$\langle f, \omega_t \rangle = \langle f, \omega_0 \rangle - \theta \int_0^t \langle f, \omega_s \rangle ds + \int_0^t : \langle H_f, \omega_s \otimes \omega_s \rangle : ds + \langle f, \Pi_t \rangle, \quad (2.23)$$

as one can verify integrating by parts in time. Both versions will be useful in what follows, but we deem (2.22) more suggestive.

Remark 2.15. The nonlinear term of (2.22) is well-defined thanks to the isometry properties of Gaussian and Poissonian double integral (see section 2): indeed, the integrand is bounded in $L^2(\mathbb{P})$ uniformly in time, so that, in particular, $\int_0^t : \langle H_f, \omega_s \otimes \omega_s \rangle : ds$ is a continuous function of time.

We are now able to state our main result.

Theorem 2.16. There exist weak solutions of (2.20) in all the outlined cases, stationary (as $H^{-1-\delta}$ -valued stochastic processes) in the cases (Ps) and (Gs).

As already remarked, equation (2.20) is difficult to deal with directly in the Gaussian (or even the stationary Poisson) case: for instance it does not seem possible to treat it with fixed point or semigroup techniques. We prove existence of stationary solutions by taking limits of point vortices solutions, corresponding to the case (P). We begin with a solution ω_M of the equation (2.20) with noise Σ_t^λ starting from finitely many vortices distributed as $\Xi_M^{\theta, \lambda}$. Well-posedness in this case is ensured by a generalisation of Theorem 2.6, whose proof is the content of section 3. The first limit we consider is $M \rightarrow \infty$, so to build a stationary solution with invariant measure $\Xi^{\theta, \lambda}$ and thus obtain existence in case (Ps). Scaling intensities $\sigma \rightarrow \frac{\sigma}{\sqrt{N}}$ and generation rate $\lambda \rightarrow N\lambda$, we prove that as $N \rightarrow \infty$ the limit points are stationary solutions of (2.20) driven by space-time white noise and with invariant measure the space white noise. The nonstationary Gaussian case (G) will be derived analogously, in this sort of central limit theorem.

We are applying a *compactness method*: first, we prove probabilistic bounds on the involved distribution, in order to -second step- apply a compactness criterion ensuring tightness of the approximating processes; finally, we pass to the limit the equation satisfied by the approximants.

Remark 2.17. Consider the case when no damping or forcing are present: we noted above that the classical finite vortices system (1.3) preserves the product Lebesgue’s measure, so in particular the distributions $\Xi_M^{\theta, \lambda}$ with $M < \infty$ and $\theta, \lambda > 0$ are also invariant. The very same limiting procedure we are going to use, as $M \rightarrow \infty$, proves existence of stationary solutions to Euler’s equation in its weak formulation (2.19) with invariant measure $\Xi_\infty^{\theta, \lambda}$ (or η , the case of [11]), where the definition of solution is to be given in the fashion of Theorem 2.7. More generally, Poissonian and Gaussian stationary solutions, as suggested in [2], should be particular cases of stationary solutions with independently scattered random distributions.

3 Solutions with finitely many vortices

Even in the case of initial data distributed as $\Xi_M^{\lambda,\theta}$, that is with almost surely finitely many initial vortices, solving the nonlinear equation

$$d\omega = -\theta \omega dt + (K * \omega) \cdot \nabla \omega dt + d\Sigma_t^\lambda \quad (3.1)$$

is not a trivial task. We will build a solution describing explicitly how the initial vortices and the ones added by the noise term evolve, as a system of increasingly numerous differential equations for the positions of vortices x_i . Intuitively, the process $\omega_{M,t}$ is defined as follows: from the initial datum $\omega_M(0)$, which is sampled under $\Xi_M^{\theta,\lambda}$, we let the system evolve according to the deterministic dynamics

$$\dot{x}_i = \sum_{j \neq i} \xi_j e^{-\theta t} K(x_i, x_j)$$

until the first jump time t_1 of the driving noise Σ_t^λ , when we add the vortex corresponding to the jump, and so on. To treat the model rigorously, let us introduce the following notation: let $x_{1,0}, \dots, x_{n,0}$ and $\xi_{1,0}, \dots, \xi_{n,0}$ be the (random) positions and signs of vortices of the initial datum, and set for notational convenience $t_1 = \dots = t_n = 0$ their birth time; at time t_i it is added a vortex with intensity $\xi_{i,t_i} = \pm 1$ in the position x_{i,t_i} , but we can pretend it to actually have existed since time 0, and just come into play at the time t_i . Thus, our equations are

$$x_{i,t} = x_{i,t_i} + \mathbf{1}_{t_i \leq t} \int_{t_i}^t \sum_{j \neq i: t_j \leq s} \xi_{j,s} K(x_{i,s}, x_{j,s}) ds, \quad (3.2)$$

$$\xi_{i,t} = \begin{cases} \xi_{i,0} & t < t_i \\ e^{-\theta(t-t_i)} \xi_{i,0} & t \geq t_i \end{cases}. \quad (3.3)$$

In this formulation of the problem, part of the randomness consists in the positions and intensities of the initial vortices and the ones to be: the random jump times t_i then determine when the latter ones become part of the system. Let us thus fix the t_i 's (that is, condition the process given the distribution of the t_i 's) so to reduce us to a deterministic problem with random initial data. The existence of a solution for almost every initial condition is ensured by the following generalisation of Proposition 2.6.

Proposition 3.1. *Let $(x_{i,0})_{i \in \mathbb{N}}$ be a sequence of i.i.d uniform variables on \mathbb{T}^2 . For every locally finite sequence of jump times $0 \leq t_1 \leq \dots \leq t_i \leq \dots \leq \infty$ and initial intensities $(\xi_{i,0}) \in [-1, 1]$ the system of equations (3.2) and (3.3) possesses a unique, piecewise smooth and continuous, global in time solution, for a full probability set which does not depend on the choice of $t_i, \xi_{i,0}$. At any time, the joint law of positions x_i is the infinite product of Lebesgue measure on \mathbb{T}^2 .*

We use the hypothesis that the jump times t_i are locally finite (there are only finitely many of them in every compact $[0, T]$) so to reduce ourselves to a system of finitely many vortices. In fact, we repeat the proof of [10, 23] adapting it to our context. The issue is the possibility of collapsing vortices, which is ruled out as follows. We define an approximating system with interaction kernel smoothed in a ball around 0: the smooth interaction readily gives well-posedness of the approximants, on which we evaluate a Lyapunov functional measuring how close the vortices can get. Bounding the Lyapunov function then ensures that as the regularisation parameter goes to 0, the approximant vortices in fact perform the same motion prescribed by the non-smoothed equation.

Proof. Let $\delta > 0$, and consider smooth functions G_δ coinciding with G outside the fattened diagonal $\{(x, y) \in \mathbb{T}^{2 \times 2} : d(x, y) < \delta\}$ (d being the distance on the torus \mathbb{T}^2), and

such that

$$|G_\delta(x, y)| \leq C|G(x, y)|, \quad |\nabla G_\delta(x, y)| \leq \frac{C}{d(x, y)} \quad \forall x, y \in \mathbb{T}^2. \quad (3.4)$$

Note in particular that the latter inequality was already true for G . Let us first restrict ourselves to a time interval $[0, T]$: in particular, we can consider only the finitely many vortices with $t_i \leq T$, let them be x_1, \dots, x_n . Notice that smoothing K does not effect the evolution of the intensities $\xi_{i,t}$.

Thanks to Cauchy-Lipschitz theorem, the system with smoothed interaction kernel $K_\delta = \nabla^\perp G_\delta$ has a unique smooth solution $x_{i,t}^\delta$ for $t \in [0, t_1]$. The time derivative $\dot{x}_{i,t}^\delta$ is not right-continuous at $t = t_1$, but on $(t_1, t_2]$ is again smooth, so we can extend the unique solution applying Cauchy-Lipschitz in $[t_1, t_2]$ starting from x_{i,t_1}^δ ; notice that the resulting solution on $[0, t_2]$ is continuous, although not differentiable at t_1 . Proceeding as such we extend well-posedness to all $t \geq 0$.

Because of the Hamiltonian structure of the equations, that is, since $K_\delta = \nabla^\perp G_\delta$, it holds $\operatorname{div} \dot{x}_{i,t}^\delta = 0$ for any $t \neq t_1, \dots, t_n$. As a consequence, by Liouville's theorem (see for instance [9, Section 2.2]) the flow is measure preserving on intervals $(t_i, t_{i+1}]$, where it is smooth. But we have seen that the solution $x_{i,t}^\delta$ is given by a composition of such transformations, so that the product Lebesgue measure is preserved at all times.

Let us now introduce a Lyapunov function measuring how close the existing vortices are by means of G_δ :

$$L_\delta(t) = L_\delta(t, x_{1,t}^\delta, \dots, x_{n,t}^\delta) = - \sum_{i \neq j: t_i, t_j \leq t} G_\delta(x_{i,t}^\delta, x_{j,t}^\delta).$$

By replacing G_δ with $G_\delta - k$ for a large enough $k > 0$ in the definition of L_δ we can assume that L_δ is nonnegative. Observe that, because of (3.4), $\int_{\mathbb{T}^{2 \times n}} L_\delta(0) dx_1, \dots, dx_n \leq C$ for a constant C independent of δ . Upon differentiating, and keeping in mind that

$$\dot{x}_{i,t}^\delta = \sum_{j \neq i: t_j < t} \xi_{j,t} \nabla^\perp G_\delta(x_{i,t}^\delta, x_{j,t}^\delta), \quad \forall t > t_i, t \neq t_1, \dots, t_n,$$

(again, the flow is continuous but only differentiable away from jump times), we obtain, for all $t \neq t_1, \dots, t_n$,

$$\begin{aligned} \frac{d}{dt} L_\delta(t) &= - \sum_{i \neq j: t_i, t_j \leq t} \nabla G_\delta(x_{i,t}^\delta, x_{j,t}^\delta) \cdot (\dot{x}_{i,t}^\delta + \dot{x}_{j,t}^\delta) \\ &= \sum_{i,j,k \leq n} \tilde{a}_{ijk}(t) \nabla G_\delta(x_{i,t}^\delta, x_{j,t}^\delta) \cdot \nabla^\perp G_\delta(x_{i,t}^\delta, x_{k,t}^\delta), \end{aligned}$$

where $\tilde{a}_{ijk}(t)$ depend on time t as functions of the intensities $\xi_{i,t}$, $\tilde{a}_{ijk} = 0$ whenever two indices are equal, since $\nabla G_\delta(x_{i,t}^\delta - x_{j,t}^\delta) \cdot \nabla^\perp G_\delta(x_{i,t}^\delta - x_{j,t}^\delta) = 0$ and it always holds $|\tilde{a}_{ijk}(t)| \leq 1$. As a consequence, and using the fact that the solution $x_{i,t}^\delta$ is continuous, we have

$$L_\delta(t) = L_\delta(0) + \sum_{i,j,k \leq n} \int_0^t \tilde{a}_{ijk}(s) \nabla G_\delta(x_{i,s}^\delta, x_{j,s}^\delta) \cdot \nabla^\perp G_\delta(x_{i,s}^\delta, x_{k,s}^\delta) ds.$$

We can use this to prove the following integral bound on L_δ : denoting by dx^n the n -fold

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Lebesgue measure of the distribution of initial position,

$$\begin{aligned}
\int_{\mathbb{T}^{2 \times n}} \sup_{t \in [0, T]} L_\delta(t) dx^n &\leq \int_{\mathbb{T}^{2 \times n}} L_\delta(0) dx^n \\
&+ \sum'_{i,j,k} \int_0^T \int_{\mathbb{T}^{2 \times n}} |\nabla G_\delta(x_{i,s}^\delta, x_{j,s}^\delta) \cdot \nabla^\perp G_\delta(x_{i,s}^\delta, x_{k,s}^\delta)| dx^n ds \\
&\leq \int_{\mathbb{T}^{2 \times n}} L_\delta(0) dx^n + T \sum'_{i,j,k} \int_{\mathbb{T}^{2 \times n}} |\nabla G_\delta(x_i, x_j) \cdot \nabla^\perp G_\delta(x_i, x_k)| dx^n \\
&\leq \int_{\mathbb{T}^{2 \times n}} L_\delta(0) dx^n + TC_n \int_{\mathbb{T}^{2 \times 3}} |\nabla G_\delta(x, y) \cdot \nabla^\perp G_\delta(x, z)| dx dy dz \leq C_T,
\end{aligned}$$

C_T being a constant depending only on T (n depends on T). In the second and third lines, \sum' denotes summation over indices $i, j, k = 1, \dots, n$ such that no pair of them coincide. In the second inequality we have used the invariance of Lebesgue's measure, in the third one the fact that summation is over distinct indices and in the last step the aforementioned integrability of $L_\delta(0)$ and the fact that, because of (3.4), the integrands in the second term are bounded by

$$|\nabla G_\delta(x - y) \cdot \nabla^\perp G_\delta(x - z)| \leq \frac{C}{|x - y||x - z|}.$$

With these estimates at hand, we can now pass to the limit as $\delta \rightarrow 0$: let

$$d_{\delta,T}(x^n) = \min_{t \in [0, T]} \min_{i \neq j} d(x_{i,t}^\delta - x_{j,t}^\delta),$$

so that

$$d_{\delta,T}(x^n) < \delta \Rightarrow \sup_{t \in [0, T]} L_\delta(t) > -C \log(\delta),$$

since when two points x, y are closer than δ , $G_\delta(x, y) \geq C \log(\delta)$ for some universal constant C . As a consequence, by Čebyšev's inequality,

$$\mathbb{P}(\Omega_{\delta,T}) := \mathbb{P}(d_{\delta,T}(x^n) < \delta) \leq C'(-\log \delta)^{-1}.$$

By construction, in the event $\Omega_{\delta,T}^c$ the solution $x_{i,t}^\delta$ is in fact a solution of the original system in $[0, T]$. Hence, the thesis holds if the event

$$\bar{\Omega} = \bigcup_{T>0} \bigcap_{\delta>0} \Omega_{\delta,T}$$

is negligible. But this is true: $\Omega_{\delta,T}$ is monotone in its arguments, so that the intersection in δ is negligible because of the above estimates, hence the increasing union in T must be negligible too. \square

The forthcoming Corollary is a direct consequence of Proposition 3.1: indeed to complete our construction we only need to randomise the jump times and intensities so that the initial conditions and driving noise have the correct distribution. Assume that

- $(x_{1,0}, \xi_{1,0}), \dots, (x_{n,0}, \xi_{n,0})$ are positions and intensities of vortices sampled under $\Xi_M^{\theta, \lambda}$,
- $(t_{n+m}, x_{n+m,0}, \xi_{n+m,0} = \sigma_{n+m})_{m \geq 1}$ is sampled under π^λ ,

both in the sense of Theorem 2.1, with variables defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$. Then there exists a piecewise smooth, *cadlag* solution of the system of equations (3.2)

and (3.3) for all $t \in [0, \infty)$, \mathbb{P} -almost surely. Moreover, the positions of vortices at any time t , $x_{i,t}$, are i.i.d. uniform variables on the torus \mathbb{T}^2 .

Corollary 3.2. *In the outlined setting, the process $\omega_{M,t} = \sum_{i:t_i \leq t} \xi_{i,t} \delta_{x_{i,t}}$ is a \mathcal{M} -valued cadlag Markov process with fixed time marginals $\omega_{M,t} \sim \Xi_{M+t}^{\theta,\lambda}$ for all $t \geq 0$. It is a strong solution of*

$$d\omega_M = -\theta \omega_M dt + (K * \omega_M) \cdot \nabla \omega_M dt + d\Sigma_t^\lambda,$$

in the sense of Theorem 2.13

Proof. Fix $s < t$: by construction, given the positions $x_{i,0}$, the initial intensities $\xi_{i,0}$ and the jump times t_i (in a \mathbb{P} -full measure event), $\omega_{M,t}$ is given by a deterministic function of $(x_{i,s}, \xi_{i,s})_{i:t_i < s}$ and $(t_i, x_{i,0}, \xi_{i,0})_{i:s \leq t_i < t}$. As a consequence, $\omega_{M,t}$ is a function of $\omega_{M,s}$ and of the driving noise $(\Sigma_r^\lambda)_{s \leq r < t}$, which is independent from $\omega_{M,s}$: this implies the Markov property. Since the trajectories of positions $x_{i,t}$ and the evolution of intensities $\xi_{i,t}$ are smooth in time, $\omega_{M,t}$ is also smooth in time, save for the jump times t_i when a new Dirac's delta is added.

As for the marginal distributions, let us first evaluate:

$$\begin{aligned} \mathbb{E} \left[e^{i\alpha \langle \omega_{M,t}, f \rangle} \right] &= \mathbb{E} \left[\exp \left(i\alpha \sum_{i:t_i \leq t} \xi_{i,t} f(x_{i,t}) \right) \right] \\ &= \mathbb{E} \left[\mathbb{E} \left[\exp \left(i\alpha \sum_{i:t_i \leq t} \xi_{i,t} f(x_{i,t}) \right) \middle| (t_i)_{i \geq 0} \right] \right] \\ &= \mathbb{E} \left[\prod_{i:t_i \leq t} \int_{\mathbb{T}^2} e^{i\alpha \xi_{i,t} f(x)} dx \right] =: \mathbb{E} \left[\prod_{i:t_i \leq t} F(\xi_{i,t}) \right]. \end{aligned}$$

Using the definition of $\xi_{i,t}$, and distinguishing the cases $i \leq n$ and $i > n$ (which correspond to two independent groups of random variables), we can write

$$\begin{aligned} \mathbb{E} \left[e^{i\alpha \langle \omega_{M,t}, f \rangle} \right] &= \mathbb{E}_N \left[\prod_{s_i \in [0, M]} F(e^{-\theta s_i}) \right] \cdot \mathbb{E}_N \left[\prod_{s_i \in [0, t]} F(e^{-\theta(t-s_i)}) \right] \\ &= \mathbb{E}_N \left[\prod_{s_i \in [0, M+t]} F(e^{-\theta s_i}) \right] \end{aligned}$$

where N is a Poisson point process of parameter λ on \mathbb{R} whose points are denoted by s_i , and the second passage follows from the fact that the points N in disjoint intervals are independent and their distribution does not change if we reverse the parametrisation of the interval. Comparing to the characteristic function of $\Xi_{M+t}^{\theta,\lambda}$ given in (2.8), we conclude that $\omega_{M,t} \sim \Xi_{M+t}^{\theta,\lambda}$.

Observe now that in this case it holds, for any $f \in C^\infty(\mathbb{T}^2)$, \mathbb{P} -almost surely for all $t \geq 0$,

$$:\langle H_f, \omega_{M,t} \otimes \omega_{M,t} \rangle: = \sum_{\substack{i,j: t_i, t_j \leq t \\ i \neq j}} \xi_{i,t} \xi_{j,t} H_f(x_{i,t}, x_{j,t}),$$

(cf. with subsection 2.1). Given this, it is straightforward to show that we do have built

solutions of (2.22): for $f \in C^\infty(\mathbb{T}^2)$, by (3.2) and (3.3),

$$\begin{aligned}
 \langle f, \omega_{M,t} \rangle &= \sum_{i:t_i \leq t} \xi_{i,t} f(x_{i,t}) \\
 &= \sum_{i:t_i \leq t} \xi_{i,t} \left(f(x_{i,t_i}) + \int_{t_i}^t \sum_{j \neq i: t_j \leq s} \xi_{j,s} \nabla f(x_{i,s}) \cdot K(x_{i,s}, x_{j,s}) ds \right) \\
 &= \left(\sum_{i=1}^n + \sum_{i>n: t_i \leq t} \right) \xi_{i,t} f(x_{i,t_i}) + \sum_{i:t_i \leq t} \xi_{i,t} \int_{t_i}^t \sum_{j \neq i: t_j \leq s} \xi_{j,s} \nabla f(x_{i,s}) \cdot K(x_{i,s}, x_{j,s}) ds \\
 &= \sum_{i=1}^n e^{-\theta t} f(x_{i,0}) + \sum_{i>n: t_i \leq t} e^{-\theta(t-t_i)} f(x_{i,t_i}) \\
 &\quad + \int_0^t \sum_{i,j: t_i, t_j \leq s, i \neq j} e^{-\theta(t-s)} \xi_{i,s} \xi_{j,s} \nabla f(x_{i,s}) \cdot K(x_{i,s}, x_{j,s}) ds \\
 &= e^{-\theta t} \langle f, \omega_{M,0} \rangle + \int_0^t e^{-\theta(t-s)} \langle f, d\Sigma_s \rangle + \int_0^t : \langle H_f, \omega_{M,s} \otimes \omega_{M,s} \rangle : ds.
 \end{aligned}$$

The latter equation holds regardless of the choice of initial positions, intensities and jump times (as soon as the dynamics is defined) so in particular it holds \mathbb{P} -almost surely uniformly in t , and this concludes the proof. \square

The method of [23] thus provides, quite remarkably, existence and pathwise uniqueness of measure-valued strong solutions. Unfortunately, it only seems to apply to systems of finitely many vortices, since it relies on the very particular, discrete nature of the measures involved to control the “diagonal collapse” issue. We refer to [17] for further uniqueness results for point vortices systems obtained by means of refinements of the above techniques.

Let us conclude this section noting that we have obtained the first piece of Theorem 2.16, namely we have built solutions in the case (P) for all $M < \infty$.

4 Proof of the main result

In section 3 we built the point vortices processes $\omega_{M,t} = \sum_{i:t_i \leq t} \xi_{i,t} \delta_{x_{i,t}}$. Let us introduce the scaling in $N \geq 1$: we will denote $\omega_{M,N,t} = \sum_{i:t_i \leq t} \frac{\xi_{i,t}}{\sqrt{N}} \delta_{x_{i,t}}$ where $x_{i,t}, \xi_{i,t}$ solve equations (3.2) and (3.3), and where the t_i 's are the jump times of a real valued Poisson process of intensity $N\lambda$. In other words, by Theorem 3.2, $\omega_{M,N,t}$ is a strong solution of

$$d\omega_{M,N} = -\theta \omega_{M,N} dt + (K * \omega_{M,N}) \cdot \nabla \omega_{M,N} dt + \frac{1}{\sqrt{N}} d\Sigma_t^{N\lambda}, \quad (4.1)$$

(in the sense of Theorem 2.13) with fixed time marginals $\omega_{M,N,t} \sim \frac{1}{\sqrt{N}} \Xi_{M+t}^{\theta, N\lambda}$. It is worth to note here that, by construction of $\omega_{M,N,t}$, its natural filtration \mathcal{F}_t coincides with the one generated by the driving noise $\Sigma_t^{N\lambda}$ and the initial datum.

The forthcoming paragraphs deal with, respectively: a recollection of some compactness criterions, the bounds proving that the laws of $\omega_{M,N}$ are tight, the proof of the fact that limit points of our family of processes are indeed solutions in the sense of Theorem 2.13, that is, the main result.

4.1 Compactness results

Let us first review a deterministic compactness criterion due to Simon (we refer to [32] for the result and the required generalities on Banach-valued Sobolev spaces).

Proposition 4.1 (Simon). Assume that

- $X \hookrightarrow B \hookrightarrow Y$ are Banach spaces such that the embedding $X \hookrightarrow Y$ is compact and there exists $0 < \theta < 1$ such that for all $v \in X \cap Y$

$$\|v\|_B \leq M \|v\|_X^{1-\theta} \|v\|_Y^\theta;$$

- $s_0, s_1 \in \mathbb{R}$ are such that $s_\theta = (1-\theta)s_0 + \theta s_1 > 0$.

If $\mathcal{F} \subset W$ is a bounded family in

$$W = W^{s_0, r_0}([0, T], X) \cap W^{s_1, r_1}([0, T], Y)$$

with $r_0, r_1 \in [0, \infty]$, and we define

$$\frac{1}{r_\theta} = \frac{1-\theta}{r_0} + \frac{\theta}{r_1}, \quad s_* = s_\theta - \frac{1}{r_\theta},$$

then if $s_* \leq 0$, \mathcal{F} is relatively compact in $L^p([0, T], B)$ for all $p < -\frac{1}{s_*}$. In the case $s_* > 0$, \mathcal{F} is moreover relatively compact in $C([0, T], B)$.

Let us specialise this result to our framework. Take

$$X = H^{-1-\delta/2}(\mathbb{T}^2), \quad B = H^{-1-\delta}(\mathbb{T}^2), \quad Y = H^{-3-\delta}(\mathbb{T}^2),$$

with $\delta > 0$: by Gagliardo-Nirenberg estimates the interpolation inequality is satisfied with $\theta = \delta/2$. Let us take moreover $s_0 = 0$, $s_1 = 1/2 - \gamma$ with $\gamma > 0$, $r_1 = 2$ and $r_0 = q \geq 1$, so that the discriminating parameter is

$$s_* = -\gamma\theta - \frac{1-\theta}{q}.$$

Note that as we take δ smaller and smaller, and q bigger and bigger, we can get $s_* < 0$ arbitrarily close to 0, but not 0. We have thus derived:

Corollary 4.2. If the sequence

$$\{v_n\} \subset L^p([0, T], H^{-1-\delta}) \cap W^{1/2-\gamma, 2}([0, T], H^{-3-\delta})$$

is bounded for any choice of $\delta > 0$ and $p \geq 1$, and for some $\gamma > 0$, then it is relatively compact in $L^q([0, T], H^{-1-\delta})$ for any $1 \leq q < \infty$. As a consequence, if a sequence of stochastic processes $u^n : [0, T] \rightarrow H^{-1-\delta}$ defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ is such that, for any $\delta > 0$, $p \geq 1$ and some $\gamma > 0$, there exists a constant $C_{\delta, \gamma, p}$ for which

$$\sup_n \mathbb{E} \left[\|u^n(t)\|_{L^p([0, T], H^{-1-\delta})} + \|u^n\|_{W^{1/2-\gamma, 1}([0, T], H^{-3-\delta})} \right] \leq C_{\delta, \gamma, p}, \quad (4.2)$$

then the laws of u_n on $L^q([0, T], H^{-1-\delta})$ are tight for any $1 \leq q < \infty$.

The processes we will consider are discontinuous in time: this is why we consider only fractional Sobolev regularity in time. However, as we have just observed, this prevents us to use Simon's criterion to prove any time regularity beyond L^q . This is why we will combine the latter result with a compactness criterion for *cadlag* functions. We refer to [24] for both the forthcoming result and the necessary preliminaries on the space $\mathbb{D}([0, T], S)$ of *cadlag* functions taking values in a complete separable metric space S .

Theorem 4.3 (Aldous' Criterion). Consider a sequence of stochastic processes $u^n : [0, T] \rightarrow S$ defined on probability spaces $(\Omega^n, \mathcal{F}^n, \mathbb{P}^n)$ and adapted to filtrations \mathcal{F}_t^n . The laws of u^n are tight on $\mathbb{D}([0, T], S)$ if:

1. for any $t \in [0, T]$ (a dense subset suffices) the laws of the variables u_t^n are tight;
2. for all $\varepsilon, \varepsilon' > 0$ there exists $R > 0$ such that for any sequence of \mathcal{F}^n -stopping times $\tau_n \leq T$ it holds

$$\sup_n \sup_{0 \leq r \leq R} \mathbb{P}^n (d(u_{\tau_n}^n, u_{\tau_n+r}^n) \geq \varepsilon') \leq \varepsilon.$$

4.2 Tightness of point vortices processes

The following estimate on our Poissonian random measures is the crux in all the forthcoming bounds; it is essentially a Poissonian analogue of the ones in Section 3 of [11].

Proposition 4.4. *Let $\omega_{M,N} \sim \frac{1}{\sqrt{N}} \Xi_M^{\theta, N\lambda}$. For any $1 \leq p < \infty$ there exists a constant $C_p > 0$ such that for any measurable bounded functions $h : \mathbb{T}^2 \rightarrow \mathbb{R}$ and $f : \mathbb{T}^{2 \times 2} \rightarrow \mathbb{R}$ it holds*

$$\mathbb{E} [\langle h, \omega_{M,N} \rangle^{2p}] \leq C_p \|h\|_\infty^{2p}, \quad \mathbb{E} [\langle f, \omega_{M,N} \otimes \omega_{M,N} \rangle^p] \leq C_p \|f\|_\infty^p, \quad (4.3)$$

uniformly in $N \geq 0$ and $M \in [0, \infty]$. As a consequence, since for $\delta > 0$ the Green function $\Delta^{-1-\delta}$ is smooth,

$$\mathbb{E} [\|\omega_{M,N}\|_{H^{-1-\delta}}^{2p}] = \mathbb{E} [\langle \Delta^{-1-\delta}, \omega_{M,N} \otimes \omega_{M,N} \rangle^p] \leq C_{p,\delta}, \quad (4.4)$$

uniformly in M, N .

Proof. Since

$$\langle f, \omega_{M,N} \otimes \omega_{M,N} \rangle = \left\langle \tilde{f}, \omega_{M,N} \otimes \omega_{M,N} \right\rangle, \quad \tilde{f}(x, y) = \frac{1}{2}(f(x, y) + f(y, x)),$$

we reduce ourselves to symmetric functions. Moreover, without loss of generality we can check (4.3) for functions with separate variables $f(x, y) = h(x)h(y)$, $h : \mathbb{T}^2 \rightarrow \mathbb{R}$ measurable and bounded, for which it holds

$$\mathbb{E} [\langle f, \omega_{M,N} \otimes \omega_{M,N} \rangle^p] = \mathbb{E} [\langle h, \omega_{M,N} \rangle^{2p}].$$

Moments of the random variable $\langle h, \omega_{M,N} \rangle$ can be evaluated by differentiating the moment generating function (2.8): using Faà di Bruno's formula to take $2p$ derivatives we get

$$\begin{aligned} \mathbb{E} [\langle h, \omega_{M,N} \rangle^{2p}] &= \\ &= (2p!) \sum_{\substack{r_1, \dots, r_{2p} \geq 0 \\ r_1 + 2r_2 + \dots + 2pr_{2p} = 2p}} \prod_{k=1}^{2p} \frac{1}{(k!)^{r_k} r_k!} \left(N\lambda \int_{[0,M] \times \{\pm 1\} \times \mathbb{T}^2} \frac{\sigma^k}{N^{k/2}} e^{-\theta t k} h(x) d\sigma dx dt \right)^{r_k} \\ &\leq (2p!) \sum_{\substack{r_1, \dots, r_{2p} \geq 0 \\ r_1 + 2r_2 + \dots + 2pr_{2p} = 2p}} \prod_{k=1}^{2p} \frac{(N\lambda)^{r_k} \|h\|_k^{kr_k} \mathbf{1}_{2|k}^{r_k}}{(\theta k)^{r_k} N^{kr_k/2} (k!)^{r_k} r_k!} \\ &= \frac{(2p!) \|h\|_\infty^{2p}}{N^p} \sum_{\substack{r_1, \dots, r_{2p} \geq 0 \\ r_1 + 2r_2 + \dots + 2pr_{2p} = 2p}} \prod_{k=1}^{2p} \frac{(N\lambda)^{r_k} \mathbf{1}_{2|k}^{r_k}}{(\theta k)^{r_k} (k!)^{r_k} r_k!} \end{aligned}$$

(see [27, 28] for similar classical computations). Let us stress that when an integral in the latter formula is null, its 0-th power is to be interpreted as $0^0 = 1$. The contribution of $\mathbf{1}_{2|k} = \int \sigma^k d\sigma$ is crucial: when k is odd, $\mathbf{1}_{2|k}$ is null, so only terms with $m_k = 0$ survive in the sum (again, $0^0 = 1$). Thus, the highest power of N appearing is $N^{r_2} \leq N^{2p/2} = N^p$, which is compensated by the N^{-p} we factored out, and this concludes the proof. \square

We can now discuss convergence at fixed times.

Proposition 4.5. *The laws of a family of variables $\omega_{M,N} \sim \frac{1}{\sqrt{N}} \Xi_M^{\theta, N\lambda}$, defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ and taking values in $H^{-1-\delta}$ are tight, for any fixed $\delta > 0$. Moreover,*

- the limit as $M \rightarrow \infty$ at fixed N , say $N = 1$, is the law of $\Xi_\infty^{\theta,\lambda}$;
- the limit as $N \rightarrow \infty$ at fixed M (any $M \in (0, \infty]$) is the law of $\sqrt{\frac{(1-e^{-2\theta M})\lambda}{2\theta}}\eta$;

and if the variables converge almost surely, they do so also in $L^p(\Omega, H^{-1-\delta})$ for any $1 \leq p < \infty, \delta > 0$.

Proof. The embedding $H^\alpha \hookrightarrow H^\beta$ is compact as soon as $\alpha > \beta$, and we know that the variables are uniformly bounded elements of $L^p(\Omega, H^{-1-\delta})$ for any $p \geq 1$ by (4.3), so by Čebyšëv's inequality their laws are tight.

Identification of limit laws is yet another consequence of (2.8): by Theorem 2 of [16] (an infinite-dimensional Lévy theorem) we only need to check that characteristic functions $\mathbb{E}[e^{i\langle \omega_{M,N}, h \rangle}]$ converge to the ones of the announced limits for any $h \in H^{1+\delta}$. Since (2.8) is valid for all $M \in [0, \infty]$, the limit for $M \rightarrow \infty$ poses no problem. As for the limit $N \rightarrow \infty$, for any test function $h \in H^{1+\delta}$,

$$\begin{aligned} \mathbb{E}[\exp(i\langle h, \omega_{M,N} \rangle)] &= e^{-N\lambda} \exp\left(N\lambda \int_{[0,M] \times \{\pm 1\} \times \mathbb{T}^2} \exp\left(\frac{i\sigma}{\sqrt{N}} h(x)e^{-\theta t}\right) dx d\sigma dt\right) \\ &= e^{-N\lambda} \exp\left(N\lambda \int_0^M \frac{1}{N} \|h\|_2^2 e^{-2\theta t} dt + O_h\left(\frac{1}{N}\right)\right) \\ &\xrightarrow{N \rightarrow \infty} \exp\left(\frac{\lambda}{2\theta} \|h\|_2^2 (1 - e^{-2\theta M})\right), \end{aligned}$$

where in the second step we used the following elementary expansion: for $\phi \in C(\mathbb{T}^2)$,

$$\left| \frac{1}{2} \int_{\mathbb{T}^2} \left(\exp\left(\frac{\phi(x)}{\sqrt{N}}\right) + \exp\left(-\frac{\phi(x)}{\sqrt{N}}\right) \right) dx - 1 - \frac{\|\phi\|_2^2}{2N} \right| \leq \frac{\|\phi\|_4^4}{24N^2}. \quad (4.5)$$

Since $\mathbb{E}[\exp(i\langle h, \eta \rangle)] = \exp(-\|h\|_2^2)$, this concludes the proof. \square

The latter result provides compactness “in space” (“equi-boundedness”): in order to apply Theorem 4.2 and Theorem 4.3, we also need to obtain a control on the regularity “in time” (“equi-continuity”). We will obtain it by exploiting the equation satisfied by $\omega_{M,N}$, which we derived in Theorem 3.2, which allows us to prove the forthcoming estimate on increments.

Proposition 4.6. *Let $\omega_{M,N} : [0, T] \rightarrow H^{-1-\delta}$ be the stochastic process defined at the beginning of this Section. For any \mathcal{F}_t -stopping time $\tau \leq T$ (possibly constant), $r, \delta > 0$, there exists a constant $C_{\delta,T}$ independent of M, N, τ, r such that*

$$\mathbb{E}\left[\|\omega_{M,N,\tau+r} - \omega_{M,N,\tau}\|_{H^{-3-\delta}}^2\right] \leq C_{\delta,T} \cdot r. \quad (4.6)$$

Proof. In order to lighten notation, and since the final result must not depend on M, N , let us drop them when writing $\omega_{M,N,t} = \omega_t$. By its definition in 4.1 and Theorem 2.14 we know that the process satisfies the integral equation

$$\langle f, \omega_{t+r} \rangle - \langle f, \omega_t \rangle = -\theta \int_t^{t+r} \langle f, \omega_s \rangle ds + \int_t^{t+r} : \langle H_f, \omega_s \otimes \omega_s \rangle : ds + \left\langle f, \frac{1}{\sqrt{N}} (\Sigma_{t+r}^{N\lambda} - \Sigma_t^{N\lambda}) \right\rangle, \quad (4.7)$$

for any smooth $f \in C^\infty(\mathbb{T}^2)$. Since this equation holds \mathbb{P} -almost surely uniformly in $s, t \in [0, T]$, it is also true when we replace t with the stopping time τ . It is convenient to recall that

$$\|u\|_{H^{-3-\delta}}^2 = \sum_{k \in \mathbb{Z}^2} (1 + |k|^2)^{-3-\delta} |\hat{u}_k|^2,$$

so we can use the weak integral equation against the orthonormal functions e_k to control the full norm:

$$\mathbb{E} \left[\|\omega_{\tau+r} - \omega_\tau\|_{H^{-3-\delta}}^2 \right] = \sum_{k \in \mathbb{Z}^2} (1 + |k|^2)^{-3-\delta} \mathbb{E} \left[|\langle \omega_{\tau+r} - \omega_\tau, e_k \rangle|^2 \right]. \quad (4.8)$$

We estimate increments by bounding separately the terms in the equation, let us start from the linear one:

$$\mathbb{E} \left[\left| \int_\tau^{\tau+r} \langle f, \omega_s \rangle ds \right|^2 \right] \leq r \mathbb{E} \left[\int_0^T |\langle f, \omega_s \rangle|^2 ds \right] = r \int_0^T \mathbb{E} \left[|\langle f, \omega_s \rangle|^2 \right] ds \leq CTr \|f\|_\infty^2, \quad (4.9)$$

where the last passage makes use of the uniform estimate (4.3). The nonlinearity is the harder one, and its singularity is the reason why we can not obtain space regularity beyond $H^{-3-\delta}$,

$$\mathbb{E} \left[\left| \int_\tau^{\tau+r} : \langle H_f, \omega_s \otimes \omega_s \rangle : ds \right|^2 \right] \leq r \int_0^T \mathbb{E} \left[|\langle H_f, \omega_s \otimes \omega_s \rangle|^2 \right] ds \quad (4.10)$$

$$\leq CTr \|H_f\|_\infty^2 \leq CTr \|f\|_{C^2(\mathbb{T}^2)}^2, \quad (4.11)$$

where the second passage uses (4.3), and the third is due to the fact that by Taylor expansion

$$|H_f(x, y)| = \frac{1}{2} |K(x, y)(\nabla f(x) - \nabla f(y))| \leq C \frac{|\nabla f(x) - \nabla f(y)|}{d(x, y)} \leq C \|f\|_{C^2(\mathbb{T}^2)}.$$

By (2.5), the martingale $(\langle f, N^{-1/2}(\Sigma_{t+r}^{N\lambda} - \Sigma_t^{N\lambda}) \rangle)_{t \in [0, T]}$ has constant quadratic variation $\lambda r \|f\|_{L^2}^2$, so Burkholder-Davis-Gundy inequality gives

$$\mathbb{E} \left[\left| \langle f, N^{-1/2}(\Sigma_{t+r}^{N\lambda} - \Sigma_t^{N\lambda}) \rangle \right|^2 \right] \leq \mathbb{E} \left[\sup_{t \in [0, T]} \left| \langle f, N^{-1/2}(\Sigma_{t+r}^{N\lambda} - \Sigma_t^{N\lambda}) \rangle \right|^2 \right] \leq C\lambda r \|f\|_{L^2}^2. \quad (4.12)$$

Applying estimates (4.9, 4.10, 4.12) to the functions e_k , from (4.7) and Cauchy-Schwarz inequality we get

$$\mathbb{E} \left[|\langle \omega_{\tau+r} - \omega_\tau, e_k \rangle|^2 \right] \leq C_{\theta, \lambda, T} r |k|^4,$$

so that (4.8) gives us

$$\mathbb{E} \left[\|\omega_{\tau+r} - \omega_\tau\|_{H^{-3-\delta}}^2 \right] \leq \sum_{k \in \mathbb{Z}^2} (1 + |k|^2)^{-3-\delta} Cr (T + |k|^4 T + \lambda) \leq C_{\theta, \lambda, T, \delta} r,$$

which concludes the proof. \square

Proposition 4.7. *The laws of the processes $\omega_{M,N} : [0, T] \rightarrow H^{-1-\delta}$ are tight in*

$$L^q([0, T], H^{-1-\delta}) \cap \mathbb{D}([0, T], H^{-3-\delta})$$

for any $\delta > 0, 1 \leq q < \infty$.

Proof. Since $\omega_{M,N,t} \sim \frac{1}{\sqrt{N}} \Xi_{M+t}^{\theta, N\lambda}$, they are bounded in $L^p(\Omega, H^{-1-\delta})$ for any $\delta > 0, 1 \leq p < \infty$ uniformly in M, N, t as shown in Theorem 4.5, and as a consequence the processes $\omega_{M,N}$ are uniformly bounded in $L^p(\Omega \times [0, T], H^{-1-\delta})$, for any $\delta > 0, 1 \leq p < \infty$. Moreover, we have proved fixed-time tightness. We are thus left to prove Aldous' condition in $H^{-3-\delta}$

and to control a fractional Sobolev norm in time in order to apply Theorem 4.2 and Theorem 4.3, concluding the proof. As in the previous proof, we denote $\omega_{M,N,t} = \omega_t$.

We only need to apply the uniform bound on increments (4.6). Starting from the fractional Sobolev norm, we evaluate

$$\begin{aligned}\mathbb{E} [\|\omega\|_{W^{\alpha,1}([0,T],H^{-3-\delta})}] &= \mathbb{E} \left[\int_0^T \int_0^T \frac{\|\omega_t - \omega_s\|_{H^{-3-\delta}}}{|t-s|^{1+\alpha}} dt ds \right] \\ &\leq \int_0^T \int_0^T \frac{\mathbb{E} [\|\omega_t - \omega_s\|_{H^{-3-\delta}}]}{|t-s|^{1+\alpha}} dt ds \\ &\leq C \int_0^T \int_0^T |t-s|^{-1/2-\alpha},\end{aligned}$$

which converges as soon as $\alpha < 1/2$. Aldous's condition follows from Čebyšev's inequality: if τ is a stopping time for ω_t , then

$$\sup_{0 \leq r \leq R} \mathbb{P} (\|\omega_{\tau+r} - \omega_\tau\|_{H^{-3-\delta}} \geq \varepsilon) \leq \varepsilon^{-1} \sup_{0 \leq r \leq R} \mathbb{E} [\|\omega_{\tau+r} - \omega_\tau\|_{H^{-3-\delta}}] \leq C\varepsilon^{-1}R^{1/2},$$

where the right-hand side is smaller than $\varepsilon' > 0$ as soon as R , which we can choose, is small enough. \square

Let us conclude this paragraph with a martingale central limit theorem concerning the driving noise of our approximant processes.

Proposition 4.8. *Let $(\Pi_t^N)_{t \in [0,T], N \in \mathbb{N}}$ be a sequence of $H^{-1-\delta}$ -valued martingale with laws $\Pi^N \sim \frac{1}{\sqrt{N}} \Sigma^{N\lambda}$ (fix $\delta > 0$). The laws of Π^N are tight in*

$$L^q([0,T], H^{-1-\delta}) \cap \mathbb{D}([0,T], H^{-1-\delta}) \quad (4.13)$$

for any $\delta > 0$, $1 \leq q < \infty$, and limit points have the law of the Wiener process $\sqrt{\lambda}W_t$ on $H^{-1-\delta}$ with covariance

$$\mathbb{E} [\langle W_t, f \rangle, \langle W_s, g \rangle] = t \wedge s \langle f, g \rangle_{L^2(\mathbb{T}^2)}.$$

Proof. By (4.12) we readily get

$$\mathbb{E} [\|\Pi_{\tau+r}^N - \Pi_\tau^N\|_{H^{-1-\delta}}^2] \leq C_{\delta,\lambda} r$$

for any $N \in \mathbb{N}, \delta, r > 0$ and any τ stopping time for Π^N , uniformly in N . The very same argument of the last proposition (here with a better space regularity) proves then the claimed tightness. The martingale property (with respect to the processes own filtrations) carries on to limit points since it can be expressed by means of the following integral formulation: for any $s, t \in [0, T]$,

$$\mathbb{E} [(\Pi_t^N - \Pi_s^N) \Phi(\Pi^N|_{[0,s]})] = 0$$

for all the real bounded measurable functions Φ on $(H^{-1-\delta})^{[0,s]}$. Limit points are Gaussian processes, since at any fixed time

$$\frac{1}{\sqrt{N}} \Sigma_t^{N\lambda} \sim \frac{1}{\sqrt{N}} \Xi_t^{\theta=0, N\lambda} \xrightarrow{N \rightarrow \infty} \sqrt{\lambda}t\eta \sim \sqrt{\lambda}W_t,$$

as one can show by repeating the computations on characteristic functions in Theorem 4.5 with $\theta = 0, M = t$. It now suffices to recall the covariance formulas (2.2) and (2.4),

$$\mathbb{E} \left[\left\langle \frac{1}{\sqrt{N}} \Sigma_t^{N\lambda}, f \right\rangle \left\langle \frac{1}{\sqrt{N}} \Sigma_s^{N\lambda}, g \right\rangle \right] = \lambda(t \wedge s) \langle f, g \rangle_{L^2}^2 = \mathbb{E} [\langle \sqrt{\lambda}W_t, f \rangle, \langle \sqrt{\lambda}W_s, g \rangle],$$

to conclude that any limit point has the law of $\sqrt{\lambda}W$. \square

4.3 Identifying limits

The last step is to prove that limit points of the family of processes $\omega_{M,N}$ satisfy Theorem 2.13. First, let us recall once again our setup for the sake of clarity:

- $\lambda, \theta > 0$ are fixed throughout;
- there is a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ on which the stochastic processes $\Sigma_t^{N\lambda}$ and the random variables $\Xi_M^{\theta, N\lambda}$ are defined, for $M \geq 0, N \in \mathbb{N}$, their laws being as in section 2;
- the processes $(\omega_{M,N,t})_{t \in [0,T]}$ are defined as at the beginning of this section: strong solutions of (4.1) with initial datum $\frac{1}{\sqrt{N}}\Xi_M^{\theta, N\lambda}$ and driving noise $\frac{1}{\sqrt{N}}\Sigma_t^{N\lambda}$, built as in Theorem 3.2.

To fix notation, let us consider separately the following three cases: by Theorem 4.7, we can consider converging sequences

- (Ps) $(\omega_{M_n, N=1})_{n \in \mathbb{N}}$, with $M_n \rightarrow \infty$ as $n \rightarrow \infty$, the limit being ω_t^P ;
 (G) $(\omega_{M, N_n})_{n \in \mathbb{N}}$, with $N_n \rightarrow \infty$ as $n \rightarrow \infty$ and fixed $M < \infty$, the limit being $\omega_{M,t}^G$;
 (Gs) $(\omega_{M_n, N_n})_{n \in \mathbb{N}}$, with $M_n, N_n \rightarrow \infty$ as $n \rightarrow \infty$, the limit being ω_t^G ;

the convergence in law takes place in $L^q([0, T], H^{-1-\delta}) \cap \mathbb{D}([0, T], H^{-3-\delta})$, for any fixed $\delta > 0, 1 \leq q < \infty$. By Theorem 4.5, the Poissonian limit (Ps) has marginals $\omega_t^P \sim \Xi_\infty^{\theta, \lambda}$, and the Gaussian ones $\omega_{M,t}^G \sim \sqrt{\frac{\lambda}{2\theta}(1 - e^{-2\theta(M+t)})}\eta$ for all $t \in [0, T]$, $M \in [0, \infty)$, and $\omega_t^G \sim \sqrt{\frac{\lambda}{2\theta}}\eta$ (the labels are given so to match the ones in Theorem 2.13). Notice that $(\omega_m^P)_{m \in \mathbb{N}}$ have all the same driving noise Σ_t^λ , but different initial data, while in the Gaussian limiting sequences the driving noise also varies. Let us show that the limit laws in the cases where $M \rightarrow \infty$ are stationary.

Proposition 4.9. *The processes ω_t^P and ω_t^G are stationary.*

Proof. As the intuition suggests, the key is the fact that M is a time-like parameter, and taking $M \rightarrow \infty$ corresponds to the infinite time limit. Formally, we observe that for all $r > 0$, $0 \leq t_1 \leq \dots \leq t_k < \infty$, and M, N ,

$$(\omega_{M,N,t_1+r}, \dots, \omega_{M,N,t_k+r}) \sim (\omega_{M+r,N,t_1}, \dots, \omega_{M+r,N,t_k}). \quad (4.14)$$

Indeed, by construction (see section 3), for all $s < t$, $\omega_{M,N,t}$ is given as a measurable function of $\omega_{M,N,s}$ and the driving noise,

$$\omega_{M,N,t} = F_{s,t}(\omega_{M,N,s}, \Sigma^{N\lambda}|_{[s,t]}) \quad (4.15)$$

this, combined with the fact that $\omega_{M,N,t} \sim \omega_{M+t,N,0}$ and the invariance of $\Sigma^{N\lambda}$ by time shifts proves (4.14). Passing (4.14) to the limits (Ps) and (Gs) concludes the proof, since the dependence on r of the right-hand side disappears. \square

Remark 4.10. Equation (4.15) is equivalent to the Markov property, *cf.* the beginning of the proof to Theorem 3.2. Equation (4.14) is the time homogeneity property. The Markov property is a consequence of uniqueness for the system (3.2), (3.3). Since uniqueness result in cases (Ps), (G) and (Gs) of Theorem 2.13 seem to be out of reach by now, we can not hope to derive the Markov property as well.

We are only left to show that our limits do produce the sought solutions of Theorem 2.16. First, we apply Skorokhod's theorem to obtain almost sure convergence.

Proposition 4.11. *There exist stochastic processes $(\tilde{\omega}_n^P)_{n \in \mathbb{N}}, \tilde{\Sigma}_t^\lambda$, defined on a probability space $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$, such that their joint distribution coincides with the one of the original objects and with $\tilde{\omega}_m^P$ converging to a limit $\tilde{\omega}^P$ almost surely in $L^q([0, T], H^{-1-\delta}) \cap \mathbb{D}([0, T], H^{-3-\delta})$ for any fixed $\delta > 0, 1 \leq q < \infty$.*

Analogously, there exist $(\tilde{\omega}_{M,n}^G, \tilde{\omega}_n^G, \tilde{\Sigma}_t^{N_n \lambda})_{n \in \mathbb{N}}$, defined on $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$, such that their joint distribution coincides with the one of the original objects and with $\tilde{\omega}_{M,n}^G, \tilde{\omega}_n^G$ converging respectively to limits $\tilde{\omega}_M^G, \tilde{\omega}^G$ almost surely in $L^q([0, T], H^{-1-\delta}) \cap \mathbb{D}([0, T], H^{-3-\delta})$ for any fixed $\delta > 0, 1 \leq q < \infty$.

The proof is a straightforward application of the following version of Skorokhod's theorem, which we borrow from [25] (see references therein). The required tightness is provided by Theorem 4.7 and Theorem 4.8.

Theorem 4.12 (Skorokhod Representation). *Let $X_1 \times X_2$ be the product of two Polish spaces, $\chi^n = (\chi_n^1, \chi_n^2)$ be a sequence of $X_1 \times X_2$ -valued random variables, defined on a probability space $(\Omega, \mathcal{F}, \mathbb{P})$, converging in law and such that χ_n^1 have all the same law ρ . Then there exist a sequence $\tilde{\chi}^n = (\tilde{\chi}_n^1, \tilde{\chi}_n^2)$ of $X_1 \times X_2$ -valued random variables, defined on a probability space $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$, such that*

- χ^n and $\tilde{\chi}^n$ have the same law for all n ;
- $\tilde{\chi}^n$ converge almost surely to a $X_1 \times X_2$ -valued random variable $\tilde{\chi} = (\tilde{\chi}_1^n, \tilde{\chi}_2^n)$ on $(\tilde{\Omega}, \tilde{\mathcal{F}}, \tilde{\mathbb{P}})$;
- the variable $\tilde{\chi}_1^n$ and $\tilde{\chi}_1$ coincide almost surely.

Proof of Theorem 4.11. In the case (P) we apply the above result with $X_1 = X_2 = X = L^q([0, T], H^{-1-\delta}) \cap \mathbb{D}([0, T], H^{-3-\delta})$ and $\chi_1^m = \Sigma_t^\lambda, \chi_2^m = \omega_m^P$, while for the case (G) we take $X_1 = \{0\}$ and $X_2 = X \times X$, with $\chi_2^n = (\omega_n^G, \Sigma_t^{N_n \lambda})$. \square

The new processes still are weak solutions of (4.1) in the sense of Theorem 2.13. Consider for instance the $\tilde{\omega}_n^G$ (the other case being identical): clearly their trajectories have the same regularity as ω_n^G , and they have the same fixed time distributions. As for the equation, it holds, for any $f \in C^\infty(\mathbb{T}^2)$ and $t \in [0, T]$, \mathbb{P} -almost surely,

$$\langle f, \tilde{\omega}_{n,t}^G \rangle - \langle f, \tilde{\omega}_{n,0}^G \rangle + \theta \int_0^t \langle f, \tilde{\omega}_{n,s}^G \rangle ds - \int_0^t : \langle H_f, \tilde{\omega}_{n,s}^G \otimes \tilde{\omega}_{n,s}^G \rangle : ds - \left\langle f, \frac{1}{\sqrt{N_n}} \Sigma_t^{N_n \lambda} \right\rangle = 0,$$

since taking the expectation of the absolute value (capped by 1) of the right-hand side gives a functional of the law of $\tilde{\omega}_n^G, \tilde{\Sigma}_t^{N_n \lambda}$, which is the same of the original ones. Moreover, since all the terms in the last equation are *cadlag* functions in time (in fact they are all continuous but the noise term), one can choose the $\tilde{\mathbb{P}}$ -full set on which the equation holds uniformly in $t \in [0, T]$.

Remark 4.13. In fact, one can prove more. Following the proof of Lemma 28 in [11], it is possible to show that the new Skorokhod process have in fact the same point vortices structure of $\omega_{M,N}$, namely it is possible to represent $\tilde{\omega}_{m,t}^P$ and $\tilde{\omega}_{M,n,t}^G, \tilde{\omega}_{n,t}^G$ as sums of vortices satisfying equations (3.2) and (3.3) of section 3. The argument would be quite long, and we feel that it would not add much to our discussion, so we refrain to go into details, contenting us with our analytically weak notion of solution.

To ease notation we will now drop all tilde symbols, implying that we are going to work only with the new processes and noise terms. We are finally ready to pass to the limit the stochastic equations satisfied by our approximating processes, thus concluding the proof of our main result.

Proof of Theorem 2.16. The limits of $\omega_n^P, \omega_{M,n}^G$ and ω_n^G provide respectively the sought solutions in the cases (Ps), (G) and (Gs) of Theorem 2.13. We focus again our attention on ω_n^G , case (Gs), the other ones being analogous.

Since ω_n^G converges almost surely in the spaces (4.13), we immediately deduce that, for any $f \in C^\infty(\mathbb{T}^2)$ and $t \in [0, T]$, \mathbb{P} -almost surely,

$$\langle f, \omega_{n,t}^G \rangle \rightarrow \langle f, \omega_t^G \rangle, \quad (4.16)$$

$$\int_0^t \langle f, \omega_{n,s}^G \rangle ds \rightarrow \int_0^t \langle f, \omega_s^G \rangle ds. \quad (4.17)$$

The nonlinear term is only slightly more difficult. Let $H_k \in C^\infty(\mathbb{T}^{2 \times 2})$, $k \in \mathbb{N}$, be symmetric functions vanishing on the diagonal converging to H_f as $k \rightarrow \infty$ (it is yet another equivalent of the approximation procedure (2.11)). Then

$$:\langle H_k, \omega_{n,t}^G \otimes \omega_{n,t}^G \rangle: = \langle H_k, \omega_{n,t}^G \otimes \omega_{n,t}^G \rangle \rightarrow \langle H_k, \omega_t^G \otimes \omega_t^G \rangle = :\langle H_k, \omega_t^G \otimes \omega_t^G \rangle:$$

in $L^2(\Omega \times [0, T])$ (the last passage is due to (2.13)). Almost sure convergence of the noise terms is ensured by Theorem 4.11, and the limiting law has been determined in Theorem 4.8, hence, summing up, it holds \mathbb{P} -almost surely

$$\langle f, \omega_t^G \rangle - \langle f, \omega_0^G \rangle + \theta \int_0^t \langle f, \omega_s^G \rangle ds - \int_0^t :\langle H_f, \omega_s^G \otimes \omega_s^G \rangle: ds - \left\langle f, \sqrt{\lambda} W_t \right\rangle = 0.$$

As already noted above, quantifiers in \mathbb{P} and $t \in [0, T]$ can be exchanged thanks to the fact that we are dealing with *cadlag* processes in time. Stationarity of ω_t^P and ω_t^G follows from Theorem 4.9. This concludes the proof of Theorem 2.16. \square

References

- [1] Sergio Albeverio and Ana Bela Cruzeiro. Global flows with invariant (Gibbs) measures for Euler and Navier-Stokes two-dimensional fluids. *Comm. Math. Phys.*, 129(3):431–444, 1990. MR-1051499
- [2] S. Albeverio and B. Ferrario. Invariant measures of Lévy-Khintchine type for 2D fluids. In *Probabilistic methods in fluids*, pages 130–143. World Sci. Publ., River Edge, NJ, 2003. MR-2083369
- [3] S. Albeverio, M. Ribeiro de Faria, and R. Höegh-Krohn. Stationary measures for the periodic Euler flow in two dimensions. *J. Statist. Phys.*, 20(6):585–595, 1979. MR-0537263
- [4] G. Benfatto, P. Picco, and M. Pulvirenti. On the invariant measures for the two-dimensional Euler flow. *J. Statist. Phys.*, 46(3-4):729–742, 1987. MR-0883549
- [5] Hakima Bessaih and Benedetta Ferrario. Inviscid limit of stochastic damped 2D Navier-Stokes equations. *Nonlinearity*, 27(1):1–15, 2014. MR-3151089
- [6] Guido Boffetta and Robert E. Ecke. Two-dimensional turbulence. In *Annual review of fluid mechanics. Volume 44, 2012*, volume 44 of *Annu. Rev. Fluid Mech.*, pages 427–451. Annual Reviews, Palo Alto, CA, 2012. MR-2882604
- [7] Zdzisław Brzeźniak, Franco Flandoli, and Mario Maurelli. Existence and uniqueness for stochastic 2D Euler flows with bounded vorticity. *Arch. Ration. Mech. Anal.*, 221(1):107–142, 2016. MR-3483892
- [8] P. Constantin and F. Ramos. Inviscid limit for damped and driven incompressible Navier-Stokes equations in \mathbb{R}^2 . *Comm. Math. Phys.*, 275(2):529–551, 2007. MR-2335784
- [9] I. P. Cornfeld, S. V. Fomin, and Ya. G. Sinai. *Ergodic theory*, volume 245 of *Grundlehren der Mathematischen Wissenschaften [Fundamental Principles of Mathematical Sciences]*. Springer-Verlag, New York, 1982. Translated from the Russian by A. B. Sosinskiĭ. MR-0832433
- [10] D. Dürr and M. Pulvirenti. On the vortex flow in bounded domains. *Comm. Math. Phys.*, 85(2):265–273, 1982. MR-0676001
- [11] Franco Flandoli. Weak vorticity formulation of 2D Euler equations with white noise initial condition. *Comm. Partial Differential Equations*, 43(7):1102–1149, 2018. MR-3910197

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- [12] Franco Flandoli, Francesco Grotto, and Dejun Luo. Fokker-Planck equation for dissipative 2D Euler equations with cylindrical noise. *to appear on Stochastics and Dynamics*, page arXiv:1907.01994, Jul 2019.
- [13] F. Flandoli, M. Gubinelli, and E. Priola. Full well-posedness of point vortex dynamics corresponding to stochastic 2D Euler equations. *Stochastic Process. Appl.*, 121(7):1445–1463, 2011. MR-2802460
- [14] Franco Flandoli and Dejun Luo. Point vortex approximation for 2D Navier–Stokes equations driven by space-time white noise. *arXiv e-prints*, page arXiv:1902.09338, Feb 2019.
- [15] Franco Flandoli and Dejun Luo. Energy conditional measures and 2D turbulence. *J. Math. Phys.*, 61(1):013101, 22, 2020. MR-4047934
- [16] Leonard Gross. Harmonic analysis on Hilbert space. *Mem. Amer. Math. Soc.* No., 46:ii+62, 1963. MR-0161095
- [17] Francesco Grotto. Essential self-adjointness of Liouville operator for 2D Euler point vortices. *J. Funct. Anal.*, 279(6):108635, 2020. MR-4099477
- [18] Francesco Grotto and Marco Romito. A Central Limit Theorem for Gibbsian Invariant Measures of 2D Euler Equations. *Comm. Math. Phys.*, 376(3):2197–2228, 2020. MR-4104546
- [19] Eberhard Hopf. Statistical hydromechanics and functional calculus. *J. Rational Mech. Anal.*, 1:87–123, 1952. MR-0059119
- [20] Svante Janson. *Gaussian Hilbert spaces*, volume 129 of *Cambridge Tracts in Mathematics*. Cambridge University Press, Cambridge, 1997. MR-1474726
- [21] Robert H. Kraichnan. Remarks on turbulence theory. *Advances in Math.*, 16:305–331, 1975. MR-0371254
- [22] Andrew J. Majda and Andrea L. Bertozzi. *Vorticity and incompressible flow*, volume 27 of *Cambridge Texts in Applied Mathematics*. Cambridge University Press, Cambridge, 2002. MR-1867882
- [23] Carlo Marchioro and Mario Pulvirenti. *Mathematical theory of incompressible nonviscous fluids*, volume 96 of *Applied Mathematical Sciences*. Springer-Verlag, New York, 1994. MR-1245492
- [24] M. Métivier. *Stochastic partial differential equations in infinite-dimensional spaces*. Scuola Normale Superiore di Pisa. Quaderni. [Publications of the Scuola Normale Superiore of Pisa]. Scuola Normale Superiore, Pisa, 1988. With a preface by G. Da Prato. MR-0982268
- [25] Elżbieta Motyl. Stochastic Navier-Stokes equations driven by Lévy noise in unbounded 3D domains. *Potential Anal.*, 38(3):863–912, 2013. MR-3034603
- [26] L. Onsager. Statistical hydrodynamics. *Nuovo Cimento (9)*, 6(Supplemento, 2 (Convegno Internazionale di Meccanica Statistica)):279–287, 1949. MR-0036116
- [27] S. Peszat and J. Zabczyk. *Stochastic partial differential equations with Lévy noise*, volume 113 of *Encyclopedia of Mathematics and its Applications*. Cambridge University Press, Cambridge, 2007. An evolution equation approach. MR-2356959
- [28] Nicolas Privault. Combinatorics of Poisson stochastic integrals with random integrands. In *Stochastic analysis for Poisson point processes*, volume 7 of *Bocconi Springer Ser.*, pages 37–80. Bocconi Univ. Press, [place of publication not identified], 2016. MR-3585397
- [29] Ken-iti Sato. *Lévy processes and infinitely divisible distributions*, volume 68 of *Cambridge Studies in Advanced Mathematics*. Cambridge University Press, Cambridge, 2013. Translated from the 1990 Japanese original, Revised edition of the 1999 English translation. MR-3185174
- [30] Steven Schochet. The weak vorticity formulation of the 2-D Euler equations and concentration-cancellation. *Comm. Partial Differential Equations*, 20(5-6):1077–1104, 1995. MR-1326916
- [31] Steven Schochet. The point-vortex method for periodic weak solutions of the 2-D Euler equations. *Comm. Pure Appl. Math.*, 49(9):911–965, 1996. MR-1399201
- [32] Jacques Simon. Compact sets in the space $L^p(0, T; B)$. *Ann. Mat. Pura Appl. (4)*, 146:65–96, 1987. MR-0916688

Acknowledgments. This article was completed while the author was a Ph.D. student at Scuola Normale Superiore, under the supervision of Franco Flandoli, whom the former wishes to thank for many of the ideas here exposed.