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# Diffusion-approximation for a kinetic spray-like system with markovian forcing 

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#### Abstract

We study a kinetic toy model for a spray of particles immersed in an ambient fluid, subject to some additional random forcing given by a mixing, space-dependent Markov process. Using the perturbed test function method, we derive the hydrodynamic limit of the kinetic system. The law of the limiting density satisfies a stochastic conservation equation in Stratonovich form, whose drift and diffusion coefficients are completely determined by the law of the stationary process associated with the Markovian perturbation.


Keywords. kinetic equations, stochastic partial differential equations, diffusion-approximation, hydrodynamic limit, perturbed test functions.

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## 1 Introduction

### 1.1 Overview of the kinetic model

We consider a toy model for spray-like behavior, describing the motion of some particles immersed in an ambient fluid, given by the following conservation equation

$$
\begin{equation*}
\partial_{t} f+\varepsilon v \cdot \nabla_{x} f+\nabla_{v} \cdot(F(t, x, v) f)=0 \tag{1.1}
\end{equation*}
$$

The forcing term

$$
F(t, x, v)=\left(u_{t}(x)-v\right)+m_{t}(x)
$$

expresses how the speed of a particle tends to align with the local speed $u(x)$ of the ambient fluid. It is perturbed by some random forcing naturally modeled by a space-dependent, centered Markov process $\left(m_{t}(x)\right)_{t \geq 0}$. Assuming that the particles have very little impact on the behavior of the ambient fluid, $u \equiv u_{t}(x)$ evolves according to

$$
\begin{equation*}
\partial_{t} u-\Delta u=\varepsilon^{2} \int_{v}(v-u) f d v \tag{1.2}
\end{equation*}
$$

We could easily add a nonlinear term of Burger's type $u \cdot \nabla_{x} u$ in equation (1.2). This would not yield any mathematical difficulty but would unnecessarily give heavier proofs.

In order for the system (1.1)-(1.2) to be well-posed and admit a solution $\left(f_{t}(x, v), u_{t}(x)\right)$, we need, in fact, to restrict to the one-dimensional case:

$$
x \in \mathbb{T}=\mathbb{R} / \mathbb{Z}, \quad v \in \mathbb{R}
$$

One can refer to section 8 below (and in particular Remark 8.2) for details regarding the wellposedness of this system. The rescaled variables

$$
f_{t}^{\varepsilon}(x, v):=f_{\varepsilon^{-2} t}(x, v), \quad u_{t}^{\varepsilon}(x):=\varepsilon^{-1} u_{\varepsilon^{-2} t}(x), \quad m_{t}^{\varepsilon}(x):=m_{\varepsilon^{-2} t}(x)
$$

then satisfy

$$
\left\{\begin{array}{l}
\partial_{t} f^{\varepsilon}+\varepsilon^{-1} v \partial_{x} f^{\varepsilon}+\varepsilon^{-2} \partial_{v}\left[\left(\varepsilon u^{\varepsilon}(x)-v+m^{\varepsilon}(x)\right) f^{\varepsilon}\right]=0  \tag{1.3}\\
\partial_{t} u^{\varepsilon}-\partial_{x}^{2} u^{\varepsilon}=\int\left(v-\varepsilon u^{\varepsilon}\right) f^{\varepsilon} d v
\end{array}\right.
$$

Introducing the integrated quantity

$$
\begin{equation*}
\rho^{\varepsilon}(x)=\int f^{\varepsilon}(x, v) d v \tag{1.4}
\end{equation*}
$$

we are interested in the behavior of $\left(\rho^{\varepsilon}, u^{\varepsilon}\right)$ as $\varepsilon$ goes to zero. In essence, we wish to describe how the random forcing $m$ added on a "mesoscopic level" translates on a "macroscopic level" in
the scaling limit. More precisely, we shall determine the stochastic partial differential equation satisfied by the limiting law of $\rho^{\varepsilon}$ as $\varepsilon$ goes to 0 .

This work is part of a program on stochastic diffusion-approximation in the context of kinetic equation and fluid limits (see $[5,6,7]$ ). In the case considered here, it is to be noted that due to the absence of collision term in the kinetic equation of (1.3), the limiting SPDE satisfied by $\rho^{\varepsilon}$ will be of order one: it is in a fact a stochastic conservation equation when written in Stratonovich form.

This latter stochastic equation is linear but our proof provides a solution which has very low regularity and uniqueness in this context is not trivial. We need to use a duality argument and backward SPDEs to prove uniqueness.

The study of the scaling limit relies essentially on the perturbed test function method, first introduced in [10] and developed further in [9]. In contrast with previous works on the matter, particular attention will be paid to the rigorous construction of the correctors involved in the perturbed test function method. This will be achieved by handling infinitesimal generators and their domains in an overall more cautious way, and by exploiting the mixing assumptions made on the driving Markov process to properly solve Poisson equations associated with these generators.

### 1.2 Main results

Throughout our study, we will be led to work with the following spaces:

$$
\begin{aligned}
& L_{x, v}^{p}:=L^{p}(\mathbb{T} \times \mathbb{R}), \quad L_{x}^{p}:=L^{p}(\mathbb{T}) \\
& W_{x}^{s, p}:=W^{s, p}(\mathbb{T}), \quad H_{x}^{\sigma}:=W^{2, \sigma}(\mathbb{T})
\end{aligned}
$$

We will denote by $\mathcal{P}(E)$ the set of probability measures on a metric space $E$. When it is well defined, we will denote

$$
\langle f, g\rangle=\int f(x, v) g(x, v) d x d v \text { or } \int f(x) g(x) d x
$$

depending on the context. Finally, we shall sometimes use the handy notation $A(z) \lesssim B(z)$ to signify that there exists a constant $C>0$ such that $A(z) \leq C B(z)$ for all $z$.

Let us now state our diffusion-approximation result.
Theorem 1. Let $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}, \mathbb{P}\right)$ be a filtered probability space equipped with a Markov process $m=\left(m_{t}\right)_{t \geq 0}$ satisfying Assumptions 1 through 4, described in section 2.2 below. Let $T>0$. Assume that the initial data $\left(f_{0}^{\varepsilon}, u_{0}^{\varepsilon}\right)_{\varepsilon>0}$ satisfy

$$
\begin{aligned}
& f_{0}^{\varepsilon} \geq 0 \text { and } \iint f_{0}^{\varepsilon}(x, v) d x d v=1 \\
& \rho_{0}^{\varepsilon} \rightarrow \rho_{0} \text { in } L_{x}^{1}, \quad u_{0}^{\varepsilon} \rightarrow u_{0} \text { in } H_{x}^{\alpha} \text { for some } \alpha \in(1 / 2,3 / 2]
\end{aligned}
$$

and, for all fixed $\varepsilon>0$,

$$
\begin{equation*}
\int_{x} \int_{v}\left(1+|v|^{4}\right)\left|f_{0}^{\varepsilon}\right|^{2} d x d v+\int_{x} \int_{v}(1+|v|)\left|\nabla_{x, v} f_{0}^{\varepsilon}\right| d x d v<\infty, \quad u_{0}^{\varepsilon} \in H_{x}^{\eta} \text { for } \eta(\varepsilon)>3 / 2 \tag{1.5}
\end{equation*}
$$

Assume the following uniform bound:

$$
\begin{equation*}
\sup _{\varepsilon>0} \int_{x} \int_{v}|v|^{3} f_{0}^{\varepsilon} d x d v<\infty \tag{1.6}
\end{equation*}
$$

Let $\left(f^{\varepsilon}, u^{\varepsilon}\right)$ be the unique solution of (1.3) with initial data $\left(f_{0}^{\varepsilon}, u_{0}^{\varepsilon}\right)$ in the sense given in Proposition 3.1 below. Then for all $\gamma>1 / 2$ and $\beta<\alpha$, there exists a subsequence $\left(\rho^{\varepsilon_{k}}, u^{\varepsilon_{k}}\right)_{k \geq 1}$ such that

$$
\left(\rho^{\varepsilon_{k}}, u^{\varepsilon_{k}}\right) \rightarrow(\rho, u) \text { in law in } C\left([0, T] ; H_{x}^{-\gamma} \times H_{x}^{\beta}\right)
$$

where $u$ satisfies the heat equation

$$
\partial_{t} u-\partial_{x}^{2} u=0, \quad \text { with } \quad u(0)=u_{0}
$$

and $\rho$ is a martingale solution (see [4], Chapter 8, for the definition) of

$$
\begin{equation*}
d \rho=\partial_{x}[(a-u) \rho] d t+\partial_{x}\left[\rho Q^{1 / 2} \circ d W_{t}\right], \quad \text { with } \rho(0)=\rho_{0}, \tag{1.7}
\end{equation*}
$$

satisfying $\sup _{t \in[0, T]}\left\|\rho_{t}\right\|_{L_{x}^{1}} \leq 1$ a.s. In (1.7), $\left(W_{t}\right)_{t \geq 0}$ denotes a cylindrical Wiener process on $L_{x}^{2}$. The coefficient $a \equiv a(x)$ and the trace operator $Q$ on $L_{x}^{2}$ are completely determined by the law of the stationary process $\left(\widetilde{m}_{t}\right)_{t \geq 0}$ associated with the random driving term $m$ (see section 2.2 below for a definition of $\widetilde{m}$ ):

$$
\begin{gather*}
a(x)=\frac{1}{2} \int_{0}^{\infty} \mathbb{E}\left[\partial_{x} \widetilde{m}_{0}(x) \widetilde{m}_{t}(x)-\widetilde{m}_{0}(x) \partial_{x} \widetilde{m}_{t}(x)\right] d t+\int_{0}^{\infty} e^{-t} \mathbb{E}\left[\widetilde{m}_{0}(x) \partial_{x} \widetilde{m}_{t}(x)\right] d t,  \tag{1.8}\\
Q f(x)=\int_{y} k(x, y) f(y) d y, \text { with } k(x, y)=\int_{\mathbb{R}} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right] d t . \tag{1.9}
\end{gather*}
$$

Note that assumption (1.5) is only necessary to ensure that the system (1.3) can indeed be uniquely solved for fixed $\varepsilon>0$, as will be stated later on, in Proposition 3.1. The result of Theorem 1 holds when the underlying probability space $\Omega$ is equipped with the probability $\mathbb{P}_{\mu}$ - so that $\mathbb{P}_{\mu}\left[m_{0} \in A\right]=\mu(A)$ - for any initial law $\mu \in \mathcal{P}(E)$. For instance, one may consider the case $\mu=\delta_{n}$ where the process $\left(m_{t}\right)_{t \geq 0}$ starts at some fixed $n \in E$; or one may consider the case where $\mu=\nu$ is the invariant measure of $m$, so that $\left(m_{t}\right)_{t \geq 0}:=\left(\widetilde{m}_{t}\right)_{t \geq 0}$ is a stationary process (see again section 2.2 below).

The weak existence result established in Theorem 1 is completed by the following strong uniqueness result.
Theorem 2. Let $\rho_{0} \in L_{x}^{1}$ and $u_{0} \in W_{x}^{4, \infty}$. Assume furthermore that the coefficient $a$ and the kernel $k$ defined in (1.8), (1.9) satisfy

$$
\begin{equation*}
a \in W_{x}^{4, \infty}, \quad \sup _{y \in \mathbb{T}}\|k(\cdot, y)\|_{W_{x}^{10+\delta, \infty}}<\infty \text { for some } \delta>0 . \tag{1.10}
\end{equation*}
$$

Then a solution $\rho \in C\left([0, T] ; H_{x}^{-\gamma}\right)$ of (1.7) satisfying the bound

$$
\sup _{t \in[0, T]}\left\|\rho_{t}\right\|_{L_{x}^{1}} \leq 1 \text { a.s }
$$

is path-wise unique.
As a consequence of Theorem 2, Yamada-Watanabe's Theorem guarantees that uniqueness in law (i.e. weak uniqueness in the probabilistic sense) holds for (1.7) and the whole sequence converges in Theorem 1. The regularity properties (1.10) will be satisfied under some adequate assumptions on the driving process $\left(m_{t}\right)_{t \geq 0}$ (see the end of section 2.2).

In the limit where $\varepsilon$ goes to 0 , the rather complex random process $\left(\rho_{t}^{\varepsilon}\right)_{t \geq 0}$ driven by a general Markov process $m=\left(m_{t}\right)_{t \geq 0}$ can therefore be approximated by a much simpler diffusion $\left(\rho_{t}\right)_{t \geq 0}$ whose drift and covariance are explicitely determined by the law of the associated stationary process $\widetilde{m}=\left(\widetilde{m}_{t}\right)_{t \geq 0}$. Let us give a quick overview of the proof of Theorems 1 and 2 .

In section 2 we start by introducing a rigorous framework for solving Poisson equations associated with generators of Markov processes, which will be of great use in section 5 in particular. The tools and results presented here are all very classical, however the authors could not find a reference really suited for the specific setup adopted in the present work. We then describe the assumptions made on the driving process $m$ and their consequences in section 2.2. In section 3, we establish some estimates regarding the solutions of (1.3) and rigorously determine the generator of the process $\left(f^{\varepsilon}, u^{\varepsilon}, m^{\varepsilon}\right)$. The essential work is developed in section 5 where we use the perturbed test function method to determine the limiting equation (1.7) and pave the way towards the convergence of the martingale problem. Section 6 concludes the proof by proving the convergence announced in Theorem 1. Finally, the path-wise uniqueness of (1.7) stated in Theorem 2 is looked upon in section 7 .

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## 2 Generators and Poisson equations

### 2.1 Generalities on generators

On a filtered probability space $\left(\Omega, \mathcal{F},\left(\mathcal{F}_{t}\right)_{t \geq 0}\right)$, let us consider a càdlàg, time-homogeneous Markov process $X=\left(X_{t}\right)_{t \geq 0}$ taking values in some separable metric space $\mathcal{X} \equiv\left(\mathcal{X}, d_{\mathcal{X}}\right)$. Throughout this paper, we will naturally use the handy notation $\left(X_{t}(x)\right)_{t \geq 0}$ to denote $\left(X_{t}\right)_{t \geq 0}$ under the probability $\mathbb{P}_{\delta_{x}}$ (that is $X_{0}(x)=x$ a.s). All the processes we shall encounter in this work satisfy the following assumptions:

- $X$ is stochastically continuous:

$$
\begin{equation*}
\forall x \in \mathcal{X}, \quad \forall t \geq 0, \quad \forall \varepsilon>0, \quad \mathbb{P}\left[d_{\mathcal{X}}\left(X_{s}(x), X_{t}(x)\right)>\varepsilon\right] \underset{s \rightarrow t}{\longrightarrow} 0 \tag{2.1}
\end{equation*}
$$

- $X$ is locally bounded (in space and time):

$$
\begin{equation*}
\forall r>0, \forall T>0, \exists R>0, \forall x \in B(0, r), \mathbb{P}\left[\forall t \in[0, T], \quad X_{t}(x) \in B(0, R)\right]=1 \tag{2.2}
\end{equation*}
$$

Let us denote by $\left(Q_{t}\right)_{t \geq 0}$ the semigroup associated to $\left(X_{t}\right)_{t \geq 0}$. The local boundedness assumption (2.2) allows to define it over locally bounded test functions: let us introduce the spaces

$$
\begin{align*}
& \mathcal{B}^{l o c}(\mathcal{X})=\{\psi: \mathcal{X} \rightarrow \mathbb{R} \text { measurable, } \psi \text { maps bounded sets onto bounded sets }\}  \tag{2.3}\\
& C_{b}^{l o c}(\mathcal{X})=\{\psi: \mathcal{X} \rightarrow \mathbb{R} \text { continuous, } \psi \text { maps bounded sets onto bounded sets }\} \tag{2.4}
\end{align*}
$$

and define the semigroup $Q_{t}: \mathcal{B}^{l o c}(\mathcal{X}) \rightarrow \mathcal{B}^{l o c}(\mathcal{X})$ as

$$
\forall \psi \in \mathcal{B}^{l o c}(\mathcal{X}), \quad Q_{t} \psi(x)=\mathbb{E}\left[\psi\left(X_{t}(x)\right)\right] .
$$

Note that, using assumption (2.1), the map $t \mapsto Q_{t} \psi(x)$ is continuous whenever $\psi \in C_{b}^{\text {loc }}(\mathcal{X})$. The Markov property

$$
\forall \psi \in \mathcal{B}^{l o c}(\mathcal{X}), \quad \mathbb{E}\left[\psi\left(X_{t+s}\right) \mid \mathcal{F}_{s}\right]=Q_{t} \psi\left(X_{s}\right) \text { a.s }
$$

yields in particular the semigroup property $Q_{t+s}=Q_{t} Q_{s}$ on $\mathcal{B}^{l o c}(\mathcal{X})$.

Definition 2.1 (Infinitesimal generator). The infinitesimal generator of $X$ is the operator

$$
A: D(A) \subset \mathcal{B}^{l o c}(\mathcal{X}) \rightarrow C_{b}^{l o c}(\mathcal{X})
$$

defined as follows: for all $\psi \in \mathcal{B}^{l o c}(\mathcal{X})$ and $\theta \in C_{b}^{l o c}(\mathcal{X})$,

$$
\psi \in D(A) \text { and } A \psi=\theta \quad \text { if and only if } \quad \frac{Q_{t} \psi-\psi}{t} \xrightarrow[t \rightarrow 0^{+}]{\text {b.p.c }} \theta \in C_{b}^{l o c}(\mathcal{X})
$$

where b.p.c denotes the (locally) bounded pointwise convergence:

$$
\theta_{t} \xrightarrow{\text { b.p.c }} \theta \Longleftrightarrow\left\{\begin{array}{l}
\forall x \in \mathcal{X}, \quad \theta_{t}(x) \rightarrow \theta(x) \\
\forall B \subset \mathcal{X} \text { bounded, } \sup _{t>0} \sup _{x \in B}\left|\theta_{t}(x)\right|<\infty
\end{array}\right.
$$

We then have the following classical result.
Proposition 2.1 (Kolmogorov's equation). For all $\psi \in D(A)$, for all $t \geq 0$,

$$
\frac{Q_{t+s} \psi-Q_{t} \psi}{s} \xrightarrow[s \rightarrow 0]{b . p . c} Q_{t} A \psi
$$

In particular, for all $x \in \mathcal{X}$,

$$
\frac{d}{d t} Q_{t} \psi(x)=Q_{t} A \psi(x)
$$

Remark 2.1. Provided that $Q_{t} A \psi \in C_{b}^{l o c}(\mathcal{X})$, this states that $Q_{t} \psi \in D(A)$ and $A Q_{t} \psi=Q_{t} A \psi$.
Proof. For $s>0$, we may write

$$
\frac{Q_{t+s} \psi-Q_{t} \psi}{s}(x)=Q_{t}\left[\frac{Q_{s} \psi-\psi}{s}\right](x)=\mathbb{E}\left[\frac{Q_{s} \psi\left(X_{t}(x)\right)-\psi\left(X_{t}(x)\right)}{s}\right] .
$$

Making use of (2.1) and (2.2), we may use dominated convergence as $s$ goes to 0 to conclude for the right derivative. As for the left derivative, we simply note that

$$
\frac{Q_{t} \psi-Q_{t-s} \psi}{s}=Q_{t-s}\left[\frac{Q_{s} \psi-\psi}{s}-A \psi\right]+Q_{t-s} A \psi
$$

and conclude in the same way, keeping in mind that $A \psi \in C_{b}^{l o c}(\mathcal{X})$ so that $u \mapsto Q_{u} A \psi(x)$ is continuous. .

The following property will be of great interest, and could in fact be an alternative way of defining the generator $A$.

Proposition 2.2 (Martingale problem). For all $\psi \in D(A)$, for all $x \in \mathcal{X}$,

$$
M_{\psi}(t)=\psi\left(X_{t}(x)\right)-\int_{0}^{t} A \psi\left(X_{s}(x)\right) d s, \quad t \geq 0
$$

defines an $\left(\mathcal{F}_{t}\right)_{t \geq 0}$-martingale.
Proof. Let us write $X_{t}:=X_{t}(x)$ for simplicity. We need to check that, for all $h$ bounded, $\mathcal{F}_{s}$-measurable,

$$
\mathbb{E}\left[\left(\psi\left(X_{t}\right)-\psi\left(X_{s}\right)\right) h\right]=\mathbb{E}\left[\left(\int_{s}^{t} A \psi\left(X_{u}\right) d u\right) h\right]
$$

Using the Markov property (twice) and Proposition 2.1, the right hand side can be written as

$$
\begin{aligned}
\mathbb{E}\left[\mathbb{E}\left[\int_{s}^{t} A \psi\left(X_{u}\right) d u \mid \mathcal{F}_{s}\right] h\right] & =\mathbb{E}\left[\left(\int_{s}^{t}\left[\mathbb{E} A \psi\left(X_{u}\right) \mid \mathcal{F}_{s}\right] d u\right) h\right]=\mathbb{E}\left[\left(\int_{s}^{t} Q_{u-s} A \psi\left(X_{s}\right) d u\right) h\right] \\
& =\mathbb{E}\left[\left(\int_{0}^{t-s} \frac{d}{d u} Q_{u} \psi\left(X_{s}\right) d u\right) h\right]=\mathbb{E}\left[\left(Q_{t-s} \psi\left(X_{s}\right) d u-\psi\left(X_{s}\right)\right) h\right] \\
& =\mathbb{E}\left[\left(\psi\left(X_{t}\right)-\psi\left(X_{s}\right)\right) h\right]
\end{aligned}
$$

Again, we have used the fact $u \mapsto Q_{u} A \psi(x)$ is continous since $A \psi$ is continuous.
Let us now give a natural context in which one can interchange generator and integral.
Proposition 2.3 (Generator under the integral sign). Let $(t, x) \in \mathbb{R}^{+} \times \mathcal{X} \mapsto \psi_{t}(x) \in \mathbb{R}$ be $a$ measurable map, such that $\psi_{t} \in D(A)$ for all $t \geq 0$. Assume that, for all $B \subset \mathcal{X}$ bounded,

$$
\begin{align*}
& \int_{0}^{\infty} \sup _{x \in B}\left|\psi_{t}(x)\right| d t<\infty  \tag{2.5}\\
& \int_{0}^{\infty} \sup _{x \in B}\left|A \psi_{t}(x)\right| d t<\infty \tag{2.6}
\end{align*}
$$

Then $\psi=\int_{0}^{\infty} \psi_{t} d t \in D(A)$ and $A \psi=\int_{0}^{\infty} A \psi_{t} d t$.
Proof. Let us not get into the questions of measurability. First, $\psi$ is well defined, in $\mathcal{B}^{\text {loc }}(X)$, and we have $Q_{s} \psi(x)=\int_{0}^{\infty} Q_{s} \psi_{t}(x) d t$. The interversion is justified by Fubini's theorem since, using (2.5),

$$
\mathbb{E}\left[\int_{0}^{\infty}\left|\psi_{t}\right|\left(X_{s}(x)\right) d t\right] \leq \int_{0}^{\infty} \sup _{y \in B}\left|\psi_{t}(y)\right| d t<\infty
$$

where $B \subset \mathcal{X}$ is a bounded set such that $\forall \sigma \in[0, s], \quad X_{\sigma}(x) \in B, \mathbb{P}$ a.s. It follows that

$$
\frac{Q_{s} \psi(x)-\psi(x)}{s}=\int_{0}^{\infty} \frac{Q_{s} \psi_{t}(x)-\psi_{t}(x)}{s} d t
$$

For $x \in B(0, r)$ for some $r>0$, for $s \leq 1$, we have the bound

$$
\left|\frac{Q_{s} \psi_{t}(x)-\psi_{t}(x)}{s}\right| \leq \sup _{\sigma \in[0, s]}\left|\frac{d}{d \sigma} Q_{\sigma} \psi_{t}(x)\right|=\sup _{\sigma \in[0, s]}\left|Q_{\sigma} A \psi_{t}(x)\right| \leq \sup _{y \in B}\left|A \psi_{t}(y)\right|
$$

where $B \subset \mathcal{X}$ is bounded and such that $\forall x \in B(0, r), \forall \sigma \in[0,1], \quad X_{\sigma}(x) \in B$, $\mathbb{P}$-a.s. Dominated convergence then gives the b.p.c of $\int_{0}^{\infty} s^{-1}\left(Q_{s} \psi_{t}(x)-\psi_{t}(x)\right) d t$ to $\int_{0}^{\infty} A \psi_{t}(x) d t$. Assumption (2.6) and the fact that $A \psi_{t} \in C_{b}^{l o c}(\mathcal{X})$ for all $t \geq 0$ guarantee that $\int_{0}^{\infty} A \psi_{t} d t$ is in $C_{b}^{l o c}(\mathcal{X})$, so that we can indeed state $\psi \in D(A)$.

Proposition 2.4 (Resolution of a Poisson equation $A \psi=\theta$ ). Let $\theta \in D(A) \cap C_{b}^{\text {loc }}(\mathcal{X})$. Assume that, for all $B \subset \mathcal{X}$ bounded,

$$
\begin{gather*}
\int_{0}^{\infty} \sup _{x \in B}\left|Q_{t} \theta(x)\right| d t<\infty  \tag{2.7}\\
\int_{0}^{\infty} \sup _{x \in B}\left|Q_{t} A \theta(x)\right| d t<\infty \tag{2.8}
\end{gather*}
$$

Then $\psi=-\int_{0}^{\infty} Q_{t} \theta d t \in D(A)$ and $A \psi=\theta$.

Proof. Since $Q_{t} A \theta$ is not necessarily continuous, we cannot exactly state that $\psi_{t}=Q_{t} \theta \in D(A)$. Let us hence be a little more careful: using the same arguments as the previous proof, we can show that

$$
\frac{Q_{s} \psi(x)-\psi(x)}{s} \xrightarrow[s \rightarrow 0]{\text { b.p.c }}-\int_{0}^{\infty} Q_{t} A \theta(x) d t=\lim _{T \rightarrow \infty}-\int_{0}^{T} \frac{d}{d t} Q_{t} \theta(x) d t=\theta(x)-\lim _{T \rightarrow \infty} Q_{T} \theta(x)
$$

Hence, $\lim _{T \rightarrow \infty} Q_{T} \theta(x)$ exists and assumption (2.7) implies $\lim _{T \rightarrow \infty} Q_{T} \theta(x)=0$. We have shown

$$
\frac{Q_{s} \psi-\psi}{s} \xrightarrow[s \rightarrow 0]{\text { b.p.c }} \theta
$$

and $\theta$ is in $C_{b}^{l o c}(\mathcal{X})$ by assumption, so that we have indeed $\psi \in D(A)$.
Remark 2.2. Assuming (2.7) only, one can still define the $\mathcal{B}^{\text {loc }}(\mathcal{X})$ function $\psi=-\int_{0}^{\infty} Q_{t} \theta d t$. If for some reason we know in advance that $\psi \in D(A)$, then we can conclude immediately that $A \psi=\theta$ without assumption (2.8). Indeed,

$$
Q_{s} \psi(x)=-\int_{0}^{\infty} Q_{t+s} \theta(x) d t=-\int_{s}^{\infty} Q_{t} \theta(x) d t
$$

which gives, since $t \mapsto Q_{t} \theta(x)$ is continuous,

$$
\left.\frac{d}{d s^{+}} Q_{s} \psi(x)\right|_{s=0}=Q_{0} \theta(x)=\theta(x)
$$

and using Kolmogorov's equation, we know that

$$
\left.\frac{d}{d s^{+}} Q_{s} \psi(x)\right|_{s=0}=A \psi(x)
$$

### 2.2 The random driving term

We now describe the assumptions made on the process $m=\left(m_{t}\right)_{t \geq 0}$ driving equation (1.3). Let us define its canonical filtration

$$
\mathcal{F}_{t}=\sigma\left(m_{s}, 0 \leq s \leq t\right), \quad t \geq 0
$$

We assume that $m=\left(m_{t}\right)_{t \geq 0}$ is a càdlàg, stochastically continuous Markov process taking values in some separable, complete normed space $\left(E,\|\cdot\|_{E}\right)$, satisfying the continuous embedding

$$
\begin{equation*}
m_{t} \in E \subset W_{x}^{2, \infty}, \text { with }\|\cdot\|_{W_{x}^{2, \infty}} \lesssim\|\cdot\|_{E} \tag{2.9}
\end{equation*}
$$

This enables us to consider linear test functions of the form

$$
\psi_{h}: n \in E \mapsto\langle h, n\rangle, \quad h \in\left(W_{x}^{2, \infty}\right)^{\prime}
$$

Assumption 1 (Bounded state space). For some $C^{*}>0$,

$$
\begin{equation*}
\|n\|_{E} \leq C^{*} \text { for all } n \in E \tag{2.10}
\end{equation*}
$$

We denote by $\left(P_{t}\right)_{t \geq 0}$ the semigroup associated to $m$ and by $(M, D(M))$ its infinitesimal generator, in the sense of the b.p.c convergence introduced in section 2.1.

Assumption 2 (Stationary measure). There exists some measure $\nu \in \mathcal{P}(E)$ satisfying:

$$
\begin{align*}
& \forall t \geq 0, \forall \psi \in \mathcal{B}(E), \quad \int_{E} P_{t} \psi(n) d \nu(n)=\int_{E} \psi(n) d \nu(n)  \tag{Stationary}\\
& \forall h \in\left(W_{x}^{2, \infty}\right)^{\prime}, \quad \int \psi_{h}(n) d \nu(n)=0 \tag{Centered}
\end{align*}
$$

We shall denote by $\left(\widetilde{m}_{t}\right)_{t \geq 0}$ the process with initial law $\nu$.

Note that it is possible to build a càdlàg stationary process $\left(\widehat{m}_{t}\right)_{t \in \mathbb{R}}$ indexed by $t \in \mathbb{R}$ : one may indeed use the Kolmogorov extension theorem with the finite dimensional distributions

$$
\mathbb{P}\left[\widehat{m}_{s} \in A_{0}, \widehat{m}_{s+t_{1}} \in A_{1}, \ldots \widehat{m}_{s+t_{n}} \in A_{n}\right]=\mathbb{P}\left[\widetilde{m}_{0} \in A_{0}, \widetilde{m}_{t_{1}} \in A_{1}, \ldots \widetilde{m}_{t_{n}} \in A_{n}\right]
$$

for $s \in \mathbb{R}$ and $t_{1}, \ldots, t_{n} \geq 0$. This extended stationary process will still be denoted by $\left(\widetilde{m}_{t}\right)_{t \in \mathbb{R}}$.
Furthermore, we require that $m$ satisfies the following mixing property.
Assumption 3 (Mixing). There exists a non-increasing, integrable function $\gamma_{m i x} \in L^{1}\left(\mathbb{R}^{+}\right)$such that, for any initial state $n \in E$, there exists a càdlàg coupling $\left(m_{t}^{*}(n), \widetilde{m}_{t}^{*}\right)_{t \geq 0}$, whose marginal laws are respectively those of $\left(m_{t}(n)\right)_{t \geq 0}$ and $\left(\widetilde{m}_{t}\right)_{t \geq 0}$, satisfying

$$
\begin{equation*}
\mathbb{E}\left\|m_{t}^{*}(n)-\widetilde{m}_{t}^{*}\right\|_{E} \leq \gamma_{m i x}(t), \quad t \geq 0 \tag{2.11}
\end{equation*}
$$

An immediate consequence is the following: for any Lipschitz-continuous function $\psi: E \rightarrow \mathbb{R}$ such that $\int_{E} \psi(n) d \nu(n)=0$,

$$
\begin{equation*}
\left|P_{t} \psi(n)\right|=\left|\mathbb{E}\left[\psi\left(m_{t}^{*}(n)\right)-\psi\left(\widetilde{m}_{t}^{*}\right)\right]\right| \leq\|\psi\|_{L i p} \gamma_{m i x}(t) \in L_{t}^{1}\left(\mathbb{R}^{+}\right) . \tag{2.12}
\end{equation*}
$$

As a result, the coefficient $a(x)$ and the kernel $k(x, y)$ introduced in (1.8), (1.9) are well-defined: for $k(x, y)$ we may for instance write (using the stationarity of $\widetilde{m}$ )

$$
\int_{\mathbb{R}}\left|\mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right]\right| d t=\int_{0}^{\infty}\left|\mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right]\right| d t+\int_{0}^{\infty}\left|\mathbb{E}\left[\widetilde{m}_{0}(y) \widetilde{m}_{t}(x)\right]\right| d t
$$

and note that, using the Markov property,

$$
\int_{0}^{\infty}\left|\mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right]\right| d t=\int_{0}^{\infty}\left|\int_{E} n(x) \mathbb{E}\left[m_{t}(n)(y)\right] d \nu(n)\right| d t \leq C^{*} \int_{0}^{\infty} \gamma_{m i x}(t) d t<\infty .
$$

As discussed in section 2.1, the property (2.12) allows to solve a number of Poisson equations. We may indeed define, for $h \in\left(W_{x}^{2, \infty}\right)^{\prime}$,

$$
M^{-1} \psi_{h}(n)=-\int_{0}^{\infty} P_{t} \psi_{h}(n) d t=-\int_{0}^{\infty} \mathbb{E}\left\langle h, m_{t}(n)\right\rangle d t .
$$

It is then natural to introduce the function $n \in E \mapsto M^{-1} I(n)$ defined by

$$
\forall h \in\left(W_{x}^{2, \infty}\right)^{\prime}, \quad\left\langle M^{-1} I(n), h\right\rangle=M^{-1} \psi_{h}(n) .
$$

In other words,

$$
\begin{equation*}
M^{-1} I(n)=-\int_{0}^{\infty} \mathbb{E}\left[m_{t}(n)\right] d t \in W_{x}^{2, \infty} \tag{2.13}
\end{equation*}
$$

Note that the function $M^{-1} I$ is bounded: since $\mathbb{E}\left[\widetilde{m}_{t}\right]=\int_{E} n d \nu(n)=0$ in $W_{x}^{2, \infty}$, we have

$$
\begin{aligned}
\left\|M^{-1} I(n)\right\|_{W_{x}^{2, \infty}} & \leq \int_{0}^{\infty}\left\|\mathbb{E}\left[m_{t}(n)\right]\right\|_{W_{x}^{2, \infty}} d t=\int_{0}^{\infty}\left\|\mathbb{E}\left[m_{t}^{*}(n)-\widetilde{m}_{t}^{*}\right]\right\|_{W_{x}^{2, \infty}} d t \\
& \leq \int_{0}^{\infty} \mathbb{E}\left\|m_{t}^{*}(n)-\widetilde{m}_{t}^{*}\right\|_{E} d t \leq \int_{0}^{\infty} \gamma_{m i x}(t) d t<\infty .
\end{aligned}
$$

Assumption 4 (Generator). For all $h \in\left(W_{x}^{2, \infty}\right)^{\prime}, M^{-1} \psi_{h} \in D(M)$ and $M^{-1} I$ satisfies

$$
\begin{equation*}
\left\|M^{-1} I(n)-M^{-1} I\left(n^{\prime}\right)\right\|_{W_{x}^{2, \infty}} \lesssim\left\|n-n^{\prime}\right\|_{E} . \tag{2.14}
\end{equation*}
$$

Furthermore, for all $h \in L_{x}^{1}$,

$$
\begin{equation*}
\left|M^{-1} \psi_{h}\right|^{2} \in D(M) \quad \text { and }\left.\quad|M| M^{-1} \psi_{h}\right|^{2}(n) \mid \lesssim\|h\|_{L_{x}^{1}}^{2} . \tag{2.15}
\end{equation*}
$$

Note that, as stated in Remark 2.2, the assumption $M^{-1} \psi_{h} \in D(M)$ guarantees that

$$
M\left\langle M^{-1} I(n), h\right\rangle=\langle n, h\rangle .
$$

It is typically obtained under the assumptions of Proposition 2.4. Combining (2.14) with the mixing property (2.11) essentially allows to define $M^{-2} I$, which shall prove necessary in the calculations developed later on.

Let us conclude by discussing the additional assumption (1.10) made in Theorem 2. This requires of course some more regularity for the stationary process $\widetilde{m}$. A natural setup would be the following: assume that the stationary measure $\nu$ is reversible (that is $\left(\widetilde{m}_{-t}\right)_{t \geq 0}=\left(\widetilde{m}_{t}\right)_{t \geq 0}$ in law) and supported in $W_{x}^{10+\delta, \infty}$, with

$$
\mathbb{E}\left\|\widetilde{m}_{t}\right\|_{W_{x}^{10+\delta, \infty}}^{2}=\int_{E}\|n\|_{W_{x}^{10+\delta, \infty}}^{2} d \nu(n)<\infty .
$$

Making use of the time-reversibility (1.8) and (1.9) become

$$
\begin{equation*}
a(x)=\int_{0}^{\infty} e^{-t} \mathbb{E}\left[\widetilde{m}_{0}(x) \partial_{x} \widetilde{m}_{t}(x)\right] d t, \quad k(x, y)=2 \int_{0}^{\infty} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right] d t \tag{2.16}
\end{equation*}
$$

and it follows that

$$
\|a\|_{W_{x}^{4, \infty}} \lesssim\left(\int_{0}^{\infty} e^{-t} d t\right)\left(\int_{E}\|n\|_{W_{x}^{5, \infty}}^{2} d \nu(n)\right)<\infty
$$

and

$$
\begin{aligned}
\|k(\cdot, y)\|_{W_{x}^{10+\delta, \infty}} & \leq 2 \int_{0}^{\infty}\left\|\int_{E} n \mathbb{E}\left[m_{t}(n)(y)\right] d \nu(n)\right\|_{W_{x}^{10+\delta, \infty}} d t \\
& \leq 2\left(\int_{0}^{\infty} \gamma_{m i x}(t) d t\right)\left(\int_{E}\|n\|_{W_{x}^{10+\delta, \infty}} d \nu(n)\right)<\infty .
\end{aligned}
$$

Alternatively, one could drop the reversibility assumption on $\nu$ and assume that the whole process $\left(m_{t}\right)_{t \geq 0}$ takes values in $E \subset W_{x}^{10+\delta, \infty}$ with $\|\cdot\|_{W_{x}^{10+\delta, \infty}} \lesssim\|\cdot\|_{E}$.

### 2.3 An elementary example

Let us give a simple example of a Markov process $m=\left(m_{t}\right)_{t \geq 0}$ satisfying the assumptions detailed above, in the form of a jump process. Consider a discrete-time Markov chain $\left(M_{i}\right)_{i \in \mathbb{N}}$ whose (bounded) countable state space $E$ consists of sufficiently smooth functions:

$$
E=\left\{n^{j}, j \in \mathbb{N}^{*}\right\} \subset C^{k}\left(\mathbb{T}_{x}\right), \quad\|\cdot\|_{E}:=\|\cdot\|_{C_{x}^{k}},
$$

for some $k$ large enough, say $k>10$. Let us denote by $P=\left(P\left(n, n^{\prime}\right)\right)_{n, n^{\prime} \in E}$ its transition matrix. Let us consider a probability space $(\Omega, \mathcal{F}, \mathbb{P})$ where an i.i.d sequence $\left(U_{i}\right)_{i \in \mathbb{N}^{*}}$ of random variables uniformly distributed in $[0,1)$ is given. The trajectory of $\left(M_{i}(n)\right)$ for any initial data $n \in E$ can then classically be simulated in the following manner: $M_{0}(n)=n$ and, for all $i \in \mathbb{N}$,

$$
M_{i+1}(n)=T\left(M_{i}(n) ; U_{i+1}\right),
$$

where the jump function is given by

$$
T(M ; U)=n^{j} \text { when } \sum_{\ell=1}^{j-1} P\left(M, n^{\ell}\right) \leq U<\sum_{\ell=1}^{j} P\left(M, n^{\ell}\right) .
$$

We may define the time-continuous, constant-rate jump process associated with the discrete Markov chain $\left(M_{i}\right)$ : considering a Poisson process $(N(t))_{t \geq 0}$ of rate 1 with $N(0)=0$, independent of $\left(U_{i}\right)$, we simply define

$$
m_{t}(n)=M_{N(t)}(n), \quad t \geq 0
$$

Let us assume that the transition matrix $P=\left(P\left(n, n^{\prime}\right)\right)$ is irreducible, aperiodic and positive recurrent. As a result, let us introduce the unique invariant law $\nu \in \mathcal{P}(E)$ of the chain $\left(N_{i}\right)_{i \in \mathbb{N}}$ (which we may assumed to be centered, in accordance with Assumption 2).

Let us prove that the mixing property given in Assumption 3 is satisfied. To this intent, we introduce a coupling similar to the one developed in the classical proof of the ergodic theorem for Markov chains. Let us consider a random variable $\widetilde{n} \in E$ of distribution $\nu$ and an i.i.d sequence $\left(\widetilde{U}_{i}\right)_{i \in \mathbb{N}^{*}}$ of uniform random variables, mutually independent and independent of $\left(U_{i}\right)_{i \in \mathbb{N}^{*}}$. One can then introduce the Markov chain $\left(\widetilde{M}_{i}(\widetilde{n})\right)_{i \in \mathbb{N}^{*}}$ defined by $\widetilde{M}_{0}(\widetilde{n})=\widetilde{n}$ and, for all $i \in \mathbb{N}$,

$$
\widetilde{M}_{i+1}(\widetilde{n})=T\left(\widetilde{M}_{i}(n) ; \widetilde{U}_{i+1}\right)
$$

As a result, $\left(\widetilde{M}_{i}(\widetilde{n})\right)$ defines a stationary Markov chain of transition matrix $P$, independent of $\left(M_{i}(n)\right)$. For fixed $n \in E$, we then define the stopping time

$$
\tau=\inf \left\{i \geq 0, \quad M_{i}(n)=\widetilde{M}_{i}(\widetilde{n})\right\}
$$

and the process

$$
M_{i}^{*}(n)=\left\{\begin{array}{l}
M_{i}(n) \text { if } i \leq \tau \\
\widetilde{M}_{i}(\widetilde{n}) \text { if } i>\tau
\end{array}\right.
$$

One can easily see that $\left(M_{i}^{*}(n)\right)$ defines a Markov chain of transition matrix $P$, so that

$$
\left(M_{i}^{*}(n), i \geq 0\right) \sim\left(M_{i}(n), i \geq 0\right) \text { in law. }
$$

Moreover, $\left(C_{i}\right)_{i \in \mathbb{N}}=\left(M_{i}(n), \widetilde{M}_{i}(\widetilde{n})\right)_{i \mathbb{N}}$ defines a Markov chain on $E \times E$, of transition matrix $P \otimes P$. The assumptions made on $P$ guarantee that $P \otimes P$ is irreducible, and since $\nu \otimes \nu$ is an invariant probability, it is also positive recurrent. As a result, since $\tau$ is the first hitting time of the diagonal $D=\{(n, n), \quad n \in E\}$ by $\left(C_{i}\right)$, we deduce that $\mathbb{E}[\tau]<\infty$. Now, defining the continuous-time processes

$$
m_{t}^{*}(n)=M_{N(t)}^{*}(n), \quad \widetilde{m}_{t}^{*}=\widetilde{M}_{N(t)}(\widetilde{n}), \quad t \geq 0
$$

the coupling $\left(m_{t}^{*}(n), \widetilde{m}_{t}^{*}\right)_{t \geq 0}$ has the expected marginal laws and

$$
\mathbb{P}\left[m_{t}^{*}(n) \neq \widetilde{m}_{t}^{*}\right]=\mathbb{P}[N(t)<\tau]=\mathbb{E}\left[e^{-t} \sum_{k=0}^{\tau-1} \frac{t^{k}}{k!}\right]
$$

This tends to zero as $t$ goes to infinity by dominated convergence. Regarding the integrability in time, simple calculations give

$$
\int_{0}^{\infty} \mathbb{P}\left[m_{t}^{*}(n) \neq \widetilde{m}_{t}^{*}\right] d t=\sum_{k=0}^{\infty}(k+1) \mathbb{P}(\tau>k)=\frac{1}{2} \mathbb{E}[\tau(\tau+1)]
$$

The mixing Assumption 3 is therefore satisfied as soon as the hitting time $\tau$ admits a second moment. This holds for instance if the state space $E$ is finite: it is indeed a well-known fact that hitting times of Markov chains with finite state space have exponential tails (see e.g [12], Lemma 3.25). Another possible setup is the one where the transition matrix is given by $P\left(n, n^{\prime}\right)=\nu\left(n^{\prime}\right)$,
that is when $M_{i+1}$ is drawn on $E$ according to the law $\nu$, independently of $M_{i}$. In the general case, the generator of $\left(m_{t}\right)_{t \geq 0}$ is given by

$$
M \psi(n)=\sum_{n^{\prime} \in E} P\left(n, n^{\prime}\right) \psi\left(n^{\prime}\right)-\psi(n)
$$

and we derive that $M^{-1} I=\left(M^{-1} I(n)\right)_{n \in E}$ satisfies the linear system

$$
\forall n \in E, \quad \sum_{n^{\prime} \in E} P\left(n, n^{\prime}\right) M^{-1} I\left(n^{\prime}\right)-M^{-1} I(n)=n .
$$

When the state space $E$ is finite, it is clear that $M^{-1} I$ satisfies the Assumption 4. In the case where $P\left(n, n^{\prime}\right)=\nu\left(n^{\prime}\right)$, since $\nu$ is centered, we get simply

$$
M^{-1} I(n)=-n
$$

## 3 Solutions of the kinetic system

For $k \geq 0$, let us introduce the Banach space ( $G^{k},\|\cdot\|_{G^{k}}$ ) defined by

$$
\begin{equation*}
G^{k}:=\left\{f \in L_{x, v}^{1},\|f\|_{G^{k}}:=\iint\left(1+|v|^{k}\right)|f(x, v)| d x d v<\infty\right\} . \tag{3.1}
\end{equation*}
$$

When $f \in G_{2}$, we shall denote

$$
\rho(f)(x)=\int f(x) d v, \quad J(f)(x)=\int v f(x) d v, \quad K(f)=\int v^{2} f(x) d v .
$$

Since solutions $f^{\varepsilon}$ of (1.3) are of course expected to be probability densities, we also introduce the complete normed space $\left(G_{0}^{k},\|\cdot\|_{G^{k}}\right)$ defined by

$$
\begin{equation*}
G_{0}^{k}:=G^{k} \cap\left\{f \geq 0, \quad \iint f(x, v) d x d v=1\right\} . \tag{3.2}
\end{equation*}
$$

### 3.1 Path-wise weak solutions

Let us start by stating the well-posedness of the PDE system (1.3). Taking $\varepsilon=1$ for simplicity, it may be written as

$$
\left\{\begin{array}{l}
\partial_{t} f+v \partial_{x} f+\partial_{v}[(m+u-v) f]=0,  \tag{3.3}\\
\partial_{t} u-\partial_{x}^{2} u=J(f)-\rho u .
\end{array}\right.
$$

We consider a fixed $\omega \in \Omega$, so that a trajectory $m \equiv m(\omega)$ belongs to $D([0, T] ; E)$, space of càdlàg functions taking values in the separable complete space $E$.
Whenever $w \in D([0, T] ; E)$ is given, the solution $f[w]$ to the linear conservation equation

$$
\left\{\begin{array}{l}
\partial_{t} f+v \partial_{x} f+\partial_{v}[(w-v) f]=0 \\
f_{0} \in L_{x, v}^{1}
\end{array}\right.
$$

is naturally expressed as

$$
\begin{equation*}
f_{t}[w](x, v)=e^{t} f_{0} \circ \Phi_{0}^{t}(x, v), \tag{3.4}
\end{equation*}
$$

where $\Phi_{s}^{t}(x, v)=\left(X_{s}^{t}(x, v), V_{s}^{t}(x, v)\right)$ is the flow associated to the characteristics:

$$
\begin{cases}\frac{d}{d s} X_{s}^{t}=V_{s}^{t}, & X_{t}^{t}(x, v)=x  \tag{3.5}\\ \frac{d}{d s} V_{s}^{t}=w_{s}\left(X_{s}^{t}\right)-V_{s}^{t}, & V_{t}^{t}(x, v)=v\end{cases}
$$

which is clearly well defined since $w$ is globally Lipschitz-continuous in the space variable. Note that the fact that $w$ is only càdlàg in time causes no issue since this differential problem should be thought of in integral form. It is in fact easy to check that, for $x \in D([0, T] ; \mathbb{R})$,

$$
\forall t \in[0, T], \quad y_{t}=\int_{0}^{t} x_{s} d s \quad \text { if and only if } \forall t \in[0, T], \frac{d y}{d t^{+}}=x_{t}, \frac{d y}{d t^{-}}=x_{t^{-}}
$$

In the following, we will sometimes write "full" derivatives $\frac{d}{d t}$ for simplicity, but rigorously speaking one should rather consider left and right derivatives. We may now state our wellposedness result.
Proposition 3.1. Assume that

$$
\begin{equation*}
\int_{x} \int_{v}\left(1+|v|^{4}\right)\left|f_{0}\right|^{2} d x d v+\int_{x} \int_{v}(1+|v|)\left|\nabla_{x, v} f_{0}\right| d x d v<\infty, \quad u_{0} \in H_{x}^{\eta} \text { for } \eta \in(3 / 2,2) . \tag{3.6}
\end{equation*}
$$

Then, for all $\beta \in(3 / 2, \eta)$, there exists a unique couple $(f, u)$ with $u \in C\left([0, T] ; H_{x}^{\beta}\right)$ solution of the system (3.3) in the sense that

$$
\left\{\begin{array}{l}
f_{t}=f_{t}[m+u] \quad \text { as defined in (3.4) }  \tag{3.7}\\
u_{t}=S(t) u_{0}+\int_{0}^{t} S(t-s)\left[J\left(f_{s}\right)-\rho_{s} u_{s}\right] d s
\end{array}\right.
$$

where $S(t)=e^{t \partial_{x}^{2}}$ denotes the semigroup associated to the heat equation. For any $k \geq 0$, assuming $f_{0} \in G_{0}^{k}$, we have additionally $f \in C\left([0, T] ; G_{0}^{k}\right)$.

The proof of Proposition 3.1 is based on a rather classical fixed point argument. Since it is not a central aspect of the present work, it is postponed to section 8 .
Remark 3.1. The uniqueness of the solutions via an explicit iterative construction method classically guarantees that $\omega \mapsto\left(f_{t}(\omega), u_{t}(\omega)\right) \in G_{0}^{k} \times H_{x}^{\beta}$ is $\mathcal{F}_{t}=\sigma\left(m_{s}, 0 \leq s \leq t\right)$-measurable. It is of course to be expected that $\left(f_{t}, u_{t}, m_{t}\right)_{t \geq 0}$ defines a Markov process. We will not however attempt to rigorously prove it, since it is in fact not necessary for the purpose of this work.

Let us simply justify here that the unique solution $(f, u)$ in the sense of (3.7) satisfies some weak formulation of (3.3). In particular, we derive some estimates that will be useful to rigorously determine the generator of the process $\left(f_{t}, u_{t}, m_{t}\right)_{t \geq 0}$ later on.

In order to properly express the weak formulation of the equation satisfied by $f_{t} \in G^{k}$, we now introduce the space $\mathcal{C}_{k}$ which should be thought of as the natural dual of $G^{k}$ "along" equation (3.3).
Definition 3.1. For $k \geq 0$ and $\xi \in C^{1}(\mathbb{T} \times \mathbb{R})$ let us introduce the norm

$$
\begin{equation*}
N_{k}(\xi)=\left\|\frac{\xi}{1+|v|^{k}}\right\|_{L_{x, v}^{\infty}}+\left\|\frac{(1+|v|) \nabla_{x, v} \xi}{1+|v|^{k}}\right\|_{L_{x, v}^{\infty}} \tag{3.8}
\end{equation*}
$$

and the associated normed space, given by the closure

$$
\mathcal{C}_{k}=\overline{\left\{\xi \in C^{1}\left(\mathbb{T}^{d} \times \mathbb{R}^{d}\right), \quad N_{k}(\xi)<\infty\right\}} .
$$

Proposition 3.2. Let $k \geq 0$ and assume that $f_{0} \in G^{k}$. Then for all $\xi \in \mathcal{C}_{k}$,

$$
\left\langle f_{t}, \xi\right\rangle-\left\langle f_{0}, \xi\right\rangle=\int_{0}^{t}\left\langle f_{s}, v \partial_{x} \xi+\left(u_{s}+m_{s}-v\right) \partial_{v} \xi\right\rangle d s
$$

and the following estimate holds: for all $t, s \in[0, T]$,

$$
\begin{equation*}
\left|\frac{\left(f_{t+s}, \xi\right)-\left(f_{t}, \xi\right)}{s}\right| \leq C\left(k, C^{*},\|u\|_{L_{t, x}^{\infty}}, T\right) N_{k}(\xi)\left\|f_{0}\right\|_{G^{k}} \tag{3.9}
\end{equation*}
$$

Proof. Let us denote $\Phi_{t}:=\Phi_{t}^{0}$ for simplicity. Since

$$
\frac{\left\langle f_{t+s}, \xi\right\rangle-\left\langle f_{t}, \xi\right\rangle}{s}=\iint f_{0}(x, v) \frac{\xi \circ \Phi_{t+s}(x, v)-\xi \circ \Phi_{t}(x, v)}{s} d x d v
$$

we want to dominate the integrand uniformly for $t, t+s \in[0, T]$. For all $t \geq 0$,

$$
\frac{\xi \circ \Phi_{t+s}(x, v)-\xi \circ \Phi_{t}(x, v)}{s}=\frac{1}{s} \int_{t}^{t+s} \frac{d}{d \sigma} \xi\left(\Phi_{\sigma}(x, v)\right) d \sigma .
$$

with

$$
\frac{d}{d t} \xi\left(\Phi_{t}(x, v)\right)=V_{t}(x, v) \partial_{x} \xi \circ \Phi_{t}(x, v)+\left(\left(u_{t}+m_{t}\right)\left(X_{t}(x, v)\right)-V_{t}(x, v)\right) \partial_{v} \xi \circ \Phi_{t}(x, v)
$$

which, by design, can be bounded by

$$
\left(1+C^{*}+\|u\|_{L_{t, x}^{\infty}}\right) N_{k}(\xi)\left(1+\left|V_{t}(x, v)\right|^{k}\right) .
$$

From the sub-linearity of the equation of characteristics (3.5), we easily derive

$$
\forall t \in[0, T],\left|V_{t}(x, v)\right|^{k} \leq C\left(C^{*},\|u\|_{L_{t, x}^{\infty}}, T\right)\left(1+|v|^{k}\right) .
$$

Combining these bounds, we have dominated

$$
\left|f(x, v) \frac{\xi \circ \Phi_{t+s}(x, v)-\xi \circ \Phi_{t}(x, v)}{s}\right|
$$

uniformly for $t, t+s \in[0, T]$ by an expression of the form

$$
C\left(k, C^{*},\|u\|_{L_{t, x}^{\infty}}, T\right)\left|f_{0}(x, v)\right|\left(1+|v|^{k}\right)
$$

This is exactly the expected estimate, and the dominated convergence theorem leads to

$$
\frac{d}{d t}\left\langle f_{t}, \xi\right\rangle=\int_{x, v} f(x, v) \frac{d}{d t} \xi \circ \Phi_{t}(x, v) d x d v=\left\langle f_{t}, v \partial_{x} \xi+\left(u_{t}+m_{t}-v\right) \partial_{v} \xi\right\rangle
$$

Similarly, $u$ is a weak solution of the corresponding equation:
Proposition 3.3. For all $\xi \in C^{2}(\mathbb{T})$,

$$
\left\langle u_{t}, \xi\right\rangle-\left\langle u_{0}, \xi\right\rangle=\int_{0}^{t}\left\langle u_{s}, \partial_{x}^{2} \xi\right\rangle+\left\langle J\left(f_{s}\right)-\rho_{s} u_{s}, \xi\right\rangle d s
$$

We deduce the estimate: for all $t, s \in[0, T]$,

$$
\begin{equation*}
\left|\frac{\left\langle u_{t+s}, \xi\right\rangle-\left\langle u_{t}, \xi\right\rangle}{s}\right| \leq C\left(\|u\|_{L_{t, x}^{\infty}},\left\|f_{0}\right\|_{G_{1}}, T\right)\left(\|\xi\|_{\infty}+\left\|\partial_{x}^{2} \xi\right\|_{\infty}\right) \tag{3.10}
\end{equation*}
$$

Proof. The weak formulation is classically derived from the mild form (3.7) (cf. [1] for example). The estimate is easily deduced.

### 3.2 Estimates uniform in $\varepsilon$

We now establish some uniform bounds for the solution $\left(f^{\varepsilon}, u^{\varepsilon}\right)$ of (1.3).
Proposition 3.4. For all $\varepsilon \in(0,1)$, we have (almost surely)

$$
\forall t \in[0, T],\left\|f_{t}^{\varepsilon}\right\|_{L_{x, v}^{1}}=\left\|f_{0}^{\varepsilon}\right\|_{L_{x, v}^{1}}=1
$$

and

$$
\begin{aligned}
& \sup _{t \in[0, T]}\left\|f_{t}^{\varepsilon}\right\|_{G_{1}} \leq C\left(T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}},\left\|f_{0}^{\varepsilon}\right\|_{G_{1}}\right), \\
& \sup _{t \in[0, T]}\left\|u_{t}^{\varepsilon}\right\|_{L_{x}^{\infty}} \leq C\left(T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}},\left\|f_{0}^{\varepsilon}\right\|_{G_{1}}\right) .
\end{aligned}
$$

Proof. The conservation of the $L^{1}$ norm of $f^{\varepsilon}$ is clear from the transport form (3.7). Assuming $f_{0}^{\varepsilon} \geq 0, f_{0}^{\varepsilon} \in C_{c}^{\infty}$ (we conclude in the general case by density), we get

$$
\frac{d}{d t} \int_{x} \int_{v}|v| f_{t}^{\varepsilon}(x, v) d x d v=-\varepsilon^{-2} \int_{x} \int_{v}|v| \partial_{v}\left[\left(m_{t}^{\varepsilon}+\varepsilon u_{t}^{\varepsilon}-v\right) f_{t}^{\varepsilon}\right] d x d v
$$

hence denoting $\bar{J}(f)=\int_{x} \int_{v}|v| f d x d v$, (so that $\|f\|_{G_{1}}=\|f\|_{L_{x, v}^{1}}+\bar{J}(f)$ ) we get

$$
\frac{d}{d t} \bar{J}\left(f_{t}^{\varepsilon}\right) \leq \varepsilon^{-2}\left(\left\|m_{t}^{\varepsilon}\right\|_{L_{x}^{\infty}}+\varepsilon\left\|u_{t}^{\varepsilon}\right\|_{L_{x}^{\infty}}\right)-\varepsilon^{-2} \bar{J}\left(f_{t}^{\varepsilon}\right)
$$

and Grönwall's lemma gives

$$
\bar{J}\left(f_{t}^{\varepsilon}\right) \leq e^{-\varepsilon^{-2} t} \bar{J}\left(f_{0}^{\varepsilon}\right)+\int_{0}^{t} \varepsilon^{-2} e^{-\varepsilon^{-2}(t-s)}\left(\left\|m_{s}^{\varepsilon}\right\|_{L_{x}^{\infty}}+\varepsilon\left\|u_{s}^{\varepsilon}\right\|_{L_{x}^{\infty}}\right) d s
$$

We may note that $\int_{0}^{t} \varepsilon^{-2} e^{-\varepsilon^{-2} s} d s=1-e^{-\varepsilon^{-2} t} \leq 1$ so that we get

$$
\begin{equation*}
\bar{J}\left(f_{t}^{\varepsilon}\right) \leq \bar{J}\left(f_{0}^{\varepsilon}\right)+\left(C^{*}+\varepsilon \sup _{s \in[0, t]}\left\|u_{s}^{\varepsilon}\right\|_{L_{x}^{\infty}}\right) . \tag{3.11}
\end{equation*}
$$

On the other hand, using the mild form (3.7) for $u^{\varepsilon}$, usual heat-semigroup estimates, and making use of the embedding $H_{x}^{1 / 2+\delta} \subset L_{x}^{\infty}$ for any $\delta>0$, we have

$$
\begin{aligned}
\left\|u_{t}^{\varepsilon}\right\|_{L_{x}^{\infty}} & \leq\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}}+C \int_{0}^{t}|t-s|^{-1 / 2-\delta}\left\|J\left(f_{s}^{\varepsilon}\right)-\varepsilon \rho_{s}^{\varepsilon} u_{s}^{\varepsilon}\right\|_{L_{x}^{1}} d s \\
& \leq\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}}+C \int_{0}^{t}|t-s|^{-1 / 2-\delta}\left\|J\left(f_{s}^{\varepsilon}\right)\right\|_{L_{x}^{1}}+C \varepsilon \int_{0}^{t}|t-s|^{-1 / 2-\delta}\left\|u_{s}^{\varepsilon}\right\|_{L_{x}^{\infty}} d s \\
& \leq\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}}+C\left(\left\|J\left(f_{0}^{\varepsilon}\right)\right\|_{L_{x}^{1}}+C^{*}\right)+C \varepsilon \int_{0}^{t}|t-s|^{-1 / 2-\delta} \sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{\varepsilon}\right\|_{L_{x}^{\infty}} d s,
\end{aligned}
$$

using (3.11). We finally get an estimate of the form

$$
\sup _{s \in[0, t]}\left\|u^{\varepsilon}\right\|_{L_{x}^{\infty}} \leq C\left(\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}},\left\|f_{0}^{\varepsilon}\right\|_{G_{1}}, C^{*}\right)\left(1+\varepsilon \int_{0}^{t}|t-s|^{-1 / 2-\delta} \sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{\varepsilon}\right\|_{L_{x}^{\infty}} d s\right)
$$

and deduce the desired estimates using a singular Grönwall inequality.
The higher $v$-moments for $f$ are obtained by induction using the same arguments. For the sake of brevity, we omit the proof.

Proposition 3.5 (Higher moments for $f$ ). For all $k \geq 1$, we have (almost surely)

$$
\sup _{t \in[0, T]}\left\|f_{t}^{\varepsilon}\right\|_{G^{k}} \leq C\left(k, T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}},\left\|f_{0}^{\varepsilon}\right\|_{G^{k}}\right)
$$

Let us finally state some compactness estimate for $u^{\varepsilon}$.
Proposition 3.6 (Compactness for $u$ ). For $\alpha \in(1 / 2,3 / 2)$, we have (almost surely)

$$
\begin{aligned}
& \left\|u_{t}^{\varepsilon}\right\|_{H_{x}^{\alpha}} \leq C\left(\alpha, T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{H_{x}^{\alpha}},\left\|f_{0}^{\varepsilon}\right\|_{G_{1}}\right) \\
& \forall \beta<\alpha\left\|u_{t}^{\varepsilon}-u_{s}^{\varepsilon}\right\|_{H_{x}^{\beta}} \leq C\left(\alpha, \beta, T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{H_{x}^{\alpha}},\left\|f_{0}^{\varepsilon}\right\|_{G_{1}}\right) \int_{s}^{t} \theta(\sigma) d \sigma,
\end{aligned}
$$

where $\theta \in L^{1}([0, T])$ is independent of $\varepsilon$.
Proof. Consider the mild form

$$
\begin{equation*}
u_{t}^{\varepsilon}=S(t) u_{0}^{\varepsilon}+\int_{0}^{t} S(t-\sigma) g_{\sigma}^{\varepsilon} d \sigma, \quad \text { where } g_{\sigma}^{\varepsilon}=J\left(f_{\sigma}^{\varepsilon}\right)-\varepsilon \rho_{\sigma}^{\varepsilon} u_{\sigma}^{\varepsilon} \tag{3.12}
\end{equation*}
$$

The embedding $L_{x}^{1} \subset H_{x}^{-1 / 2-\delta}$ for all $\delta>0$, and a classical regularity estimate for the heat kernel $S(t)$ lead to

$$
\begin{equation*}
\|S(t) w\|_{H_{x}^{\alpha}} \lesssim\left(1+t^{-(\alpha+1 / 2+\delta) / 2}\right)\|w\|_{L_{x}^{1}} \tag{3.13}
\end{equation*}
$$

from which we derive

$$
\left\|u_{t}^{\varepsilon}\right\|_{H_{x}^{\alpha}} \lesssim\left\|u_{0}^{\varepsilon}\right\|_{H^{\alpha}}+\int_{0}^{t}\left(1+|t-\sigma|^{-(\alpha+1 / 2+\delta) / 2}\right)\left\|g_{\sigma}^{\varepsilon}\right\|_{L_{x}^{1}} d \sigma
$$

Using Proposition 3.4, we have $\left\|g_{t}^{\varepsilon}\right\|_{L_{x, v}^{1}} \leq C\left(T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{L^{\infty}},\left\|f_{0}^{\varepsilon}\right\|_{G_{1}}\right)$, which leads to the first estimate, since $(\alpha+1 / 2+\delta) / 2<1$ for $\delta$ small enough.
For the second estimate, we first write, for all $\xi \in H_{x}^{-\beta}$,

$$
\begin{aligned}
\left|\left\langle S(t) u_{0}^{\varepsilon}-S(s) u_{0}^{\varepsilon}, \xi\right\rangle\right| & =\left|\int_{s}^{t}\left\langle S(\sigma) u_{0}^{\varepsilon}, \partial_{x}^{2} \xi\right\rangle d \sigma\right| \leq\|\xi\|_{H_{x}^{-\beta}} \int_{s}^{t}\left\|S(\sigma) u_{0}^{\varepsilon}\right\|_{H_{x}^{2+\beta}} d \sigma \\
& \leq C\|\xi\|_{H_{x}^{-\beta}}\left\|u_{0}^{\varepsilon}\right\|_{H_{x}^{\alpha}} \int_{s}^{t}\left(1+\sigma^{-(2+\beta-\alpha) / 2}\right) d \sigma
\end{aligned}
$$

which gives $\left\|S(t) u_{0}^{\varepsilon}-S(s) u_{0}^{\varepsilon}\right\|_{H_{x}^{\beta}} \leq\left\|u_{0}^{\varepsilon}\right\|_{H_{x}^{\alpha}} \int_{s}^{t} \theta_{1}(\sigma) d \sigma$ with $\theta_{1} \in L^{1}$ since $(2+\beta-\alpha) / 2<1$. Finally,

$$
\left\|\int_{s}^{t} S(t-\sigma) g_{\sigma}^{\varepsilon} d \sigma\right\|_{H_{x}^{\beta}} \leq \int_{s}^{t}\left\|S(t-\sigma) g_{\sigma}^{\varepsilon} d \sigma\right\|_{H_{x}^{\beta}} d \sigma
$$

and we use again (3.13) with $\beta$ instead of $\alpha$ to conclude from the mild form (3.12).

### 3.3 Generator of the process

As previously mentioned, $\left(f_{t}^{\varepsilon}, u^{\varepsilon}\right)_{t \geq 0}$ is not a Markov process in itself, however the process $\left(f_{t}^{\varepsilon}, u^{\varepsilon}, m_{t}^{\varepsilon}\right)_{t \geq 0}$ is. In view of (1.3), its generator is naturally expected to be

$$
\begin{equation*}
\mathcal{L}^{\varepsilon} \psi=\frac{1}{\varepsilon^{2}} \mathcal{L}_{\#} \psi+\frac{1}{\varepsilon} \mathcal{L}_{b} \psi+\mathcal{L}_{0} \psi+\varepsilon \mathcal{L}_{1} \psi \tag{3.14}
\end{equation*}
$$

where

$$
\left\{\begin{array}{l}
\mathcal{L}_{\#} \psi(f, u, n)=-\left\langle\partial_{v}[(n-v) f], D_{f} \psi(f, u, n)\right\rangle+M \psi(f, u, n)  \tag{3.15}\\
\mathcal{L}_{b} \psi(f, u, n)=-\left\langle v \partial_{x} f+u \partial_{v} f, D_{f} \psi(f, u, n)\right\rangle \\
\mathcal{L}_{0} \psi(f, u, n)=\left\langle\partial_{x}^{2} u+J(f), D_{u} \psi(f, u, n)\right\rangle \\
\mathcal{L}_{1} \psi(f, u, n)=-\left\langle\rho u, D_{u} \psi(f, u, n)\right\rangle
\end{array}\right.
$$

Let us quickly get into the details of the notation " $D$ " used here. For $\psi: G^{k} \rightarrow \mathbb{R}$, let $\operatorname{Diff}_{f} \psi$ denote the usual differential in $f \in G^{k}$ of $\psi$, that is

$$
\forall f \in G^{k}, \quad \operatorname{Diff}_{f} \psi(f) \in\left(G^{k}\right)^{\prime}=L^{1}\left(\mathbb{T}^{d} \times \mathbb{R}^{d},\left(1+|v|^{k}\right) d x d v\right)^{\prime} \simeq L_{x, v}^{\infty}
$$

There exists $h(f) \in L_{x, v}^{\infty}$ such that

$$
\forall g \in G^{k}, \quad \operatorname{Diff}_{f} \psi(f)(g)=\langle g, h(f)\rangle_{G^{k},\left(G^{k}\right)^{\prime}}=\left\langle g,\left(1+|v|^{k}\right) h(f)\right\rangle
$$

The differential $\operatorname{Diff}_{f} \psi(f)$ is naturally identified with the $L_{x, v}^{\infty}$ function $h(f)$. For convenience, we shall rather denote by $D_{f} \psi(f)$ the function $\left(1+|v|^{k}\right) h(f)$, so that

$$
\forall g \in G^{k}, \quad \operatorname{Diff}_{f} \psi(f)(g)=\left\langle g, D_{f} \psi(f)\right\rangle
$$

To avoid cluttering notation, we do not indicate the dependence of $D_{f} \psi(f)$ on the choice of the parameter $k \geq 1$. In section 5 , where the perturbed test-function method is developed, the value $k=3$ is fixed for good. The natural bound on $D_{f} \psi$ is consequently

$$
\left\|\frac{D_{f} \psi(f)}{1+|v|^{k}}\right\|_{L_{x, v}^{\infty}}=\left\|\operatorname{Diff}_{f} \psi(f)\right\|_{G^{k^{\prime}}}<\infty
$$

Similarly, for $\psi: H_{x}^{\beta} \rightarrow \mathbb{R}, D_{u} \psi$ is defined so that

$$
\forall w \in H_{x}^{\beta}, \quad \operatorname{Diff}_{u} \psi(u)(w)=\left\langle w, D_{u} \psi(u)\right\rangle
$$

We will not properly identify the complete domain of the infinitesimal generator $\mathcal{L}^{\varepsilon}$ : for the purpose of this paper, it will in fact be sufficient to show that $\mathcal{L}^{\varepsilon}$ satisfies the martingale problem associated with $\left(f^{\varepsilon}, u^{\varepsilon}, m^{\varepsilon}\right)$ for a large enough class of test functions.

Definition 3.2 (Good test function). Let $k \geq 1$ and $\alpha \in(1 / 2,3 / 2)$ be given. Recalling the space $G_{0}^{k}$ introduced (3.2), define the state space for the process $\left(f_{t}^{\varepsilon}, u_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)_{t \geq 0}$ :

$$
\begin{equation*}
\mathcal{X}_{k}^{\alpha}=G_{0}^{k} \times H_{x}^{\alpha} \times E \tag{3.16}
\end{equation*}
$$

A function $\psi:(f, u, n) \in G^{k} \times H_{x}^{\alpha} \times E \rightarrow \mathbb{R}$ is said to be a "good test function" on $\mathcal{X}_{k}^{\alpha}$ when:

1. $\psi \in \mathcal{B}^{l o c}\left(\mathcal{X}_{k}^{\alpha}\right)$
2. For all $(f, u) \in G_{0}^{k} \times H_{x}^{\alpha}, \psi(f, u, \cdot) \in D(M)$ and $M \psi \in C_{b}^{l o c}\left(\mathcal{X}_{k}^{\alpha}\right)$
3. For all $(u, n) \in H_{x}^{\alpha} \times E, \psi(\cdot, u, n)$ is $G^{k}$-differentiable and

$$
\begin{array}{cccc}
D_{f} \psi: & \mathcal{X}_{k}^{\alpha} & \longrightarrow & \left(\mathcal{C}_{k}, N_{k}\right) \\
(f, u, n) & \longmapsto & D_{f} \psi(f, u, n)
\end{array} \quad \text { is } C_{b}^{l o c}\left(\mathcal{X}_{k}^{\alpha} ; \mathcal{C}_{k}\right)
$$

where the normed space $\left(\mathcal{C}_{k}, N_{k}\right)$ has been introduced in Definition (3.1)
4. For all $(f, n) \in G_{0}^{k} \times E, \psi(f, \cdot, n)$ is $H_{x}^{\alpha}$-differentiable and

$$
\begin{array}{ccc}
D_{u} \psi: & \mathcal{X}_{k}^{\alpha} & \longrightarrow C^{2}(\mathbb{T}) \\
(f, u, n) & \longmapsto & D_{u} \psi(f, u, n)
\end{array} \quad \text { is } C_{b}^{l o c}\left(\mathcal{X}_{k}^{\alpha} ; C^{2}(\mathbb{T})\right) .
$$

Note that the local boundedness required in point 3. can be expressed as

$$
\left\{\begin{array}{l}
\left|D_{f} \psi(f, u, n)\right| \leq C\left(\|f\|_{G^{k}},\|u\|_{H_{x}^{\alpha}}\right)\left(1+|v|^{k}\right) \\
\left|v \partial_{x} D_{f} \psi(f, u, n)\right| \leq C\left(\|f\|_{G^{k}},\|u\|_{H_{x}^{\alpha}}\right)\left(1+|v|^{k}\right) \\
\left|(1+|v|) \partial_{v} D \psi(f, n)\right| \leq C\left(\|f\|_{G^{k}},\|u\|_{H_{x}^{\alpha}}\right)\left(1+|v|^{k}\right)
\end{array}\right.
$$

which is a natural condition to ensure that $\mathcal{L}^{\varepsilon} \psi$ is well-defined and $C_{b}^{\text {loc }}\left(\mathcal{X}_{k}^{\alpha}\right)$. This suggests that, when $f \in G_{0}^{k}$, typical good test functions involve moments of order at most $k-1$ for $f$. Later on, we will for instance consider the test function

$$
\psi(f, n)=\left\langle J_{k-1}(f),\left(\partial_{x} M^{-1} I(n)\right) \xi\right\rangle, \quad \xi \in W_{x}^{1, \infty}
$$

and linear combinations of such functions. Introducing the canonical filtration of $m^{\varepsilon}$

$$
\mathcal{F}_{t}^{\varepsilon}=\sigma\left(m_{s}^{\varepsilon}, 0 \leq s \leq t\right)=\mathcal{F}_{\varepsilon^{-2} t}, \quad t \geq 0
$$

the process $t \mapsto\left(f_{t}^{\varepsilon}, u_{t}^{\varepsilon}\right) \in G_{0}^{k} \times H_{x}^{\alpha}$ is naturally $\left(\mathcal{F}_{t}^{\varepsilon}\right)$-adapted, as mentioned in Remark 3.1. We can now state the following result:

Proposition 3.7. Let $k \geq 1$, and $\alpha \in(1 / 2,3 / 2)$ be given. Let $\psi \equiv \psi(f, u, n)$ be a good test function on $\mathcal{X}_{k}^{\alpha}$. Then $\psi$ satisfies the martingale problem for $\mathcal{L}^{\varepsilon}$ in the sense that

$$
\forall(f, u, n) \in \mathcal{X}_{k}^{\alpha}, \quad \mathcal{L}^{\varepsilon} \psi(f, u, n) \text { is well-defined, } \quad \mathcal{L}^{\varepsilon} \psi \in C_{b}^{l o c}\left(\mathcal{X}_{k}^{\alpha}\right)
$$

and for any $\left(f_{0}^{\varepsilon}, u_{0}^{\varepsilon}, n\right) \in \mathcal{X}_{k}^{\alpha}$ (with $f_{0}^{\varepsilon}, u_{0}^{\varepsilon}$ satisfying the assumption (3.6)),

$$
M_{\psi}^{\varepsilon}(t)=\psi\left(f_{t}^{\varepsilon}, u_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)-\int_{0}^{t} \mathcal{L}^{\varepsilon} \psi\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s, \quad t \geq 0
$$

defines a càdlàg, $L^{2}$ martingale with respect to the filtration $\left(\mathcal{F}_{t}^{\varepsilon}\right)_{t \geq 0}$ under $\mathbb{P}_{\left(f_{0}^{\varepsilon}, u_{0}^{\varepsilon}, n\right)}$. Moreover, if $|\psi|^{2}$ is also a good test function, the predictable quadratic variation is given by:

$$
\left[M_{\psi}^{\varepsilon}\right](t)=\frac{1}{\varepsilon^{2}} \int_{0}^{t}\left(M|\psi|^{2}-2 \psi M \psi\right)\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m^{\varepsilon}, s\right) d s
$$

Remark 3.2. The predictable quadratic variation $[M](t)$ is the unique càdlàg, predictable, increasing process such that $|M(t)|^{2}-[M](t)$ defines a martingale. When $M(t)$ is moreover continuous, this coincides with the standard quadratic variation (see [11], p.38).
Proof. The assumptions on $\psi$ clearly guarantee that $\mathcal{L}^{\varepsilon} \psi$ is well defined and $C_{b}^{\text {loc }}\left(\mathcal{X}_{k}^{\alpha}\right)$. To simplify, we shall consider the case $\varepsilon=1$ and write $\left(f_{t}, u_{t}, m_{t}\right)$ instead of $\left(f_{t}^{\varepsilon}, u_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)$. The estimates established in the previous subsections show that, $\left(f_{0}, u_{0}\right)$ being fixed,

$$
\forall T>0, \forall t \in[0, T], \quad\left\|f_{t}\right\|_{G^{k}}+\left\|u_{t}\right\|_{H_{x}^{\alpha}} \leq C(T), \quad\left\|m_{t}\right\|_{E} \leq C^{*}
$$

almost surely under the probability $\mathbb{P}_{\left(f_{0}^{\varepsilon}, u_{0}^{\varepsilon}, n\right)}$. If follows that $M_{\psi}^{\varepsilon}(t)$ is well defined and in $L^{\infty}(\Omega)$ for $t \geq 0$. We now need to prove that, for $s<t$ and $h$ bounded and $\mathcal{F}^{\varepsilon}{ }_{s}$-measurable,

$$
\begin{equation*}
\mathbb{E}\left[\left(\psi\left(f_{t}, u_{t}, m_{t}\right)-\psi\left(f_{s}, u_{s}, m_{s}\right)-\int_{s}^{t} \mathcal{L}^{\varepsilon} \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right) d \sigma\right) h\right]=0 . \tag{3.17}
\end{equation*}
$$

To this purpose, let us formulate two natural lemmas:

Lemma 3.1. We have

$$
\begin{aligned}
\frac{d}{d t} \psi\left(f_{t}, u, n\right) & ="\left\langle D_{f} \psi\left(f_{t}, u, n\right), \partial_{t} f_{t}\right\rangle " \\
& :=\left\langle f_{t}, v \partial_{x} D_{f} \psi\left(f_{t}, u, n\right)+\left(u_{t}+m_{t}-v\right) \partial_{v} D_{f} \psi\left(f_{t}, u, n\right)\right\rangle .
\end{aligned}
$$

Proof of the lemma. One can simply write

$$
\frac{\psi\left(f_{t+s}, u, n\right)-\psi\left(f_{t}, u, n\right)}{s}=\left\langle D_{f} \psi\left(f_{t}, u, n\right), \frac{f_{t+s}-f_{t}}{s}\right\rangle+R_{s}
$$

with

$$
R_{s}=\int_{0}^{1}\left\langle D_{f} \psi\left(f_{t}+\theta\left(f_{t+s}-f_{t}\right), u, n\right)-D_{f} \psi\left(f_{t}, u, n\right), \frac{f_{t+s}-f_{t}}{s}\right\rangle d \theta
$$

Since $D_{f} \psi\left(f_{t}, u, n\right) \in \mathcal{C}_{k}$, Proposition 3.2 guarantees that

$$
\left\langle D_{f} \psi\left(f_{t}, u, n\right), \frac{f_{t+s}-f_{t}}{s}\right\rangle \underset{s \rightarrow 0}{ } "\left\langle D_{f} \psi\left(f_{t}, u, n\right), \partial_{t} f_{t}\right\rangle "
$$

and that

$$
\left|R_{s}\right| \lesssim \sup _{\theta \in[0,1]} N_{k}\left(D_{f} \psi\left(f_{t}+\theta\left(f_{t+s}-f_{t}\right), u, n\right)-D_{f} \psi\left(f_{t}, u, n\right)\right)
$$

which tends to 0 as $s \rightarrow 0$, since $f_{t+s} \rightarrow f_{t}$ in $G_{0}^{k}$ and $D_{f} \psi$ is continuous.
This second lemma is proved in a similar fashion.
Lemma 3.2. We have

$$
\begin{aligned}
\frac{d}{d t} \psi\left(f, u_{t}, n\right) & ="\left\langle D_{u} \psi\left(f, u_{t}, n\right), \partial_{t} u_{t}\right\rangle " \\
& :=\left\langle u_{s}, \partial_{x}^{2} D_{u} \psi\left(f, u_{t}, n\right)\right\rangle+\left\langle J\left(f_{t}\right)-\rho_{t} u_{t}, D_{u} \psi\left(f, u_{t}, n\right)\right\rangle .
\end{aligned}
$$

Introducing a subdivision $t_{0}=s<t_{1}<\ldots<t_{N}=t$ of $[s, t]$ of step $\delta=\max _{i}\left|t_{i+1}-t_{i}\right|$, we may split (3.17) into 3 terms:

$$
\begin{aligned}
A & =\mathbb{E}\left[\left(\sum_{i} \psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i+1}}\right)-\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i}}\right)-\int_{t_{i}}^{t_{i+1}} M \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right) d \sigma\right) h\right], \\
B & =\mathbb{E}\left[\left(\sum_{i} \psi\left(f_{t_{i+1}}, u_{t_{i+1}}, m_{t_{i+1}}\right)-\psi\left(f_{t_{i}}, u_{t_{i+1}}, m_{t_{i+1}}\right)-\int_{t_{i}}^{t_{i+1}}\left\langle D_{f} \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right), \partial_{\sigma} f_{\sigma}\right\rangle d \sigma\right) h\right], \\
C & =\mathbb{E}\left[\left(\sum_{i} \psi\left(f_{t_{i}}, u_{t_{i+1}}, m_{t_{i+1}}\right)-\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i+1}}\right)-\int_{t_{i}}^{t_{i+1}}\left\langle D_{u} \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right), \partial_{\sigma} u_{\sigma}\right\rangle d \sigma\right) h\right] .
\end{aligned}
$$

Let us treat the first term $A$ : conditioning with respect to $\mathcal{F}_{t_{i}}^{\varepsilon}$ and using the Markov property,

$$
\mathbb{E}\left[\left(\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i+1}}\right)-\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i}}\right)\right) h\right]=\mathbb{E}\left[\left(P_{t_{i+1}-t_{i}}\left[\psi\left(f_{t_{i}}, u_{t_{i}}, \cdot\right)\right]\left(m_{t_{i}}\right)-\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i}}\right)\right) h\right] .
$$

We note that, for all $t \geq 0$ and $(f, u, n) \in \mathcal{X}_{k}^{\alpha}$,

$$
P_{t}[\psi(f, u, .)](n)-\psi(f, u, n)=\int_{0}^{t} P_{\sigma} M \psi(f, u, n) d \sigma=\mathbb{E}_{n} \int_{0}^{t} M \psi\left(f, u, m_{\sigma}\right) d \sigma
$$

where we have used Kolmogorov's forward equation from Proposition 2.1 and the interversion $\mathbb{E} \int=\int \mathbb{E}$ is justified since $M \psi$ is locally bounded. We are hence led to

$$
\begin{aligned}
\mathbb{E}\left[\left(\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i+1}}\right)-\psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i}}\right)\right) h\right] & =\mathbb{E}\left[\left(\mathbb{E}_{m_{t_{i}}} \int_{0}^{t_{i+1}-t_{i}} M \psi\left(f_{t_{i}}, u_{t_{i}}, m_{\sigma}\right) d \sigma\right) h\right] \\
& =\mathbb{E}\left[\left(\int_{0}^{t_{i+1}-t_{i}} M \psi\left(f_{t_{i}}, u_{t_{i}}, m_{t_{i}+\sigma}\right) d \sigma\right) h\right] \\
& =\mathbb{E}\left[\left(\int_{t_{i}}^{t_{i+1}} M \psi\left(f_{t_{i}}, u_{t_{i}}, m_{\sigma}\right) d \sigma\right) h\right]
\end{aligned}
$$

where we have once again used the Markov property. It follows that $A$ can be rewritten as

$$
A=\mathbb{E}\left[\left(\int_{s}^{t} a_{\delta}(\sigma) d \sigma\right) h\right]
$$

with

$$
a_{\delta}(\sigma)=\sum_{i} \mathbb{1}_{\left[t_{i}, t_{i+1}\right]}(\sigma)\left(M \psi\left(f_{t_{i}}, u_{t_{i}}, m_{\sigma}\right)-M \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right)\right) .
$$

For fixed $\sigma \in[s, t]$, since $t \mapsto\left(f_{t}, u_{t}\right) \in G_{0}^{k} \times H_{x}^{\alpha}$ and $M \psi$ are continuous, it is clear that, as $\delta \rightarrow 0$, we have $a_{\delta}(\sigma) \rightarrow 0$ almost surely (hence in probability). Moreover, $M \psi$ being locally bounded, we get $a_{\delta} \in L^{\infty}([s, t] \times \Omega)$ with a bound uniform in $\delta$. In particular, $\mathbb{E}\left[a_{\delta}(\sigma) h\right]$ tends to 0 and is uniformly integrable in $\sigma \in[s, t]$, so that we may conclude $A \underset{\delta \rightarrow 0}{\longrightarrow} 0$.

Using the two lemmas introduced earlier, one can write in a similar fashion

$$
B=\mathbb{E}\left[\left(\int_{s}^{t} b_{\delta}(\sigma) d \sigma\right) h\right], \quad C=\mathbb{E}\left[\left(\int_{s}^{t} c_{\delta}(\sigma) d \sigma\right) h\right],
$$

with

$$
\begin{aligned}
& b_{\delta}(\sigma)=\sum_{i} \mathbb{1}_{\left[t_{i}, t_{i+1}\right]}(\sigma)\left\langle D_{f} \psi\left(f_{\sigma}, u_{t_{i+1}}, m_{t_{i+1}}\right)-D_{f} \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right), \partial_{\sigma} f_{\sigma}\right\rangle \\
& c_{\delta}(\sigma)=\sum_{i} \mathbb{1}_{\left[t_{i}, t_{i+1}\right]}(\sigma)\left\langle D_{u} \psi\left(f_{t_{i}}, u_{\sigma}, m_{t_{i+1}}\right)-D_{u} \psi\left(f_{\sigma}, u_{\sigma}, m_{\sigma}\right), \partial_{\sigma} u_{\sigma}\right\rangle
\end{aligned}
$$

and conclude in the same way using the estimates from Propositions 3.2 and 3.3 and the fact that $\left(m_{t}\right)_{t \geq 0}$ is stochastically continuous and $D_{f} \psi, D_{u} \psi$ are $C_{b}^{\text {loc }}\left(\mathcal{X}_{k}^{\alpha}\right)$.

If $|\psi|^{2}$ is also a good test function, then $\mathcal{L}^{\varepsilon}|\psi|^{2}$ is locally bounded, and $M_{\left.|\psi|\right|^{2}}^{\varepsilon}(t)$ defines a martingale. Using these facts, (we may refer to the proof of Theorem B. 3 in [7] for details) one can show that the predictable quadratic variation of $M_{\psi}^{\varepsilon}$ is given by

$$
d\left[M_{\psi}^{\varepsilon}\right](t)=\left(\mathcal{L}^{\varepsilon}|\psi|^{2}-2 \psi \mathcal{L}^{\varepsilon} \psi\right)\left(f_{t}^{\varepsilon}, u_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right) d t
$$

which, in our case, boils down to $\varepsilon^{-2}\left(M|\psi|^{2}-2 \psi M \psi\right)$ since $\mathcal{D}|\psi|^{2}=2 \psi \mathcal{D} \psi$ for any first order differential operator $\mathcal{D}$.

## 4 The auxiliary process

In section 5, we will be led to solve Poisson equations relative to the partial generator $\mathcal{L}_{\#}$ defined in (3.15). To this purpose, it is natural to study the auxiliary Markov process $\left(g_{t}, m_{t}\right)_{t \geq 0}$ whose generator is expected to be $\mathcal{L}_{\#}$. We hence define $g_{t} \equiv g_{t}(f, n)$ as the solution of

$$
\left\{\begin{array}{l}
\partial_{t} g_{t}+\partial_{v}\left[\left(m_{t}(n)-v\right) g_{t}\right]=0  \tag{4.1}\\
g_{0}(f, n)=f, \quad m_{0}(n)=n
\end{array}\right.
$$

and let $\left(Q_{t}\right)_{t \geq 0}$ be the corresponding semigroup

$$
\forall \psi \in \mathcal{B}^{l o c}\left(G_{0}^{k} \times E\right), \quad Q_{t} \psi(f, n)=\mathbb{E}\left[\psi\left(g_{t}(f, n), m_{t}(n)\right)\right] .
$$

The characteristics associated with the conservation equation (4.1) solve

$$
\dot{V}_{t}(v)=m_{t}(x)-V_{t}(v),
$$

leading to the explicit expression

$$
V_{t}(v)=e^{-t}\left(v+\int_{0}^{t} e^{u} m_{u}(x) d u\right)
$$

Definition 4.1. The solution of (4.1) is defined by

$$
\begin{equation*}
g_{t}(f, n)(x, v)=e^{t} f\left(x, e^{t}\left[v-w_{t}(n)(x)\right]\right), \quad t \geq 0 \tag{4.2}
\end{equation*}
$$

where

$$
\begin{equation*}
w_{t}(n)=\int_{0}^{t} e^{-(t-s)} m_{s}(n) d s=\int_{0}^{t} e^{-s} m_{t-s}(n) d s \in W_{x}^{2, \infty} . \tag{4.3}
\end{equation*}
$$

It is easy to verify that, for all $f \in G_{0}^{k}$, we have $g(f, n) \in C\left([0, T] ; G_{0}^{k}\right)$ with a bound

$$
\left\|g_{t}(f, n)\right\|_{G^{k}} \leq C\left(T, C^{*},\|f\|_{G^{k}}\right) .
$$

Additionally, $g_{t}(f, n)$ defines a weak solution of (4.1) just as in Proposition 3.2.

### 4.1 Generator of the auxiliary process

With the explicit expression (4.2), one could easily prove that the auxiliary process $\left(g_{t}, m_{t}\right)_{t \geq 0}$ rigorously defines a (locally bounded) Markov process on the state space $G_{0}^{k} \times E$. In order to solve Poisson equations later, it is of interest this time to truly characterize the generator in terms of b.p.c convergence.
Proposition 4.1. Let $k \geq 1$ be given, and let $\psi \equiv \psi(f, n)$ be a good test function on $G_{0}^{k} \times E$. Then $\psi \in D\left(\mathcal{L}_{\#}\right)$ in the sense that $\mathcal{L}_{\#} \psi \in C_{b}^{\text {loc }}\left(G_{0}^{k} \times E\right)$ and

$$
\frac{Q_{t} \psi(f, n)-\psi(f, n)}{t} \xrightarrow[t \rightarrow 0^{+}]{\text {b.p.c }} \mathcal{L}_{\#} \psi(f, n) .
$$

Proof. The properties satisfied by the good test function $\psi$ clearly guarantee that $\mathcal{L}_{\#} \psi$ is well defined and $C_{b}^{l o c}\left(G_{0}^{k} \times E\right)$. Let us denote $\left(g_{t}(f, n), m_{t}(n)\right)$ by $\left(g_{t}, m_{t}\right)$ for simplicity. We split the ratio into two terms

$$
\frac{Q_{t} \psi(f, n)-\psi(f, n)}{t}=\mathbb{E}\left[\frac{\psi\left(f, m_{t}\right)-\psi(f, n)}{t}\right]+\mathbb{E}\left[\frac{\psi\left(g_{t}, m_{t}\right)-\psi\left(f, m_{t}\right)}{t}\right] .
$$

The first term is simply $t^{-1}\left(P_{t} \psi(f, n)-\psi(f, n)\right)$ which converges to $M \psi(f, n)$ and can be bounded as $t \rightarrow 0$ by

$$
\sup _{t \in[0,1]}\left|\frac{d}{d s} P_{s} \psi(f, n)\right|=\sup _{t \in[0,1]}\left|P_{t} M \psi(f, n)\right|
$$

which is locally bounded in $(f, n)$. For the second term, we write

$$
\begin{equation*}
\frac{\psi\left(g_{t}, m_{t}\right)-\psi\left(f, m_{t}\right)}{t}=\left\langle D_{f} \psi(f, n), \frac{g_{t}-f}{t}\right\rangle+R_{t} \tag{4.4}
\end{equation*}
$$

with

$$
R_{t}=\int_{0}^{1}\left\langle D_{f} \psi\left(f+\theta\left(g_{t}-f\right), m_{t}\right)-D_{f} \psi(f, n), \frac{g_{t}-f}{t}\right\rangle d \theta .
$$

As in the proof of Lemma 3.1, we see that (4.4) converges to $\left(D_{f} \psi(f, n),-\partial_{v}[(n-v) f]\right)$ and is uniformly locally bounded in $(f, n)$ as $t \rightarrow 0$. We conclude by dominated convergence.

### 4.2 Limiting measures and mixing properties

We will now determine to some extent the limiting law of $\left(g_{t}, m_{t}\right)(f, n)$ as $t$ goes to infinity, which naturally corresponds to an invariant measure for the semi-group $\left(Q_{t}\right)_{t \geq 0}$. More precisely, we shall describe how the mixing assumption (2.11) on the driving process $\left(m_{t}\right)_{t \geq 0}$ impacts the speed of convergence of $Q_{t}$ to this invariant measure, along a certain type of test functions. From the explicit form (4.2) for example, we see that

$$
\rho\left(g_{t}(f, n)\right)=\rho(f),
$$

for all $t \geq 0$. The density $\rho$ being invariant along the trajectory $\left(g_{t}, m_{t}\right)(f, n)$, the limiting law (which we define in (4.6)) will therefore be naturally parametrized by $\rho$.

Definition 4.2. For $\rho: \mathbb{T} \rightarrow \mathbb{R}_{+}$with $\int \rho(x) d x=1$ and $w: \mathbb{\mathbb { X }} \rightarrow \mathbb{R}$, we define the probability measure $\rho \otimes \delta_{w}$ on $\mathbb{T} \times \mathbb{R}$ as

$$
\forall \xi \in C_{b}(\mathbb{T} \times \mathbb{R}), \quad\left\langle\xi, \rho \otimes \delta_{w}\right\rangle=\int_{x} \xi(x, w(x)) \rho(x) d x,
$$

Keeping in mind expression (4.2), we easily have

$$
e^{t} f\left(x, e^{t}[v-w(x)]\right) \underset{t \rightarrow \infty}{\longrightarrow} \rho(x) \otimes \delta_{w(x)}
$$

in the distributions sense. Moreover, we will see that $w_{t}(n)=\int_{0}^{t} e^{-s} m_{t-s}(n) d s \underset{t \rightarrow \infty}{\longrightarrow} \widetilde{w}$ in law, where

$$
\begin{equation*}
\widetilde{w}=\int_{-\infty}^{0} e^{s} \widetilde{m}_{s} d s \tag{4.5}
\end{equation*}
$$

We therefore expect the limiting measures to have the following form.
Definition 4.3 (Invariant measures). For $\rho: \mathbb{T} \rightarrow \mathbb{R}_{+}$with $\int \rho(x) d x=1$, we (partially) define the probability measure $\mu_{\rho}$ on $\mathcal{P}(\mathbb{T} \times \mathbb{R}) \times E$ by

$$
\begin{equation*}
\forall \psi \equiv \psi(f, n), \quad\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle:=\int \psi(f, n) d \mu_{\rho}(f, n)=\mathbb{E}\left[\psi\left(\rho \otimes \delta_{\widetilde{w}}, \widetilde{m}_{0}\right)\right] . \tag{4.6}
\end{equation*}
$$

when this expression makes sense.
Remaining on a formal level, one may compute the first moments

$$
\begin{aligned}
& \rho\left(\rho \otimes \delta_{\widetilde{w}}\right):=\int \rho \otimes \delta_{\widetilde{w}} d v=\rho, \\
& J\left(\rho \otimes \delta_{\widetilde{w}}\right):=\int v \rho \otimes \delta_{\widetilde{w}} d v=\widetilde{w} \rho, \\
& K\left(\rho \otimes \delta_{\widetilde{w}}\right):=\int v^{2} \rho \otimes \delta_{\widetilde{w}} d v=\widetilde{w}^{2} \rho .
\end{aligned}
$$

We are now ready to state the following result.
Proposition 4.2 (Mixing speed for the auxiliary process). Let us work with $f \in G_{0}^{k}$ for $k \geq 3$. Let $\psi$ be a function of the form

$$
\begin{equation*}
\psi(f, n)=\left\langle\xi_{1}(n) \rho, \xi_{2}(n)\right\rangle \text { or }\left\langle J(f), \xi_{1}(n)\right\rangle \text { or }\langle K(f), \xi\rangle \tag{4.7}
\end{equation*}
$$

with $\xi_{1}, \xi_{2}: E \rightarrow W_{x}^{1, \infty}$ Lipschitz-continuous functions:

$$
\left\|\xi_{i}(n)-\xi_{i}\left(n^{\prime}\right)\right\|_{W_{x}^{1, \infty}} \lesssim\left\|n-n^{\prime}\right\|_{E}
$$

or $\xi \in W_{x}^{1, \infty}$. Then for all $(f, n) \in G_{0}^{k} \times E$, letting $\rho=\rho(f)$, we have

$$
Q_{t} \psi(f, n) \underset{t \rightarrow \infty}{ }\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle
$$

More precisely, for all $B \subset G_{0}^{k} \times E$ bounded,

$$
\begin{aligned}
& \sup _{(f, n) \in B}\left|Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right| \in L^{1}\left(\mathbb{R}^{+}\right) \\
& \sup _{(f, n) \in B} N_{k}\left(D_{f}\left[Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right]\right) \in L^{1}\left(\mathbb{R}^{+}\right)
\end{aligned}
$$

Remark 4.1. - Note that in particular, when $\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle=0$, we have

$$
\int_{0}^{\infty} \sup _{(f, n) \in B}\left|Q_{t} \psi(f, n)\right| d t<\infty, \quad \int_{0}^{\infty} \sup _{(f, n) \in B} N_{k}\left(D_{f} Q_{t} \psi(f, n)\right) d t<\infty
$$

This allows to define $\psi_{1}(f, n)=\int_{0}^{\infty} Q_{t} \psi(f, n) d t$ which then satisfies

$$
\psi_{1} \in \mathcal{B}^{l o c}\left(G_{0}^{k} \times E\right) \text { and } D_{f} \psi_{1} \in C_{b}^{l o c}\left(G_{0}^{k} \times E ; \mathcal{C}_{k}^{1}\right)
$$

as in Definition 3.2 of a good test function.

- Later on, we will typically consider $\xi_{i}(n)=n, M^{-1} I(n), \partial_{x} M^{-1} I(n)$ which are indeed Lipschitz-continuous functions according to the assumptions made on the driving process.
- The proof of this result will in fact emphasize that it holds for any test function $\psi$ such that $\psi\left(g_{t}(f, n), m_{t}(n)\right)$ depends "at most quadratically on $n$ " with the general rule of thumb:

$$
\begin{aligned}
\xi(n)(\text { with } \xi \text { Lipschitz }) & \longleftrightarrow n \\
\rho & \longleftrightarrow \\
J(f) & \longleftrightarrow n \\
K(f) & \longleftrightarrow n^{2}
\end{aligned}
$$

For instance, the test function $\psi(f, n)=\langle\rho, \xi(n)\rangle\langle J(f), \xi\rangle$ can be considered.
Proof. Recalling (4.2), simple computations lead to

$$
\begin{aligned}
& \rho\left(g_{t}(f, n)\right)=\rho(f)=\rho \\
& J\left(g_{t}(f, n)\right)=e^{-t} J(f)+w_{t}(n) \rho(f) \\
& K\left(g_{t}(f, n)\right)=e^{-2 t} K(f)+2 e^{-t} w_{t}(n) J(f)+w_{t}(n)^{2} \rho(f) .
\end{aligned}
$$

Let us introduce the coupling $\left(m_{t}^{*}(n), \widetilde{m}_{t}^{*}\right)_{t \geq 0}$ given by the mixing Assumption 3, and define

$$
w_{t}^{*}(n):=\int_{0}^{t} e^{-s} m_{t-s}^{*}(n) d s, \quad \widetilde{w}_{t}^{*}:=\int_{-\infty}^{0} e^{s} \widetilde{m}_{t+s}^{*} d s
$$

Recalling (4.3) and (4.5), one can observe that, since $\left(\widetilde{m}_{t}\right)_{t \in \mathbb{R}}$ is stationary,

$$
\left(w_{t}(n), m_{t}(n)\right) \sim\left(w_{t}^{*}(n), m_{t}^{*}(n)\right) \text { in law, }\left(\widetilde{w}, \widetilde{m}_{0}\right) \sim\left(\widetilde{w}_{t}^{*}, \widetilde{m}_{t}^{*}\right) \text { in law. }
$$

The mixing assumption gives

$$
\begin{aligned}
\mathbb{E}\left\|w_{t}^{*}(n)-\widetilde{w}_{t}^{*}\right\|_{E} & \leq \int_{0}^{t} e^{-s} \mathbb{E}\left\|m_{t-s}^{*}(n)-\widetilde{m}_{t-s}^{*}\right\|_{E} d s+\int_{t}^{\infty} e^{-s} \mathbb{E}\left\|\widetilde{m}_{t-s}^{*}\right\|_{E} d s \\
& \lesssim \int_{0}^{t} e^{-s} \gamma_{\text {mix }}(t-s) d s+e^{-t}
\end{aligned}
$$

Noting that

$$
\begin{aligned}
& \int_{0}^{t} e^{-s} \gamma_{\text {mix }}(t-s) d s \leq \int_{0}^{\infty} e^{-s} \gamma_{\text {mix }}(t-s) d s \xrightarrow[t \rightarrow \infty]{ } 0 \\
& \int_{t=0}^{\infty} \int_{s=0}^{t} e^{-s} \gamma_{\text {mix }}(t-s) d s d t=\int_{s=0}^{\infty} e^{-s} \int_{t=s}^{\infty} \gamma_{\text {mix }}(t-s) d t d s=\int_{0}^{\infty} \gamma_{\text {mix }}(t) d t<\infty
\end{aligned}
$$

we have shown that

$$
t \mapsto \mathbb{E}\left\|w_{t}^{*}(n)-\widetilde{w}_{t}^{*}\right\|_{E} \in L^{1}\left(\mathbb{R}^{+}\right)
$$

and tends to 0 as $t \rightarrow \infty$. Let us first consider $\psi(f, n)=\left\langle\xi_{1}(n) \rho, \xi_{2}(n)\right\rangle$. In this case,

$$
\begin{aligned}
& Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle=\mathbb{E}\left\langle\xi_{1}\left(m_{t}^{*}(n)\right) \rho, \xi_{2}\left(m_{t}^{*}(n)\right)\right\rangle-\mathbb{E}\left\langle\xi_{1}\left(\widetilde{m}_{t}^{*}\right) \rho, \xi_{2}\left(\widetilde{m}_{t}^{*}\right)\right\rangle, \\
& D_{f}\left[Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right]=\mathbb{E}\left[\xi_{1}\left(m_{t}^{*}(n)\right) \xi_{2}\left(m_{t}^{*}(n)\right)-\xi_{1}\left(\widetilde{m}_{t}^{*}\right) \xi_{2}\left(\widetilde{m}_{t}^{*}\right)\right] .
\end{aligned}
$$

The Lipschitz assumption on $\xi_{i}$ and the inequality

$$
\begin{equation*}
|\langle h, k\rangle-\langle\widetilde{h}, \widetilde{k}\rangle| \leq\|h\|\|k-\widetilde{k}\|+\|k\|\|h-\widetilde{h}\| \tag{4.8}
\end{equation*}
$$

lead to

$$
\left|Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right| \lesssim\|f\|_{G_{0}} \gamma_{\operatorname{mix}}(t), \quad N_{k}\left(D_{f}\left[Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right]\right) \lesssim \gamma_{\operatorname{mix}}(t)
$$

Let us now consider $\psi(f, n)=\left\langle J(f), \xi_{1}(n)\right\rangle$. In this case,
$Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle=e^{-t} \mathbb{E}\left\langle J(f), \xi_{1}\left(m_{t}(n)\right)\right\rangle+\left[\mathbb{E}\left\langle w_{t}^{*}(n) \rho, \xi_{1}\left(m_{t}^{*}(n)\right)\right\rangle-\mathbb{E}\left\langle\widetilde{w}_{t}^{*}(n) \rho, \xi_{1}\left(\widetilde{m}_{t}^{*}\right)\right\rangle\right]$
$D_{f}\left[Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right]=e^{-t} \mathbb{E}\left[v \xi_{1}\left(m_{t}(n)\right)\right]+\mathbb{E}\left[w_{t}^{*}(n) \xi_{1}\left(m_{t}^{*}(n)\right)-\widetilde{w}_{t}^{*}(n) \xi_{1}\left(\widetilde{m}_{t}^{*}(n)\right)\right]$.
The first terms can be bounded by some $C e^{-t}$, and we can use (4.8) again to bound the second terms. Finally, let us consider $\psi(f, n)=\langle K(f), \xi\rangle$. In this case,

$$
\begin{aligned}
& Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle=e^{-2 t} \mathbb{E}\langle K(f), \xi\rangle+2 e^{-t} \mathbb{E}\left\langle w_{t}(n) J(f), \xi\right\rangle \\
&+\left[\mathbb{E}\left\langle w_{t}^{*}(n)^{2} \rho, \xi\right\rangle-\mathbb{E}\left\langle\left(\widetilde{w}_{t}^{*}\right)^{2} \rho, \xi\right\rangle\right] \\
& D_{f}\left[Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right]= e^{-2 t} v^{2} \xi+2 e^{-t} \mathbb{E}\left[w_{t}(n)\right] v \xi+\mathbb{E}\left[w_{t}^{*}(n)^{2}-\left(\widetilde{w}_{t}^{*}\right)^{2}\right] \xi
\end{aligned}
$$

We conclude in the same way.
Let us complete this section by pointing out some identities regarding the laws of $\widetilde{w}, \widetilde{m}_{0}$, and $M^{-1} I\left(\widetilde{m}_{0}\right)$ which will be useful for calculations later on.

Proposition 4.3. We have the following identities:

$$
\begin{align*}
& \mathbb{E}[\widetilde{w}(x) \widetilde{w}(y)]=\frac{1}{2} \int_{\mathbb{R}} e^{-|t|} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right] d t  \tag{4.9}\\
& \begin{aligned}
\mathbb{E}\left[\widetilde{w}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right] & =-\int_{t=0}^{\infty}\left(1-e^{-t}\right) \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t}(y)\right] d t \\
& =\mathbb{E}\left[\widetilde{m}_{0}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)+\widetilde{w}(x) \widetilde{m}_{0}(y)\right]
\end{aligned} \tag{4.10}
\end{align*}
$$

Proof. Recalling Definition 4.5 and using the fact that $\tilde{m}$ is stationary, the quantity $\mathbb{E}[\widetilde{w}(x) \widetilde{w}(y)]$ can be written

$$
\int_{s=-\infty}^{0} \int_{t=-\infty}^{0} e^{t+s} \mathbb{E}\left[\widetilde{m}_{s}(x) \widetilde{m}_{t}(y)\right] d t d s=\int_{s=-\infty}^{0} \int_{t=-\infty}^{0} e^{t+s} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t-s}(y)\right] d t d s
$$

We set $r=t-s$ and use Fubini's Theorem to obtain the expression

$$
\int_{r=-\infty}^{0}\left(\int_{s=-\infty}^{0} e^{2 s} d s\right) e^{r} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{r}(y)\right] d r+\int_{r=0}^{\infty}\left(\int_{s=-\infty}^{-r} e^{2 s} d s\right) e^{r} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{r}(y)\right] d r
$$

which gives (4.9). As for the second relation (4.10), we have

$$
\begin{equation*}
\mathbb{E}\left[\widetilde{w}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right]=\int_{-\infty}^{0} e^{s} \mathbb{E}\left[\widetilde{m}_{s}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right] d s \tag{4.12}
\end{equation*}
$$

For $s \leq 0$, recalling definition (2.13) and using the Markov property,

$$
\begin{align*}
\mathbb{E}\left[\widetilde{m}_{s}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right] & =\int_{E} \mathbb{E}\left[m_{s}(n)(x) M^{-1} I(n)\right] d \nu(n) \\
& =-\int_{E} \int_{0}^{\infty} \mathbb{E}\left[m_{s}(n)(x) \mathbb{E}\left[m_{t}(n)(y) \mid \mathcal{F}_{0}\right]\right] d t d \nu(n) \\
& =-\int_{E} \int_{0}^{\infty} \mathbb{E}\left[m_{s}(n)(x) m_{t}(n)(y)\right] d t d \nu(n) \\
& =-\int_{0}^{\infty} \mathbb{E}\left[\widetilde{m}_{s}(x) \widetilde{m}_{t}(y)\right] d t \tag{4.13}
\end{align*}
$$

It follows that (4.12) is equal to

$$
-\int_{s=-\infty}^{0} \int_{t=0}^{\infty} e^{s} \mathbb{E}\left[\widetilde{m}_{s}(x) \widetilde{m}_{t}(y)\right] d s d t=-\int_{s=-\infty}^{0} \int_{t=0}^{\infty} e^{s} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{t-s}(y)\right] d s d t
$$

Setting again $r=t-s$ and using Fubini's Theorem leads to the expression

$$
-\int_{r=0}^{\infty}\left(1-e^{-r}\right) \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{m}_{r}(y)\right] d r
$$

for (4.12). Using (4.13), we also recognize this to be (4.11)

## 5 The perturbed test function method

From this point on, we shall work in the state space $(f, u, n) \in \mathcal{X}_{k}^{\alpha}=G_{0}^{k} \times H_{x}^{\alpha} \times E$, for $k=3$ and $\alpha \in(1 / 2,3 / 2)$. The aim of this section is to identify a "limiting generator" $\mathcal{L}$ satisfying, for a large class of test functions $\phi$,

$$
\begin{equation*}
\mathcal{L}^{\varepsilon} \phi^{\varepsilon}(f, u, n)=\mathcal{L} \phi(\rho, u)+O(\varepsilon) \tag{5.1}
\end{equation*}
$$

where $\phi^{\varepsilon}$ is a perturbed version of $\phi$ of the form

$$
\phi^{\varepsilon}(f, u, n)=\phi(f, u)+\varepsilon \phi_{1}(f, u, n)+\varepsilon^{2} \phi_{2}(f, u, n)
$$

and the correctors $\phi_{1}, \phi_{2}$ are good test functions to be identified. The $O(\varepsilon)$ is to be understood as some $R^{\varepsilon}(f, u, n)$ such that $\sup _{B}\left|R^{\varepsilon}(f, u, n)\right|=O(\varepsilon)$ as $\varepsilon \rightarrow 0$ for all bounded $B \subset \mathcal{X}_{k}^{\alpha}$. Let us start with a good test function $\phi \equiv \phi(f, u) \in \mathcal{B}^{l o c}\left(\mathcal{X}_{k}^{\alpha}\right)$ of the form

$$
\begin{equation*}
\phi(f, u)=\Phi(\langle f, \xi\rangle ;\langle u, \zeta\rangle)=\Phi(\langle\rho, \xi\rangle ;\langle u, \zeta\rangle) \tag{5.2}
\end{equation*}
$$

with $\Phi \in C^{3}\left(\mathbb{R}^{2}\right)$, and $\xi \equiv \xi(x), \zeta \equiv \zeta(x) \in C^{2}(\mathbb{T})$. Since, for fixed initial data $f$, the quantity $\rho(x)$ is conserved along the auxiliary equation (4.1), in the following we shall often write $\Phi, \partial_{i} \Phi$ (etc...) instead of specifying $\Phi(\langle\rho, \xi\rangle ;\langle u, \zeta\rangle), \partial_{i} \Phi(\langle\rho, \xi\rangle ;\langle u, \zeta\rangle)$ (etc...). Identifying the different orders in $\varepsilon$ in (5.1), we get the set of equations

$$
\begin{cases}\varepsilon^{-2}: & \mathcal{L}_{\#} \phi(f, u, n)=0 \\ \varepsilon^{-1}: & \mathcal{L}_{\#} \phi_{1}(f, u, n)+\mathcal{L}_{b} \phi(f, u, n)=0, \\ \varepsilon^{0}: & \mathcal{L}_{\#} \phi_{2}(f, u, n)+\mathcal{L}_{b} \phi_{1}(f, u, n)+\mathcal{L}_{0} \phi(f, u, n)=\mathcal{L} \phi(\rho, u),\end{cases}
$$

as well as

$$
\begin{cases}\varepsilon^{1}: & \mathcal{L}_{1} \phi(f, u, n)+\mathcal{L}_{0} \phi_{1}(f, u, n)+\mathcal{L}_{b} \phi_{2}(f, u, n)=O(1)  \tag{5.3}\\ \varepsilon^{2}: & \mathcal{L}_{1} \phi_{1}(f, u, n)+\mathcal{L}_{0} \phi_{2}(f, u, n)=O\left(\varepsilon^{-1}\right) \\ \varepsilon^{3}: & \mathcal{L}_{1} \phi_{2}(f, u, n)=O\left(\varepsilon^{-2}\right)\end{cases}
$$

### 5.1 Order $\varepsilon^{-2}$

Given the form (5.2) for $\phi$, we have $\phi(f, u)=\phi(\rho(f), u)$, which is invariant along the trajectories of the auxiliary process $\left(g_{t}, m_{t}\right)_{t \geq 0}$. It follows that, indeed $\mathcal{L}_{\#} \phi(f, u, n)=0$. Furthermore, one may observe that

$$
\mathcal{L}_{b} \phi(f, u)=\left\langle J(f), \partial_{x} D_{f} \phi(f, u)\right\rangle=\left\langle J(f), \partial_{x} \xi\right\rangle \partial_{1} \Phi
$$

is a good test function satisfying the assumptions of Proposition 4.2.

### 5.2 Order $\varepsilon^{-1}$ : first corrector

We wish to solve the Poisson equation in $\phi_{1}$ :

$$
\mathcal{L}_{\#} \phi_{1}(f, u, n)=-\mathcal{L}_{b} \phi(f, u)
$$

The solution is expected to be given by

$$
\phi_{1}(f, u, n)=\int_{0}^{\infty} Q_{t}\left[\mathcal{L}_{b} \phi\right](f, u, n) d t
$$

Proposition 4.2 guarantees that this corrector is indeed well-defined and locally bounded as long as $\left\langle\left\langle\mathcal{L}_{b} \phi, \mu_{\rho}\right\rangle\right\rangle=0$, and indeed, since $\widetilde{m}$ is centered,

$$
\left\langle\left\langle\mathcal{L}_{b} \phi, \mu_{\rho}\right\rangle\right\rangle=\mathbb{E}\left\langle\rho \tilde{w}, \partial_{x} D_{f} \phi(\rho, u)\right\rangle=0
$$

The first corrector $\phi_{1}$ can in fact be calculated explicitly:

$$
\phi_{1}(f, u, n)=\int_{0}^{\infty} \mathbb{E}\left\langle f, w_{t}(n) \partial_{x} D_{f} \phi(f, u)\right\rangle d t=\left\langle f, \int_{0}^{\infty} \mathbb{E}\left[w_{t}(n)\right] d t \partial_{x} D_{f} \phi(f, u)\right\rangle
$$

with

$$
\begin{aligned}
\int_{0}^{\infty} \mathbb{E}\left[w_{t}(n)\right] d t & =\int_{t=0}^{\infty} \int_{s=0}^{t} e^{-(t-s)} \mathbb{E}\left[m_{s}(n)\right] d s d t=\int_{s=0}^{\infty} e^{s} \mathbb{E}\left[m_{s}(n)\right] \int_{t=s}^{\infty} e^{-t} d t d s \\
& =\int_{s=0}^{\infty} \mathbb{E}\left[m_{s}(n)\right] d s=-M^{-1} I(n)
\end{aligned}
$$

We hence get the expression

$$
\begin{equation*}
\phi_{1}(f, u, n)=\left\langle\partial_{x}\left[M^{-1} I(n) \rho-J(f)\right], D_{f} \phi(f, u)\right\rangle \tag{5.4}
\end{equation*}
$$

that is, with (5.2),

$$
\phi_{1}(f, u, n)=\left\langle J(f)-M^{-1} I(n) \rho, \partial_{x} \xi\right\rangle \partial_{1} \Phi
$$

which is clearly a good test function, thus in the domain of generator $\mathcal{L}_{\#}$. Making use of Remark 2.2, we have hence rigorously proven the following.

Proposition 5.1. The first corrector given by (5.4) is a good test function which satisfies

$$
\mathcal{L}_{\#} \phi_{1}(f, u, n)=-\mathcal{L}_{b} \phi(f, u)
$$

A simple calculation then leads to

$$
\begin{align*}
\left\langle D_{f} \phi_{1}(f, u, n), h\right\rangle= & \left\langle\partial_{x}\left[M^{-1} I(n) \rho(h)-J(h)\right], D_{f} \phi(\rho, u)\right\rangle  \tag{5.5}\\
& +D_{f}^{2} \phi(\rho, u)\left(\rho(h), \partial_{x}\left[M^{-1} I(n) \rho-J(f)\right]\right)
\end{align*}
$$

### 5.3 Order $\varepsilon^{0}$ : Limiting generator and second corrector

We now wish to solve the Poisson equation in $\phi_{2}$ :

$$
\mathcal{L}_{\#} \phi_{2}(f, u, n)=-\mathcal{L}_{b} \phi_{1}(f, u, n)-\mathcal{L}_{0} \phi(f, u, n)+\mathcal{L} \phi(\rho, u) .
$$

Again, the solution is expected to be given by

$$
\phi_{2}(f, u, n)=\int_{0}^{\infty} Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi(\rho, u)\right](f, u, n) d t
$$

which should converge whenever

$$
\begin{equation*}
\mathcal{L} \phi(\rho, u)=\left\langle\left\langle\mathcal{L}_{b} \phi_{1}, \mu_{\rho}\right\rangle\right\rangle+\left\langle\left\langle\mathcal{L}_{0} \phi, \mu_{\rho}\right\rangle\right\rangle . \tag{5.6}
\end{equation*}
$$

This last condition determines the form of the limiting generator $\mathcal{L} \phi$.

### 5.3.1 Limiting generator

We first compute $\mathcal{L}_{b} \phi_{1}(f, u, n)$ : letting $h=-\left(v \partial_{x} f+u \partial_{v} f\right)$ in (5.5), noting that

$$
\rho(h)=-\partial_{x} J(f), \quad J(h)=-\partial_{x} K(f)+u \rho
$$

leads to the expression

$$
\begin{align*}
\mathcal{L}_{b} \phi_{1}(f, u, n)= & \left\langle\partial_{x}\left[-M^{-1} I(n) \partial_{x} J(f)+\partial_{x} K(f)-u \rho\right], D_{f} \phi(\rho, u)\right\rangle \\
& +D_{f}^{2} \phi(\rho, u)\left(-\partial_{x} J(f), \partial_{x}\left[M^{-1} I(n) \rho-J(f)\right]\right) \tag{5.7}
\end{align*}
$$

According to (5.6), the first contribution $\mathcal{L}_{\rho}$ to the limiting generator is hence given by

$$
\begin{align*}
\mathcal{L}_{\rho} \phi(\rho, u):=\left\langle\left\langle\mathcal{L}_{b} \phi_{1}, \mu_{\rho}\right\rangle\right\rangle= & \mathbb{E}\left\langle\partial_{x}\left[\partial_{x}\left[(\widetilde{w})^{2} \rho\right]-M^{-1} I\left(\widetilde{m}_{0}\right) \partial_{x}[\widetilde{w} \rho]-u \rho\right], D_{f} \phi(\rho, u)\right\rangle \\
& +\mathbb{E} D_{f}^{2} \phi(\rho, u)\left(\partial_{x}[\widetilde{w} \rho], \partial_{x}\left[\left(\widetilde{w}-M^{-1} I\left(\widetilde{m}_{0}\right)\right) \rho\right]\right) . \tag{5.8}
\end{align*}
$$

Similarly, recalling that

$$
\begin{equation*}
\mathcal{L}_{0} \phi(f, u)=\left\langle\partial_{x}^{2} u+J(f), D_{u} \phi(f, u)\right\rangle \tag{5.9}
\end{equation*}
$$

the second contribution $\mathcal{L}_{u}$ to the limiting generator is given by

$$
\begin{equation*}
\mathcal{L}_{u} \phi(\rho, u):=\left\langle\left\langle\mathcal{L}_{0} \phi, \mu_{\rho}\right\rangle\right\rangle=\left\langle\partial_{x}^{2} u, D_{u} \phi(\rho, u)\right\rangle . \tag{5.10}
\end{equation*}
$$

Defining the operator $\mathcal{L}$ as

$$
\begin{equation*}
\mathcal{L} \phi(\rho, u)=\mathcal{L}_{\rho} \phi(\rho, u)+\mathcal{L}_{u} \phi(\rho, u) \tag{5.11}
\end{equation*}
$$

let us now describe the corresponding process.

### 5.3.2 Limiting SPDE

The limiting process corresponding to the generator $\mathcal{L}$, will be denoted by $\left(\rho_{t}, u_{t}\right)_{t \geq 0}$ for simplicity in this subsection only. Since it appears to be a diffusion process, we are naturally led to determine the associated SPDE.

Firstly, note that this limiting process $\left(\rho_{t}, u_{t}\right)_{t \geq 0}$ is decoupled: the contribution $\mathcal{L}_{u} \phi$ given by (5.10) corresponds to the simple deterministic PDE

$$
\partial_{t} u=\partial_{x}^{2} u
$$

Let us hence focus on the contribution $\mathcal{L}_{\rho} \phi$ given by (5.8) which describes the evolution of $\left(\rho_{t}\right)_{t \geq 0}$. The second order differential in $\phi$ is associated to to the diffusive part of the SPDE: we must identify the corresponding covariance operator. Let us assume that $D_{f}^{2} \phi(\rho, u)$ is associated to some symmetric kernel $\Psi$. Then,

$$
\begin{align*}
& \mathbb{E}\left[D_{f}^{2} \phi(\rho, u)\left(\partial_{x}[\widetilde{w} \rho], \partial_{x}\left[\left(\widetilde{w}-M^{-1} I\left(\widetilde{m}_{0}\right)\right) \rho\right]\right)\right] \\
& =\mathbb{E} \int_{x} \int_{y} \Psi(x, y) \partial_{x}[\widetilde{w} \rho](x) \partial_{y}\left[\left(\widetilde{w}-M^{-1} I\left(\widetilde{m}_{0}\right)\right) \rho\right](y) d x d y \\
& =\mathbb{E} \int_{x} \int_{y} \Psi(x, y) \partial_{x}\left[\partial_{y}\left[\widetilde{w}(x)\left(\widetilde{w}-M^{-1} I\left(\widetilde{m}_{0}\right)\right)(y) \rho(x) \rho(y)\right]\right] d x d y \tag{5.12}
\end{align*}
$$

We are thus naturally led to study the kernel $\mathbb{E}\left[\widetilde{w}(x)\left(\widetilde{w}(y)-M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right)\right]$. Note that since $\Psi$ is symmetric in (5.12), we may as well consider the symmetrized version

$$
\mathbb{E}[\widetilde{w}(x) \widetilde{w}(y)]-\frac{1}{2}\left(\mathbb{E}\left[\widetilde{w}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right]+\mathbb{E}\left[\widetilde{w}(y) M^{-1} I\left(\widetilde{m}_{0}\right)(x)\right]\right)
$$

Making use of the expressions obtained in Proposition 4.3 as well as the stationary character of $\left(\tilde{m}_{t}\right)_{t \in \mathbb{R}}$, we easily derive

$$
\begin{equation*}
\mathbb{E}[\widetilde{w}(x) \widetilde{w}(y)]-\frac{1}{2}\left(\mathbb{E}\left[\widetilde{w}(x) M^{-1} I\left(\widetilde{m}_{0}\right)(y)\right]+\mathbb{E}\left[\widetilde{w}(y) M^{-1} I\left(\widetilde{m}_{0}\right)(x)\right]\right)=\frac{1}{2} k(x, y) \tag{5.13}
\end{equation*}
$$

where $k(x, y)$ is the kernel defined in (1.9). Let us denote by $Q$ the linear operator on $L_{x}^{2}$ associated to this kernel:

$$
\begin{equation*}
\forall f \in L_{x}^{2}, \quad Q f(x)=\int_{y} k(x, y) f(y) d y \tag{5.14}
\end{equation*}
$$

Proposition 5.2. The operator $Q$ is self-adjoint, compact and non-negative:

$$
\forall f \in L^{2}(\mathbb{T} ; \mathbb{R}) \quad(Q f, f) \geq 0
$$

We refer to [6], Lemma 1 for a proof. As a result, we can introduce the square root $Q^{1 / 2}$ of the operator $Q$, and denote by $q$ the associated kernel:

$$
Q^{1 / 2} f(x)=\int_{y} q(x, y) f(y) d y, \quad k(x, y)=\int_{z} q(x, z) q(z, y) d z
$$

We may now define the operator on $L_{x}^{2}$

$$
\begin{equation*}
\partial_{x}\left[\rho Q^{1 / 2}\right](f)(x):=\int_{y} \partial_{x}[\rho(x) q(x, y) f(y)] d y \tag{5.15}
\end{equation*}
$$

whose adjoint is given by

$$
\partial_{x}\left[\rho Q^{1 / 2}\right]^{*}(g)(y)=\int_{x} \partial_{x}[\rho(x) q(x, y)] g(x) d x
$$

A straight-forward calculation then leads to

$$
\partial_{x}\left[\rho Q^{1 / 2}\right]^{*} D^{2} \phi(\rho, u) \partial_{x}\left[\rho Q^{1 / 2}\right](f)(t)=\int_{z} N(t, z) f(z) d z
$$

where

$$
N(t, z)=\int_{x} \int_{y} \partial_{x}[\rho(x) q(x, t)] \Psi(x, y) \partial_{y}[\rho(y) q(y, z)] d x d y .
$$

We can now get back to (5.12) and derive

$$
\begin{aligned}
& \mathbb{E}\left[D_{f}^{2} \phi(\rho, u)\left(\partial_{x}[\widetilde{w} \rho], \partial_{x}\left[\left(\widetilde{w}-M^{-1} I\left(\widetilde{m}_{0}\right)\right) \rho\right]\right)\right]=\frac{1}{2} \int_{x} \int_{y} \Psi(x, y) \partial_{x}\left[\partial_{y}[\rho(x) \rho(y) k(x, y)]\right] d x d y \\
& =\frac{1}{2} \int_{x} \int_{y} \int_{z} \Psi(x, y) \partial_{x}\left[\partial_{y}[\rho(x) \rho(y) q(x, z) q(z, y)]\right] d x d y d z=\frac{1}{2} \int_{z} N(z, z) d z \\
& =\frac{1}{2} \operatorname{Tr}\left(\partial_{x}\left[\rho Q^{1 / 2}\right]^{*} D_{f}^{2} \phi(\rho) \partial_{x}\left[\rho Q^{1 / 2}\right]\right) .
\end{aligned}
$$

Consequently, the contribution $\mathcal{L}_{\rho}$ to the limiting generator $\mathcal{L}$ may finally be rewritten as

$$
\begin{equation*}
\mathcal{L}_{\rho} \phi(\rho, u)=\left\langle\mathcal{A}^{I} \rho, D_{f} \phi(\rho, u)\right\rangle+\frac{1}{2} \operatorname{Tr}\left(\partial_{x}\left[\rho Q^{1 / 2}\right]^{*} D_{f}^{2} \phi(\rho, u) \partial_{x}\left[\rho Q^{1 / 2}\right]\right), \tag{5.16}
\end{equation*}
$$

where $\mathcal{A}^{I}$ denotes the second order differential operator

$$
\begin{equation*}
\mathcal{A}^{I} \rho=\mathbb{E}\left[\partial_{x}\left(\partial_{x}\left[\widetilde{w}^{2} \rho\right]-M^{-1} I\left(\widetilde{m}_{0}\right) \partial_{x}[\widetilde{w} \rho]-\rho u\right)\right] . \tag{5.17}
\end{equation*}
$$

This naturally corresponds to the following SPDE, given in Itô form

$$
\begin{equation*}
d \rho_{t}=\mathcal{A}^{I} \rho_{t} d t+\partial_{x}\left[\rho_{t} Q^{1 / 2} d W_{t}\right] . \tag{5.18}
\end{equation*}
$$

Finally, let us determine the Stratonovich form associated with (5.18). Recalling (5.13), the Itô operator $\mathcal{A}^{I}$ may be rewritten as

$$
\mathcal{A}^{I} \rho=\frac{1}{2} \partial_{x}^{2}[k(x, x) \rho]+\partial_{x}\left[\mathbb{E}\left(\partial_{x}\left[M^{-1} I\left(\widetilde{m}_{0}\right)\right] \widetilde{w}\right) \rho\right]-\partial_{x}[u \rho] .
$$

The correction from the Itô to the Stratonovich form is given by the formula

$$
\partial_{x}\left[\rho Q^{1 / 2} \circ d W_{t}\right]=\partial_{x}\left[\rho Q^{1 / 2} d W_{t}\right]+\mathcal{A}^{I \rightarrow S} \rho
$$

where

$$
\mathcal{A}^{I \rightarrow S} \rho:=\frac{1}{2} \partial_{x}\left[\int_{y} q(x, y) \partial_{x}(\rho(x) q(x, y)) d y\right]=\frac{1}{2} \partial_{x}^{2}[k(x, x) \rho]-\frac{1}{4} \partial_{x}\left[\left(\partial_{x} k(x, x)\right) \rho\right]
$$

since $k(x, x)=\int_{y} q(x, y)^{2} d y$. Hence, (5.18) can be expressed in Stratonovich form as

$$
d \rho_{t}=\mathcal{A}^{S} \rho_{t} d t+\partial_{x}\left[\rho Q^{1 / 2} \circ d W_{t}\right]
$$

where the differential operator $\mathcal{A}^{S}$ is of order one:

$$
\mathcal{A}^{S} \rho=\mathcal{A}^{I} \rho-\mathcal{A}^{I \rightarrow S} \rho=\partial_{x}\left[\left(\mathbb{E} \partial_{x}\left[M^{-1} I\left(\widetilde{m}_{0}\right)\right] \widetilde{w}+\frac{1}{4} \partial_{x} k(x, x)-u\right) \rho\right] .
$$

Using once again expression (5.13) and the identities from Proposition 4.3, straight-forward calculations give

$$
\mathcal{A}^{S} \rho=\partial_{x}[(a-u) \rho]
$$

where $a \equiv a(x)$ is given by (1.8).

### 5.3.3 Second corrector

Using the form (5.2) for $\phi$, recalling (5.7) and (5.9), calculations lead to the expressions

$$
\begin{align*}
\mathcal{L}_{b} \phi_{1}(f, u, n)= & -\left\langle J(f), \partial_{x}\left(M^{-1} I(n) \partial_{x} \xi\right)\right\rangle \partial_{1} \Phi+\left\langle K(f), \partial_{x}^{2} \xi\right\rangle \partial_{1} \Phi+\left\langle\rho, u \partial_{x} \xi\right\rangle \partial_{1} \Phi  \tag{5.19}\\
& -\left\langle J(f), \partial_{x} \xi\right\rangle\left\langle\rho, M^{-1} I(n) \partial_{x} \xi\right\rangle \partial_{1}^{2} \Phi+\left\langle J(f), \partial_{x} \xi\right\rangle^{2} \partial_{1}^{2} \Phi
\end{align*}
$$

and

$$
\begin{equation*}
\mathcal{L}_{0} \phi(f, u)=\left\langle u, \partial_{x}^{2} \zeta\right\rangle \partial_{2} \Phi+\langle J(f), \zeta\rangle \partial_{2} \Phi . \tag{5.20}
\end{equation*}
$$

We now prove the following result.
Proposition 5.3. The second corrector given by

$$
\phi_{2}(f, u, n)=\int_{0}^{\infty} Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right](f, u, n) d t
$$

defines a good test function which satisfies

$$
\begin{equation*}
\mathcal{L}_{\#} \phi_{2}(f, u, n)=-\mathcal{L}_{b} \phi_{1}(f, u, n)-\mathcal{L}_{0} \phi(f, u, n)+\mathcal{L} \phi(\rho, u) \tag{5.21}
\end{equation*}
$$

where $\mathcal{L}$ is the limiting generator given by (5.11).
Proof. Firstly, note that $\theta:=\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi \in C_{b}^{\text {loc }}\left(\mathcal{X}_{k}^{\alpha}\right)$ (recall that $k=3$ ) has the form of a good test function. In particular, it is in the domain of $\mathcal{L}_{\#}$. For $B \subset \mathcal{X}_{k}^{\alpha}$ bounded, we shall prove the following points:

1. First point: $\int_{0}^{\infty} \sup _{(f, u, n) \in B}\left|Q_{t}\left[\mathcal{L}_{b} \phi_{1} \mathcal{L}_{0} \phi-\mathcal{L} \phi\right](f, u, n)\right| d t<\infty$,
2. Second point: $\int_{0}^{\infty} \sup _{(f, u, n) \in B}\left|Q_{t} \mathcal{L}_{\#}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right](f, u, n)\right| d t<\infty$,
3. Third point: $\int_{0}^{\infty} \sup _{(f, u, n) \in B} N_{k}\left(D_{f} Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right](f, u, n)\right) d t<\infty$,
4. Fourth point: $\int_{0}^{\infty} \sup _{(f, u, n) \in B}\left\|D_{u} Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right](f, u, n)\right\|_{C^{2}(\mathbb{T})} d t<\infty$,

Point 1 shows that $\phi_{2}$ is well defined and in $\mathcal{B}^{l o c}\left(\mathcal{X}_{k}^{\alpha}\right)$. Once point 2 is obtained, we can use Corollary 2.4 to conclude immediately that $\phi_{2} \in D\left(\mathcal{L}_{\#}\right)$ and satisfies (5.21). It then remains to prove that $\phi_{2}$ is a good test function.

Points 3 and 4 guarantee that $\phi_{2}$ meets the last two requirements of Definition 3.2. Moreover, using the same decomposition as in the proof of Proposition 4.1, we may write

$$
\frac{P_{t} \phi_{2}(f, u, n)-\phi_{2}(f, u, n)}{t}=\frac{Q_{t} \phi_{2}(f, u, n)-\phi_{2}(f, u, n)}{t}-\mathbb{E}\left[\frac{\phi_{2}\left(g_{t}, u, m_{t}\right)-\phi_{2}\left(f, u, m_{t}\right)}{t}\right]
$$

Considering the b.p.c convergence of each term of the right hand side as $t$ goes to zero, we obtain $\phi_{2}(f, u, \cdot) \in D(M)$ and

$$
M \phi_{2}(f, u, n)=\mathcal{L}_{\#} \phi_{2}(f, u, n)-\left\langle\partial_{v}[(v-n) f], D_{f} \phi_{2}(f, u, n)\right\rangle
$$

which is in $C_{b}^{l o c}\left(\mathcal{X}_{k}^{\alpha}\right)$, concluding the proof. Let us now review the different points.

First point It is clear from the explicit expressions (5.19) and (5.20) that $\mathcal{L}_{b} \phi_{1}$ and $\mathcal{L}_{0} \phi$ are linear combinations of functions satisfying the assumptions of Proposition 4.2, i.e essentially "quadratic in $n$ ". By definition, $\mathcal{L} \phi(f, u)=\left\langle\left\langle\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi, \mu_{\rho}\right\rangle\right\rangle$, so that

$$
Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right]=Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi\right]-\left\langle\left\langle\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi, \mu_{\rho}\right\rangle\right\rangle
$$

One can easily check that the estimates obtained from Proposition 4.2 are also locally bounded in $u \in L_{x}^{\infty} \subset H_{x}^{\alpha}($ since $\alpha>1 / 2)$.

Second point Let us compute explicitely $\mathcal{L}_{\#}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right]$, recalling that

$$
\mathcal{L}_{\#} \psi(f, n)=\left\langle h, D_{f} \psi(f, n)\right\rangle+M \psi(f, n)
$$

where $h=\partial_{v}[(v-n) f]$. Since $\mathcal{L} \phi \equiv \mathcal{L} \phi(\rho, u)$, we have $\mathcal{L}_{\#}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right]=\mathcal{L}_{\#} \mathcal{L}_{b} \phi_{1}+\mathcal{L}_{\#} \mathcal{L}_{0} \phi$. Using Proposition 4.2 once again, we simply need to show that the test functions involved have the right form and are centered with respect to the invariant measure $\mu_{\rho}$ of the auxiliary process. Calculations give, for $\chi \equiv \chi(x)$,

$$
\begin{aligned}
& \langle h, \chi\rangle=\langle\rho(h), \chi\rangle=0 \\
& \langle J(h), \chi\rangle=\langle n \rho-J(f), \chi\rangle \\
& \langle K(h), \chi\rangle=2\langle n J(f)-K(f), \chi\rangle
\end{aligned}
$$

Let us review the different terms in (5.19) and express the corresponding parts in $\mathcal{L}_{\#} \mathcal{L}_{b} \phi_{1}$.
The first term is $\psi_{1}(f, n)=-\left\langle J(f), \partial_{x}\left(M^{-1} I(n) \partial_{x} \xi\right)\right\rangle \partial_{1} \Phi$. We get

$$
\begin{aligned}
& M \psi_{1}(f, n)=-\left\langle J(f), \partial_{x}\left(n . \partial_{x} \xi\right)\right\rangle \partial_{1} \Phi \\
& \left\langle h, D \psi_{1}(f, n)\right\rangle=\left\langle J(f)-n \rho, \partial_{x}\left(M^{-1} I(n) . \partial_{x} \xi\right)\right\rangle \partial_{1} \Phi
\end{aligned}
$$

so that $\mathcal{L}_{\#} \psi_{1}$ has the form required in Proposition 4.2. Calculations give the following expression for $\left\langle\left\langle\mathcal{L}_{\#} \psi_{1}, \mu_{\rho}\right\rangle\right\rangle$ :

$$
\mathbb{E}\left[-\left\langle\widetilde{w} \rho, \partial_{x}\left(\widetilde{m}_{0} \cdot \partial_{x} \xi\right)\right\rangle+\left\langle\widetilde{w} \rho, \partial_{x}\left(M^{-1} I\left(\widetilde{m}_{0}\right) \cdot \partial_{x} \xi\right)\right\rangle-\left\langle\widetilde{m}_{0} \rho, \partial_{x}\left(M^{-1} I\left(\widetilde{m}_{0}\right) \cdot \partial_{x} \xi\right)\right\rangle\right] \partial_{1} \Phi
$$

Using the identity (4.11) from Proposition 4.3 , this is indeed 0.
The next term is $\psi_{2}(f, n)=\left\langle K(f), \partial_{x}^{2} \xi\right\rangle \partial_{1} \Phi$. We get

$$
\begin{aligned}
& M \psi_{2}(f, n)=0 \\
& \left\langle h, D \psi_{2}(f, n)\right\rangle=2\left\langle n J(f)-K(f), \partial_{x}^{2} \xi\right\rangle \partial_{1} \Phi
\end{aligned}
$$

so that $\mathcal{L}_{\#} \psi_{2}$ satisfies the assumption of Proposition 4.2. Here,

$$
\left\langle\left\langle\mathcal{L}_{\#} \psi_{2}, \mu_{\rho}\right\rangle\right\rangle=2 \mathbb{E}\left\langle\left(\widetilde{m}_{0} \widetilde{w}-\widetilde{w}^{2}\right) \rho, \partial_{x}^{2} \xi\right\rangle \partial_{1} \Phi .
$$

Using the identity (4.9) from Proposition 4.3 as well as the stationarity of $\tilde{m}$, we see that $\mathbb{E}\left(\widetilde{m}_{0} \widetilde{w}\right)=\mathbb{E}\left(\widetilde{w}^{2}\right)$, so that this is indeed 0 .
The next term $\left\langle\rho, u \partial_{x} \xi\right\rangle \partial_{1} \Phi$ is in fact constant along the trajectories of $\left(g_{t}, m_{t}\right)_{t \geq 0}$, giving no contribution. The following term is $\psi_{3}(f, n)=-\left\langle J(f), \partial_{x} \xi\right\rangle\left\langle\rho, M^{-1} I(n) \partial_{x} \xi\right\rangle \partial_{1}^{2} \Phi$. We get

$$
\begin{aligned}
& M \psi_{3}(f, n)=-\left\langle J(f), \partial_{x} \xi\right\rangle\left\langle\rho, n \partial_{x} \xi\right\rangle \partial_{1}^{2} \Phi \\
& \left\langle h, D \psi_{3}(f, n)\right\rangle=-\left\langle n \rho-J(f), \partial_{x} \xi\right\rangle\left\langle\rho, M^{-1} I(n) . \partial_{x} \xi\right\rangle \partial_{1}^{2} \Phi
\end{aligned}
$$

so that $\mathcal{L} \# \psi_{3}$ satisfies the assumption of Proposition 4.2. Here,

$$
\begin{aligned}
\left\langle\left\langle\mathcal{L}_{\#} \psi_{3}, \mu_{\rho}\right\rangle\right\rangle=\mathbb{E}[ & -\left\langle\widetilde{w} \rho, \partial_{x} \xi\right\rangle\left\langle\rho, \widetilde{m}_{0} \cdot \partial_{x} \xi\right\rangle-\left\langle\tilde{m}_{0} \rho, \partial_{x} \xi\right\rangle\left\langle\rho, M^{-1} I\left(\widetilde{m}_{0}\right) \cdot \partial_{x} \xi\right\rangle \\
& \left.+\left\langle\widetilde{w} \rho, \partial_{x} \xi\right\rangle\left\langle\rho, M^{-1} I\left(\widetilde{m}_{0}\right) \cdot \partial_{x} \xi\right\rangle\right] \partial_{1}^{2} \Phi .
\end{aligned}
$$

Again, using the identity (4.11) from Proposition 4.3, this is indeed 0. Finally, the last term is $\psi_{4}(f, n)=\left\langle J(f), \partial_{x} \xi\right\rangle^{2} \partial_{1}^{2} \Phi$. We get

$$
\begin{aligned}
& M \psi_{4}(f, n)=0 \\
& \left\langle h, D \psi_{4}(f, n)\right\rangle=2\left\langle n \rho-J(f), \partial_{x} \xi\right\rangle\left\langle J(f), \partial_{x} \xi\right\rangle \partial_{1}^{2} \Phi
\end{aligned}
$$

so that $\mathcal{L}_{\#} \psi_{4}$ satisfies the assumption of Proposition 4.2. Here,

$$
\begin{equation*}
\left\langle\left\langle\mathcal{L}_{\#} \psi_{4}, \mu_{\rho}\right\rangle\right\rangle=2 \mathbb{E}\left[\left\langle\widetilde{m}_{0} \rho, \partial_{x} \xi\right\rangle\left\langle\widetilde{w} \rho, \partial_{x} \xi\right\rangle-\left\langle\widetilde{w} \rho, \partial_{x} \xi\right\rangle\left\langle\widetilde{w} \rho, \partial_{x} \xi\right\rangle\right] \partial_{1}^{2} \Phi . \tag{5.22}
\end{equation*}
$$

Using the identity (4.9) as well as the stationarity of $\widetilde{m}$, we see that

$$
\frac{1}{2} \mathbb{E}\left[\widetilde{m}_{0}(x) \widetilde{w}(y)+\widetilde{m}_{0}(y) \widetilde{w}(x)\right]=\mathbb{E}[\widetilde{w}(x) \widetilde{w}(y)]
$$

hence, making use of the symmetry of (5.22), it is indeed 0 . It remains to study $\mathcal{L}_{\#} \mathcal{L}_{0} \phi(f, u)$. The first term $\left\langle u, \partial_{x}^{2} \zeta\right\rangle \partial_{2} \Phi$ is constant. As for the second term $\langle J(f), \zeta\rangle \partial_{2} \Phi$, we get

$$
\mathcal{L}_{\#}\left[\langle J(f), \zeta\rangle \partial_{2} \Phi\right]=\langle J(h), \zeta\rangle \partial_{2} \Phi=\langle n \rho-J(f), \zeta\rangle \partial_{2} \Phi
$$

which satisfies the assumption of Proposition 4.2 and

$$
\left\langle\left\langle\langle n \rho-J(f), \zeta\rangle \partial_{2} \Phi, \mu_{\rho}\right\rangle\right\rangle=\mathbb{E}\left\langle\widetilde{m}_{0} \rho-\widetilde{w} \rho, \zeta\right\rangle \partial_{2} \Phi=0 .
$$

This concludes the proof of the second point.
Third point Once again, we simply use Proposition 4.2 (the part regarding the mixing speed of $\left.D_{f} Q_{t} \psi\right)$ since, as previously mentioned,

$$
Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right]=\left(Q_{t}\left[\mathcal{L}_{b} \phi_{1}\right]-\left\langle\left\langle\mathcal{L}_{b} \phi_{1}, \mu_{\rho}\right\rangle\right\rangle\right)+\left(Q_{t}\left[\mathcal{L}_{0} \phi\right]-\left\langle\left\langle\mathcal{L}_{0} \phi, \mu_{\rho}\right\rangle\right\rangle\right)
$$

and the different terms in (5.19), (5.20) are good test functions satisfying the requested assumptions.

Fourth point Reviewing $Q_{t}\left[\mathcal{L}_{b} \phi_{1}+\mathcal{L}_{0} \phi-\mathcal{L} \phi\right]$, we see that it is a sum of terms of the form

$$
\left(Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right) \partial_{j}^{i} \Phi((f, \xi) ;(u, \zeta))
$$

with $i, j=1$ or 2 and $\psi$ of the form required in Proposition 4.2 (the terms involving $u$ in (5.19) and (5.20) are constant along the trajectories of $\left(g_{t}, m_{t}\right)_{t \geq 0}$ and are hence canceled out). The differential with respect to $u$ of such a term is hence simply

$$
\left(Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right) \partial_{2} \partial_{j}^{i} \Phi((f, \xi) ;(u, \zeta)) \zeta .
$$

Since we assumed $\Phi \in C^{3}$, its $\|\cdot\|_{C^{2}(\mathbb{T})}$ norm over $B \subset \mathcal{X}_{k}^{\alpha}$ can be bounded by

$$
C \sup _{(f, u, n) \in B}\left|Q_{t} \psi(f, n)-\left\langle\left\langle\psi, \mu_{\rho}\right\rangle\right\rangle\right|\|\zeta\|_{C^{2}(\mathbb{T})}
$$

which is integrable over $\mathbb{R}^{+}$according to Proposition 4.2.

### 5.4 Consequence

Note that since $\phi, \phi_{1}, \phi_{2}$ are good test functions, the equations with a positive order (5.3) are automatically satisfied since $\mathcal{L}_{\star} \phi, \mathcal{L}_{\star} \phi_{i}$ are locally bounded for $i=1,2$, and $\star=0,1, b$. Let us summarize the results we have obtained: for all $\phi$ good test-function of the form

$$
\phi(f)=\Phi((\rho, \xi) ;(u, \zeta))
$$

with $\Phi \in C^{3}\left(\mathbb{R}^{2}\right), \xi \equiv \xi(x), \zeta \equiv \zeta(x) \in C^{2}(\mathbb{T})$, there exist correctors $\phi_{1}$, $\phi_{2}$ which are good test functions, such that the perturbed test function $\phi^{\varepsilon}=\phi+\varepsilon \phi_{1}+\varepsilon^{2} \phi_{2}$ satisfies

$$
\mathcal{L}^{\varepsilon} \phi^{\varepsilon}(f, u, n)=\mathcal{L} \phi(\rho, u)+R^{\varepsilon}(f, u, n)
$$

where $\mathcal{L}$ is the limiting generator identified in (5.11), and $R^{\varepsilon} / \varepsilon$ is locally bounded. As a result, we know from Proposition 3.7 that

$$
M_{\phi}^{\varepsilon}(t)=\phi^{\varepsilon}\left(f_{t}^{\varepsilon}, u_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)-\int_{0}^{t} \mathcal{L}^{\varepsilon} \phi^{\varepsilon}\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s, \quad t \geq 0
$$

defines a $\left(\mathcal{F}_{t}^{\varepsilon}\right)_{t \geq 0}$-martingale, so that

$$
N_{\phi}^{\varepsilon}(t)=\phi\left(\rho_{t}^{\varepsilon}, u_{t}^{\varepsilon}\right)-\int_{0}^{t} \mathcal{L} \phi\left(\rho_{s}^{\varepsilon}, u_{s}^{\varepsilon}\right) d s, \quad t \geq 0
$$

is "almost" a martingale, up to some $O(\varepsilon)$. This is an essential step towards our diffusionapproximation result.

## 6 Diffusion-approximation

### 6.1 Tightness

We now prove the following result.
Proposition 6.1. Let $\alpha \in(1 / 2,3 / 2]$ and assume that

$$
\sup _{\varepsilon}\left\|u_{0}^{\varepsilon}\right\|_{H_{x}^{\alpha}}+\sup _{\varepsilon}\left\|f_{0}^{\varepsilon}\right\|_{G_{3}}<\infty .
$$

Then the family of random variables $\left(\rho^{\varepsilon}, u^{\varepsilon}\right)_{\varepsilon>0}$ is tight in the space

$$
C\left([0, T] ; H_{x}^{-\gamma}\right) \times C\left([0, T] ; H_{x}^{\beta}\right), \quad \gamma>1 / 2, \quad \beta<\alpha .
$$

Proof. The estimates from Proposition 3.6 guarantee that $\left(u^{\varepsilon}\right)_{\varepsilon}$ lies in a (deterministic) compact subset of $C\left([0, T] ; H_{x}^{\beta}\right)$. We only need to prove the tightness of $\left(\rho^{\varepsilon}\right)_{\varepsilon}$ in $C\left([0, T] ; H_{x}^{-\gamma}\right)$. To this intent, we will once again use the perturbed test function method, but this time only up to the first order. Firstly, note that since the injection $L_{x}^{1} \subset H_{x}^{-\gamma}$ is compact for $\gamma>1 / 2$, it is sufficient to prove

$$
\begin{align*}
& \limsup _{\varepsilon}^{\sup } \mathbb{P}\left(\left\|\rho^{\varepsilon}\right\|_{L_{t}^{\infty} L_{x}^{1}}>M\right) \xrightarrow[M \rightarrow \infty]{ } 0,  \tag{6.1}\\
& \forall r>0, \underset{\varepsilon}{\lim \sup } \mathbb{P}\left(w_{\gamma}\left(\rho^{\varepsilon} ; \delta\right)>r\right) \xrightarrow[\delta \rightarrow 0]{\longrightarrow} 0, \tag{6.2}
\end{align*}
$$

where $w_{\gamma}(\rho ; \delta)$ denote the modulus of continuity

$$
w_{\gamma}(\rho ; \delta)=\sup _{|t-s| \leq \delta}\left\|\rho_{t}-\rho_{s}\right\|_{H_{x}^{-\gamma}} .
$$

One may refer to [2], Chapter 2, for details on this tightness criterion. The first point (6.1) is obvious since $\left\|\rho_{t}^{\varepsilon}\right\|_{L_{x}^{1}}=\left\|\rho_{0}^{\varepsilon}\right\|_{L_{x}^{1}}=1$. Note that for any $\eta>\gamma$ an interpolation inequality gives, for some $\theta \in(0,1)$,

$$
w_{\gamma}(\rho ; \delta) \leq\left(2 \sup _{t \in[0, T]}\left\|\rho_{t}\right\|_{L_{x}^{1}}\right)^{\theta} w_{\eta}(\rho ; \delta)^{1-\theta}
$$

so that it is in fact sufficient to establish (6.2) for $\gamma$ large enough. Let us define the diagonal operator

$$
J=\left(I d-\Delta_{x}\right)^{-1 / 2}
$$

that is, considering the $L^{2}(\mathbb{T})$ Hilbert-base $\left\{e_{j}(x)=(2 \pi)^{-1 / 2} e^{i j x}, j \in \mathbb{Z}\right\}$,

$$
J \xi=\sum_{j}\left(1+|j|^{2}\right)^{-1 / 2}\left\langle\xi, e_{j}\right\rangle e_{j}
$$

so that

$$
\left\|J^{\gamma} \rho\right\|_{L^{2}}^{2}=\sum_{j}\left|\left\langle\rho, J^{\gamma} e_{j}\right\rangle\right|^{2}=\sum_{j}\left(1+|j|^{2}\right)^{-\gamma}\left|\left\langle\rho, e_{j}\right\rangle\right|^{2}=\|\rho\|_{H_{x}^{-\gamma}}^{2}
$$

For fixed $j \in \mathbb{Z}$, we consider the good test function

$$
\begin{equation*}
\phi_{j}(f)=\left\langle\rho, J^{\gamma} e_{j}\right\rangle \tag{6.3}
\end{equation*}
$$

which is indeed of the form (5.2) used in the perturbed test function method, with

$$
\left\|J^{\gamma} e_{j}\right\|_{C^{2}}=(2 \pi)^{-1 / 2}|j|^{2}\left(1+|j|^{2}\right)^{-\gamma / 2}
$$

The associated first corrector

$$
\begin{equation*}
\phi_{j}^{1}(f, n)=\left\langle J(f)-M^{-1} I(n) \rho, \partial_{x} J^{\gamma} e_{j}\right\rangle \tag{6.4}
\end{equation*}
$$

is designed so that $\mathcal{L}_{\#} \phi_{j}^{1}(f, n)=-\mathcal{L}_{b} \phi_{j}(f)$. Defining the first order perturbation

$$
\phi_{j}^{\varepsilon}(f, n)=\phi_{j}(\rho)+\varepsilon \phi_{j}^{1}(f, n)
$$

it results that

$$
\begin{equation*}
\mathcal{L}^{\varepsilon} \phi_{j}^{\varepsilon}(f, u, n)=\mathcal{L}_{b} \phi_{j}^{1}(f, u, n)=O(1) \tag{6.5}
\end{equation*}
$$

in the sense that it can be bounded by $C\left\|J^{\gamma} e_{j}\right\|_{C^{2}}$ when $(f, n)$ lies in a bounded $B \subset \mathcal{X}_{k}^{\alpha}$. Let us choose $\gamma$ large enough so that

$$
\sum_{j}\left\|J^{\gamma} e_{j}\right\|_{C^{2}}^{2}<\infty
$$

that is in fact, $\gamma>5 / 2$. We now define

$$
\theta^{\varepsilon}(t)=J^{-\gamma} \sum_{j} \phi_{j}^{\varepsilon}\left(f_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right) e_{j} \in H_{x}^{-\gamma}
$$

and notice that, using the form (6.4) and Proposition 3.4, for all $t \in[0, T]$,

$$
\begin{aligned}
\left\|\rho_{t}^{\varepsilon}-\theta^{\varepsilon}(t)\right\|_{H_{x}^{-\gamma}} & =\left\|J^{\gamma} \rho_{t}^{\varepsilon}-J^{\gamma} \theta^{\varepsilon}(t)\right\|_{L^{2}}=\sum_{j}\left|\phi_{j}\left(\rho_{t}^{\varepsilon}\right)-\phi_{j}^{\varepsilon}\left(f_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)\right|^{2}=\varepsilon \sum_{j}\left|\phi_{j}^{1}\left(f_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)\right|^{2} \\
& \leq \varepsilon C\left(T, C^{*},\left\|u_{0}^{\varepsilon}\right\|_{L_{x}^{\infty}},\left\|f_{0}^{\varepsilon}\right\|_{G^{1}}\right) \sum_{j}\left\|J^{\gamma} e_{j}\right\|_{C^{1}}^{2} \lesssim \varepsilon
\end{aligned}
$$

We deduce that (6.2) can be derived from the tightness of $\left(\theta^{\varepsilon}\right)_{\varepsilon}$ in $D\left([0, T] ; H_{x}^{-\gamma}\right)$ : indeed, $\rho^{\varepsilon}$ being continuous,

$$
w_{\gamma}\left(\rho^{\varepsilon} ; \delta\right) \lesssim w_{\gamma}^{\prime}\left(\rho^{\varepsilon} ; \delta\right) \lesssim w_{\gamma}^{\prime}\left(\theta^{\varepsilon} ; \delta\right)+\varepsilon
$$

where $w_{\gamma}^{\prime}$ denotes the modulus of continuity in $D\left([0, T] ; H_{x}^{-\gamma}\right)$ (see again [2], Chapter 3, for details). We apply Aldou's criterion to prove the tightness of $\left(\theta^{\varepsilon}\right)_{\varepsilon}$ and conclude: let $\delta>0$ and $\tau_{1}, \tau_{2}$ be two $\left(\mathcal{F}_{t}^{\epsilon}\right)_{t \geq 0}$-stopping time satisfying

$$
\begin{equation*}
\tau_{1} \leq \tau_{2} \leq \tau_{1}+\delta, \quad \tau_{2} \leq T \tag{6.6}
\end{equation*}
$$

Introducing the square integrable martingale

$$
M_{j}^{\varepsilon}(t)=\phi_{j}^{\varepsilon}\left(f_{t}^{\varepsilon}, m_{t}^{\varepsilon}\right)-\int_{0}^{t} \mathcal{L}^{\varepsilon} \phi_{j}^{\varepsilon}\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s
$$

we have

$$
\begin{aligned}
\mathbb{E}\left\|\theta^{\varepsilon}\left(\tau_{2}\right)-\theta^{\varepsilon}\left(\tau_{1}\right)\right\|_{H_{x}^{-\gamma}}^{2} & =\sum_{j}\left|\phi_{j}^{\varepsilon}\left(f_{\tau_{2}}^{\varepsilon}, m_{\tau_{2}}^{\varepsilon}\right)-\phi_{j}^{\varepsilon}\left(f_{\tau_{1}}^{\varepsilon}, m_{\tau_{1}}^{\varepsilon}\right)\right|^{2} \\
& =\sum_{j} \mathbb{E}\left|M_{j}^{\varepsilon}\left(\tau_{2}\right)-M_{j}^{\varepsilon}\left(\tau_{1}\right)-\int_{\tau_{1}}^{\tau_{2}} \mathcal{L}^{\varepsilon} \phi_{j}^{\varepsilon}\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s\right|^{2} .
\end{aligned}
$$

On one hand, using Doob's optional sampling theorem, we have

$$
\mathbb{E}\left|M_{j}^{\varepsilon}\left(\tau_{2}\right)-M_{j}^{\varepsilon}\left(\tau_{1}\right)\right|^{2}=\mathbb{E}\left(\left|M_{j}^{\varepsilon}\left(\tau_{2}\right)\right|^{2}-\left|M_{j}^{\varepsilon}\left(\tau_{1}\right)\right|^{2}\right)=\mathbb{E}\left(\left[M_{j}^{\varepsilon}\right]\left(\tau_{2}\right)-\left[M_{j}^{\varepsilon}\right]\left(\tau_{1}\right)\right)
$$

where $\left[M_{j}^{\varepsilon}\right](t)$ denotes the predictable quadratic variation of $M_{j}^{\varepsilon}(t)$. We know from Proposition 3.7 that it is given by

$$
\left[M_{j}^{\varepsilon}(t)\right]=\frac{1}{\varepsilon^{2}} \int_{0}^{t}\left(M\left|\phi_{j}^{\varepsilon}\right|^{2}-2 \phi_{j}^{\varepsilon} M \phi_{j}^{\varepsilon}\right)\left(f_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s=\int_{0}^{t}\left(M\left|\phi_{j}^{1}\right|^{2}-2 \phi_{j}^{1} M \phi_{j}^{1}\right)\left(f_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s
$$

after some straight-forward calculation. Using expression (6.4) and the quadratic inequality (2.15), we deduce that the integrand can be bounded over $[0, T]$ by $C\left\|J^{\gamma} e_{j}\right\|_{C^{1}}^{2}$, so that

$$
\begin{equation*}
\sum_{j} \mathbb{E}\left|M_{j}^{\varepsilon}\left(\tau_{2}\right)-M_{j}^{\varepsilon}\left(\tau_{1}\right)\right|^{2} \lesssim \delta . \tag{6.7}
\end{equation*}
$$

On the other hand, making use of (6.5),

$$
\left|\int_{s}^{t} \mathcal{L}^{\varepsilon} \phi_{j}^{\varepsilon}\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s\right|^{2} \lesssim \int_{s}^{t}\left|\mathcal{L}^{\varepsilon} \phi_{j}^{\varepsilon}\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right)\right|^{2} d s \lesssim\left\|J^{\gamma} e_{j}\right\|_{C^{2}}^{2}|t-s|
$$

so that the choice of $\gamma$ guarantees that

$$
\begin{equation*}
\sum_{j} \mathbb{E}\left|\int_{\tau_{1}}^{\tau_{2}} \mathcal{L}^{\varepsilon} \phi_{j}^{\varepsilon}\left(f_{s}^{\varepsilon}, u_{s}^{\varepsilon}, m_{s}^{\varepsilon}\right) d s\right|^{2} \lesssim \delta . \tag{6.8}
\end{equation*}
$$

Combining (6.7) and (6.8) and using Markov's inequality, we conclude that

$$
\forall r>0, \quad \lim \sup _{\varepsilon} \sup _{\tau_{1}, \tau_{2}} \mathbb{P}\left(\left|\theta^{\varepsilon}\left(\tau_{2}\right)-\theta^{\varepsilon}\left(\tau_{1}\right)\right|>r\right) \underset{\delta \rightarrow 0}{\longrightarrow} 0
$$

where the supremum is taken over all stopping times $\tau_{1}, \tau_{2}$ satisfying (6.6). Aldou's criterion is met, which concludes the proof.

In the following, let us fix $\gamma>1 / 2$ and $\beta<\alpha$ and a subsequence $\left(\rho^{\varepsilon_{j}}, u^{\varepsilon_{j}}\right)_{j \geq 1}$ which converges in law to some $(\rho, u)$. Using Skorokhod's representation theorem, we may introduce some probability space $\left(\Omega^{*}, \mathbb{P}^{*}\right)$ hosting some random variables $\left(\rho_{*}^{\varepsilon_{j}}, u_{*}^{\varepsilon_{j}}\right)_{j \geq 1},\left(\rho_{*}, u_{*}\right)$ satisfying

$$
\left(\rho_{*}^{\varepsilon_{j}}, u_{*}^{\varepsilon_{j}}\right) \sim\left(\rho^{\varepsilon_{j}}, u^{\varepsilon_{j}}\right), \quad\left(\rho_{*}, u_{*}\right) \sim(\rho, u) \quad \text { in law in } C\left([0, T] ; H_{x}^{-\gamma}\right) \times C\left([0, T] ; H_{x}^{\beta}\right)
$$

and

$$
\left(\rho_{*}^{\varepsilon_{j}}, u_{*}^{\varepsilon_{j}}\right) \xrightarrow[j \rightarrow \infty]{\longrightarrow}\left(\rho_{*}, u_{*}\right) \text { a.s in } C\left([0, T] ; H_{x}^{-\gamma}\right) \times C\left([0, T] ; H_{x}^{\beta}\right)
$$

Additionally, note that passing to the limit in the equality $\left\|\rho_{*}^{\varepsilon}(t)\right\|_{L_{x}^{1}}=1$ leads to the inequality

$$
\sup _{[0, T]}\left\|\rho_{*}(t)\right\|_{L_{x}^{1}} \leq 1 \quad \text { a.s }
$$

or equivalently, regarding the limiting law,

$$
\begin{equation*}
\mathbb{P}\left[\sup _{t \in[0, T]}\left\|\rho_{t}\right\|_{L_{x}^{1}} \leq 1\right]=1 \tag{6.9}
\end{equation*}
$$

### 6.2 Convergence of the martingale problem

Let us consider a good test-function $\phi$ of the form (5.2). The final result obtained in section 5 can be summarized as follows: for any $s_{1}<s_{2}<\ldots<s_{n}=s<t$ and any $\theta:\left(H_{x}^{-\gamma} \times H_{x}^{\beta}\right)^{n} \rightarrow \mathbb{R}$ continuous and bounded,

$$
\begin{equation*}
\mathbb{E}\left[\left(\phi\left(\rho_{t}^{\varepsilon}, u_{t}^{\varepsilon}\right)-\phi\left(\rho_{s}^{\varepsilon}, u_{s}^{\varepsilon}\right)-\int_{s}^{t} \mathcal{L} \phi\left(\rho_{\sigma}^{\varepsilon}, u_{\sigma}^{\varepsilon}\right) d \sigma\right) \theta\left(\rho_{s_{1}}^{\varepsilon}, u_{s_{1}}^{\varepsilon}, \ldots, \rho_{s_{n}}^{\varepsilon}, u_{s_{n}}^{\varepsilon}\right)\right]=O(\varepsilon) \tag{6.10}
\end{equation*}
$$

Regarding the adaptedness of the process, note that since $\omega \mapsto\left(\rho_{t}^{\varepsilon}, u_{t}^{\varepsilon}\right)(\omega) \in L_{x}^{1} \times H_{x}^{\alpha}$ is $\mathcal{F}_{t}^{\varepsilon}-$ measurable, and since the injection $L_{x}^{1} \times H_{x}^{\alpha} \subset H_{x}^{-\gamma} \times H_{x}^{\beta}$ is continuous, the random variable $\omega \mapsto\left(\rho_{t}^{\varepsilon}, u_{t}^{\varepsilon}\right)(\omega) \in H_{x}^{-\gamma} \times H_{x}^{\beta}$ is also $\mathcal{F}_{t}^{\varepsilon}$-measurable. The martingale property (6.10) only depends on the law of $\left(\rho^{\varepsilon}, u^{\varepsilon}\right) \in C\left([0, T] ; H_{x}^{-\gamma} \times H_{x}^{\beta}\right)$, so that we may in fact write, on the new probability space $\left(\Omega^{*}, \mathbb{P}^{*}\right)$,

$$
\mathbb{E}^{*}\left[\left(\phi\left(\rho_{*}^{\varepsilon_{j}}(t), u_{*}^{\varepsilon_{j}}(t)\right)-\phi\left(\rho_{*}^{\varepsilon_{j}}(s), u_{*}^{\varepsilon_{j}}(s)\right)-\int_{s}^{t} \mathcal{L} \phi\left(\rho_{*}^{\varepsilon_{j}}(\sigma), u_{*}^{\varepsilon_{j}}(\sigma)\right) d \sigma\right) \theta_{*}\right]=O(\varepsilon)
$$

where $\theta_{*}=\theta\left(\rho_{*}^{\varepsilon_{j}}\left(s_{1}\right), \ldots, u_{*}^{\varepsilon_{j}}\left(s_{n}\right)\right)$. From the expressions (5.8), (5.10) and (5.11), we easily deduce that $\mathcal{L} \phi$ is continuous and locally bounded on $H_{x}^{-\gamma} \times H_{x}^{\beta}$ as soon as the functions $\xi$ and $\zeta$ involved in (5.2) have enough regularity (namely $\xi \in H_{x}^{\gamma+2}, \zeta \in H_{x}^{2-\beta}$ ). We may then apply the dominated convergence theorem on the probability space $\Omega^{*}$ to send $\varepsilon_{j}$ to zero. Using once again the invariance of the law, we obtain

$$
\mathbb{E}\left[\left(\phi\left(\rho_{t}, u_{t}\right)-\phi\left(\rho_{s}, u_{s}\right)-\int_{s}^{t} \mathcal{L} \phi\left(\rho_{\sigma}, u_{\sigma}\right) d \sigma\right) \theta\left(\rho_{s_{1}}, u_{s_{1}}, \ldots, \rho_{s_{n}}, u_{s_{n}}\right)\right]=0
$$

This equality holding true for any $s_{1}<\ldots<s_{n}=s<t$ and $\theta$, we deduce that

$$
N_{\phi}(t)=\phi\left(\rho_{t}, u_{t}\right)-\phi\left(\rho_{0}, u_{0}\right)-\int_{0}^{t} \mathcal{L} \phi\left(\rho_{s}, u_{s}\right) d s, \quad t \geq 0
$$

defines a continuous martingale, starting at 0 , with respect to the natural filtration

$$
\mathcal{F}_{t}=\sigma\left(\left(\rho_{s}, u_{s}\right) \in H_{x}^{-\gamma} \times H_{x}^{\beta}, s \in[0, t]\right), \quad t \geq 0
$$

This reasoning also holds for the good test function $|\phi|^{2}$ which is still of the form (5.2), so that $N_{|\phi|^{2}}(t)$ defines a continuous martingale as well. As mentioned in the proof of Proposition 3.7 (see again [7], Theorem B.3) we derive that the quadratic variation of $N_{\phi}(t)$ is given by

$$
\left[N_{\phi}\right](t)=\int_{0}^{t}\left(\mathcal{L}|\phi|^{2}-2 \phi \mathcal{L} \phi\right)\left(\rho_{s}, u_{s}\right) d s
$$

Let us consider $\phi(\rho, u)=(\rho, \xi)+(u, \zeta)$. Simple calculations from the expression (5.16) give

$$
\begin{aligned}
& \mathcal{L} \phi(\rho, u)=\left(\mathcal{A}^{I} \rho, \xi\right)+\left(\partial_{x}^{2} u, \zeta\right) \\
& \left(\mathcal{L}|\phi|^{2}-2 \phi \mathcal{L} \phi\right)(\rho, u)=\int_{x} \int_{y} k(x, y) \rho(x) \partial_{x} \xi(x) \rho(y) \partial_{y} \xi(y) d x d y
\end{aligned}
$$

where the operator $\mathcal{A}^{I}$ and the kernel $k(x, y)$ are defined in (5.17) and (1.9). Let us introduce the $H_{x}^{-\gamma-2} \times H_{x}^{\beta-2}$-valued process

$$
N(t)=\left(\rho_{t}-\rho_{0}-\int_{0}^{t} \mathcal{A}^{I} \rho_{s} d s, u_{t}-u_{0}-\int_{0}^{t} \partial_{x}^{2} u_{s} d s\right), \quad t \geq 0
$$

Then for all $(\xi, \zeta) \in H_{x}^{\gamma+2} \times H_{x}^{2-\beta}$,

$$
\langle N(t),(\xi, \zeta)\rangle, \quad t \geq 0
$$

defines a continuous real-valued $\left(\mathcal{F}_{t}\right)_{t \geq 0}$-martingale of quadratic variation

$$
\begin{equation*}
[\langle N,(\xi, \zeta)\rangle](t)=\int_{0}^{t} \int_{x} \int_{y} k(x, y) \rho_{t}(x) \partial_{x} \xi(x) \rho_{t}(y) \partial_{y} \xi(y) d x d y \tag{6.11}
\end{equation*}
$$

Using the polarization formula $\left\langle N, h_{1}\right\rangle\left\langle N, h_{2}\right\rangle=\frac{1}{4}\left(\left\langle N, h_{1}+h_{2}\right\rangle-\left\langle N, h_{1}-h_{2}\right\rangle\right)$, we derive from (6.11) that

$$
\left\langle N(t),\left(\xi_{1} n \zeta_{1}\right)\right\rangle\left\langle N(t),\left(\xi_{2}, \zeta_{2}\right)\right\rangle-\left\langle V(t) \xi_{1}, \xi_{2}\right\rangle, \quad t \geq 0
$$

is a continuous real-valued $\left(\mathcal{F}_{t}\right)_{t \geq 0}$-martingale, where the operator $V(t)$ is given by

$$
\begin{aligned}
\left\langle V(t) \xi_{1}, \xi_{2}\right\rangle & =\int_{0}^{t} \int_{x} \int_{y} k(x, y) \rho_{t}(x) \partial_{x} \xi_{1}(x) \rho_{t}(y) \partial_{y} \xi_{2}(y) d x d y \\
& =\int_{z}\left(\int_{x} \partial_{x}\left(\rho_{t}(x) q(x, z)\right) \xi_{1}(x) d x\right)\left(\int_{x} \partial_{x}\left(\rho_{t}(y) q(y, z)\right) \xi_{2}(y) d y\right) d z \\
& =\int_{z}\left(\partial_{x}\left[\rho_{t} Q^{1 / 2}\right]^{*} \xi_{1}(z)\right)\left(\partial_{x}\left[\rho_{t} Q^{1 / 2}\right]^{*} \xi_{2}(z)\right) d z
\end{aligned}
$$

We have used the notation introduced in (5.15). That is in fact

$$
V(t)=\partial_{x}\left[\rho_{t} Q^{1 / 2}\right] \partial_{x}\left[\rho_{t} Q^{1 / 2}\right]^{*}
$$

We may therefore apply the martingale representation theorem [4, Theorem 8.2]: there exists some filtered probability space $\left(\widehat{\Omega},\left(\widehat{\mathcal{F}}_{t}\right)_{t \geq 0}\right)$ equipped with a Wiener process $\hat{W}$ and a process ( $\widehat{\rho}, \widehat{u}$ ) defined on $\widehat{\Omega}$ whose law is that of $(\rho, u)$, such that the process

$$
\widehat{N}(t)=\left(\widehat{\rho}_{t}-\rho_{0}-\int_{0}^{t} \mathcal{A}^{I} \widehat{\rho}_{s} d s, \widehat{u}_{t}-u_{0}-\int_{0}^{t} \partial_{x}^{2} \widehat{u}_{s} d s\right), \quad t \geq 0
$$

can be represented as the Wiener integral

$$
\widehat{N}(t)=\left(\int_{0}^{t} \partial_{x}\left[\widehat{\rho}_{t} Q^{1 / 2} d \widehat{W}_{t}\right] ; 0\right)
$$

This shows that the limiting process $\left(\rho_{t}, u_{t}\right)_{t \geq 0}$ is a weak solution of the expected SPDE (1.7), which concludes the proof of Theorem 1.

## 7 Uniqueness for the limiting SPDE

We now prove the path-wise uniqueness for the limiting SPDE stated in Theorem 2. For conveniency, let us drop the hats on the quantities $\widehat{\rho}, \widehat{u}$ introduced earlier and consider the limiting process $(\rho, u) \in C\left([0, T] ; H_{x}^{-\gamma} \times H_{x}^{\beta}\right)$, where $\rho$ satisfies the SPDE of (1.7) given on a filtered probability space $\left(\Omega,\left(\mathcal{F}_{t}\right)_{t \geq 0}\right)$ equipped with a Wiener process $W=\left(W_{t}\right)_{t \geq 0}$. Of course, $u=\left(u_{t}\right)_{t \in[0, T]}$ is completely determined as the solution the deterministic heat equation with initial data $u_{0}$. Let us denote

$$
b_{t}(x)=a(x)-u_{t}(x)
$$

with $a$ given by (1.8). Diagonalizing the compact operator $Q^{1 / 2}$ involved in (1.7), we may introduce functions $\left(\phi_{k}\right)_{k \geq 0}$ such that

$$
Q^{1 / 2} \circ d W_{t}=\sum_{k} \phi_{k} \circ d \beta_{t}^{k}, \quad\left(\beta^{k}\right) \text { independent Brownian motions. }
$$

As a consequence, equation (1.7) can be re-written as

$$
\left\{\begin{array}{l}
d \rho_{t}+\partial_{x}\left(b_{t} \rho_{t}\right) d t=\partial_{x}\left(\rho_{t} \circ \sum_{k} \phi_{k} d \beta_{t}^{k}\right) \\
\rho(0)=\rho_{0} \in L_{x}^{1}
\end{array}\right.
$$

or, equivalently, in Itô form,

$$
\left\{\begin{array}{l}
d \rho_{t}+\left(\partial_{x}\left(b_{t} \rho_{t}\right)-\frac{1}{2} \sum_{k} \partial_{x}\left(\phi_{k} \partial_{x}\left(\phi_{k} \rho_{t}\right)\right)\right) d t=\sum_{k} \partial_{x}\left(\phi_{k} \rho_{t}\right) d \beta_{t}^{k}  \tag{7.1}\\
\rho(0)=\rho_{0} \in L_{x}^{1}
\end{array}\right.
$$

Recall that, as mentioned in (6.9), the solution $\rho$ constructed as a weak limit of $\left(\rho^{\varepsilon}\right)_{\varepsilon}$ in Theorem 1 satisfies the additional bound

$$
\begin{equation*}
\sup _{t \in[0, T]}\left\|\rho_{t}\right\|_{L_{x}^{1}} \leq 1, \text { a.s. } \tag{7.2}
\end{equation*}
$$

Under the additional regularity assumptions of Theorem 2 , it is clear that $b \in L^{\infty}\left([0, T] ; W_{x}^{4, \infty}\right)$. Regarding the regularity of $\phi_{k}$, one may note that, for $\sigma>4+1 / 2$,

$$
\sum_{k}\left\|\phi_{k}\right\|_{W_{x}^{4, \infty}}^{2} \leq \sum_{k}\left\|\phi_{k}\right\|_{H_{x}^{\sigma}}^{2}=\left\|Q^{1 / 2}\right\|_{\mathcal{L}_{2}\left(L^{2} ; H_{x}^{\sigma}\right)}^{2}=\operatorname{Tr}_{L^{2}}\left(J^{-2 \sigma} Q\right)
$$

where $J$ is the positive regularizing diagonal operator introduced in (6.1). Applying Corollary C. 2 and Proposition C. 3 from appendix C of [4], we deduce, for any $\theta>1$,

$$
\operatorname{Tr}_{L_{x}^{2}}\left(J^{-2 \sigma} Q\right)=\operatorname{Tr}_{L_{x}^{2}}\left(J^{\theta} J^{-\theta} J^{-2 \sigma} Q\right) \leq \operatorname{Tr}_{L_{x}^{2}}\left(J^{\theta}\right)\left\|J^{-(\theta+2 \sigma)} Q\right\|_{\mathcal{L}\left(L_{x}^{2}\right)}
$$

On one hand, $\operatorname{Tr}_{L_{x}^{2}}\left(J^{\theta}\right)=\sum_{k}\left(1+|k|^{2}\right)^{-\theta / 2}<\infty$ since $\theta>1$. On the other hand, for $f \in L_{x}^{2}$, we have

$$
J^{-(\theta+2 \sigma)} Q f(x)=\int_{y}\left(I d-\partial_{x}^{2}\right)^{(\theta+2 \sigma) / 2} k(x, y) f(y) d y
$$

so that the bound $\left\|J^{-(\theta+2 \sigma)} Q\right\|_{\mathcal{L}\left(L_{x}^{2}\right)}<\infty$ is ensured as soon as

$$
\sup _{y \in \mathbb{T}}\|k(\cdot, y)\|_{W_{x}^{\theta+2 \sigma, \infty}}<\infty
$$

The number $\theta+2 \sigma$ can be chosen as close to $1+2 \times(4+1 / 2)=10$ as desired. Hence, the regularity assumption (1.10) of Theorem 2 leads to

$$
\sum_{k}\left\|\phi_{k}\right\|_{W_{x}^{4, \infty}}^{2}<\infty
$$

It is now clear that the proof of Theorem 2 boils down to that of the following result.

Proposition 7.1. Assume that $\rho_{0} \in L_{x}^{1}$ and

$$
\begin{equation*}
b \in L^{\infty}\left([0, T] ; W_{x}^{4, \infty}\right), \quad \sum_{k}\left\|\phi_{k}\right\|_{W_{x}^{4, \infty}}^{2}<\infty \tag{7.3}
\end{equation*}
$$

Then a solution $\rho \in C\left([0, T] ; H_{x}^{-\gamma}\right)$ a.s of (7.1) satisfying (7.2) is path-wise unique.
We are naturally led to consider the dual BSPDE (Backward SPDE) associated with (7.1), of unknown $(\psi, Z)$ :

$$
\left\{\begin{array}{l}
d \psi_{t}+\left(-b_{t} \partial_{x} \psi_{t}+\frac{1}{2} \sum_{k} \phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)-\sum_{k} \phi_{k} \partial_{x} Z_{t}^{k}\right) d t=\sum_{k} Z_{t}^{k} d \beta_{t}^{k}  \tag{7.4}\\
\psi(T)=\psi_{T}
\end{array}\right.
$$

For $\sigma>0$, let us define the spaces (where "prog." stands for "progressively measurable")

$$
\begin{aligned}
& \mathbb{H}_{\psi}^{\sigma}=\left\{\psi: \Omega \times[0, T] \rightarrow H_{x}^{\sigma} \text { prog., }\|\psi\|_{H^{\sigma}}^{2}:=\mathbb{E} \int_{0}^{T}\left\|\psi_{t}\right\|_{H_{x}^{\sigma}}^{2} d t<\infty\right\} \\
& \mathbb{H}_{Z}^{\sigma}=\left\{Z=\left(Z^{k}\right)_{k}: \Omega \times[0, T] \rightarrow H_{x}^{\sigma} \text { prog., }\|Z\|_{H^{\sigma}}^{2}:=\sum_{k} \mathbb{E} \int_{0}^{T}\left\|Z_{t}\right\|_{H_{x}^{\sigma}}^{2} d t<\infty\right\} .
\end{aligned}
$$

The existence theory for (7.4) is linked to the uniqueness for (7.1) in the following way.
Lemma 7.1. Assume that for some $s>1 / 2$, for all $\mathcal{F}_{T}$-measurable $\psi_{T}: \Omega \rightarrow C^{\infty}\left(\mathbb{T}_{x}\right)$, (7.4) has a solution $(\psi, Z) \in \uplus_{\psi}^{s+2} \times \Vdash_{Z}^{s+1}$. Then the conclusion of Proposition 7.1 holds.
Proof of the lemma. Considering two solutions $\rho^{1}, \rho^{2}$ of (7.1), let us define $\rho=\rho^{1}-\rho^{2}$, so that $\rho(0)=0$. We start by regularizing (7.1) using an appropriate kernel $k_{\delta} \in C^{\infty}\left(\mathbb{T}_{x}\right), \delta>0$ (one may for instance consider the Féjer kernel). Denoting $\rho_{\delta}=k_{\delta} * \rho$ and $\partial_{\delta}=k_{\delta} * \partial_{x}$, we have

$$
\left\{\begin{array}{l}
d \rho_{t}^{\delta}+\left(\partial_{\delta}\left(b_{t} \rho_{t}\right)-\frac{1}{2} \sum_{k} \partial_{\delta}\left(\phi_{k} \partial_{x}\left(\phi_{k} \rho_{t}\right)\right)\right) d t=\sum_{k} \partial_{\delta}\left(\phi_{k} \rho_{t}\right) d \beta_{t}^{k} \\
\rho^{\delta}(0)=k_{\delta} * \rho_{0}:=\rho_{0}^{\delta}
\end{array}\right.
$$

in the "strong sense" that $\left(\rho_{t}^{\delta}\right)_{t \geq 0}$ is an $L_{x}^{2}$-valued semi-martingale. Itô's formula is now licit (one may verify that the assumptions of [4] Theorem 4.17 are satisfied) and gives

$$
\begin{aligned}
d\left\langle\rho_{t}^{\delta}, \psi_{t}\right\rangle & =\left[-\left\langle\partial_{\delta}\left(b_{t} \rho_{t}\right), \psi_{t}\right\rangle-\left\langle\rho_{t}^{\delta}, b \partial_{x} \psi_{t}\right\rangle\right] d t \\
& \left.+\left[\frac{1}{2} \sum_{k}\left\langle\partial_{\delta}\left(\phi_{k} \partial_{x} \phi_{k} \rho_{t}\right)\right), \psi_{t}\right\rangle-\left\langle\rho_{t}^{\delta}, \phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right\rangle\right] d t \\
& +\left[\sum_{k}\left\langle\rho_{t}^{\delta}, \phi_{k} \partial_{x} Z_{t}^{k}\right\rangle+\left\langle\partial_{\delta}\left(\phi_{k} \rho_{t}\right), Z_{t}^{k}\right\rangle\right] d t \\
& +\sum_{k}\left[\left\langle\partial_{\delta}\left(\phi_{k} \rho_{t}\right), \psi_{t}\right\rangle+\left\langle\rho_{t}^{\delta}, Z_{t}^{k}\right\rangle\right] d \beta_{t}^{k}
\end{aligned}
$$

The assumptions easily guarantee that the last term defines a square-integrable martingale. We may hence write

$$
\begin{align*}
\mathbb{E}\left[\left\langle\rho_{T}, k_{\delta} * \psi_{T}\right\rangle\right] & =-\mathbb{E} \int_{0}^{T}\left\langle\rho_{t}, b_{t} \partial_{x}\left(k_{\delta} * \psi_{t}\right)-k_{\delta} *\left[b_{t} \partial_{x} \psi_{t}\right]\right\rangle d t \\
& -\frac{1}{2} \mathbb{E} \sum_{k} \int_{0}^{T}\left\langle\rho_{t}, \phi_{k} \partial_{x}\left(\phi_{k} \partial_{x}\left(k_{\delta} * \psi_{t}\right)\right)-k_{\delta} *\left[\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right]\right\rangle d t  \tag{7.5}\\
& -\mathbb{E} \sum_{k} \int_{0}^{T}\left\langle\rho_{t}, k_{\delta} *\left[\phi_{k} \partial_{x} Z_{t}^{k}\right]-\phi_{k} \partial_{x}\left(k_{\delta} * Z_{t}^{k}\right)\right\rangle d t \tag{7.6}
\end{align*}
$$

The left-hand side easily goes to $\mathbb{E}\left[\left\langle\rho_{T}, \psi_{T}\right\rangle\right]$, while the right-hand side vanishes as $\delta \rightarrow 0$. For instance, using (7.2), one may control the second term (7.5) with

$$
\sum_{k} \mathbb{E} \int_{0}^{T}\left\|\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x}\left(k_{\delta} * \psi_{t}\right)\right)-k_{\delta} *\left[\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right]\right\|_{L_{x}^{\infty}} d t
$$

The assumptions guarantee that $\psi_{t} \in H_{x}^{s+2} \subset W_{x}^{2, \infty}$ a.s, so that the integrand tends to 0 for all fixed $k, t$ and almost every $\omega$. With the uniform bound

$$
\left\|\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x}\left(k_{\delta} * \psi_{t}\right)\right)-k_{\delta} *\left[\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right]\right\|_{L_{x}^{\infty}} \leq C\left\|\phi_{k}\right\|_{W_{x}^{1, \infty}}^{2}\left\|\psi_{t}\right\|_{W_{x}^{2, \infty}} \leq C\left\|\phi_{k}\right\|_{W_{x}^{1, \infty}}^{2}\left\|\psi_{t}\right\|_{H_{x}^{s+2}}
$$

and the regularity assumption (7.3), we may apply dominated convergence to conclude. Note that the bound for the third term (7.6) is

$$
\left\|k_{\delta} *\left[\phi_{k} \partial_{x} Z_{t}^{k}\right]-\phi_{k} \partial_{x}\left(k_{\delta} * Z_{t}^{k}\right)\right\|_{L_{x}^{\infty}} \leq\left\|\phi_{k}\right\|_{L^{\infty}}\left\|Z_{t}^{k}\right\|_{W_{x}^{1, \infty}} \leq\left\|\phi_{k}\right\|_{L^{\infty}}^{2}+\left\|Z_{t}^{k}\right\|_{H_{x}^{s+1}}^{2} .
$$

Since $\left\|\rho_{T}\right\|_{L_{x}^{1}} \leq 1$ a.s, the equality $\mathbb{E}\left[\left\langle\rho_{T}, \psi_{T}\right\rangle\right]=0$ for all $\mathcal{F}_{T}$-measurable $\psi_{T}: \Omega \rightarrow C^{\infty}(\mathbb{T})$ leads to $\rho_{T}(x)=0$ for almost every $(x, \omega) \in \mathbb{T} \times \Omega$. This holds for all $T>0$ and since $\rho \in C\left([0, T] ; H_{x}^{-\gamma}\right)$ almost surely, we conclude $\rho=0$ a.s.

It now remains to prove the existence of sufficiently regular solutions to the BSPDE (7.4). Let us consider the equation perturbed with an additional $\frac{\varepsilon}{2} \partial_{x}^{2} \psi$ :

$$
\left\{\begin{array}{l}
d \psi_{t}+\left(\left(-b_{t}+\frac{1}{2} \sum_{k} \phi_{k} \partial_{x} \phi_{k}\right) \partial_{x} \psi_{t}-\sum_{k} \phi_{k} \partial_{x} Z_{t}^{k}\right) d t+\frac{1}{2}\left(\sum_{k} \phi_{k}^{2}+\varepsilon\right) \partial_{x}^{2} \psi_{t} d t=\sum_{k} Z_{t}^{k} d \beta_{t}^{k}  \tag{7.7}\\
\psi(T)=\psi_{T} .
\end{array}\right.
$$

In this super-parabolic setup, in view of the assumptions (7.3) on the coefficients, assuming furthermore that

$$
\psi_{T} \in L^{2}\left(\left(\Omega, \mathcal{F}_{T}\right) ; H_{x}^{4}\right)
$$

Theorem 2.3 of [8] guarantees the existence of a solution $\left(\psi^{\varepsilon}, Z^{\varepsilon}\right)$ of (7.7) such that

$$
\begin{equation*}
\left(\psi^{\varepsilon}, Z^{\varepsilon}\right) \in \mathbb{H}^{5} \times \mathbb{H}^{4} \quad \text { and } \quad \mathbb{E}\left[\sup _{t \in[0, T]}\left\|\psi_{t}^{\varepsilon}\right\|_{H_{x}^{4}}^{2}\right]<\infty . \tag{7.8}
\end{equation*}
$$

Of course, the corresponding estimates in these spaces may depend on $\varepsilon>0$.
Proposition 7.2. For all $0 \leq m \leq 4$, the solution $\left(\psi^{\varepsilon}, Z^{\varepsilon}\right) \equiv(\psi, Z)$ of (7.7) satisfies the estimate

$$
\begin{align*}
& \frac{1}{2}\left\|\partial_{x}^{m} \psi_{t}\right\|_{L_{x}^{2}}^{2}+\frac{1}{2} \int_{t}^{T}\left(\varepsilon\left\|\partial_{x}^{m+1} \psi_{s}\right\|_{L_{x}^{2}}^{2}+\sum_{k}\left\|\partial_{x}^{m} Z_{s}^{k}+m\left(\partial_{x}^{m} \psi_{s}\right)\left(\partial_{x} \phi_{k}\right)-\partial_{x}\left(\partial_{x}^{m} \psi_{s} \phi_{k}\right)\right\|_{L_{x}^{2}}^{2}\right) d s \\
& \leq \frac{1}{2}\left\|\partial_{x}^{m} \psi_{T}\right\|_{L_{x}^{2}}^{2}+C \int_{t}^{T}\left(\left\|\psi_{s}\right\|_{H_{x}^{m}}^{2}+\sum_{k}\left\|Z_{s}^{k}\right\|_{H_{x}^{m-1}}^{2}\right) d s-\sum_{k} \int_{t}^{T}\left\langle\partial_{x}^{m} \psi_{s}, \partial_{x}^{m} Z_{s}^{k}\right\rangle d \beta_{s}^{k}, \tag{7.9}
\end{align*}
$$

with the convention $\|Z\|_{H_{x}^{m-1}}=0$ when $m=0$, where the constant $C$ depends only on the quantities involved in (7.3).

Remark 7.1. Note that (7.8) guarantees that, for $0 \leq m \leq 4$,

$$
M_{t}=\sum_{k} \int_{0}^{t}\left\langle\partial_{x}^{m} \psi_{s}, \partial_{x}^{m} Z_{s}^{k}\right\rangle d \beta_{s}^{k}, \quad t \geq 0
$$

defines a real-valued $L^{1}$ martingale. Indeed,

$$
V_{T}:=\sum_{k} \int_{0}^{T}\left|\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m} Z_{t}^{k}\right\rangle\right|^{2} d t \leq \sup _{t \in[0, T]}\left\|\psi_{t}\right\|_{H_{x}^{m}}^{2} \sum_{k} \int_{0}^{T}\left\|Z_{t}^{k}\right\|_{H_{x}^{m}}^{2} d t<\infty \text { a.s }
$$

ensures that $M_{t}$ is a local martingale. Burkholder-Davis-Gundy's inequality [3] then gives

$$
\mathbb{E}\left[\sup _{t \in[0, T]}\left|M_{t}\right|\right] \lesssim \mathbb{E}\left[V_{T}^{1 / 2}\right] \lesssim \mathbb{E}\left[\sup _{t \in[0, T]}\left\|\psi_{t}\right\|_{H_{x}^{m}}^{2}\right]+\sum_{k} \mathbb{E} \int_{0}^{T}\left\|Z_{t}^{k}\right\|_{H_{x}^{m}}^{2} d t<\infty
$$

Let us now give a proof of Proposition 7.2.
Proof. We apply Itô's formula to $\left\|\partial_{x}^{m} \psi_{t}\right\|_{L_{x}^{2}}^{2}$. To make it rigorous (so that the assumptions of [4] Theorem 4.17 are again satisfied), one may start by regularizing the coefficients $b$ and $\left(\phi_{k}\right)_{k}$ involved in (7.7), so that [8] provides the regularity $(\psi, Z) \in \mathbb{H}^{6} \times \Vdash^{5}$. We obtain

$$
\frac{1}{2} d\left\|\partial_{x}^{m} \psi_{t}\right\|_{L_{x}^{2}}^{2}=\left(A_{t}^{1}+A_{t}^{2}+A_{t}^{3}+A_{t}^{4}+A_{t}^{5}\right) d t+\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m} Z_{t} d W_{t}\right\rangle
$$

with

$$
\begin{aligned}
& A_{t}^{1}:=-\frac{\varepsilon}{2}\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m+2} \psi_{t}\right\rangle, \quad A_{t}^{2}:=\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m}\left(b_{t} \partial_{x} \psi_{t}\right)\right\rangle, \quad A_{t}^{3}:=\sum_{k}\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m}\left(\phi_{k} \partial_{x} Z_{t}^{k}\right)\right\rangle, \\
& A_{t}^{4}:=\frac{1}{2} \sum_{k}\left\|\partial_{x}^{m} Z_{t}^{k}\right\|_{L_{x}^{2}}^{2}, \quad A_{t}^{5}:=-\frac{1}{2} \sum_{k}\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m}\left(\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right)\right\rangle d t .
\end{aligned}
$$

Let us review the different terms. The first term $A_{t}^{1}$ appears untouched in (7.9). We may expand the second term into

$$
\begin{aligned}
A_{t}^{2} & =\left\langle\partial_{x}^{m} \psi_{t}, b_{t} \partial_{x}^{m+1} \psi_{t}\right\rangle+\sum_{l=0}^{m-1}\binom{m}{l}\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m-l} b_{t} \partial_{x}^{l+1} \psi_{s}\right\rangle \\
& \left.=-\left.\frac{1}{2}\langle | \partial_{x}^{m} \psi_{t}\right|^{2}, \partial_{x} b_{t}\right\rangle+\sum_{l=0}^{m-1}\binom{m}{l}\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x}^{m-l} b_{t} \partial_{x}^{l+1} \psi_{s}\right\rangle
\end{aligned}
$$

which can be bounded by $C\|b\|_{L_{t}^{\infty} W_{x}^{m, \infty}}\left\|\psi_{t}\right\|_{H_{x}^{m}}^{2}$. Similarly, the third term can be expanded into

$$
\begin{aligned}
A_{t}^{3}= & {\left[\sum_{k}\left\langle\partial_{x}^{m} \psi_{t}, \phi_{k} \partial_{x}^{m+1} Z_{t}^{k}\right\rangle+m\left\langle\partial_{x}^{m} \psi_{t}, \partial_{x} \phi_{k} \partial_{x}^{m} Z_{t}^{k}\right\rangle\right] } \\
& +\left[\sum_{k} \sum_{l=0}^{m-2}\binom{m}{l}\left\langle\partial_{x}^{m} \psi_{t} \partial_{x}^{m-l} \phi_{k}, \partial_{x}^{l+1} Z_{t}^{k}\right\rangle\right] \quad:=B_{t}^{1}+B_{t}^{2}
\end{aligned}
$$

The term $B_{t}^{2}$ can be bounded by

$$
C \sum_{k}\left\|\psi_{t}\right\|_{H_{x}^{m}}\left\|\phi_{k}\right\|_{W_{x}^{m, \infty}}\left\|Z_{t}^{k}\right\|_{H_{x}^{m-1}} \lesssim\left\|\psi_{t}\right\|_{H_{x}^{m}}^{2} \sum_{k}\left\|\phi_{k}\right\|_{W_{x}^{m, \infty}}^{2}+\sum_{k}\left\|Z_{t}^{k}\right\|_{H_{x}^{m-1}}^{2}
$$

Using a polarizing identity, the term $B_{t}^{1}$ may be rewritten as

$$
\begin{aligned}
B_{t}^{1}= & -\sum_{k}\left\langle\partial_{x}\left(\partial_{x}^{m} \psi_{t} \phi_{k}\right)-m \partial_{x}^{m} \psi_{t} \partial_{x} \phi_{k}, \partial_{x}^{m} Z_{t}^{k}\right\rangle \\
= & -\frac{1}{2} \sum_{k}\left\|\partial_{x}^{m} Z_{t}^{k}\right\|_{L_{x}^{2}}^{2}-\frac{1}{2} \sum_{k}\left\|\partial_{x}\left(\partial_{x}^{m} \psi_{t} \phi_{k}\right)-m \partial_{x}^{m} \psi_{t} \partial_{x} \phi_{k}\right\|_{L_{x}^{2}}^{2} \\
& +\frac{1}{2} \sum_{k}\left\|\partial_{x}^{m} Z_{t}^{k}+m \partial_{x}^{m} \psi_{t} \partial_{x} \phi_{k}-\partial_{x}\left(\partial_{x}^{m} \psi_{t} \phi_{k}\right)\right\|_{L_{x}^{2}}^{2} \\
= & C_{t}^{1}+C_{t}^{2}+C_{t}^{3}
\end{aligned}
$$

The term $C_{t}^{1}$ cancels out $A_{t}^{4}$. The term $C_{t}^{3}$ appears untouched in (7.9). We are hence left to study

$$
D_{t}:=A_{t}^{5}+C_{t}^{2}=-\frac{1}{2} \sum_{k}\left(\partial_{x}^{m} \psi_{t}, \partial_{x}^{m}\left(\phi_{k} \partial_{x}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right)\right)-\frac{1}{2} \sum_{k}\left\|\partial_{x}\left(\partial_{x}^{m} \psi_{t} \phi_{k}\right)-m \partial_{x}^{m} \psi_{t} \partial_{x} \phi_{k}\right\|_{L_{x}^{2}}^{2} .
$$

By expanding the derivatives and integrating by parts as done previously, we get the expression

$$
-\frac{1}{2} \sum_{k}\left[\left(\left|\partial_{x}^{m} \psi_{t}\right|^{2}, \partial_{x}\left(\phi_{k} \partial_{x} \phi_{k}\right)\right)-\left(\left|\partial_{x}^{m+1} \psi_{t}\right|^{2},\left|\phi_{k}\right|^{2}\right)+\sum_{l=0}^{m-2}\binom{m}{l}\left(\partial_{x}^{m} \psi_{t}, \partial_{x}^{m-l} \phi_{k} \partial_{x}^{l+1}\left(\phi_{k} \partial_{x} \psi_{t}\right)\right)\right]
$$

for the term $A_{t}^{5}$, and

$$
C_{t}^{2}=-\frac{1}{2} \sum_{k}\left[\left(\left|\partial_{x}^{m} \psi_{t}\right|^{2},(m-1)^{2}\left|\partial_{x} \phi_{k}\right|^{2}+(m-1) \partial_{x}\left(\phi_{k} \partial_{x} \phi_{k}\right)\right)+\left(\left|\partial_{x}^{m+1} \psi_{t}\right|^{2},\left|\phi_{k}\right|^{2}\right)\right] .
$$

We see that the unwanted terms involving $\partial_{x}^{m+1} \psi_{t}$ therefore cancel out in $D_{t}$, and the remaining terms can once again be bounded by $C\left\|\psi_{t}\right\|_{H_{x}^{m}}^{2} \sum_{k}\left\|\phi_{k}\right\|_{W_{x}^{m, \infty}}^{2}$.

From Proposition 7.2, we derive the following uniform bound.
Corollary 7.1. The solution $\left(\psi^{\varepsilon}, Z^{\varepsilon}\right)$ of (7.7) satisfies the bound

$$
\left\|\psi^{\varepsilon}\right\|_{\mathbb{H}^{4}}^{2}+\left\|Z^{\varepsilon}\right\|_{\mathbb{H}^{3}}^{2} \lesssim 1 .
$$

The constant involved in $\lesssim$ depends only on $T$ and the quantities of (7.3).
Proof. Let us simply denote $(\psi, Z)=\left(\psi^{\varepsilon}, Z^{\varepsilon}\right)$. As mentioned in Remark 7.1, we may take the expectation in (7.9), which gives, in particular, for $0 \leq m \leq 4$, the estimate

$$
\begin{equation*}
\mathbb{E}\left\|\psi_{t}\right\|_{H_{x}^{m}}^{2} \leq \mathbb{E}\left\|\psi_{T}\right\|_{H_{x}^{m}}^{2}+C\left(\int_{t}^{T} \mathbb{E}\left\|\psi_{s}\right\|_{H_{x}^{m}}^{2} d s+\sum_{k} \int_{t}^{T} \mathbb{E}\left\|Z_{s}^{k}\right\|_{H_{x}^{m-1}}^{2} d s\right) . \tag{7.10}
\end{equation*}
$$

Writing

$$
\left\|\partial_{x}^{m} Z_{s}^{k}\right\|_{L_{x}^{2}} \leq\left\|\partial_{x}^{m} Z_{s}^{k}+m\left(\partial_{x}^{m} \psi_{s}\right)\left(\partial_{x} \phi_{k}\right)-\partial_{x}\left(\partial_{x}^{m} \psi_{s} \phi_{k}\right)\right\|_{L_{x}^{2}}+\left\|m\left(\partial_{x}^{m} \psi_{s}\right)\left(\partial_{x} \phi_{k}\right)-\partial_{x}\left(\partial_{x}^{m} \psi_{s} \phi_{k}\right)\right\|_{L_{x}^{2}},
$$ and using (7.9) again, we also obtain

$$
\begin{equation*}
\sum_{k} \int_{t}^{T} \mathbb{E}\left\|Z_{s}^{k}\right\|_{H_{x}^{m}}^{2} d s \leq \mathbb{E}\left\|\psi_{T}\right\|_{H_{x}^{m}}^{2}+C\left(\int_{t}^{T} \mathbb{E}\left\|\psi_{s}\right\|_{H_{x}^{m+1}}^{2} d s+\sum_{k} \int_{t}^{T} \mathbb{E}\left\|Z_{s}^{k}\right\|_{H_{x}^{m-1}}^{2} d s\right) . \tag{7.11}
\end{equation*}
$$

Recalling the convention $\|Z\|_{H_{x}^{-1}}=0$, in Proposition 7.2, inequality (7.11) yields, for $0 \leq m \leq 3$,

$$
\begin{equation*}
\sum_{k} \int_{t}^{T} \mathbb{E}\left\|Z_{s}^{k}\right\|_{H_{x}^{m}}^{2} d s \lesssim \mathbb{E}\left\|\psi_{T}\right\|_{H_{x}^{3}}^{2}+\int_{t}^{T} \mathbb{E}\left\|\psi_{s}\right\|_{H_{x}^{4}}^{2} . \tag{7.12}
\end{equation*}
$$

Coming back to (7.10), we get, for $0 \leq m \leq 4$,

$$
\mathbb{E}\left\|\psi_{t}\right\|_{H_{x}^{4}}^{2} \lesssim \mathbb{E}\left\|\psi_{T}\right\|_{H_{x}^{4}}^{2}+\int_{t}^{T} \mathbb{E}\left\|\psi_{s}\right\|_{H_{x}^{4}}^{2} d s
$$

and Grönwall's Lemma provides the desired estimate for $\psi$. The estimate for $Z$ is obtained by taking $t=0$ in (7.12).

Now, let us consider $\varepsilon_{1}, \varepsilon_{2}>0$ and denote $(\psi, Z)=\left(\psi^{\varepsilon_{1}}-\psi^{\varepsilon_{2}}, Z^{\varepsilon_{1}}-Z^{\varepsilon_{2}}\right)$. In the same fashion as Proposition 7.2, one can prove the following estimate, for $0 \leq m \leq 3$ :

$$
\begin{aligned}
& \frac{1}{2}\left\|\partial_{x}^{m} \psi_{t}\right\|_{L_{x}^{2}}^{2}+\frac{1}{2} \int_{t}^{T} \sum_{k}\left\|\partial_{x}^{m} Z_{s}^{k}+m\left(\partial_{x}^{m} \psi_{s}\right)\left(\partial_{x} \phi_{k}\right)-\partial_{x}\left(\partial_{x}^{m} \psi_{s} \phi_{k}\right)\right\|_{L_{x}^{2}}^{2} d s \\
& \leq C\left|\varepsilon_{2}-\varepsilon_{1}\right| \int_{0}^{T}\left(\left\|\psi_{s}^{\varepsilon_{1}}\right\|_{H_{x}^{m+1}}^{2}+\left\|\psi_{s}^{\varepsilon_{2}}\right\|_{H_{x}^{m+1}}^{2}\right) d s+C \int_{t}^{T}\left(\left\|\psi_{s}\right\|_{H_{x}^{m}}^{2}+\sum_{k}\left\|Z_{s}^{k}\right\|_{H_{x}^{m-1}}^{2}\right) d s \\
& \quad-\sum_{k} \int_{t}^{T}\left(\partial_{x}^{m} \psi_{s}, \partial_{x}^{m} Z_{s}^{k}\right) d \beta_{s}^{k}
\end{aligned}
$$

Relying on the bound from Corollary 7.1 and applying similar calculations as in its proof, we deduce

$$
\left\|\psi^{\varepsilon_{1}}-\psi^{\varepsilon_{2}}\right\|_{\mathbb{H}^{3}}^{2}+\left\|Z^{\varepsilon_{1}}-Z^{\varepsilon_{2}}\right\|_{\mathbb{H}^{2}}^{2} \lesssim\left|\varepsilon_{1}-\varepsilon_{2}\right| .
$$

As a consequence, $\left(\psi^{\varepsilon}, Z^{\varepsilon}\right)$ converges to some $(\psi, Z)$ in the Banach space $\mathbb{H}^{3} \times \uplus^{2}$ as $\varepsilon$ goes to zero. It is now easy to send $\varepsilon$ to zero in (7.7) and verify that $(\psi, Z) \in \Vdash^{3} \times \Vdash^{3}$ is indeed a solution of (7.4). Lemma 7.1 may be applied, which concludes the proof of Proposition 7.1, and therefore that of Theorem 2.

## 8 Appendix: Well-posedness of the path-wise kinetic system

For the sake of completeness, we wish to prove the well-posedness of the deterministic PDE system

$$
\left\{\begin{array}{l}
\partial_{t} f+v \partial_{x} f+\partial_{v}[(m+u-v) f]=0  \tag{8.1}\\
\partial_{t} u-\partial_{x}^{2} u=J(f)-\rho u
\end{array}\right.
$$

Here, we consider a fixed $\omega \in \Omega$, so that $m \equiv m(\omega)$ belongs to $D([0, T] ; E)$, space of càdlàg functions taking values in the separable complete space $E$.

## 8.1 $L^{1}$ and $L^{2}$ estimates for a conservation equation with Lipschitz coefficients

Recall that, whenever $w \in D([0, T] ; E)$ is given, the solution $f[w]$ to the linear conservation equation

$$
\left\{\begin{array}{l}
\partial_{t} f+v \partial_{x} f+\partial_{v}[(w-v) f]=0 \\
f_{0} \in L_{x, v}^{1}
\end{array}\right.
$$

is naturally expressed as

$$
f_{t}[w](x, v)=e^{t} f_{0} \circ \Phi_{0}^{t}(x, v)
$$

where $\Phi_{s}^{t}(x, v)=\left(X_{s}^{t}(x, v), V_{s}^{t}(x, v)\right)$ is the flow associated to the characteristics:

$$
\left\{\begin{align*}
\frac{d}{d s} X_{s}^{t} & =V_{s}^{t}, & X_{t}^{t}(x, v) & =x  \tag{8.2}\\
\frac{d}{d s} V_{s}^{t} & =w_{s}\left(X_{s}^{t}\right)-V_{s}^{t}, & V_{t}^{t}(x, v) & =v
\end{align*}\right.
$$

Grönwall's lemma immediately gives the following result.
Lemma 8.1. Let $w, \widetilde{w} \in D([0, T] ; E)$ and let $\Phi[w], \Phi[\widetilde{w}]$ denote the respective flows. Then,

$$
\forall s, t \in[0, T], \quad\left\|\Phi_{s}^{t}[w]-\Phi_{s}^{t}[\widetilde{w}]\right\|_{L_{x, v}^{\infty}} \lesssim\|w-\widetilde{w}\|_{L_{t}^{\infty} L_{x}^{\infty}}
$$

The constant involved in $\lesssim d e p e n d s$ on $T,\left\|\partial_{x} w\right\|_{L_{t, x}^{\infty}},\left\|\partial_{x} \widetilde{w}\right\|_{L_{t, x}^{\infty}}$ only.

Let us now fix some $w \in D([0, T] ; E)$ and give some estimates regarding the $v$-moments of $f_{t}[w]$. Recall the notation, for $f \in G^{1}$,

$$
\rho(f)(x)=\int f(x, v) d v, \quad J(f)(x)=\int v f(x, v) d v
$$

Lemma 8.2 ( $L^{1}$ moments estimates). Let $f_{0} \in L_{x, v}^{1}$. Then for all $t \in[0, T]$,

$$
\left\|f_{t}[w]\right\|_{L_{x, v}^{1}}=\left\|f_{0}\right\|_{L_{x, v}^{1}}
$$

Assume furthermore that $f_{0} \in G^{k}, k \geq 1$. Then for all $t \in[0, T]$,

$$
\left\|f_{t}[w]\right\|_{G^{k}} \lesssim\left\|f_{0}\right\|_{G^{k}}
$$

The constant involved in $\lesssim d e p e n d s$ on $T, k$ and $\|w\|_{L_{t, x}^{\infty}}$ only.
The estimate for the first moment is precisely:

$$
\left\|J\left(f_{t}[w]\right)\right\|_{L_{x, v}^{1}} \leq\left(\left\|J\left(f_{0}\right)\right\|_{L_{x}^{1}}+t \sup _{s \in[0, t]}\left\|w_{s}\right\|_{L_{x}^{\infty}}\left\|f_{0}\right\|_{L_{x, v}^{1}}\right) e^{-t}
$$

Proof. Without loss of generality, we may assume $f_{0} \geq 0$. When $w \in D([0,1] ; E)$ is smooth, it is clear that the flow $\Phi[\widetilde{w}]$ is a globally-defined $C^{\infty}$-diffeomorphism, whose Jacobian determinant is given by $J \Phi_{0}^{t}[\widetilde{w}]=\operatorname{det}\left(D \Phi_{0}^{t}[\widetilde{w}]\right)=e^{t}$. The conservation of the $L^{1}$-norm is then obtained by a change of variable using the form (3.4). The result is generalized to $w \in D([0,1] ; E)$ by density using Lemma 8.1.

For simplicity, we only prove the estimate in $G^{1}$ : the estimates on higher moments can be obtained by induction using the same process. Assuming $f_{0} \in C_{c}^{\infty}(\mathbb{T} \times \mathbb{R})$, the following calculations are justified:

$$
\begin{aligned}
\frac{d}{d t}\left\|J\left(f_{t}\right)\right\|_{L_{x}^{1}} & =\frac{d}{d t} \int_{x} \int_{v}|v| f_{t}(x, v) d x d v=-\int_{x} \int_{v}|v|\left(v \partial_{x} f+\partial_{v}\left[\left(w_{t}(x)-v\right) f_{t}(x, v)\right]\right) d x d v \\
& =-\int_{v} \operatorname{sign}(v) w_{t}(x) f_{t}(x, v) d v+\int_{v}|v| f_{t}(x, v) d v
\end{aligned}
$$

so that

$$
\frac{d}{d t}\left\|J\left(f_{t}\right)\right\|_{L_{x}^{1}} \leq\left\|w_{t}\right\|_{L_{x}^{\infty}}\left\|f_{0}\right\|_{L_{x}^{1}}-\left\|J\left(f_{t}\right)\right\|_{L_{x}^{1}} \leq \sup _{s \in[0, t]}\left\|w_{s}\right\|_{L_{x}^{\infty}}\left\|f_{0}\right\|_{L_{x}^{1}}-\left\|J\left(f_{t}\right)\right\|_{L_{x}^{1}}
$$

that is, integrating over $[0, t]$,

$$
\left\|J\left(f_{t}\right)\right\|_{L_{x}^{1}} \leq\left(\left\|J\left(f_{0}\right)\right\|_{L_{x}^{1}}+t \sup _{s \in[0, t]}\left\|w_{s}\right\|_{L_{x}^{\infty}}\left\|f_{0}\right\|_{L_{x}^{1}}\right)-\int_{0}^{t}\left\|J\left(f_{s}\right)\right\|_{L_{x}^{1}} d s
$$

Grönwall's lemma gives the expected result. Again, the inequality holds for a general initial data $f_{0} \geq 0$ by a density argument, using Fatou's Lemma.

Corollary 8.1. Let $f_{0} \in L_{x, v}^{2}$. Then for all $t \in[0, T]$,

$$
\left\|f_{t}[w]\right\|_{L_{x, v}^{2}}=e^{t / 2}\left\|f_{0}\right\|_{L_{x, v}^{2}} .
$$

Assume furthermore that $\left|f_{0}\right|^{2} \in G^{k}, k \geq 1$. Then for all $t \in[0, T]$,

$$
\left\|f_{t}[w]^{2}\right\|_{G^{k}} \lesssim\left\|f_{0}^{2}\right\|_{G^{k}}
$$

The constant involved in $\lesssim d e p e n d s$ on $T, k$ and $\|w\|_{L_{t, x}^{\infty}}$ only.

Proof. Simply note that $h_{t}:=e^{-t}\left(f_{t}[w]\right)^{2}=e^{t} f_{0}^{2} \circ \Phi_{0}^{t}$ and apply Lemma 8.2 to $h$.
Corollary 8.2 ( $L^{2}$ moments estimates). Let $f_{0} \in L_{x, v}^{2}$ and $f_{0}^{2} \in G_{4}$. Then, for all $t \in[0, T]$,

$$
\begin{aligned}
& \left\|\rho\left(f_{t}[w]\right)\right\|_{L_{x}^{2}} \lesssim\left\|f_{0}^{2}\right\|_{G_{2}}^{1 / 2}, \\
& \left\|J\left(f_{t}[w]\right)\right\|_{L_{x}^{2}} \lesssim\left\|f_{0}^{2}\right\|_{G_{4}}^{1 / 2} .
\end{aligned}
$$

The constants involved in $\lesssim$ depends on $T$ and $\|w\|_{L_{t, x}^{\infty}}$ only.
Proof. For instance, for the second estimate, Hölder's inequality gives

$$
\left(\int_{v}|v| f d v\right)^{2}=\left(\int_{v}(1+|v|)|v| f \times \frac{1}{1+|v|} d v\right)^{2} \leq\left(\int_{v}(1+|v|)^{2}|v|^{2} f^{2} d v\right)\left(\int_{v} \frac{1}{(1+|v|)^{2}} d v\right)
$$

from which we deduce $\int_{x} J(f)^{2} d x \leq C\left\|f^{2}\right\|_{G_{4}}$, and we use the estimate from Corollary 8.1.
We will now establish some kind of stability estimate for $f$ as a function of $w$. Let us start with a lemma regarding the gradient.

Lemma 8.3 ( $L^{1}$ gradient estimates). Assume that $\left\|\nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1}}<\infty$. Then,

$$
\left\|\nabla_{x, v} f_{t}[w]\right\|_{L_{x, v}^{1},} \leq C\left(T,\left\|\partial_{x} w\right\|_{L_{t, x}^{\infty}}\right)\left\|\nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1}, v}
$$

Assume furthermore that $\left\|v \nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1}}<\infty$. Then,

$$
\left\|v \nabla_{x, v} f_{t}[w]\right\|_{L_{x, v}^{1}} \leq C\left(T,\|w\|_{L_{t, x}^{\infty}},\left\|\partial_{x} w\right\|_{L_{t, x}^{\infty}},\left\|\nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1},}\left\|v \nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1}}\right)
$$

Proof. Let us assume that $f_{0} \in C_{c}^{\infty}$, so that we are dealing with smooth, compact-supported, strong solutions. The general case is once again deduced by density arguments. From

$$
\partial_{t} f+v \partial_{x} f+(w-v) \partial_{v} f-f=0
$$

we derive that $g_{1}=\partial_{x} f, g_{2}=\partial_{v} f$ satisfy

$$
\begin{align*}
& \partial_{t} g_{1}+v \partial_{x} g_{1}+\partial_{x} w g_{2}+(w-v) \partial_{v} g_{1}-g_{1}=0,  \tag{8.3}\\
& \partial_{t} g_{2}+g 1+v \partial_{x} g_{2}+(w-v) \partial_{v} g_{2}-2 g_{2}=0 \tag{8.4}
\end{align*}
$$

Multiplying (8.3) by $\operatorname{sign}\left(g_{1}\right)$ and integrating over time (to be rigorous, one should replace $\operatorname{sign}(x)$ by $s_{\eta}(x)=x /(|x|+\eta)$ and carefully let $\eta$ go to zero), we are easily led to

$$
\frac{d}{d t} \int_{x} \int_{v}\left|g_{1}\right| d x d v \leq\left\|\partial_{x} w\right\|_{L_{t, x}^{\infty}} \int_{x} \int_{v}\left|g_{2}\right| d x d v
$$

Similarly, multiplying (8.4) by $\operatorname{sign}\left(g_{2}\right)$ we get

$$
\frac{d}{d t} \int_{x} \int_{v}\left|g_{2}\right| d x d v \leq \int_{x} \int_{v}\left|g_{1}\right| d x d v+\int_{x} \int_{v}\left|g_{s}\right| d x d v
$$

Summing these and applying Grönwall's inequality gives the expected estimate. Let us now establish the estimate on $\left\|v \nabla_{x, v} f\right\|_{L_{x, v}^{1}}$. From (8.3), (8.4), we derive that $h_{1}:=v \partial_{x} f$ and $h_{2}:=v \partial_{v} f$ satisfy

$$
\begin{aligned}
& \partial_{t} h_{1}+v \partial_{x} h_{1}+\left(\partial_{x} w\right) h_{2}+v(w-v) \partial_{v} g_{1}-h_{1}=0, \\
& \partial_{t} h_{2}+h_{1}+v \partial_{x} h_{2}+v(w-v) \partial_{v} g_{2}-2 h_{2}=0
\end{aligned}
$$

Now since $(w-v) \partial_{v} h_{i}=(w-v) g_{i}+v(w-v) \partial_{v} g_{i}$, we can rewrite

$$
v(w-v) \partial_{v} g_{i}=(w-v) \partial_{v} h_{i}+h_{i}-w g_{i}
$$

and get

$$
\begin{aligned}
& \partial_{t} h_{1}+v \partial_{x} h_{1}+\left(\partial_{x} w\right) h_{2}+(w-v) \partial_{v} h_{1}-w g_{1}=0 \\
& \partial_{t} h_{2}+h_{1}+v \partial_{x} h_{2}+(w-v) \partial_{v} h_{2}-h_{2}-w g_{2}=0
\end{aligned}
$$

We then use the same approach to obtain

$$
\begin{aligned}
& \frac{d}{d t}\left\|h_{1}(t)\right\|_{L_{x, v}^{1}} \leq\left\|\partial_{x} w\right\|_{L_{t, x}^{\infty}}\left\|h_{2}(t)\right\|_{L_{x, v}^{1}}-\left\|h_{1}(t)\right\|_{L_{x, v}^{1}}+\|w\|_{L_{t, x}^{\infty}}\left\|g_{1}(t)\right\|_{L_{x, v}^{1}} \\
& \frac{d}{d t}\left\|h_{2}(t)\right\|_{L_{x, v}^{1}} \leq\left\|h_{1}(t)\right\|_{L_{x, v}^{1}}+\|w\|_{L_{t, x}^{\infty}}\left\|g_{2}(t)\right\|_{L_{x, v}^{1}}
\end{aligned}
$$

and conclude by applying the previous estimates and Grönwall's inequality.
Corollary 8.3 ( $L^{1}$ stability). Let $w, \widetilde{w} \in D([0, T] ; E)$. Assume

$$
\left\|(1+|v|) \nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1}}<+\infty
$$

Then the following estimates hold:

$$
\begin{aligned}
& \left\|\rho_{t}[w]-\rho_{t}[\widetilde{w}]\right\|_{L_{x}^{1}} \leq\left\|f_{t}[w]-f_{t}[\widetilde{w}]\right\|_{L_{x, v}^{1}} \lesssim\|w-\widetilde{w}\|_{L_{t, x}^{\infty}}, \\
& \left\|J\left(f_{t}[w]\right)-J\left(f_{t}[\widetilde{w}]\right)\right\|_{L_{x}^{1}} \lesssim\|w-\widetilde{w}\|_{L_{t, x}^{\infty}} .
\end{aligned}
$$

The constants involved in $\lesssim$ depend on $T,\left\|(1+|v|) \nabla_{x, v} f_{0}\right\|_{L_{x, v}^{1}},\|w\|_{L_{t, x}^{\infty}},\|\widetilde{w}\|_{L_{t, x}^{\infty}},\left\|\partial_{x} w\right\|_{L_{t, x}^{\infty}}$, and $\left\|\partial_{x} \widetilde{w}\right\|_{L_{t, x}^{\infty}}$ only.
Proof. Let us once again assume that $f_{0} \in C_{c}^{\infty}$. Let us denote $f_{t}=f_{t}[w], \widetilde{f_{t}}=f_{t}[\widetilde{w}], F=f-\widetilde{f}$ and $W=w-\widetilde{w}$. These quantities satisfy

$$
\partial_{t} F+v \partial_{x} F+\partial_{v}[(w-v) F+W \tilde{f}]=0
$$

which can be rewritten as

$$
\partial_{t} F+v \partial_{x} F+(w-v) \partial_{v} F-F+W \partial_{v} \tilde{f}=0
$$

Multiplying by $\operatorname{sign}(F)$, we get

$$
\partial_{t}|F|+v \partial_{x}|F|+(w-v) \partial_{v}|F|-|F|+\operatorname{sign}(F) W \partial_{v} \widetilde{f}=0
$$

and we deduce

$$
\frac{d}{d t}\|F(t)\|_{L_{x, v}^{1}} \leq\|W\|_{L_{t, x}^{\infty}}\left\|\nabla_{x, v} \tilde{f}\right\|_{L_{x, v}^{1}}
$$

Using the $L^{1}$ gradient estimates on $\tilde{f}$ from Lemma 8.3, we get the first inequality. Now let $G=v F$. Multiplying (8.1) by $v$, we get

$$
\partial_{t} G+v \partial_{x} G+v(w-v) \partial_{v} F-G+v W \partial_{v} \tilde{f}=0
$$

Since $(w-v) \partial_{v} G=(w-v) F+v(w-v) \partial_{v} F$ this can be rewritten as

$$
\partial_{t} G+v \partial_{x} G+(w-v) \partial_{v} G-w F+v W \partial_{v} \tilde{f}=0
$$

Multiplying by $\operatorname{sign}(G)$, we get

$$
\partial_{t}|G|+v \partial_{x}|G|+(w-v) \partial_{v}|G|-\operatorname{sign}(G) w F+\operatorname{sign}(G) W v \partial_{v} \widetilde{f}=0
$$

and we deduce

$$
\frac{d}{d t}\|G(t)\|_{L_{x, v}^{1}} \leq-\|G(t)\|_{L_{x, v}^{1}}+\|w\|_{L_{t, x}^{\infty}}\|F(t)\|_{L_{x, v}^{1}}+\|W\|_{L_{t, x}^{\infty}}\left\|v \nabla_{x, v} \widetilde{f}(t)\right\|_{L_{x, v}^{1}}
$$

Using the $L^{1}$ gradient estimates on $\tilde{f}$ and the inequality we just showed, Grönwall's lemma gives the expected estimate.

### 8.2 Existence and uniqueness

We are now ready to prove the following result.
Proposition 8.1. Assume that

$$
\int_{x} \int_{v}\left(1+|v|^{4}\right)\left|f_{0}\right|^{2} d x d v+\int_{x} \int_{v}(1+|v|)\left|\nabla_{x, v} f_{0}\right| d x d v<\infty, \quad u_{0} \in H_{x}^{\eta} \text { for } \eta \in(3 / 2,2) .
$$

Then, for all $\beta \in(3 / 2, \eta)$, there exists a unique couple $(f, u)$ with $u \in C\left([0, T] ; H_{x}^{\beta}\right)$ solution of the system (3.3) in the sense that

$$
\left\{\begin{array}{l}
f_{t}=f_{t}[m+u] \quad \text { as defined in (8.1) } \\
u_{t}=S(t) u_{0}+\int_{0}^{t} S(t-s)\left[J\left(f_{s}\right)-\rho_{s} u_{s}\right] d s
\end{array}\right.
$$

where $S(t)=e^{t \partial_{x}^{2}}$ denotes the semigroup associated to the heat equation. For any $k \geq 0$, assuming $f_{0} \in G^{k}$, we have additionally $f \in C\left([0, T] ; G^{k}\right)$.

Remark 8.1. The "numerical constants" used in the following proof may all depend on the quantity $\sup _{t \in[0, T]}\left\|m_{t}(\omega)\right\|_{W_{x}^{1, \infty}}$, which is in any case assumed to be bounded in Assumption 1.

Let us define the recursive sequences

$$
\begin{aligned}
f^{n} \in C\left([0, T] ; L_{x, v}^{1}\right), & n \geq 0 \\
u^{n} \in C\left([0, T] ; H_{x}^{\eta}\right), & n \geq 0
\end{aligned}
$$

as follows: $\left(f_{t}^{0}, u_{t}^{0}\right)=\left(f_{0}, u_{0}\right)$ and

$$
\begin{align*}
f_{t}^{n+1} & =f_{t}\left[m+u^{n}\right] \quad \text { as defined in (8.1) },  \tag{8.5}\\
u_{t}^{n+1} & =S(t) u_{0}+\int_{0}^{t} S(t-s)\left[J\left(f_{s}^{n+1}\right)-\rho_{s}^{n+1} u_{s}^{n}\right] d s \tag{8.6}
\end{align*}
$$

## Claim 1: these sequences are well-defined

Since $\eta>3 / 2$, and we work in dimension 1 , the embedding $H_{x}^{\eta} \subset W_{x}^{1, \infty}$ guarantees that $f_{t}\left[m+u^{n}\right]$ can be properly defined as long as $u^{n}$ takes values in $H_{x}^{\eta}$. We have

$$
\left\|u_{t}^{n+1}\right\|_{H_{x}^{n}} \leq\left\|u_{0}\right\|_{H_{x}^{n}}+\int_{0}^{t}\left\|S(t-s)\left[J\left(f_{s}^{n+1}\right)-\rho_{s}^{n+1} u_{s}^{n}\right]\right\|_{H_{x}^{\eta}} d s
$$

with the classical estimate

$$
\forall w \in L^{2}, \quad\|S(t) w\|_{H_{x}^{\eta}} \lesssim\left(1+t^{-\eta / 2}\right)\|w\|_{L_{x}^{2}} .
$$

The $L^{2}$ estimates for $f_{t}^{n+1}=f_{t}\left[m+u^{n}\right]$ from Lemma 8.2 give

$$
\left.\left\|J\left(f_{s}^{n+1}\right)-\rho_{s}^{n+1} u_{s}^{n}\right\|_{L_{x}^{2}} \leq C\left(\sup _{t \in[0, T]}\left\|u_{t}^{n}\right\|_{L_{x}^{\infty}}\right)\left(1+\sup _{t \in[0, T]}\left\|u_{t}^{n}\right\|_{L_{x}^{\infty}}\right)\right)
$$

from which we deduce

$$
\begin{equation*}
\left.\left\|u_{t}^{n+1}\right\|_{H_{x}^{\eta}} \leq\left\|u_{0}\right\|_{H_{x}^{\eta}}+C\left(\sup _{t \in[0, T]}\left\|u_{t}^{n}\right\|_{L_{x}^{\infty}}\right)\right) \int_{0}^{t}\left(1+s^{-\eta / 2}\right) d s \tag{8.7}
\end{equation*}
$$

Since $\eta<2$, it follows that, for all $n \geq 0, u^{n} \in C\left([0, T] ; H_{x}^{\eta}\right)$ and $f^{n}$ is well defined.
Remark 8.2. This calculation explains why we have to restrict ourselves to the one-dimensional case $x \in \mathbb{T}$.

Claim 2: $\left(u^{n}\right)_{n}$ is bounded in $C\left([0, T] ; L_{x}^{\infty}\right)$
We have

$$
\left\|u_{t}^{n+1}\right\|_{L_{x}^{\infty}} \leq\left\|u_{0}\right\|_{L_{x}^{\infty}}+\int_{0}^{t}\left\|S(t-s)\left[J\left(f_{s}^{n+1}\right)-\rho_{s}^{n+1} u_{s}^{n}\right]\right\|_{L_{x}^{\infty}} d s
$$

and since $H_{x}^{1 / 2+\delta} \subset L_{x}^{\infty}$ for $\delta>0$, we get $L_{x}^{1} \subset H_{x}^{-1 / 2-\delta}$ and the estimate

$$
\forall w \in L_{x}^{1}, \quad\|S(t) w\|_{L_{x}^{\infty}} \leq\|S(t) w\|_{H_{x}^{1 / 2+\delta}} \lesssim\left(1+t^{-1 / 2-\delta}\right)\|w\|_{H_{x}^{-1 / 2-\delta}} \lesssim\left(1+t^{-1 / 2-\delta}\right)\|w\|_{L_{x}^{1}} .
$$

The $L^{1}$ estimates for $f_{t}^{n+1}=f_{t}\left[m+u^{n}\right]$ from Lemma 8.2 give

$$
\left\|J\left(f_{s}^{n+1}\right)-\rho_{s}^{n+1} u_{s}^{n}\right\|_{L_{x}^{1}} \lesssim 1+\sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{n}\right\|_{L_{x}^{\infty}}+\left\|u_{s}^{n}\right\|_{L_{x}^{\infty}} \lesssim 1+\sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{n}\right\|_{L_{x}^{\infty}} .
$$

It follows that (forgetting the $\delta>0$ for simplicity)

$$
\left\|u_{t}^{n+1}\right\|_{L_{x}^{\infty}} \leq\left\|u_{0}\right\|_{L_{x}^{\infty}}+C \int_{0}^{t}\left(1+(t-s)^{-1 / 2}\right)\left(1+\sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{n}\right\|_{L_{x}^{\infty}}\right) d s
$$

Introducing, for $\mu>0$, the norm

$$
\begin{equation*}
\forall w \in C\left([0, T], L_{x}^{\infty}\right), \quad\|w\|_{\mu}=\sup _{t \in[0, T]} e^{-\mu t}\left\|w_{t}\right\|_{L_{x}^{\infty}}, \tag{8.8}
\end{equation*}
$$

is is easy to check that $\left\|u^{n+1}\right\|_{\mu} \leq a_{\mu}\left\|u^{n}\right\|_{\mu}+b$ with $a_{\mu} \in[0,1)$ for $\mu$ chosen large enough. It follows that $\sup _{n}\left\|u^{n}\right\|_{\mu}<\infty$, hence the result claimed.

Claim 3: $\left(u^{n}\right)_{n}$ is bounded in $C\left([0, T] ; H_{x}^{\eta}\right)$
Since $\sup _{t \in[0, T]}\left\|u^{n}\right\|_{L_{x}^{\infty}} \lesssim 1$, this follows directly from (8.7).
Claim 4: $\left(u^{n}\right)_{n}$ is compact in $C\left([0, T] ; H_{x}^{\beta}\right)$
Noting that, with $g^{n}=J\left(f^{n}\right)-\rho^{n} u^{n-1}$,

$$
\forall \xi \in C_{c}^{\infty}(\mathbb{T}),\left\langle u_{t}^{n}, \xi\right\rangle-\left\langle u_{0}^{n}, \xi\right\rangle=\int_{0}^{t}\left\langle u_{s}^{n}, \partial_{x}^{2} \xi\right\rangle+\left\langle g_{s}^{n}, \xi\right\rangle d s
$$

we are led to

$$
\begin{aligned}
\left|\left\langle u_{t_{2}}^{n}-u_{t_{1}}^{n}, \xi\right\rangle\right| & \leq \int_{t_{1}}^{t_{2}}\left|\left\langle u_{s}^{n}, \partial_{x}^{2} \xi\right\rangle\right|+\left|\left\langle g_{s}^{n}, \xi\right\rangle\right| d s \\
& \left.\leq\left|t_{2}-t_{1}\right|\left(\left\|u^{n}\right\|_{C\left([0, T] ; H_{x}^{\eta}\right)}\|\xi\|_{H_{x}^{2-\eta}}+\left\|g^{n}\right\|_{C\left([0, T] ; L_{x, v}^{2}, v\right.}\right)\|\xi\|_{L_{x}^{2}}\right) \\
& \lesssim\left|t_{2}-t_{1}\right|\|\xi\|_{H_{x}^{2-\eta}},
\end{aligned}
$$

that is in fact $u^{n} \in C^{0,1}\left([0, T] ; H_{x}^{\eta-2}\right)$ with a Lipschitz constant uniform in $n$. Ascoli's theorem therefore guarantees the compactness of $\left(u^{n}\right)_{n} \in C\left([0, T] ; H_{x}^{\eta-2}\right)$. Recalling the bound from Claim 3, we get the compactness of $\left(u^{n}\right)_{n} \in C\left([0, T] ; H_{x}^{\beta}\right)$ by interpolation.

## Conclusion

We simply note that the mapping $u_{n} \mapsto u_{n+1}$ is a contraction on $C\left([0, T] ; L_{x}^{\infty}\right)$ :

$$
\begin{align*}
& \left\|u_{t}^{n+1}-u_{t}^{n}\right\|_{L_{x}^{\infty}}=\left\|\int_{0}^{t} S(t-s)\left[J\left(f_{s}^{n+1}\right)-J\left(f_{s}^{n}\right)-\left(\rho_{s}^{n+1} u_{s}^{n}-\rho_{s}^{n} u_{s}^{n-1}\right)\right] d s\right\|_{L_{x}^{\infty}}  \tag{8.9}\\
& \lesssim \int_{0}^{t}\left(1+(t-s)^{-1 / 2}\right)\left(\left\|J\left(f_{s}^{n+1}\right)-J\left(f_{s}^{n}\right)\right\|_{L_{x}^{1}}+\left\|\rho_{s}^{n+1} u_{s}^{n+1}-\rho_{s}^{n} u_{s}^{n-1}\right\|_{L_{x}^{1}}\right) d s .
\end{align*}
$$

Since $f^{n+1}=f\left[m+u^{n}\right]$ and $f^{n}=f\left[m+u^{n-1}\right]$, using Corollary (8.3) and the fact that $\left(u^{n}\right)_{n}$ is bounded in $C\left([0, T] ; W_{x}^{1, \infty}\right)$, we get for all $s \in[0, T]$

$$
\left\|J\left(f_{s}^{n+1}\right)-J\left(f_{s}^{n}\right)\right\|_{L_{x}^{1}} \lesssim \sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{n}-u_{\sigma}^{n-1}\right\|_{L_{x}^{\infty}},
$$

and

$$
\begin{aligned}
\left\|\rho_{s}^{n+1} u_{s}^{n}-\rho_{s}^{n} u_{s}^{n-1}\right\|_{L_{x}^{1}} & \leq\left\|\rho_{s}^{n+1}\left(u_{s}^{n}-u_{s}^{n-1}\right)\right\|_{L_{x}^{1}}+\left\|\left(\rho_{s}^{n+1}-\rho_{s}^{n}\right) u_{s}^{n-1}\right\|_{L_{x}^{1}} \\
& \lesssim \sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{n}-u_{\sigma}^{n-1}\right\|_{L_{x}^{\infty}}^{\infty} .
\end{aligned}
$$

Therefore,

$$
\left\|u_{t}^{n+1}-u_{t}^{n}\right\|_{L_{x}^{\infty}} \lesssim \int_{0}^{t}\left(1+(t-s)^{-1 / 2}\right) \sup _{\sigma \in[0, s]}\left\|u_{\sigma}^{n}-u_{\sigma}^{n-1}\right\|_{L_{x}^{\infty}} d s
$$

Using once again the norm $\|\cdot\|_{\mu}$ defined in (8.8), we are easily led to

$$
\left\|u^{n+1}-u^{n}\right\|_{\mu} \leq k_{\mu}\left\|u^{n}-u^{n-1}\right\|_{\mu}
$$

with $k_{\mu} \in[0,1)$ for $\mu$ large enough. As a result, $u^{n}$ converges to some $u$ in the Banach space $C\left([0, T] ; L_{x}^{\infty}\right)$. Since $\left(u^{n}\right)_{n}$ is also compact in $C\left([0, T] ; H_{x}^{\beta}\right)$ we deduce that $u^{n} \rightarrow u$ in $C\left([0, T] ; H_{x}^{\beta}\right)$. We may now introduce $f=f[m+u]$. Using the same estimates as before (only replacing $u^{n+1}$ by $u$ in (8.9)), we easily derive, as $n$ goes to infinity,

$$
u_{t}=S(t) u_{0}+\int_{0}^{t} S(t-s)\left[J\left(f_{s}\right)-\rho_{s} u_{s}\right] d s
$$

Hence, $(f, u)$ is indeed a solution to (8.1). If $(\widetilde{f}, \widetilde{u})$ is another solution with $\widetilde{u} \in C\left([0, T] ; H_{x}^{\beta}\right)$ and $\widetilde{f}=f[m+\widetilde{u}]$, the same estimates lead to

$$
\|u-\widetilde{u}\|_{\mu} \leq k_{\mu}\|u-\widetilde{u}\|_{\mu},
$$

so that $\widetilde{u}=u$ and $\widetilde{f}=f$. This concludes the proof of the well-posedness of (8.1).
Finally, let us assume that $f_{0} \in G^{k}$ and prove that $f \in C\left([0, T] ; G^{k}\right)$. It is clearly enough to show that $\left\|f_{t}-f_{0}\right\|_{G^{k}} \rightarrow 0$ as $t$ goes to 0 . Given the form (8.1) and Lemma 8.1, this result is easily obtained by dominated convergence when $f_{0} \in C_{c}\left(\mathbb{R}^{2 d}\right)$. It may then be generalized by density to any $f_{0} \in G^{k}$, since, given $\widetilde{f}_{0} \in C_{c}\left(\mathbb{R}^{2 d}\right)$

$$
\forall t \in[0,1], \quad\left\|f_{t}-\widetilde{f}_{t}\right\|_{G^{k}}=\int\left(1+\left|V_{0}^{t}(z)\right|^{k}\right)\left|f_{0}(z)-\widetilde{f}_{0}(z)\right| d z \lesssim\left\|f_{0}-\widetilde{f}_{0}\right\|_{G^{k}}
$$

We have used the estimate $\left|V_{0}^{t}(x, v)\right|^{k} \lesssim 1+|v|^{k}$ easily derived from the sub-linearity of (8.2), with $\|m+u\|_{L_{t, x}^{\infty}}<\infty$.

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