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Deterministic mean field games with control on the acceleration

Yves Achdou^{*}, Paola Mannucci[†],
Claudio Marchi[‡], Nicoletta Tchou[§]

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Abstract

In the present work, we study deterministic mean field games (MFGs) with finite time horizon in which the dynamics of a generic agent is controlled by the acceleration. They are described by a system of PDEs coupling a continuity equation for the density of the distribution of states (forward in time) and a Hamilton-Jacobi (HJ) equation for the optimal value of a representative agent (backward in time).

The state variable is the pair $(x, v) \in \mathbb{R}^N \times \mathbb{R}^N$ where x stands for the position and v stands for the velocity. The dynamics is often referred to as the *double integrator*. In this case, the Hamiltonian of the system is neither strictly convex nor coercive, hence the available results on MFGs cannot be applied. Moreover, we will assume that the Hamiltonian is unbounded w.r.t. the velocity variable v . We prove the existence of a weak solution of the MFG system via a vanishing viscosity method and we characterize the distribution of states as the image of the initial distribution by the flow associated with the optimal control.

Keywords: Mean field games, first order Hamilton-Jacobi equations, double integrator, non-coercive Hamiltonian.

2010 AMS Subject classification: 35F50, 35Q91, 49K20, 49L25.

1 Introduction

The theory of mean field games (MFGs for short) is more and more investigated since the pioneering works [20, 21, 22] of Lasry and Lions: it aims at studying the asymptotic behavior of differential games (Nash equilibria) as the number of agents tends to infinity. In the present work, we study deterministic mean field games with finite time horizon in which the dynamics of a generic agent is controlled by the acceleration. They are described by a system of PDEs coupling a continuity equation for the density of the distribution of states (forward in time) and a Hamilton-Jacobi (HJ) equation for the optimal value of a representative agent (backward in time). The state variable is the pair $(x, v) \in \mathbb{R}^N \times \mathbb{R}^N$ where x stands for the position and v stands for the velocity.

^{*}Laboratoire Jacques-Louis Lions, Université Paris-Diderot, achdou@ljl.univ-paris-diderot.fr

[†]Dipartimento di Matematica “Tullio Levi-Civita”, Università di Padova, mannucci@math.unipd.it

[‡]Dipartimento di Ingegneria dell’Informazione, Università di Padova, claudio.marchi@unipd.it

[§]Univ Rennes, CNRS, IRMAR - UMR 6625, F-35000 Rennes, France, nicoletta.tchou@univ-rennes1.fr

The systems of PDEs are of the form

$$(1.1) \quad \begin{cases} (i) & -\partial_t u - v \cdot D_x u + H(x, v, D_v u) - F[m(t)](x, v) = 0 & \text{in } \mathbb{R}^{2N} \times (0, T) \\ (ii) & \partial_t m + v \cdot D_x m - \operatorname{div}_v (D_{p_v} H(x, v, D_v u) m) = 0 & \text{in } \mathbb{R}^{2N} \times (0, T) \\ (iii) & m(x, v, 0) = m_0(x, v), u(x, v, T) = G[m(T)](x, v), & \text{on } \mathbb{R}^{2N} \end{cases}$$

where T is a positive real number, $u = u(x, v, t)$, $m = m(x, v, t)$, $(x, v) \in \mathbb{R}^{2N}$, $t \in (0, T)$ and H is defined by

$$(1.2) \quad H(x, v, p_v) = \max_{\alpha \in \mathbb{R}^N} (-\alpha p_v - l(x, v, \alpha)).$$

We take F and G strongly regularizing and we assume that the running cost has the form $l(x, v, \alpha) = l(x, v) + \frac{1}{2}|\alpha|^2 + \frac{1}{2}|v|^2$, where $(x, v) \mapsto l(x, v)$ is a bounded and C^2 -bounded function.

Formally, systems of this form arise when the dynamics of the generic player is described by a *double integrator*:

$$(1.3) \quad \begin{cases} \xi'(s) &= \eta(s), & s \in (t, T), \\ \eta'(s) &= \alpha(s), & s \in (t, T), \\ \xi(t) &= x, \\ \eta(t) &= v, \end{cases}$$

and when the control law belongs to the space of the measurable functions with values in \mathbb{R}^N and is chosen in order to minimize the cost

$$(1.4) \quad J_t := J_t(\xi, \eta, \alpha) = \int_t^T l(\xi(s), \eta(s), \alpha(s)) + F[m(s)](\xi(s), \eta(s)) ds + G[m(T)](\xi(T), \eta(T)).$$

To summarize, the main features of this model are:

1. The control α is only involved in the dynamics of the second component of the state variable, see (1.3).
2. The running cost has the form

$$(1.5) \quad l(\xi, \eta, \alpha) = l(\xi, \eta) + \frac{1}{2}|\eta|^2 + \frac{1}{2}|\alpha|^2,$$

where $(\xi, \eta) \mapsto l(\xi, \eta)$ is a bounded C^2 function, thus the former is unbounded w.r.t. the variable η . Note that $|\eta|^2$ stands for a kinetic energy, whereas the term $|\alpha|^2$ is a penalty for large accelerations. Note also that the results of the present paper hold for a fairly large class of generalizations of (1.5).

3. Setting $f(\xi, \eta, \alpha) = (\eta, \alpha)$, the Hamiltonian associated to the control problem of a generic player is

$$\mathcal{H}(\xi, \eta, p) = \max_{\alpha \in \mathbb{R}^N} \{-p \cdot f(\xi, \eta, \alpha) - l(\xi, \eta, \alpha)\} = -p_x \cdot \eta + H(\xi, \eta, p_v),$$

where $p = (p_x, p_v)$ and H is defined in (1.2). The Hamilton \mathcal{H} is neither strictly convex nor coercive with respect to $p = (p_x, p_v)$. Hence the available results on the regularity of the value function u of the associated optimal control problem ([9], [14], [10]) and on the existence of a solution of the MFG system ([10]) cannot be applied.

To the best of our knowledge, this type of MFG systems have not been investigated yet. The main result of the present work is the existence of a solution of (2.1) and a characterization of the distribution of states as the image of the initial distribution by the optimal flow associated with the Hamilton-Jacobi equation. In order to establish the representation formula for m , we use some ideas introduced by P-L Lions in his lectures at Collège de France (2012) (see [24, 10]), some results proved in [12, 11], and Ambrosio's superposition principle [1]. Indeed, the lack of coercivity of \mathcal{H} makes it impossible to directly apply the arguments of [10, Sect. 4.1], in particular a contraction property of the flow associated to the dynamics (see Lemma 4.13 of [10]).

A similar approach was recently proposed for a class of non-coercive MFG when the generic player has some "forbidden direction", see ([25]), more precisely when, in the two dimensional case, the dynamics is of the form: $x'_1 = \alpha_1$, $x'_2 = h(x_1)\alpha_2$ and $h(x_1)$ may vanish.

The paper is organized as follows. In Section 2, we list our assumptions, give the definition of (weak) solution to system (2.1) and state the existence result for the latter. In Section 3, we obtain some regularity properties for the solution u of the Hamilton-Jacobi equation (2.1)-(i) with m fixed. These properties, combined with the uniqueness of the optimal trajectories of the associated control problem, will be crucial for proving the main theorem. In Section 4, we study the continuity equation (2.1)-(ii). An important ingredient is the vanishing viscosity method that is used to characterize its solution. Finally, Section 5 is devoted to the proof of the main Theorem 2.1.

2 Assumptions and main results

We consider the running cost $l(x, v, \alpha)$ of the form $l(x, v, \alpha) = l(x, v) + \frac{1}{2}|\alpha|^2 + \frac{1}{2}|v|^2$.

Then system (1.1) can be written

$$(2.1) \quad \begin{cases} (i) & -\partial_t u - v \cdot D_x u + \frac{1}{2}|D_v u|^2 - \frac{1}{2}|v|^2 - l(x, v) - F[m](x, v) = 0, & \text{in } \mathbb{R}^{2N} \times (0, T), \\ (ii) & \partial_t m + v \cdot D_x m - \operatorname{div}_v(D_v u m) = 0, & \text{in } \mathbb{R}^{2N} \times (0, T), \\ (iii) & m(x, v, 0) = m_0(x, v), \quad u(x, v, T) = G[m(T)](x, v), & \text{on } \mathbb{R}^{2N}, \end{cases}$$

where $H(x, v, p_v) = \frac{1}{2}|p_v|^2 - \frac{1}{2}|v|^2 - l(x, v)$.

Let \mathcal{P}_1 and \mathcal{P}_2 denote the spaces of Borel probability measures on \mathbb{R}^{2N} with respectively finite first and second order moments, endowed with the Monge-Kantorovich distances \mathbf{d}_1 , respectively \mathbf{d}_2 .

Let $C^2(\mathbb{R}^{2N})$ denote the space of twice differentiable functions with continuous and bounded derivatives up to order two. It is endowed with the norm

$$\|f\|_{C^2} := \sup_{(x,v) \in \mathbb{R}^{2N}} [|f(x, v)| + |Df(x, v)| + |D^2 f(x, v)|].$$

Hereafter, we shall make the following hypotheses:

Assumptions (H)

(H1) The functions F and G are real-valued continuous functions defined on $\mathcal{P}_1 \times \mathbb{R}^{2N}$

(H2) The function l is a real-valued C^2 function defined on \mathbb{R}^{2N}

(H3) The map $m \mapsto F[m](\cdot, \cdot)$ is Lipschitz continuous from \mathcal{P}_1 to $C^2(\mathbb{R}^{2N})$; moreover, there exists $C \geq \|l\|_{C^2}$ such that

$$\|F[m](\cdot, \cdot)\|_{C^2}, \quad \text{and} \quad \|G[m](\cdot, \cdot)\|_{C^2} \leq C, \quad \forall m \in \mathcal{P}_1$$

(H4) the initial distribution $m_0 : \mathbb{R}^{2N} \rightarrow \mathbb{R}^+$ has a compactly supported density (still named m_0 , with a slight abuse of notation) $m_0 \in C^\delta(\mathbb{R}^{2N})$, for some $\delta \in (0, 1)$.

Definition 2.1. *The pair (u, m) is a solution of system (2.1) if:*

- 1) $u \in W_{\text{loc}}^{1,\infty}(\mathbb{R}^{2N} \times [0, T])$, $m \in C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ and for all $t \in [0, T]$, $m(t)$ is absolutely continuous with respect to Lebesgue measure on \mathbb{R}^{2N} . Let $m(\cdot, \cdot, t)$ denote the density of $m(t)$. The function $(x, v, t) \mapsto m(x, v, t)$ is bounded.
- 2) equation (2.1)-(i) is satisfied by u in the viscosity sense
- 3) equation (2.1)-(ii) is satisfied by m in the sense of distributions.

We can now state the main result of this paper:

Theorem 2.1. *Under the assumptions (H) :*

1. System (2.1) has a solution (u, m) in the sense of Definition 2.1,
2. m is the image of m_0 by the flow

$$(2.2) \quad \begin{cases} x'(s) = v(s), & x(0) = x, \\ v'(s) = -D_v u(x(s), v(s), s), & v(0) = v. \end{cases}$$

3 The optimal control problem

In this section, we tackle the optimal control problem related to equation (2.1)-(i) with a fixed $\bar{m} \in C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$. To alleviate the notations, we introduce the functions

$$(3.1) \quad \ell(x, v, t) := l(x, v) + F[\bar{m}(t)](x, v) \quad \text{and} \quad g(x, v) := G[\bar{m}(T)](x, v),$$

which, from the set assumptions (H), satisfy

$$(3.2) \quad \|\ell(\cdot, \cdot, t)\|_{C^2}, \|\ell(x, v, \cdot)\|_C, \|g\|_{C^2} \leq C \quad \forall t \in [0, T], (x, v) \in \mathbb{R}^{2N}.$$

With the new notation, the optimal control problem to be solved by a representative agent whose state at time t is (x, v) is to find the control law α in order to minimize

$$(3.3) \quad J_t(\xi, \eta, \alpha) = \int_t^T \left[\frac{|\alpha|^2}{2} + \frac{|v|^2}{2} + \ell(\xi(s), \eta(s), s) \right] ds + g(\xi(T), \eta(T)),$$

by following the trajectory (1.3). Then the Cauchy problem given by (2.1)-(i) and its terminal condition becomes

$$(3.4) \quad \begin{cases} -\partial_t u - v \cdot D_x u + \frac{1}{2}|D_v u|^2 - \frac{1}{2}|v|^2 - \ell(x, v, t) = 0 & \text{in } \mathbb{R}^N \times \mathbb{R}^N \times (0, T), \\ u(x, v, T) = g(x, v) & \text{in } \mathbb{R}^N \times \mathbb{R}^N. \end{cases}$$

From (3.3), it is obvious that the control α must be chosen in $L^2(t, T; \mathbb{R}^N)$. Therefore, we can introduce the value function as follows:

Definition 3.1. *The value function for the cost J_t defined in (3.3) with dynamics (1.3) is*

$$(3.5) \quad u(x, v, t) := \inf \{ J_t(\xi, \eta, \alpha) : (\xi, \eta, \alpha) \in \mathcal{A}(x, v, t) \}$$

where

$$(3.6) \quad \mathcal{A}(x, v, t) = \left\{ (\xi, \eta, \alpha) : \begin{array}{l} (\xi, \eta) \in AC([t, T]; \mathbb{R}^{2N}), \alpha \in L^2(t, T; \mathbb{R}^N), \\ (\xi, \eta, \alpha) \text{ satisfy (1.3) and } \xi(t) = x, \eta(t) = v \end{array} \right\}.$$

Lemma 3.1. *i) (Existence of an optimal control.) For every $(x, v, t) \in \mathbb{R}^N \times \mathbb{R}^N \times (0, T)$, there exists an optimal control α^* for $u(x, v, t)$.*

ii) (Concatenation.) Let (ξ^, η^*) be an optimal trajectory for $u(x, v, t)$ corresponding to the control law α^* . For $r \in (t, T)$, let $(\tilde{\xi}^*, \tilde{\eta}^*)$ be an optimal trajectory for $u(\xi^*(r), \eta^*(r), r)$ with control $\tilde{\alpha}^*$. Then the concatenation of α^* and $\tilde{\alpha}^*$ at time r is optimal for $u(x, v, t)$ and, moreover,*

$$u(x, v, t) = u(\xi^*(r), \eta^*(r), r) + \int_t^r \left[\frac{|\alpha^*(s)|^2}{2} + \frac{|\eta^*(s)|^2}{2} + \ell(\xi^*(s), \eta^*(s), s) \right] ds.$$

*iii) Under the same assumption as in point (ii), the control $\alpha^*_{[[r, T]}$ is optimal for $u(\xi^*(r), \eta^*(r), r)$.*

iv) (Dynamic Programming Principle.) The Dynamic Programming Principle holds, namely

$$u(x, v, t) = \min_{(\xi, \eta, \alpha) \in \mathcal{A}(x, v, t)} \left\{ u(\xi(r), \eta(r), r) + \int_t^r \frac{|\alpha(s)|^2}{2} + \frac{|\eta(s)|^2}{2} + \ell(\xi(s), \eta(s), s) ds \right\}.$$

Proof. (i): let $\{\alpha_n\}_n$ be a sequence of minimizing control laws and (ξ_n, η_n) be the solution of (1.3) corresponding to α_n . Then, the boundedness of ℓ and the definition of J_t ensure that $\|\alpha_n\|_{L^2(t, T; \mathbb{R}^N)}$ are uniformly bounded. Then, possibly after extracting a subsequence, $\alpha_n \rightharpoonup \alpha^*$ in $L^2(t, T; \mathbb{R}^N)$, $\eta_n \rightarrow \eta^*$ in $C([t, T]; \mathbb{R}^N)$ and $\xi_n \rightarrow \xi^*$ in $C^1([t, T]; \mathbb{R}^N)$. The lower semi-continuity of J_t yields that α^* is optimal.

Points (ii), (iii) and (iv) are obtained by arguing exactly as in [25, Proposition 5.1] (points (1), (2) and (4) respectively), see also [10]. \square

Lemma 3.2. *The value function u has the following properties:*

1. $u(x, v, t)$ is Lipschitz continuous with respect to the variable x . Moreover $u(x, v, t)$ is Lipschitz continuous with respect to the variable v , locally in v , with a Lipschitz constant bounded by $C(1 + |v|)$.
2. $u(x, v, t)$ is Lipschitz continuous with respect to the time variable t , locally in v with a Lipschitz constant bounded by $C(1 + |v|^2)$.

The constant C depends only on the constants in assumptions (H).

Proof. 1. Fix $t \in [0, T]$. Let α be an optimal control law for $u(x, v, t)$ i.e.,

$$(3.7) \quad u(x, v, t) = \int_t^T \frac{1}{2} |\alpha(s)|^2 + \frac{1}{2} |v(s)|^2 + \ell(x(s), v(s), s) ds + g(x(T), v(T)),$$

where $(x(\cdot), v(\cdot))$ obeys to the dynamics (1.3).

We consider the path $(y(\cdot), w(\cdot))$ starting from (y, w) , with control $\alpha(\cdot)$. Hence, we obtain

$$\begin{aligned} y(s) &= y + w(s - t) + \int_t^s \int_t^\theta \alpha(\tau) d\tau d\theta = y - x + x(s) + (w - v)(s - t), \\ w(s) &= w + \int_t^s \alpha(\tau) d\tau = w - v + v(s). \end{aligned}$$

Note that

$$(3.8) \quad v(s) - w(s) = v - w, \quad x(s) - y(s) = x - y + (v - w)(s - t).$$

The definition of the value function (3.1) and relation (3.7) imply

$$\begin{aligned} u(y, w, t) &\leq \int_t^T \frac{1}{2} |\alpha(s)|^2 + \frac{1}{2} |w(s)|^2 + \ell(y(s), w(s), s) ds + g(y(T), w(T)) \\ &\leq u(x, v, t) - \int_t^T \frac{1}{2} |v(s)|^2 - \ell(x(s), v(s), s) ds - g(x(T), v(T)) \\ &\quad + \int_t^T \frac{1}{2} |w(s)|^2 + \ell(y(s), w(s), s) ds + g(y(T), w(T)) \\ &\leq u(x, v, t) + L_\ell(|x(s) - y(s)| + |v(s) - w(s)|) \\ &\quad + L_g(|x(T) - y(T)| + |v(T) - w(T)|) + \int_t^T \frac{1}{2} (|w(s)|^2 - |v(s)|^2) ds, \end{aligned}$$

where L_ℓ and L_g denote respectively the Lipschitz constants of ℓ and g w.r.t (x, v) . Hence, by (3.8),

$$\begin{aligned} \int_t^T \frac{1}{2} (|w(s)|^2 - |v(s)|^2) ds &= \int_t^T \frac{1}{2} |w - v| |w(s) + v(s)| ds \leq \\ &|w - v| \int_t^T |w + v + 2 \int_t^s \alpha(\tau) d\tau| ds \leq C |w - v| (|w + v| + (1 + |v|)), \end{aligned}$$

where the last inequality comes from (3.33) of Corollary 3.1 below. Hence we obtain

$$(3.9) \quad u(y, w, t) \leq u(x, v, t) + C|x - y| + K(v, w)|v - w|,$$

where $K(v, w) = C(|w + v| + 1 + |v|)$. Reverting the roles of (x, v) and (y, w) , we get the first result.

2. We fix (x, v) . From the concatenation property of optimal trajectories established in Lemma 3.1, if α is optimal for $u(x, v, t)$ and $(x(s), v(s))$ is the associated optimal trajectory, then

$$u(x, v, t) = u(x(s), v(s), s) + \int_t^s \frac{1}{2} |\alpha(r)|^2 + \frac{1}{2} |v(r)|^2 + \ell(x(r), v(r), r) dr$$

for any $s \in [t, T]$. Then

$$\begin{aligned} |u(x, v, t) - u(x, v, s)| &\leq |u(x, v, t) - u(x(s), v(s), s)| + |u(x(s), v(s), s) - u(x, v, s)| \\ &\leq \int_t^s \frac{1}{2} |\alpha(r)|^2 + \frac{1}{2} |v(r)|^2 + |\ell(x(r), v(r), r)| dr + L|x(s) - x| + L(v)|v(s) - v|, \end{aligned}$$

where the last two terms come from the Lipschitz continuity of u w.r.t. (x, v) : L is the Lipschitz constant of u with respect to x and $L(v)$ is a local Lipschitz constant of u with respect to v .

From (1.3) and the bound (3.33) in Corollary 3.1 below, we get the bounds for $x(s)$, $v(s)$ and α , i.e. $|v(s) - v| + |x(s) - x| \leq C(1 + |v|)|s - t|$, hence

$$(3.10) \quad |u(x, v, t) - u(x, v, s)| \leq C(1 + |v|^2)|s - t|,$$

which ends the proof. □

Remark 3.1. Taking $\alpha \equiv 0$ in (1.3), we get $\eta(s) = v$ and $\xi(s) = x + v(s - t)$; then, thanks to the boundedness of ℓ in (3.2), the value function verifies:

$$(3.11) \quad C \leq u(x, v, t) \leq C(1 + |v|^2),$$

and, by Lemma 3.2, it is also continuous. Hence, using the DPP in Lemma 3.1-(iv) and Proposition 3.1 below, the value function is the unique viscosity solution of (3.4) in the class of functions verifying (3.11).

Proposition 3.1. *There exists a viscosity solution u of (3.4) s.t. $\exists C \in \mathbb{R}$, $C \leq u(x, v, t) \leq C(|v|^2 + 1)$ and it is unique in this class of functions.*

Proof. We apply the comparison principle stated in [15] in order to prove the existence of a solution of (2.1)(i) by Perron's method. Testing (2.1)(i) with $w(x, v, t) = C_1(t - T) + \frac{C_2}{2}|v|^2 + C_3$ yields

$$-\partial_t u - v \cdot D_x u + \frac{1}{2}|D_v u|^2 - \frac{1}{2}|v|^2 - \ell(x, v, t) = -C_1 + C_2 \frac{1}{2}|v|^2 - \frac{1}{2}|v|^2 - \ell(x, v, t).$$

If $C_2 = 0$, $C_3 \leq -\|g\|_\infty$ and $-C_1 \leq \|\ell\|_\infty$, we have a subsolution. If $C_2 = 1$, $C_3 \geq \|g\|_\infty$ and $-C_1 \geq \|\ell\|_\infty$ we have a supersolution. \square

The following lemma deals with the semi-concavity of $u(x, v, t)$ w.r.t. (x, v) :

Lemma 3.3. *Under Hypothesis (H), $u(x, v, t)$ is semi-concave w.r.t. (x, v) with a linear modulus of semi-concavity, which depends only on the constants in assumptions (H).*

Proof. For any (x, v) , (y, w) and $\lambda \in [0, 1]$, consider $x_\lambda := \lambda x + (1 - \lambda)y$, $v_\lambda := \lambda v + (1 - \lambda)w$.

Let α be an optimal control for $u(x_\lambda, v_\lambda, t)$; hence, the associated trajectory is

$$(3.12) \quad x_\lambda(s) = x_\lambda + v_\lambda(s - t) + \int_t^s \int_t^\theta \alpha(\tau) d\tau d\theta, \quad v_\lambda(s) = v_\lambda + \int_t^s \alpha(\tau) d\tau$$

and

$$u(x_\lambda, v_\lambda, t) = \int_t^T \frac{1}{2}|\alpha(s)|^2 + \frac{1}{2}|v_\lambda(s)|^2 + \ell(x_\lambda(s), v_\lambda(s), s) ds + g(x_\lambda(T), v_\lambda(T)).$$

Let $(x(s), v(s))$ be the trajectory starting at (x, v) at time t with control α and $(y(s), w(s))$ the trajectory starting at (y, w) at time t still with control α .

We have to estimate

$$\begin{aligned} & \lambda u(x, v, t) + (1 - \lambda)u(y, w, t) - u(x_\lambda, v_\lambda, t) \\ & \leq \int_t^T \frac{1}{2}\lambda|v(s)|^2 + (1 - \lambda)\frac{1}{2}|w(s)|^2 - \frac{1}{2}|v_\lambda(s)|^2 ds \\ & \quad + \int_t^T \lambda\ell(x(s), v(s), s) + (1 - \lambda)\ell(y(s), w(s), s) - \ell(x_\lambda(s), v_\lambda(s), s) ds \\ & \quad + \lambda g(x(T), v(T)) + (1 - \lambda)g(y(T), w(T)) - g(x_\lambda(T), v_\lambda(T)). \end{aligned}$$

Since

$$(3.13) \quad v(s) = v + \int_t^s \alpha(\tau) d\tau, \quad w(s) = w + \int_t^s \alpha(\tau) d\tau, \quad v_\lambda(s) = \lambda v + (1 - \lambda)w + \int_t^s \alpha(\tau) d\tau,$$

we get

$$\begin{aligned}
(3.14) \quad & \lambda \frac{1}{2} |v(s)|^2 + (1-\lambda) \frac{1}{2} |w(s)|^2 - \frac{1}{2} |v_\lambda(s)|^2 \\
&= (\lambda v + (1-\lambda)w - \lambda v - (1-\lambda)w) \int_t^s \alpha(\tau) d\tau + \lambda \frac{|v|^2}{2} + (1-\lambda) \frac{|w|^2}{2} - \frac{1}{2} |\lambda v + (1-\lambda)w|^2 \\
&= \frac{1}{2} \lambda(1-\lambda) |v|^2 + \frac{1}{2} \lambda(1-\lambda) |w|^2 - \lambda(1-\lambda) v \cdot w = \frac{1}{2} \lambda(1-\lambda) |v-w|^2.
\end{aligned}$$

Hence

$$(3.15) \quad \int_t^T \frac{1}{2} \lambda |v(s)|^2 + (1-\lambda) \frac{1}{2} |w(s)|^2 - \frac{1}{2} |v_\lambda(s)|^2 ds = \frac{1}{2} \lambda(1-\lambda) |v-w|^2 (T-t).$$

Now, we have to estimate the terms $\lambda \ell(x(s), v(s), s) + (1-\lambda) \ell(y(s), w(s), s) - \ell(x_\lambda(s), v_\lambda(s), s)$ and $\lambda g(x(T), v(T)) + (1-\lambda) g(y(T), w(T)) - g(x_\lambda(T), v_\lambda(T))$. We write the algebra for the second term, since the treatment of the first term is similar. The Taylor expansion of g centered at $(x_\lambda(T), v_\lambda(T))$ gives

$$(3.16) \quad g(x(T), v(T)) = \begin{cases} g(x_\lambda(T), v_\lambda(T)) + Dg(x_\lambda(T), v_\lambda(T))(x(T) - x_\lambda(T), v(T) - v_\lambda(T)) \\ + R_1, \end{cases}$$

where R_1 is the error term in the expansion, namely

$$(3.17) \quad R_1 = \frac{1}{2} (x(T) - x_\lambda(T), v(T) - v_\lambda(T)) D^2 g(\xi_1, \eta_1) (x(T) - x_\lambda(T), v(T) - v_\lambda(T))^T,$$

for suitable ξ_1, η_1 .

From (3.12) and (3.13), we get

$$\begin{aligned}
(3.18) \quad & x(s) - x_\lambda(s) = (1-\lambda)((x-y) + (v-w)(s-t)), \\
& v(s) - v_\lambda(s) = (1-\lambda)(v-w), \\
& y(s) - x_\lambda(s) = \lambda((y-x) + (w-v)(s-t)), \\
& w(s) - v_\lambda(s) = \lambda(w-v),
\end{aligned}$$

hence the error term can be written as

$$(3.19) \quad R_1 = \frac{1}{2} (1-\lambda)^2 (x-y + (v-w)(T-t), v-w) D^2 g(\xi_1, \eta_1) (x-y + (v-w)(T-t), v-w)^T.$$

Similarly

$$g(y(T), w(T)) = g(x_\lambda(T), v_\lambda(T)) + Dg(x_\lambda(T), v_\lambda(T))(y(T) - x_\lambda(T), w(T) - v_\lambda(T)) + R_2,$$

where

$$\begin{aligned}
R_2 &= \frac{1}{2} (y(T) - x_\lambda(T), w(T) - v_\lambda(T)) D^2 g(\xi_2, \eta_2) (y(T) - x_\lambda(T), w(T) - v_\lambda(T))^T \\
&= \frac{1}{2} \lambda^2 (y-x + (w-v)(T-t), w-v) D^2 g(\xi_2, \eta_2) (y-x + (w-v)(T-t), w-v)^T.
\end{aligned}$$

At this point, taking into account that from (3.18),

$$\begin{aligned}
(3.20) \quad & \lambda Dg(x_\lambda(T), v_\lambda(T))(x(T) - x_\lambda(T), v(T) - v_\lambda(T)) \\
&+ (1-\lambda) Dg(x_\lambda(T), v_\lambda(T))(y(T) - x_\lambda(T), w(T) - v_\lambda(T)) \\
&= Dg(x_\lambda(T), v_\lambda(T))(\lambda(x(T) - x_\lambda(T)) \\
&+ (1-\lambda)(y(T) - x_\lambda(T)), \lambda(v(T) - v_\lambda(T)) + (1-\lambda)(w(T) - v_\lambda(T))) \\
&= 0,
\end{aligned}$$

we obtain that

$$\begin{aligned}
(3.21) \quad & \lambda g(x(T), v(T)) + (1 - \lambda)g(y(T), w(T)) - g(x_\lambda(T), v_\lambda(T)) \\
& = \lambda R_1 + (1 - \lambda)R_2 \\
& \leq (1 - \lambda)\lambda C_T \|D^2 g\|_\infty (|x - y|^2 + |v - w|^2).
\end{aligned}$$

Hence from (3.15), (3.20), (3.21) we get

$$\begin{aligned}
& \lambda u(x, v, t) + (1 - \lambda)u(y, w, t) - u(x_\lambda, v_\lambda, t) \\
& \leq \frac{\lambda(1 - \lambda)}{2}(v - w)^2(T - t) + C_T(1 - \lambda)\lambda (\|D^2 g\|_\infty + \|D^2 \ell\|_\infty) (|x - y|^2 + |v - w|^2).
\end{aligned}$$

We obtain that u is semi-concave in (x, v) with a linear modulus of semi-concavity. \square

Pontryagin's maximum principle yields the following necessary optimality conditions:

Proposition 3.2 (Necessary conditions for optimality). *Let (x^*, v^*, α^*) be optimal for $u(x, v, t)$ in (3.5). There exists an arc $p = (p_x, p_v) \in AC([t, T]; \mathbb{R}^N \times \mathbb{R}^N)$, hereafter called the costate, such that*

1. (α^*, x^*, v^*, p) satisfies the adjoint equations: for a.e. $s \in [t, T]$,

$$(3.22) \quad p'_x = D_x \ell(x^*, v^*, s),$$

$$(3.23) \quad p'_v = -p_x + v^* + D_v \ell(x^*, v^*, s),$$

the transversality condition

$$(3.24) \quad p(T) = -Dg(x^*(T), v^*(T)),$$

together with the maximum condition: for almost all $s \in [t, T]$,

$$(3.25) \quad \max_{\alpha} p_x \cdot v^* + p_v \cdot \alpha - \frac{|\alpha|^2}{2} - \frac{|v^*|^2}{2} = p_x \cdot v^* + p_v \cdot \alpha^* - \frac{|\alpha^*|^2}{2} - \frac{|v^*|^2}{2}.$$

2. The optimal control α^* is given by

$$(3.26) \quad \alpha^* = p_v, \text{ a.e. in } [t, T].$$

3. The triple (x^*, v^*, p) satisfies the system of differential equations: for a.e. $s \in [t, T]$

$$(3.27) \quad x' = v,$$

$$(3.28) \quad v' = p_v,$$

$$(3.29) \quad p'_x = D_x \ell(x, v, s),$$

$$(3.30) \quad p'_v = -p_x + v + D_v \ell(x, v, s),$$

with the mixed boundary conditions $x^*(t) = x$, $v^*(t) = v$, $p(T) = -Dg(x^*(T), v^*(T))$.

Proof. 1. Hypothesis (3.2) ensures that our control problem satisfies the assumption [14, Hypothesis 22.16], so we can invoke [14, Theorem 22.17] on the maximum principle for problems with unbounded control. Moreover, since there is no constraint on the state variable at T , the same arguments as in [14, Corollary 22.3] ensure that the necessary conditions hold in normal form.

2. The maximum condition (3.25) implies that

$$D_\alpha \left(p_x \cdot v^* + p_v \cdot \alpha - \frac{|\alpha|^2}{2} - \frac{|v^*|^2}{2} - f(x^*, v^*) \right)_{\alpha=\alpha^*} = 0 \quad \text{for a.e. } s \in [t, T]$$

from which we get (3.26).

3. Conditions (3.27) – (3.28) follow directly from (1.3) and (3.26). Conditions (3.29) and (3.30) coincide with (3.22), (3.23). \square

Corollary 3.1 (Feedback control and regularity). *Let (x^*, v^*, α^*) be optimal for $u(x, v, t)$ and $p = (p_x, p_v)$ be the related costate as in Proposition 3.2. Then:*

1. *The costate p is uniquely expressed in terms of x^*, v^* for every $s \in [t, T]$ by*

$$(3.31) \quad \begin{cases} p_x(s) = -D_x g(x^*(T), v^*(T)) - \int_s^T D_x \ell(x^*(\tau), v^*(\tau), \tau) d\tau, \\ p_v(s) = -D_v g(x^*(T), v^*(T)) - \int_s^T D_v \ell(x^*(\tau), v^*(\tau), \tau) + v^*(\tau) - p_x(\tau) d\tau. \end{cases}$$

2. *The optimal control α^* is a feedback control (i.e., a function of x^*, v^*), uniquely expressed in terms of x^*, v^* for a.e. $s \in [t, T]$ by*

$$(3.32) \quad \alpha^*(s) = p_v(s).$$

3. *The optimal trajectory (x^*, v^*) and the optimal control α^* are of class C^1 . In particular the equalities (3.26) – (3.32) do hold for every $s \in [t, T]$. Moreover*

$$(3.33) \quad \begin{aligned} \|v^*\|_{C^1} + \|\alpha^*\|_{C^1} &\leq C(1 + |v|), \\ \|x^*\|_{C^1} &\leq |x| + C(1 + |v|). \end{aligned}$$

4. *Assume that, for some $k \in \mathbb{N}$, $D_x \ell(x, v, s)$, $D_v \ell(x, v, s)$ are of class C^k . Then (x^*, v^*) , p and α^* are of class C^{k+1} .*

Proof. Point 1 is an immediate consequence of (3.29) – (3.30) together with the endpoint condition $p(T) = -Dg(x^*(T))$. Point 2 follows then directly from (3.26).

3. Since x^*, v^* are continuous, the continuity of α^* follows from (3.31) and (3.32). Then (1.3) implies that $x^*, v^* \in C^1$. Relations (3.31) and (3.32) imply, respectively, that p and α^* are of class C^1 . The bounds (3.33) are obtained from the fact that

1. $(v^*)''(\tau) - v^*(\tau)$ is bounded uniformly in x and v , $\tau \in [t, T]$
2. $(v^*)(t) = v$
3. $(v^*)'(T)$ is bounded uniformly in x and v .

4. The relations (3.31) and the C^1 -regularity of x^*, v^* and p imply that, actually, $p \in C^2$. Therefore, (3.32) gives the C^2 -regularity of α^* and, finally, (2.2) yields the C^2 -regularity of x^*, v^* . Further regularity of x^*, v^*, α^* and p follows by a standard bootstrap inductive argument. \square

Remark 3.2. *Taking advantage of Corollary 3.1-(3), we will always consider the representation of the optimal control α^* which belongs to C^1 .*

Corollary 3.2 that follows implies that the optimal trajectories for $u(x, v, t)$ do not bifurcate at any time $r > t$.

Corollary 3.2. *Under Hypothesis (3.2), let (x^*, v^*) be an optimal trajectory for $u(x, v, t)$. For every $t < r < T$, there are no other optimal trajectories for $u(x^*(r), v^*(r), r)$ other than (x^*, v^*) restricted to $[r, T]$.*

Proof. 1. Let $r \in (t, T)$ and (y^*, w^*) be an optimal trajectory for $u(x^*(r), v^*(r), r)$. Lemma 3.1 ensures that (z^*, ν^*) , the concatenation of (x^*, v^*) with (y^*, w^*) at r is an optimal trajectory for $u(x, v, t)$. Let $p := (p_x, p_v), q := (q_x, q_v)$ be the costates corresponding respectively to (x^*, v^*) and to (z^*, ν^*) . Both (x^*, v^*, p) and (z^*, ν^*, q) satisfy (3.27) – (3.30) on $[t, T]$. Now, Corollary 3.1 shows that (x^*, v^*) and (z^*, ν^*) are of class C^1 . Since $x^* = z^*, v^* = \nu^*$ on $[t, r]$, we choose τ such that $t < \tau < r$. From (3.28), we get

$$p_v(\tau) = q_v(\tau).$$

Moreover, from (3.28) and (3.30), we also get that

$$p_x(\tau) = q_x(\tau).$$

Therefore, both (x^*, v^*, p) and (z^*, ν^*, q) are solutions to the same Cauchy problem on $[t, T]$ with the first order differential system (3.27)-(3.30) and Cauchy data at τ . The regularity assumptions on ℓ, g and Cauchy-Lipschitz Theorem guarantee the uniqueness of the solution. Thus $x^* = z^*, v^* = \nu^*$ on $[t, T]$, from which we obtain the desired identities $x^* = y^*$ and $v^* = w^*$ on $[r, T]$. \square

Definition 3.2. *For any $(x, v, t) \in \mathbb{R}^{2N} \times [0, T]$, let $\mathcal{U}(x, v, t)$ denote the set of optimal controls for the value function $u(x, v, t)$ defined in (3.5).*

Remark 3.3. *Lemma 3.1-(i) and Remark 3.2 ensure that $\emptyset \neq \mathcal{U}(x, v, t) \subset C^1([t, T]; \mathbb{R}^N)$.*

Lemma 3.4. *The following properties hold:*

1. *The function $u(x, \cdot, t)$ is differentiable at v if and only if the set $\{\alpha(t) : \alpha \in \mathcal{U}(x, v, t)\}$ is a singleton. Moreover $D_v u(x, v, t) = -\alpha(t)$.*
2. *In particular, if $\mathcal{U}(x, v, t)$ is a singleton, then, calling $(x(s), v(s))$ the optimal trajectory associated to the singleton $\mathcal{U}(x, v, t)$, $D_v u(x(s), v(s), s)$ exists for any $s \in [t, T]$.*
3. *If $u(\cdot, \cdot, t)$ is differentiable at (x, v) , then $\mathcal{U}(x, v, t)$ is a singleton.*

Proof. 1. We prove that if $D_v u(x, v, t)$ exists, then all $\alpha(\cdot) \in \mathcal{U}(x, v, t)$ take the same value $\alpha(t)$ at t and $D_v u(x, v, t) = -\alpha(t)$. If $\alpha(\cdot) \in \mathcal{U}(x, v, t)$, calling $(x(\cdot), v(\cdot))$ the corresponding optimal trajectory, then

$$u(x, v, t) = \int_t^T \frac{1}{2} |\alpha(s)|^2 + \frac{1}{2} |v(s)|^2 + \ell(x(s), v(s), s) ds + g(x(T), v(T)),$$

and $(x(\cdot), v(\cdot))$ and $\alpha(\cdot)$ satisfy the necessary conditions for optimality proved in Proposition 3.2. Take $h = (h_1, h_2) \in \mathbb{R}^{2N}$ and consider the solution $(y(\cdot), w(\cdot))$ of (1.3) with initial condition $(y(t), w(t)) = (x + h_1, v + h_2)$ and control α , namely

$$\begin{aligned} y(s) &= x + h_1 + (v + h_2)(s - t) + \int_t^s \int_t^\theta \alpha(\tau) d\tau d\theta = x(s) + h_1 + h_2(s - t), \\ w(s) &= v + h_2 + \int_t^s \alpha(\tau) d\tau = v(s) + h_2. \end{aligned}$$

Hence,

$$\begin{aligned} (3.34) \quad & u(x + h_1, v + h_2, t) - u(x, v, t) \\ & \leq \int_t^T \frac{1}{2} |w(s)|^2 - \frac{1}{2} |v(s)|^2 + \ell(y(s), w(s), s) - \ell(x(s), v(s), s) ds \\ & \quad + g(y(T), w(T)) - g(x(T), v(T)) \\ & = \int_t^T \frac{1}{2} |v(s) + h_2|^2 - \frac{1}{2} |v(s)|^2 + \ell(x(s) + h_1 + h_2(s - t), v(s) + h_2, s) - \ell(x(s), v(s), s) ds \\ & \quad + g(x(T) + h_1 + h_2(T - t), v(T) + h_2) - g(x(T), v(T)) \\ & = \int_t^T \frac{1}{2} h_2^2 + h_2 \cdot v(s) + \ell(x(s) + h_1 + h_2(s - t), v(s) + h_2, s) - \ell(x(s), v(s), s) ds \\ & \quad + g(x(T) + h_1 + h_2(T - t), v(T) + h_2) - g(x(T), v(T)). \end{aligned}$$

The differentiability of u w.r.t. v yields

$$\begin{aligned} D_v u(x, v, t) &= \int_t^T v(s) ds + \int_t^T D_x \ell(x(s), v(s), s)(s - t) + D_v \ell(x(s), v(s), s) ds \\ & \quad + D_x g(x(T), v(T))(T - t) + D_v g(x(T), v(T)). \end{aligned}$$

By (3.29) and (3.24), we obtain

$$\begin{aligned} \int_t^T D_x \ell(x(s), v(s), s)(s - t) ds &= \int_t^T p'_x(s)(s - t) ds = p_x(T)(T - t) - \int_t^T p_x(s) ds \\ &= -D_x g(x(T), v(T))(T - t) - \int_t^T p_x(s) ds. \end{aligned}$$

Hence

$$\begin{aligned} D_v u(x, v, t) &= \int_t^T (v(s) - p_x(s) + D_v \ell(x(s), v(s), s)) ds + D_v g(x(T), v(T)) \\ &= \int_t^T p'_v(s) ds + D_v g(x(T), v(T)) \\ &= -p_v(t) = -\alpha(t), \end{aligned}$$

where the last two inequalities are due to (3.30), (3.26) and the terminal condition for p . This uniquely determines the value of $\alpha(\cdot)$ at time t .

Conversely we prove that, if all $\alpha(\cdot) \in \mathcal{U}(x, v, t)$ take the same value $\alpha(t)$ at t , then $D_v u(x, v, t)$ exists. Fix x and t . From the semi-concavity of $u(x, \cdot, t)$, the differentiability of $u(x, \cdot, t)$ at v will follow from the fact that $D_v^* u(x, v, t)$ is a singleton (see [9, Proposition 3.3.4]). Recall that the set of reachable gradients of $u(x, \cdot, t)$ is defined by

$$D_v^* u(x, v, t) = \left\{ \chi \in \mathbb{R}^N : \exists (v_n)_{n \in \mathbb{N}} \text{ with } \begin{cases} \lim_{n \rightarrow \infty} v_n = v, \\ u(x, \cdot, t) \text{ is differentiable at } v_n, \\ \lim_{n \rightarrow \infty} D_v u(x, v_n, t) = \chi. \end{cases} \right\}.$$

Take $\chi \in D_v^* u(x, v, t)$. By definition of $D_v^* u(x, v, t)$ there exist sequences $\{v_n\}$, $\{\chi_n = D_v u(x, v_n, t)\}$ such that

$$(3.35) \quad v_n \rightarrow v \quad \text{and} \quad \chi_n \rightarrow \chi.$$

Consider $\alpha_n \in \mathcal{U}(x, v_n, t)$; by the other part of the statement (already proven), we know that

$$(3.36) \quad -\alpha_n(t) = D_v u(x, v_n, t) = \chi_n.$$

From estimate (3.33) in Corollary 3.1, we see that

$$(3.37) \quad \|\alpha'_n\|_\infty \leq C(1 + |v_n|) \leq C, \text{ for any } n.$$

Hence from Ascoli-Arzelà Theorem, we deduce that, after extracting a subsequence, α_n uniformly converge to some $\alpha \in C([t, T]; \mathbb{R}^N)$. In particular, calling $(x_n(\cdot), v_n(\cdot))$ the trajectory associated to α_n starting from (x, v_n) :

$$x_n(s) = x + v_n(s - t) + \int_t^s \int_t^\theta \alpha_n(\tau) d\tau d\theta, \quad \text{and} \quad v_n(s) = v_n + \int_t^s \alpha_n(\tau) d\tau.$$

we get:

$$\begin{aligned} x_n(s) &\rightarrow x(s) = x + v(s - t) + \int_t^s \int_t^\theta \alpha(\tau) d\tau d\theta, \quad \text{uniformly in } [t, T], \\ v_n(s) &\rightarrow v(s) = v + \int_t^s \alpha(\tau) d\tau \quad \text{uniformly in } [t, T]. \end{aligned}$$

Moreover, by classical arguments of stability, α is optimal, i.e. $\alpha \in \mathcal{U}(x, v, t)$. The uniform convergence of the α_n yields in particular that $\alpha_n(t) \rightarrow \alpha(t)$ where $\alpha(t)$ is uniquely determined by assumption. By (3.35) and (3.36), we get that $\chi_n \rightarrow \chi = \alpha(t)$. This implies that $D_v^* u(x, v, t)$ is a singleton, then $D_v u(x, v, t)$ exists. Going back to the first part of the proof, we see that $D_v u(x, v, t) = -\alpha(t)$.

2. If $\mathcal{U}(x, v, t) = \{\alpha(\cdot)\}$, then for any $s \in [t, T]$, $\alpha(s)$ is uniquely determined. Indeed, if there exists $\beta \in \mathcal{U}(x(s), v(s), s)$, then the concatenation γ of α and β (see Lemma 3.1) is also optimal, i.e. $\gamma \in \mathcal{U}(x, v, t) = \{\alpha(\cdot)\}$.

Then from point 1 with $t = s$ at $(x(s), v(s))$, we deduce that $D_v u(x(s), v(s), s)$ exists.

3. From point 1, we know that for any $\alpha(\cdot) \in \mathcal{U}(x, v, t)$, $\alpha(t)$ is unique and coincides with $-D_v u(x, v, t)$. Hence, relation (3.26) ensures $p_v(t) = -D_v u(x, v, t)$. On the other hand, note that, since $D_x u(x, v, t)$ exists, we get from (3.34) that

$$\begin{aligned} D_x u(x, v, t) &= \int_t^T D_x \ell(x(s), v(s), s) ds + D_x g(x(T), v(T)) \\ &= \int_t^T p'_x(s) ds + D_x g(x(T), v(T)) = -p_x(t); \end{aligned}$$

thus, $p_x(t)$ and $p_v(t)$ are both uniquely determined. Hence (3.27)-(3.30) is a system of differential equations with initial conditions $x(t), v(t), p_x(t)$ and $p_v(t)$ which admits a unique solution $(x(\cdot), v(\cdot), p_x(\cdot), p_v(\cdot))$ by Cauchy-Lipschitz theorem, and $(x(\cdot), v(\cdot))$ is the unique optimal trajectory starting from (x, v) , associated to the unique optimal control law $\alpha(\cdot) = p_v(\cdot)$. \square

Lemma 3.5 (optimal synthesis). *Consider $\xi \in \mathbb{R}^N$ and $\eta \in \mathbb{R}^N$.*

1. *Let $x \in C^1([t, T]; \mathbb{R}^N)$, $v \in AC([t, T]; \mathbb{R}^N)$ be such that*

$$(3.38) \quad x(t) = \xi, \quad \text{and} \quad v(t) = \eta,$$

and for almost every $s \in (t, T)$,

$$(3.39) \quad u(x(s), \cdot, s) \text{ is differentiable at } v(s),$$

and

$$(3.40) \quad \begin{aligned} x'(s) &= v(s), \\ v'(s) &= -D_v u(x(s), v(s), s), \end{aligned}$$

where u is the solution of (3.4). Under these assumptions, the control law $\alpha(s) = v'(s) = -D_v u(x(s), v(s), s)$ is optimal for $u(\xi, \eta, t)$.

2. *If $u(\cdot, \cdot, t)$ is differentiable at (ξ, η) , then problem (3.38), (3.40) has a unique solution corresponding to the optimal trajectory.*

Proof. We adapt the arguments of [10, Lemma 4.11]. Fix $(t, \xi, \eta) \in (0, T) \times \mathbb{R}^{2N}$. Let $x \in C^1([t, T]; \mathbb{R}^N)$, $v \in AC([0, T]; \mathbb{R}^N)$ be as in the statement of Lemma 3.5. Note that v is Lipschitz continuous. Therefore, from Lemma 3.2, the function $s \mapsto u(x(s), v(s), s)$ is Lipschitz continuous as well. Hence, for almost every $s \in [t, T]$,

1. (3.39) and (3.40) hold,
2. the function $u(x(\cdot), v(\cdot), \cdot)$ admits a derivative at s .

Lebourg's Theorem for Lipschitz functions (see [13, Thm 2.3.7] and [13, Thm 2.5.1]) ensures that, for any positive and sufficiently small number h , there exists (y_h, w_h, s_h) in the open line segment $((x(s), v(s), s), (x(s+h), v(s+h), s+h))$ and $(\chi_x^h, \chi_v^h, \chi_t^h) \in \text{conv}(D_{x,v,t}^* u(y_h, w_h, s_h))$ such that

$$(3.41) \quad u(x(s+h), v(s+h), s+h) - u(x(s), v(s), s) = \chi_x^h \cdot (x(s+h) - x(s)) + \chi_v^h \cdot (v(s+h) - v(s)) + \chi_t^h h.$$

Here, $\text{conv}(A)$ stands for the convex hull of a set A , and $D_{x,v,t}^* u(y_h, w_h, s_h)$ stands for the reachable gradient at (y_h, w_h, s_h) with respect to the three variables x, v and t .

By Carathéodory's theorem, see [9, Thm A.1.6], there exist $(\lambda^{h,i}, \chi_x^{h,i}, \chi_v^{h,i}, \chi_t^{h,i})_{i=1, \dots, 2N+2}$ such that $\lambda^{h,i} \geq 0$, $\sum_{i=1}^{2N+2} \lambda^{h,i} = 1$, $(\chi_x^{h,i}, \chi_v^{h,i}, \chi_t^{h,i}) \in D_{x,v,t}^* u(y_h, w_h, s_h)$ and $(\chi_x^h, \chi_v^h, \chi_t^h) = \sum_{i=1}^{2N+2} \lambda^{h,i} (\chi_x^{h,i}, \chi_v^{h,i}, \chi_t^{h,i})$. Note that, by [9, Prop 3.3.4(a)], for all $i = 1, \dots, 2N+2$, $\chi_v^{h,i}$ converges to $D_v u(x(s), v(s), s)$ as $h \rightarrow 0$; hence, χ_v^h also converges to $D_v u(x(s), v(s), s)$, as $h \rightarrow 0$.

On the other hand, since u is a viscosity solution to equation (3.4) and $(\chi_x^{h,i}, \chi_v^{h,i}, \chi_t^{h,i}) \in D_{x,v,t}^* u(y_h, w_h, s_h)$, we obtain that for all $i \in 1, \dots, 2N+2$,

$$-\chi_t^{h,i} + \frac{1}{2} |\chi_v^{h,i}|^2 - \frac{1}{2} |w_h|^2 - w_h \cdot \chi_x^{h,i} = \ell(y_h, w_h, s_h).$$

Therefore, $\chi_t^h + w_h \cdot \chi_x^h = \frac{1}{2} \sum_{i=1}^{2N+2} \lambda^{h,i} |\chi_v^{h,i}|^2 - \frac{1}{2} |w_h|^2 - \ell(y_h, w_h, s_h)$ converges to $\frac{1}{2} |D_v u(x(s), v(s), s)|^2 - \frac{1}{2} |v(s)|^2 - \ell(x(s), v(s), s)$ as $h \rightarrow 0$.

Then dividing (3.41) by h and letting h tend to 0, we get that

$$\begin{aligned} & \frac{d}{ds} (u(x(s), v(s), s)) \\ &= D_v u(x(s), v(s), s) \cdot v'(s) + \frac{1}{2} |D_v u(x(s), v(s), s)|^2 - \frac{1}{2} |v(s)|^2 - \ell(x(s), v(s), s). \end{aligned}$$

Recalling (3.40), we get

$$\frac{d}{ds} (u(x(s), v(s), s)) = -\frac{1}{2} |D_v u(x(s), v(s), s)|^2 - \frac{1}{2} |v(s)|^2 - \ell(x(s), v(s), s).$$

or in equivalent manner,

$$\frac{d}{ds} (u(x(s), v(s), s)) = -\frac{1}{2} |v'(s)|^2 - \frac{1}{2} |v(s)|^2 - \ell(x(s), v(s), s),$$

which holds for almost every s . Integrating this equality on $[t, T]$ and taking into account the terminal condition in (3.4), we obtain

$$u(x, v, t) = \int_t^T \frac{1}{2} |v'(s)|^2 + \frac{1}{2} |v(s)|^2 + \ell(x(s), v(s), s) ds + g(x(T)).$$

Therefore, the control law $\alpha(s) = v'(s) = -D_v u(x(s), v(s), s)$ is optimal. This achieves the proof of the first statement.

The second statement is a direct consequence of Lemma 3.4. \square

4 The continuity equation

In this section, our aim is to study equation (2.1)-(ii), and more precisely the well-posedness of

$$(4.1) \quad \begin{cases} \partial_t m + v \cdot D_x m - \operatorname{div}_v (m D_v u) = 0, & \text{in } \mathbb{R}^{2N} \times (0, T), \\ m(x, v, 0) = m_0(x, v), & \text{in } \mathbb{R}^{2N}, \end{cases}$$

where u is a solution to problem (3.4), which we rewrite here for the sake of clarity

$$\begin{cases} -\partial_t u - v \cdot D_x u + \frac{1}{2} |D_v u|^2 - \frac{1}{2} |v|^2 - l(x, v) = F[\overline{m}(t)](x, v), & \text{in } \mathbb{R}^N \times \mathbb{R}^N \times (0, T), \\ u(x, v, T) = G[\overline{m}(T)](x, v), & \text{on } \mathbb{R}^{2N}, \end{cases}$$

and \overline{m} is fixed and belongs to $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$.

Remark 4.1. Note that u is the value function of the optimal control problem (3.5) with J_t and ℓ respectively given by (3.3) and (3.1).

It is worth to observe that the differential equation in (4.1) can also be written

$$\partial_t m - \operatorname{div}_{x,v}(m b) = 0,$$

with $b := (-v, D_v u)$. In the present framework, the semi-concavity proved in Lemma 3.3 does not imply that the flow $\Phi(x, t, s)$ given by Lemma 3.5 has a Lipschitz continuous inverse, by contrast with [10, Lemma 4.13]. Moreover, the drift b is only locally bounded; this lack of regularity makes it impossible to apply the standard results for drifts which are Lipschitz continuous (uniqueness, existence and representation formula of m as the push-forward of m_0 through the characteristic flow; e.g., see [1, Proposition 8.1.8]). We shall overcome this difficulty by applying Ambrosio's superposition principle [1, Theorem 8.2.1]. The latter yields a representation formula of m as the push-forward of some measure on $C([0, T]; \mathbb{R}^{2N})$ through the evaluation map e_t defined by $e_t(\gamma) = \gamma(t)$ for all continuous function γ with value in \mathbb{R}^{2N} . In the following theorem, we state existence, uniqueness, and some regularity results for (4.1):

Theorem 4.1. *Under assumptions (H), for any $\overline{m} \in C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$, there is a unique $m \in C^{\frac{1}{2}}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N})) \cap L^\infty(0, T; \mathcal{P}_2(\mathbb{R}^{2N}))$ which solves problem (4.1) in the sense of Definition 2.1.*

Moreover $m(t, \cdot)$ satisfies: for any $\phi \in C_0(\mathbb{R}^{2N})$, for any $t \in [0, T]$,

$$(4.2) \quad \int_{\mathbb{R}^{2N}} \phi(x, v) m(x, v, t) dx dv = \int_{\mathbb{R}^{2N}} \phi(\overline{\gamma}_{x,v}(t)) m_0(x, v) dx dv,$$

where, for a.e. $(x, v) \in \mathbb{R}^{2N}$, $\overline{\gamma}_{x,v}$ is the solution to (2.2).

The proof of Theorem 4.1 is given in the next two subsections which are devoted respectively to existence (see Proposition 4.1) and to uniqueness and the representation formula (see Proposition 4.2).

4.1 Existence of the solution

We wish to establish the existence of a solution to the continuity equation via a vanishing viscosity method applied to the *whole* MFG system in which the viscous terms involve Laplace operators with respect to *both* x and v . This is reminiscent of [11, Appendix] (see also [10, Section 4.4]). In this way, $D_v u$ is replaced by $D_v u^\sigma$, which is regular by standard regularity theory for parabolic equations; this implies the regularity of the solution of the Fokker-Planck equation (see [12]). Note also that $D_v u$ may be unbounded; we shall overcome this issue by taking advantage of estimates similar to those in Lemma 3.2. Indeed, these estimates will allow us to apply classical results for the existence and uniqueness of the solution.

Proposition 4.1. *Under assumptions (H), for any $\overline{m} \in C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$, problem (4.1) has a solution m in the sense of Definition 2.1. Moreover $m \in C^{\frac{1}{2}}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N})) \cap L^\infty(0, T; \mathcal{P}_2(\mathbb{R}^{2N}))$.*

We consider the solution (u^σ, m^σ) to the following problem

$$(4.3) \quad \begin{cases} (i) & -\partial_t u - \sigma \Delta_{x,v} u - v \cdot D_x u + \frac{1}{2} |D_v u|^2 - \frac{1}{2} |v|^2 - l(x, v) = F[\overline{m}](x, v), & \text{in } \mathbb{R}^{2N} \times (0, T), \\ (ii) & \partial_t m - \sigma \Delta_{x,v} m - \operatorname{div}_v(m D_v u) - v \cdot D_x m = 0, & \text{in } \mathbb{R}^{2N} \times (0, T), \\ (iii) & m(x, v, 0) = m_0(x, v), \quad u(x, v, T) = G[\overline{m}(T)](x, v), & \text{on } \mathbb{R}^{2N}. \end{cases}$$

Recall that equation (4.3)-(ii) has a standard probabilistic interpretation (see relation (4.8) below). Our aim is to find a solution to problem (4.1) by letting σ tend to 0^+ . To this end, some estimates are needed.

Note that equation (4.3)-(ii) can be written in the compact form

$$(4.4) \quad \partial_t m^\sigma - \sigma \Delta_{x,v} m^\sigma - \operatorname{div}_{x,v}(m^\sigma b^\sigma) = 0, \quad \text{with } b^\sigma := (-v, D_v u^\sigma).$$

We start by establishing the well-posedness of system (4.3) and that the functions u^σ are Lipschitz continuous and semi-concave uniformly in σ .

Lemma 4.1. *Under the same assumptions as in Proposition 4.1, there exists a unique classical solution u^σ to equation (4.3)-(i) with the terminal condition contained in (4.3)-(iii). Moreover, there exists a constant $C > 0$ which depends only on the constants in assumptions (H), in particular it is independent of $\sigma \leq 1$, such that*

- (a) $|u^\sigma(x, v, t)| \leq C(1 + |v|^2),$
- (b) $\|D_x u^\sigma\|_\infty \leq C, \quad |D_v u^\sigma(x, v, t)| \leq C(1 + |v|), \quad |\partial_t u^\sigma(x, v, t)| \leq C(1 + |v|^2),$
- (c) $D_{x,v}^2 u^\sigma \leq C,$

where $D_{x,v}^2 u$ is the Hessian of u with respect to both x and v .

Proof. Following the same arguments of Proposition 3.1 (based on the comparison principle by Da Lio and Ley [15]), one can easily prove the existence of a viscosity solution to equation (4.3)-(i) with terminal condition as in (4.3)-(iii) and satisfying inequality (a). Furthermore, still by the results in [15], this solution is unique among the functions with this growth at infinity. Hence, estimate (a) is proved.

Let us now prove that this viscosity solution u^σ is a classical solution. To this end, let us assume for a moment that u^σ satisfies estimates (b) and (c). We see that u is a viscosity subsolution of

$$\partial_t u - \sigma \Delta_{x,v} u - v \cdot D_x u \leq C(1 + |v|^2).$$

Moreover, from estimate (c), we see that at any point (x, v, t) , either u^σ is twice differentiable with respect to x and v , or there does exist a smooth function that touches u^σ from below. This and estimate (b) imply that u^σ is a viscosity supersolution of

$$\partial_t u - \sigma \Delta_{x,v} u - v \cdot D_x u \geq -C(1 + |v|^2).$$

for some positive constant C . From [18], u is also a distributional subsolution (respectively supersolution) of the same linear inequalities. Therefore, both $\partial_t u^\sigma - \sigma \Delta_{x,v} u^\sigma - v \cdot D_x u^\sigma$ and $-\frac{1}{2}|D_v u^\sigma|^2 + \frac{1}{2}|v|^2 + l(x, v) + F[\overline{m}](x, v)$ are in L_{loc}^∞ . On the other hand, from (b) and (c), Alexandrov's theorem implies that u^σ is twice differentiable with respect to x and v almost everywhere, so the equation

$$\partial_t u^\sigma - \sigma \Delta_{x,v} u^\sigma - v \cdot D_x u^\sigma = -\frac{1}{2}|D_v u^\sigma|^2 + \frac{1}{2}|v|^2 + \ell(x, v, t),$$

(where ℓ and g are defined in (3.1)), holds almost everywhere, and in the sense of distributions since both the left and right hand sides are in L_{loc}^∞ .

Hence classical results on the regularity of weak solutions (including bootstrap) can be applied and yield that u is a classical solution.

Let us now prove the estimates (b) and (c), by using similar arguments to those contained in the proofs of Lemma 3.2. They use a representation formula of u arising from a stochastic optimal control problem (see, for example, [15, 5, 10]).

Let $(\Omega, \mathcal{F}, (\mathcal{F}_t), \mathbb{P})$ be a complete filtered probability space, the filtration (\mathcal{F}_t) supporting a standard $2N$ -dimensional Brownian motion $B_s = (B_{x,s}, B_{v,s})$. Let \mathcal{A}_t be the set of \mathbb{R}^N -valued (\mathcal{F}_t) -progressively measurable processes and let \mathbb{E} be the expectation with respect to the probability measure \mathbb{P} . The unique solution of (4.3)-(i) which satisfies point (a) can be written as:

$$u^\sigma(x, v, t) = \inf_{\alpha \in \mathcal{A}_t} \mathbb{E} \left(\int_t^T \left[\frac{1}{2} |\alpha(s)|^2 + \frac{1}{2} |V(s)|^2 + \ell(X(s), V(s), s) \right] ds + g(X(T), V(T)) \right)$$

where the controlled process $(X(\cdot), V(\cdot))$ satisfies

$$X(t) = x, \quad V(t) = v,$$

almost surely and is governed by the stochastic differential equations

$$(4.5) \quad \begin{cases} dX = V(s)ds + \sqrt{2\sigma}dB_{x,s}, \\ dV = \alpha(s)ds + \sqrt{2\sigma}dB_{v,s}. \end{cases}$$

Thus, almost surely,

$$(4.6) \quad \begin{cases} X(s) &= x + v(s-t) + \int_t^s \int_t^\theta \alpha(\tau) d\tau d\theta + \sqrt{2\sigma} \int_t^s \left(\int_t^\theta dB_{v,\tau} \right) d\theta + \sqrt{2\sigma} \int_t^s dB_{x,\tau}, \\ V(s) &= v + \int_t^s \alpha(\tau) d\tau + \sqrt{2\sigma} \int_t^s dB_{v,\tau}. \end{cases}$$

To prove (b), we can exactly use the same arguments as for Lemma 3.2, replacing the paths $(x(s), v(s))$ and $(y(s), w(s))$ by the processes $(X(s), V(s))$ and $(Y(s), W(s))$, and noting that, from (4.6), we get similar equalities as in (3.8).

Note that, for any σ , we get from (a) that any ϵ -optimal control α^σ for $u^\sigma(x, v, t)$ satisfies

$$\mathbb{E} \left(\int_t^T |\alpha^\sigma(s)|^2 ds \right) \leq C(1 + |v|^2),$$

hence, we get the same estimates as (3.9) and (3.10), namely estimate (b). An analytic proof of (b) is also possible, see [23, Chapter XI].

In order to prove (c), we can follow the same procedure as in the proof of Lemma 3.3, noting that:

- i) equalities (3.14) and (3.18) are still true for the stochastic processes,
- ii) if we fix $s \in [t, T]$, using a Taylor expansion of g as in (3.16), we get

$$\begin{aligned} g(X(s), V(s)) &= g(X_\lambda(s), V_\lambda(s)) + Dg(X_\lambda(s), V_\lambda(s))(X(s) - X_\lambda(s), V(s) - V_\lambda(s)) \\ &\quad + \frac{1}{2}(X(s) - X_\lambda(s), V(s) - V_\lambda(s))D^2g(\xi, \eta)(X(s) - X_\lambda(s), V(s) - V_\lambda(s))^T, \end{aligned}$$

where $\xi = X(s) + \theta_1(X(s) - X_\lambda(s))$, $\eta = V(s) + \theta_2(V(s) - V_\lambda(s))$ for suitable θ_1 and θ_2 in $[0, 1]$. For a similar proof, see [5]. \square

Lemma 4.2. *Under the same assumptions as in Proposition 4.1, there exists a unique classical solution m^σ to equation (4.3)-(ii) with initial condition as in (4.3)-(iii). Moreover, $m^\sigma > 0$.*

Proof. By Lemma 4.1, the problem for m^σ can be written

$$\partial_t m - \sigma \Delta_{x,v} m - b^\sigma \cdot D_{xv} m - (\Delta_v u^\sigma) m = 0, \quad m(0) = m_0,$$

where b^σ has been introduced in (4.4) and from the estimates contained in Lemma 4.1, $|b^\sigma| \leq C(1 + |v|)$ and $\Delta_v u^\sigma \leq C$. Using this and the results contained in [17], we get the existence and uniqueness of a classical solution m^σ of (4.3)-(ii) with initial condition as in (4.3)-(iii). From the assumptions on m_0 and Harnack inequality (see for example [19, Theorem 2.1, p.13]) we get that $m^\sigma(\cdot, t) > 0$ for $t > 0$. \square

Let us now prove some properties of the functions m^σ which will play a crucial role in the proofs of Proposition 4.1 and of Theorem 2.1.

Lemma 4.3. *Under the same assumptions of Proposition 4.1, there exists a constant $K > 0$ which depends only on the constants in assumptions (H) and on m_0 , in particular it is independent of $\sigma \leq 1$, such that:*

1. $\|m^\sigma\|_\infty \leq K$,
2. $\mathbf{d}_1(m^\sigma(t_1), m^\sigma(t_2)) \leq K(t_2 - t_1)^{1/2}, \quad \forall t_1 \leq t_2 \in [0, T]$,
3. $\int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm^\sigma(t)(x, v) \leq K \left(\int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm_0(x, v) + 1 \right), \quad \forall t \in [0, T]$.

Proof. Point 1. In order to prove this L^∞ estimate, we argue as in [11, Theorem 5.1]. We note that

$$\operatorname{div}_v(m^\sigma D_v u^\sigma) = D_v m^\sigma \cdot D_v u^\sigma + m^\sigma(\Delta_v u^\sigma) \leq D_v m^\sigma \cdot D_v u^\sigma + C m^\sigma,$$

because of the semi-concavity of u established in Lemma 4.1 and the positivity of m^σ . Therefore, from assumption (H2), the function m^σ satisfies

$$\partial_t m^\sigma - \sigma \Delta m^\sigma - v D_x m^\sigma - D_v u^\sigma \cdot D_v m^\sigma - C m^\sigma \leq 0, \quad m^\sigma(x, v, 0) \leq C.$$

Then, using $w = Ce^{Ct}$ as a supersolution (recall that C is independent of σ), we obtain that $\|m^\sigma\|_\infty \leq Ce^{CT}$, using the comparison principle proved [15].

To prove Points 2 and 3 as in the proof of [10, Lemma 3.4 and 3.5], it is convenient to introduce the stochastic differential equation

$$(4.7) \quad dY_t = b^\sigma(Y_t, t)dt + \sqrt{2\sigma}dB_t, \quad Y_0 = Z_0,$$

where $Y_t = (X_t, V_t)$, $b^\sigma(x, v, t) = (-v, D_v u^\sigma(x, v, t))$, B_t is a standard $2N$ -dimensional Brownian motion, and $\mathcal{L}(Z_0) = m_0$. By standard arguments, setting

$$(4.8) \quad m^\sigma(t) := \mathcal{L}(Y_t),$$

we know that $m^\sigma(t)$ is absolutely continuous with respect to Lebesgue measure, and that if $m^\sigma(\cdot, \cdot, t)$ is the density of $m^\sigma(t)$, then m^σ is the weak solution to (4.3)-(ii) with $m^\sigma|_{t=0} = m_0$ (from Ito's Theorem, since b^σ has at most linear growth with respect to (x, v) , Proposition 3.6 Chapter 5 [16], p.303, and the book [4]). Here again, we have used the estimate on $|D_v u^\sigma|$ given in Lemma 4.1.

Point 3: Noting that

$$\int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm^\sigma(t)(x, v) = \mathbb{E}(|Y_t|^2),$$

the desired estimate can be obtained by applying Estimate 3.17 of Problem 3.15, p. 306, (the solutions are at p. 389) of [16] with $m = 1$.

Point 2: For $t_2 \geq t_1$, it is well known that

$$\mathbf{d}_1(m^\sigma(t_1), m^\sigma(t_2)) \leq \mathbb{E}(|Y_{t_1} - Y_{t_2}|).$$

Recall also that for a suitable constant C ,

$$|b^\sigma(Y_\tau, \tau)| \leq C(|V_\tau| + 1).$$

The latter two observations imply that

$$\begin{aligned} \mathbb{E}(|Y_{t_1} - Y_{t_2}|) &\leq \mathbb{E} \left(\int_{t_1}^{t_2} |b^\sigma(Y_\tau, \tau)| d\tau + \sqrt{2\sigma} |B_{t_2} - B_{t_1}| \right) \\ &\leq \mathbb{E} \left(C \int_{t_1}^{t_2} (|V_\tau| + 1) d\tau + \sqrt{2\sigma} |B_{t_2} - B_{t_1}| \right) \\ &\leq C \left(\mathbb{E} \left(\int_{t_1}^{t_2} (|V_\tau|^2 + 1) d\tau \right) \right)^{\frac{1}{2}} \sqrt{t_2 - t_1} + \sqrt{2\sigma} \sqrt{t_2 - t_1} \\ &\leq C \left(\mathbb{E} \left(\max_{[t_1, t_2]} |Y_\tau|^2 \right) + 1 \right)^{\frac{1}{2}} (t_2 - t_1) + \sqrt{2\sigma} \sqrt{t_2 - t_1}. \end{aligned}$$

where we have used estimate [16, (3.17) p. 306].

□

Proof of Proposition 4.1. The arguments are similar to those in the proof of [11, Theorem 5.1] (see also [10, Theorem 4.20]). Lemma 4.1 imply that possibly after the extraction of a subsequence, u^σ locally uniformly converges to some function u , which is Lipschitz continuous with respect to x , locally Lipschitz continuous with respect to v , and $Du^\sigma \rightarrow Du$ a.e. (because of the semi-concavity estimate of Lemma 4.1 and [9, Theorem 3.3.3]). By standard stability result for viscosity solutions, the function u is a viscosity solution of (3.4).

On the other hand, the function m^σ satisfies the estimates stated in Lemma 4.3:

1. from point 3, $m^\sigma(t)$ is bounded in $\mathcal{P}_2(\mathbb{R}^{2N})$ uniformly in $\sigma \in [0, 1]$ and $t \in [0, T]$
2. from points 2 and 3, m^σ is bounded in $C^{1/2}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ uniformly with respect to $\sigma \in [0, 1]$.

Recalling that the subsets of $\mathcal{P}_1(\mathbb{R}^{2N})$ whose elements have uniformly bounded second moment are relatively compact in $\mathcal{P}_1(\mathbb{R}^{2N})$, see for example [10, Lemma 5.7], we can apply Ascoli-Arzelà theorem: we may extract a sequence (still called σ for simplicity) such that $\sigma \rightarrow 0^+$ and m^σ converges to some $m \in C^{1/2}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ in the $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ topology. Moreover, from point 1 in Lemma 4.3 and Banach-Alaoglu theorem, m belongs to $L_{\text{loc}}^\infty((0, T) \times \mathbb{R}^{2N})$ and the sequence m^σ converges to m in $L_{\text{loc}}^\infty((0, T) \times \mathbb{R}^{2N})$ -weak-*. Therefore, by passing to the limit, we immediately obtain that $m|_{t=0} = m_0$, $\|m\|_\infty \leq K$

and that $\mathbf{d}_1(m(t_1), m(t_2)) \leq K(t_2 - t_1)^{1/2}$, $\forall t_1 \leq t_2 \in [0, T]$.
Let us prove that for all $t \in [0, T]$,

$$(4.9) \quad \int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm(t)(x, v) \leq K \left(\int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm_0(x, v) + 1 \right).$$

For that, let us consider the increasing sequence of functions defined on \mathbb{R}_+ : $\phi_n(\rho) = 1 \wedge ((n+1 - \rho) \vee 0)$. We know from point 3 in Lemma 4.3, that for all $t \in [0, T]$,

$$(4.10) \quad \int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) \phi_n(|x|^2 + |v|^2) m^\sigma(x, v, t) dx dv \leq K \left(\int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm_0(x, v) + 1 \right).$$

For a fixed n , we can pass to the limit in (4.10) thanks to the $L_{\text{loc}}^\infty((0, T) \times \mathbb{R}^{2N})$ -weak-* convergence established above. We obtain:

$$(4.11) \quad \int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) \phi_n(|x|^2 + |v|^2) m(x, v, t) dx dv \leq K \left(\int_{\mathbb{R}^{2N}} (|x|^2 + |v|^2) dm_0(x, v) + 1 \right).$$

We then pass to the limit as $n \rightarrow +\infty$ thanks to Beppo-Levi monotone convergence theorem, and obtain (4.9).

Finally, m^σ is a solution to (4.3)-(ii),

$$\int_0^T \int_{\mathbb{R}^{2N}} m^\sigma (-\partial_t \psi - \sigma \Delta \psi + D_v \psi \cdot D_v u^\sigma - v D_x \psi) dx dv dt = 0$$

for any $\psi \in C_0^\infty((0, T) \times \mathbb{R}^{2N})$. Letting $\sigma \rightarrow 0^+$, we conclude from the L_{loc}^∞ -weak-* convergence of m^σ and the convergence $Du^\sigma \rightarrow Du$ a.e. that the function m solves (4.1) in the sense of Definition 2.1. \square

Remark 4.2. *Note that we have just proven that all the estimates on u^σ contained in Lemma 4.1 hold for u . These estimates have also been obtained directly in the proof of Lemma 3.2. Similarly, all the estimates on m^σ contained in Lemma 4.3 hold for m .*

4.2 Uniqueness of the solution

We now deal with uniqueness for (4.1).

Proposition 4.2. *Under assumptions (H), the function m found in Proposition 4.1 is the unique solution to problem (4.1) in the sense of Definition 2.1 such that*

$$m \in C^{\frac{1}{2}}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N})) \cap L^\infty(0, T; \mathcal{P}_2(\mathbb{R}^{2N})).$$

Moreover, m satisfies:

$$(4.12) \quad \int_{\mathbb{R}^{2N}} \phi(x, v) m(x, v, t) dx dv = \int_{\mathbb{R}^{2N}} \phi(\overline{\gamma}_{x,v}(t)) m_0(x, v) dx dv, \quad \forall \phi \in C_0(\mathbb{R}^{2N}), \forall t \in [0, T],$$

where, for a.e. $(x, v) \in \mathbb{R}^{2N}$, $\overline{\gamma}_{x,v}$ is the solution to (2.2).

Proof of Proposition 4.2. The proof is similar to that of [12, Proposition A.1], which relies on Ambrosio's superposition principle. Let Γ_T denote the set of continuous curves in \mathbb{R}^{2N} , namely $\Gamma_T = C([0, T]; \mathbb{R}^{2N})$. For any $t \in [0, T]$, we introduce the evaluation map:

$e_t : \Gamma_T \rightarrow \mathbb{R}^{2N}$, $e_t(\gamma) := \gamma(t)$. Hereafter, when we write “for a.e.” without specifying the measure, we intend “with respect to the Lebesgue measure”.

Let $m \in C^{1/2}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N})) \cap L^\infty(0, T; \mathcal{P}_2(\mathbb{R}^{2N}))$ be a solution of problem (4.1) in the sense of Definition 2.1. Recall the notation $b(x, v, t) = (-v, D_v u(x, v, t))$. The estimate (8.1.20) in chapter 8 of [1] is fulfilled: indeed,

$$\begin{aligned} & \int_0^T \int_{\mathbb{R}^{2N}} |b(x, v, t)|^2 dm(t)(x, v) \leq \\ & C \int_0^T \int_{\mathbb{R}^{2N}} |v|^2 dm(t)(x, v) + C \int_0^T \int_{\mathbb{R}^{2N}} |D_v u(x, v, t)|^2 dm(t)(x, v) \leq C, \end{aligned}$$

where the last inequality comes from the estimates on $D_v u$ (see Remark 4.2), Points 3 of Lemma 4.3 and the fact that $m_t := m(t, \cdot)$ is a probability measure. Therefore, the assumptions of Ambrosio’s superposition principle are fulfilled (see [1, Theorem 8.2.1] and also [1, pag. 182]). The latter and the disintegration theorem (see [1, Theorem 5.3.1]) entail that there exist a probability measure η on $\mathbb{R}^{2N} \times \Gamma_T$ and for m_0 -almost every $(x, v) \in \mathbb{R}^{2N}$, a probability measure on $\eta_{x,v}$ on Γ_T , such that

- i) $e_t \# \eta = m_t$, i.e., for every bounded and continuous real valued function ψ defined on \mathbb{R}^{2N} , for every $t \in [0, T]$,

$$\int_{\mathbb{R}^{2N}} \psi(x, v) dm_t(x, v) = \int_{\mathbb{R}^{2N} \times \Gamma_T} \psi(\zeta(t)) d\eta(x, v, \zeta).$$

In particular, $e_0 \# \eta = m_0$.

- ii)

$$\eta = \int_{\mathbb{R}^{2N}} \eta_{x,v} dm_0(x, v),$$

i.e. for every bounded Borel function $f : \mathbb{R}^{2N} \times \Gamma_T \rightarrow \mathbb{R}$,

$$\int_{\mathbb{R}^{2N} \times \Gamma_T} f(x, v, \zeta) d\eta(x, v, \zeta) = \int_{\mathbb{R}^{2N}} \left(\int_{\Gamma_T} f(x, v, \zeta) d\eta_{x,v}(\zeta) \right) dm_0(x, v).$$

- iii) For m_0 -almost every $(x, v) \in \mathbb{R}^{2N}$, the support of $\eta_{x,v}$ is contained in the set

$$(4.13) \quad \left\{ \zeta \in \text{AC}([0, T]; \mathbb{R}^{2N}) : \zeta(t) = (\xi(t), \eta(t)) : \begin{cases} \xi(0) = x, \eta(0) = v, \\ \xi'(t) = \eta(t), \\ \eta'(t) = -D_v u(\xi(t), \eta(t), t). \end{cases} \right\}.$$

Recall that in the present case, m_0 is absolutely continuous (from assumption (H4)); hence, since for all $t \in [0, T]$, $u(\cdot, \cdot, t)$ is Lipschitz continuous, the optimal synthesis in Lemma 3.5 ensures that for a.e. $(x, v) \in \mathbb{R}^{2N}$, (3.38)-(3.40) (with $t = 0$ in the present context) has a unique solution $\bar{\gamma}_{x,v}$, because it is the optimal trajectory for the control problem mentioned in Remark 4.1. Therefore, for a.e. $(x, v) \in \mathbb{R}^{2N}$, the set in (4.13) is a singleton, or in equivalent manner, $\eta_{x,v}$ coincides with $\delta_{\bar{\gamma}_{x,v}}$. In conclusion, for any

function $\phi \in C_0(\mathbb{R}^{2N})$,

$$\begin{aligned}
\int_{\mathbb{R}^{2N}} \psi(x, v) m(x, v, t) dx dv &= \int_{\Gamma_T} \psi(e_t(\zeta)) d\eta(x, v, \zeta) \\
&= \int_{\mathbb{R}^{2N}} \left(\int_{\Gamma_T} \psi(e_t(\zeta)) d\eta_{x,v}(\zeta) \right) dm_0(x, v) \\
&= \int_{\mathbb{R}^{2N}} \psi(e_t(\bar{\gamma}_{x,v})) dm_0(x, v) \\
&= \int_{\mathbb{R}^{2N}} \psi(\bar{\gamma}_{x,v}(t)) m_0(x, v) dx dv.
\end{aligned}$$

This shows that m is uniquely defined as the image of m_0 by the flow of (2.2). \square

Theorem 4.1. Existence of m comes from Proposition 4.1, uniqueness and the representation formula come from Proposition 4.2. \square

5 Proof of the main Theorem

Proof of Theorem 2.1. For point 1, we argue as in the proof of [10, Theorem 4.1]. Consider the set $\mathcal{C} := \{m \in C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N})) \mid m(0) = m_0\}$ endowed with the norm of $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ and observe that it is a closed and convex subset of $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$. We also introduce a map \mathcal{T} as follows: to any $m \in \mathcal{C}$, we associate the solution u to problem (3.4) with $\bar{m} = m$ and to this u we associate the solution $\mu =: \mathcal{T}(m)$ to problem (4.1) which, by Proposition 4.1 belongs to \mathcal{C} . Hence, \mathcal{T} maps \mathcal{C} into itself. We claim that the map \mathcal{T} has the following properties:

- (a) \mathcal{T} is a continuous map with respect to the norm of $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$
- (b) \mathcal{T} is a compact map.

Assume for the moment that these properties are true. In this case, Schauder fixed point Theorem ensures the existence of a fixed point for \mathcal{T} , namely a solution to system (2.1). Therefore it remains to prove properties (a) and (b).

Let us now prove (a). Let $(m_n)_n$ be a sequence in \mathcal{C} such that $m_n \rightarrow m$ in the $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ topology. We want to prove that $\mathcal{T}(m_n) \rightarrow \mathcal{T}(m)$ in $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$. We observe that hypothesis (H3) ensures that the functions $(x, v, t) \mapsto F[m_n(t)](x, v)$ and $(x, v) \mapsto G[m_n(T)](x, v)$ converge locally uniformly to the map $(x, v, t) \mapsto F[m(t)](x, v)$ and respectively $(x, v) \mapsto G[m(T)](x, v)$. Moreover, Lemma 3.2 entails that the solutions u_n to problem (3.4) with $\bar{m} = m_n$ are locally uniformly bounded and locally uniformly Lipschitz continuous. Therefore, by standard stability results for viscosity solutions, the sequence $(u_n)_n$ converges locally uniformly to viscosity the solution u to problem (3.4) with $\bar{m} = m$. Moreover, from Lemma 3.3, the functions u_n are uniformly semi-concave; hence, by [9, Theorem 3.3.3], Du_n converge a.e. to Du .

By Proposition 4.1 and Remark 4.2, the function $\mathcal{T}(m_n)$ verifies the bounds in Lemma 4.3 with a constant K independent of n . Hence, the sequence $(\mathcal{T}(m_n))_n$ is uniformly bounded in $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ (by Lemma 4.3-(3) and Remark 4.2, and because the subsets of $\mathcal{P}_1(\mathbb{R}^{2N})$ whose elements have uniformly bounded second moment are relatively compact in $\mathcal{P}_1(\mathbb{R}^{2N})$, and uniformly Hölder continuous in time with values in $\mathcal{P}_1(\mathbb{R}^{2N})$ (by Lemma 4.3-(2) and Remark 4.2). Therefore, by Ascoli-Arzelà and

Banach-Alaoglu theorems, there exists a subsequence $(\mathcal{T}(m_{n_k}))_k$ which converges to some $\mu \in C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ in the $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ -topology and in the $L_{\text{loc}}^\infty((0, T) \times \mathbb{R}^{2N})$ -weak-* topology. As in Remark 4.2, μ verifies the bounds in Lemma 4.3 and $\mu(0) = m_0$.

Observe that $\mathcal{T}(m_{n_k})$ solves problem (4.1) with u replaced by u_{n_k} ,

$$\int_0^T \int_{\mathbb{R}^{2N}} \mathcal{T}(m_{n_k}) (-\partial_t \psi + D_v \psi \cdot D_v u_{n_k} - v \cdot D_x \psi) dx dv dt = 0,$$

for any $\psi \in C_0^\infty((0, T) \times \mathbb{R}^{2N})$. Passing to the limit as $k \rightarrow \infty$, we get that μ is a solution to (4.1). By the uniqueness result established in Proposition 4.2, we deduce that $\mu = \mathcal{T}(m)$, and that the whole sequence $(\mathcal{T}(m_n))_n$ converges to $\mathcal{T}(m)$.

Let us now prove (b); since \mathcal{C} is closed, it is enough to prove that $\mathcal{T}(\mathcal{C})$ is a precompact subset of $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$. Let $(\mu_n)_n$ be a sequence in $\mathcal{T}(\mathcal{C})$ with $\mu_n = \mathcal{T}(m_n)$ for some $m_n \in \mathcal{C}$; we wish to prove that, possibly for a subsequence, μ_n converges to some μ in the $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ -topology as $n \rightarrow \infty$.

By Remark 4.2, the functions $\mathcal{T}(m_n)$ satisfy the estimates in Lemma 4.3 with the same constant K . Since the subsets of $\mathcal{P}_1(\mathbb{R}^{2N})$ whose elements have uniformly bounded second moment are relatively compact in $\mathcal{P}_1(\mathbb{R}^{2N})$, Lemma 4.3-(3) ensures that the sequence $(\mathcal{T}(m_n))_n$ is uniformly bounded. Moreover, Lemma 4.3-(2) yields that the sequence $(\mathcal{T}(m_n))_n$ is uniformly bounded in $C^{1/2}([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ and $L^\infty(0, T; \mathcal{P}_2(\mathbb{R}^{2N}))$. By arguing as in the proof of Proposition 4.1, we obtain that, possibly for a subsequence (still denoted by $\mathcal{T}(m_n)$), $\mathcal{T}(m_n)$ converges to some μ in the $C([0, T]; \mathcal{P}_1(\mathbb{R}^{2N}))$ -topology.

2. Theorem 4.1 ensures that, if (u, m) is a solution of (2.1), for any function $\psi \in C_0(\mathbb{R}^{2N})$,

$$(5.1) \quad \int_{\mathbb{R}^{2N}} \psi(x, v) m(x, v, t) dx dv = \int_{\mathbb{R}^{2N}} \psi(\bar{\gamma}_{x,v}(t)) m_0(x, v) dx dv$$

where $\bar{\gamma}_{x,v}$ is the solution of (2.2) (uniquely defined for a.e. $(x, v) \in \mathbb{R}^{2N}$). \square

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