# A PONTRYAGIN MAXIMUM PRINCIPLE FOR AGENT-BASED MODELS WITH CONVEX STATE SPACE 

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#### Abstract

We derive a first order optimality condition for a class of agent-based systems, as well as for their mean-field counterpart. A relevant difficulty of our analysis is that the state equation is formulated on possibly infinite-dimensional convex subsets of Banach spaces, as required by some problems in multi-population dynamics. Due to the lack of a linear structure and of local compactness, the usual tools of needle variations and linearisation procedures used to derive Pontryagin type conditions have to be generalised to the setting at hand. This is done by considering suitable notions of differentials and by a careful inspection of the underlying functional structures.


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## 1. Introduction

Presentation of the problem. Large systems of interacting agents, the so-called multiagent systems, have enjoyed growing attention from several mathematical communities in recent years. Indeed, they may encompass a broad class of applications, including, e.g.,

[^0]the modelling of autonomous vehicle ensembles [20], swarms and flocking structures in the animal kingdom [27, coordinated animal motion [8], opinion dynamics on networks [9], pedestrian flows [26], cell aggregation and motility [23, 30], cooperative robots [25], and influence of key investors in the stock market [11, Introduction]. The modeling of such systems is usually inspired from Newtonian laws of motion and is usually described by large dynamical systems based on pairwise forces accounting for repulsion/attraction, alignment, self-propulsion/friction in biological, social, or economical interactions. On the one hand, such discrete formulations often give rise to intractable problems because of their very large dimensionality. On the other hand, for many applications one is not interested in the exact pattern of the single agent, which may be indistinguishable from another one, but rather on the collective behavior of the system. This point of view is also relevant for control-theoretic purposes [2,3,33], where a policy maker aims at formulating control laws which are as generic and global as possible in order to steer a given system towards a desired goal.

To overcome the so-called curse of dimensionality one is naturally lead to consider infinite-dimensional approximations of the class of multi-agent systems at hand. They usually take the form of a continuity equation for the time-dependent distribution $\mu$ of the agent on the state space (a probability measure), driven by a possibly nonlocal velocity field accounting for particle interactions and external control. This is often referred to as mean-field approximation of large particle systems (see [32] as well as the survey $[24]$ ). These equations are studied in the space of probability measures endowed with the Wasserstein metrics of optimal transport (see Section 4 below). The wellposedness of mean-field optimal control problems, as well as their compliance with particle formulations, has been addressed in several recent contributions (see, e.g., [22, 28, 29]) hinging on a variational approach.

From an applicative standpoint, the identification of mean-field optimal controls relies on the implementation of suitable optimality conditions rather then on the direct minimisation of a cost functional. Such optimality conditions have been studied in the recent literature $[12,14,15,18,21,34$, resulting in a generalisation to Wasserstein spaces of the celebrated Pontryagin Maximum Principle (PMP). The works 12, 15, 18] are of particular interest for our purposes. Indeed, in such papers the conditions of Pontryagin type are recovered by adapting to the geometric setting of Wasserstein spaces some classical tools coming from optimal control techniques in Euclidean setting, such as the use of needle variations and linearisation of the evolution constraints. The optimal control has to maximise a suitable Hamiltonian field $\mathcal{H}(\nu, u)$, where $\nu$ is a probability measure on $\mathbb{R}^{2 d}$ accounting for both state and costate of the system, evolving through the Wasserstein Hamiltonian flow driven by $\mathcal{H}$.

In recent applications, a refined modelling comes into play, where the state space of an agent is no more the Euclidean, but is instead represented by a convex subset $C$ of a possibly infinite-dimensional Banach space $E$. This point of view has been introduced in the context of spatially inhomogeneous evolutionary games in [6]. In this case, the state variable is a pair $(x, \lambda) \in C:=\mathbb{R}^{d} \times \mathcal{P}(K)$, where $\mathcal{P}(K)$ is the space of probability measure on a compact metric space of pure strategies $K$. The agents' state is therefore described by their position $x$ and by their mixed strategy $\lambda$. A related point of view has been considered in other contributions [4,5,31] to describe the evolution of systems where each agent has a time-evolving degree of influence decoded by $\lambda$. This is useful in many applications where agents could belong to different populations, such as leaders and followers. In particular, a mean-field selective optimal control problem has been analysed in [1], where the action of the policy maker aimed at steering the whole system
through a parsimonious intervention on possibly temporary leaders, i.e., via a control law depending on the (time evolving) variable $\lambda$. Notice that also the particle formulation of this systems requires additional care, as the involved operators must satisfy a suitable set of assumptions (see, e.g., [31]), in order to guarantee existence and uniqueness of solutions for differential equations constrained to live in a given convex set (see [19]).
Our results. The goal of our paper is to extend the formulation of a PMP to the convex constrained setting. The main technical tools are (i) a notion of differentiability for functionals defined on a convex set $C$, (ii) a weak notion of local differentiability for functionals defined on $\mathcal{P}(C)$, and (iii) the introduction of a functional setting suitable for linearisation and for the derivation of the adjoint equation.

We start with the controlled particle systems, which we recast in the form

$$
\begin{equation*}
\min _{u \in \mathcal{U}}\left\{\int_{0}^{T} L(t, c(t), u(t)) \mathrm{d} t\right\} \tag{1.1}
\end{equation*}
$$

subject to

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} c(t)=A(t, c(t), u(t)) \quad \text { in }(0, T]  \tag{1.2}\\
c(0)=c_{0} \in C
\end{array}\right.
$$

In (1.1), $\mathcal{U}$ is a weakly compact space of admissible controls, while the assumptions on the operator $A$ are listed in (HAode) assure well-posedness of the state equation (1.2) constrained to $C$ and involve the notion of $C$-differentiability with respect to $c \in C$ (see also [6]). In Theorem 3.4 we prove that the an optimal pair trajectory-control $(\mathfrak{c}(t), \mathfrak{p}(t))$ solves

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t}\binom{\mathfrak{c}(t)}{\mathfrak{p}(t)}=J \mathrm{D}_{(\mathfrak{c}(t), \mathfrak{p}(t))} \mathcal{H}(t, \mathfrak{c}(t), \mathfrak{p}(t), \mathfrak{u}(t)) \quad \text { in }[0, T), \\
\mathfrak{c}(0)=c_{0} \in C, \\
\mathfrak{p}(T)=\mathbf{0} \in E_{C}^{*},
\end{array}\right.
$$

on $C \times E_{C}^{*}$, where $E_{C}:=\overline{\mathbb{R}(C-C)}$. Furthermore, the optimal control maximises the Hamiltonian $\mathcal{H}(t, c(t), p(t), u(t))=\langle p(t), A(t, c(t), u(t))\rangle-L(t, c(t), u(t))$, where the duality product is taken in $E_{C}^{*}, E_{C}$. The proof is based on needle variations and on passing to the limit in the blow-up $\left(c_{\varepsilon}-c\right) / \varepsilon$, where $c_{\varepsilon}$ denotes the solution of the perturbed problem. Here some care is required, since a linear structure is not available and can only be partially replaced by convexity. In our setting tangent directions can be approximated only at first order by the perturbed trajectories, so that a careful estimate of the behavior of rest terms is in order. We refer to the proof of Theorems 3.3 and 3.4 for full details.

In Section 4 we turn our attention to the mean-field control problem

$$
\begin{equation*}
\min _{u \in \mathcal{U}}\left\{\int_{0}^{T} L(t, \mu(t), u(t)) \mathrm{d} t\right\}, \tag{1.3}
\end{equation*}
$$

subject to

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \mu(t)+\operatorname{div}(A(t, \mu(t), \cdot, u(t, \cdot)) \mu(t))=0 \quad \text { in }(0, T]  \tag{1.4}\\
\mu(0)=\mu^{0} \in \mathcal{P}_{c}(C)
\end{array}\right.
$$

The set of admissible controls $\mathcal{U}$ has a closed-loop structure with respect to the state variable, which means that an admissible control $u$ may depend on both $t$ and $c$, and satisfies similar assumptions than the one considered in 18 to ensure compactness of controls. The operator $A$ is assumed to satisfy two sets of assumptions, namely (HA1) and (HA2). The first is similar to (HAode), while the latter is crucial for writing a linearised equation in form of a nonlocal Cauchy problem 4.19). This requires, besides
the already mentioned $C$-differentiability of $A$, a suitable notion of Wasserstein gradient with respect to $\mu$ introduced in Definition 4.2, which is tailored for the functional setting at hand. The core of our arguments is contained in Proposition 4.9, where we perform the linearisation of the flow associated to (1.4). Here is where our analysis significantly departs from 15, 18]. Indeed the states space lacks of linearity and of local compactness; furthermore we cannot in principle assume that the dependence of the velocity field $A$ on the closed-loop control $u$ is linear.

Also from a technical point of view the hardest part is contained in the proof of Proposition 4.9. A crucial technical tool is a chain-rule with respect to the notion of Wasserstein differentiability, which we prove in Proposition 4.3, allowing for a first-order expansion of suitable compositions of flow maps close to the identity. This is heavily exploited in Theorem 4.7 to analyse the behavior of the perturbed solutions to (1.4) arising from needle variations of the optimal control, as well as in the derivation of the optimality conditions for the Hamiltonian associated to the problem, which is explicitly given in (4.12). Our main result is contained in Theorem 4.10, where a PMP for (1.3)-(1.4) is derived, coupling a maximality condition for the control $u$ with a Wasserstein Hamiltonian flow in the space $C \times E_{C}^{*}$. We remark that the derivation of such principle requires some assumptions on the involved (possibly infinite-dimensional) functional spaces. In particular, we need to assume $E$ to be separable and $E_{C}$ to be reflexive. The compliance of our setting with a model example inspired by [4] is shortly discussed in Section 5 .

We conclude with the obvious remark that in the special case $C=\mathbb{R}^{d}$ one retrieves the known results obtained for instance in [18] for closed-loop controls and in 15] for open-loop controls.

Outlook. This work provides a general functional setting for identifying optimality conditions for agent-based systems with convex state space, where the evolution of the system is described by the continuity equation (1.4). Possible generalisations of the present framework may take into account evolutions driven by differential inclusions or stochastic effects, which have recently been considered in the Euclidean setting (see [16] and [13], respectively). As pointed out in [17] for systems with unconstrained state space, optimality conditions for the finite particle control problem (1.1)-(1.2) in combination with convexity assumptions on the Hamiltonian $\mathcal{H}$ may lead to Lipschitz regularity of optimal controls for the mean-field problem (1.3)-(1.4). Such topics will be the focus of future investigation.

## 2. Preliminaries and notation

In a Banach space $\left(E,\|\cdot\|_{E}\right)$, let $C$ be a closed convex subset of $E$. For $R>0$ we define $B_{R}:=\left\{e \in E:\|e\|_{E} \leq R\right\}$ and $B_{R, C}=B_{R} \cap C$. For any vectorial subspace $W$ of $E$ we denote by $\mathcal{L}(W ; E)$ the Banach space of linear bounded operator from $W$ to $E$ and by $W^{*}$ the dual space of $W$, i.e. $W^{*}=\mathcal{L}(W ; \mathbb{R})$. Recall that, by density, each operator in $\mathcal{L}(W ; E)$ uniquely extends to an operator in $\mathcal{L}(\bar{W} ; E)$. By $\langle\cdot, \cdot\rangle$ we denote the duality pairing between a vectorial space and its dual. We denote by $E_{C}$ the closed vector subspace $\overline{\mathbb{R}}(C-C) \subset E$ and, for $c \in C$, by $E_{c}$ the convex cone of directions

$$
E_{c}:=\mathbb{R}_{+}(C-c) \subset E_{C} \subset E .
$$

We have the following definition of $C$-differential.

Definition 2.1. Let $\left(F,\|\cdot\|_{F}\right)$ be a normed space and let $f: C \rightarrow F$. We say that $f$ is $C$-differentiable in $c \in C$ if there exists a linear operator $M \in \mathcal{L}\left(E_{C} ; F\right)$ such that

$$
\lim _{c^{\prime} \rightarrow c} \frac{f\left(c^{\prime}\right)-f(c)-M\left[c^{\prime}-c\right]}{\left\|c^{\prime}-c\right\|_{E}}=0 .
$$

We will denote the $C$-differential in $c \in C$ by $D_{c} f$. Notice that such map is a priori well-defined in $E_{c}$ and can be uniquely extended by linearity and density to a map in $E_{C}$.

Given $(X, d)$ a separable Radon metric space, we denote by $\mathcal{P}(X)$ the family of all Borel probability measures on $X$. For $p \geq 1$ we further consider

$$
\mathcal{P}_{p}(X):=\left\{\mu \in \mathcal{P}(X): \int_{X} d(x, \bar{x})^{p} \mathrm{~d} \mu(x)<+\infty \text { for some } \bar{x} \in X\right\}
$$

and $\mathcal{P}_{c}(X)$ the subset of $\mathcal{P}(X)$ of measures with compact support in $X$ recalling that the support is the closed set

$$
\operatorname{supp}(\mu)=\{x \in X: \mu(V)>0 \text { for each neighborhood } V \text { of } x\} .
$$

If $X$ is contained in some Banach space $Y$, we define the $p$ momentum of $\mu \in \mathcal{P}(X)$ as

$$
m_{p}(\mu):=\left(\int_{X}\|x\|_{Y}^{p} \mathrm{~d} \mu(x)\right)^{\frac{1}{p}} \quad \text { for } p \geq 1
$$

Let $X_{1}$ and $X_{2}$ be separable Radon metric spaces, we define for every $\mu^{1} \in \mathcal{P}\left(X_{1}\right)$ and $\mu^{2} \in \mathcal{P}\left(X_{2}\right)$ the transport plans with marginals $\mu^{1}$ and $\mu^{2}$

$$
\Gamma\left(\mu^{1}, \mu^{2}\right):=\left\{\gamma \in \mathcal{P}\left(X_{1} \times X_{2}\right): \pi_{\#}^{i} \gamma=\mu^{i} \text { for } i=1,2\right\},
$$

where $\pi^{i}: X_{1} \times X_{2} \rightarrow X_{i}$ is the projection on $X_{i}$ and $\pi_{\#}^{i} \boldsymbol{\gamma} \in \mathcal{P}\left(X_{i}\right)$ is the pushforward of $\boldsymbol{\gamma}$ through $\pi^{i}$. Note that $\Gamma\left(\mu^{1}, \mu^{2}\right)$ is a not empty and compact subset of $\mathcal{P}\left(X_{1} \times X_{2}\right)$ (see Remark 5.2.3 of [7]). We define the $p$-Wasserstein distance between two probability measure $\mu^{1}$ and $\mu^{2}$ in $\widehat{\mathcal{P}}_{p}(X)$ by

$$
W_{p}^{p}\left(\mu^{1}, \mu^{2}\right)=\min \left\{\int_{X \times X} d\left(x_{1}, x_{2}\right)^{p} \mathrm{~d} \boldsymbol{\gamma}\left(x_{1}, x_{2}\right): \boldsymbol{\gamma} \in \Gamma\left(\mu^{1}, \mu^{2}\right)\right\},
$$

it follows from Proposition 7.15 of $[7]$ that $\mathcal{P}_{p}(X)$ endowed with the $p$-Wasserstein distance is a separable metric space which is complete if $X$ is complete.

## 3. Pontryagin maximum principle for optimal control problems in CONVEX SPACES

Let the following assumptions hold:

- $\left(E,\|\cdot\|_{E}\right)$ is a Banach space;
- $C \subset E$ is a closed and convex subset of $E$;
- $\left(U, d_{U}\right)$ is a compact metric space of admissible controls;
- $\mathcal{U}:=\{u:[0, T] \rightarrow U$ such that $u(\cdot)$ is Lebesgue measurable $\}$.

We introduce here the main set of assumptions for the operator $A$ appearing in (1.2).
(HAode) Let $A(t, \cdot, \cdot): C \times U \rightarrow E, t \in[0, T]$, be a family of operators satisfying the following properties:
(i) for every $R>0$ there exists a constant $L_{R} \geq 0$ such that for every $c_{1}, c_{2} \in$ $B_{R, C}, t \in[0, T]$ and $u \in U$

$$
\begin{equation*}
\left\|A\left(t, c_{1}, u\right)-A\left(t, c_{2}, u\right)\right\|_{E} \leq L_{R}\left\|c_{1}-c_{2}\right\|_{E} \tag{3.1}
\end{equation*}
$$

(ii) there exists $M>0$ such that for every $c \in C, t \in[0, T]$ and $u \in U$ it holds

$$
\|A(t, c, u)\|_{E} \leq M\left(1+\|c\|_{E}\right)
$$

(iii) for every $c \in C$ and $u \in U$ the map $t \mapsto A(t, c, u)$ belongs to $L^{1}([0, T] ; E)$;
(iv) for every $t \in[0, T]$ and $c \in C$ the map $u \mapsto A(t, c, u)$ belongs to $C^{0}(U ; E)$;
(v) for every $R>0$ there exists $\theta>0$ such that for every $t \in[0, T]$ and $u \in U$

$$
c \in C,\|c\|_{E} \leq R \Rightarrow c+\theta A(t, c, u) \in C
$$

(vi) for every $c \in C, t \in[0, T]$ and $u \in U, A$ is $C$-differentiable at $c \in C$, i.e. there exists $B_{(t, c, u)} \in \mathcal{L}\left(E_{C} ; E_{C}\right)$ such that

$$
\begin{equation*}
\lim _{C \ni c^{\prime} \rightarrow c} \frac{\left\|A\left(t, c^{\prime}, u\right)-A(t, c, u)-B_{(t, c, u)}\left[c^{\prime}-c\right]\right\|_{E}}{\left\|c^{\prime}-c\right\|_{E}}=0 \tag{3.2}
\end{equation*}
$$

Note that $(v)$ implies $A(t, \cdot, u): C \rightarrow E_{c}$ and that by (3.2) we have for every $e \in E_{c}$

$$
\begin{equation*}
B_{(t, c, u)}[e]=\lim _{h \rightarrow 0^{+}} \frac{A(t, c+h e, u)-A(t, c, u)}{h} . \tag{3.3}
\end{equation*}
$$

For every $t \in[0, T]$ and $u \in U$, we denote the $C$-differential of $A$ in $c \in C$ by

$$
\mathrm{D}_{c} A:=\left\{B_{(t, c, u)} \in \mathcal{L}\left(E_{C} ; E_{C}\right): B_{(t, c, u)} \text { fullfils (3.2) }\right\}
$$

Let us give some results related to the notion of $C$-differentiability given in (vi).
Lemma 3.1. Let $A$ satisfy (vi). Then the map $\mathrm{D} A: C \rightarrow \mathcal{L}\left(E_{C} ; E_{C}\right)$ with $\mathrm{D} A(c):=\mathrm{D}_{c} A$ is single valued. Moreover, if (i) holds, $\mathrm{D}_{c(t)} A:=B_{(t, c(t), u(t))} \in L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)$ for every $c \in A C([0, T] ; C)$ and $u \in \mathcal{U}$.

Proof. First we prove that $\mathrm{D} A$ is single valued. If $e \in E_{c}$, then, by (3.3), $\mathrm{D}_{c} A[e]$ is a singleton. Moreover, if $v \in \mathbb{R}(C-C)$, then $v=\alpha\left(c_{1}-c_{2}\right)=\alpha\left(c_{1}-c\right)-\alpha\left(c_{2}-c\right)$ for some $c_{1}, c_{2} \in C$ and $\alpha \in \mathbb{R}$. Hence, by linearity, we deduce

$$
\mathrm{D}_{c} A[v]=\alpha \mathrm{D}_{c} A\left[c_{1}-c\right]-\alpha \mathrm{D}_{c} A\left[c_{2}-c\right]
$$

where $c_{1}-c, c_{2}-c \in E_{c}$. Therefore $\mathrm{D}_{c} A[v]$ too is singleton and, by density, $\mathrm{D} A$ is single valued.
Now let $c(t) \in A C([0, T] ; C)$ and $u \in \mathcal{U}$. By continuity we have that $c(t) \in B_{R, C}$ for some $R>0$. Since for $h$ small enough we have $c(t)+h e \in C \cap B_{2 R}$ with $e \in E_{c(t)}$ (recall that $C$ is convex), it follows from (3.1) and (3.3) that $\left\|\mathrm{D}_{c(t)} A[e]\right\|_{E} \leq l\|e\|_{E}$ for some $l \geq 0$ depending on $R$. This property can be extended for $v \in \mathbb{R}(C-C)$. Indeed, using the decomposition $v=\alpha\left(c_{1}-c_{2}\right)=\alpha\left(c_{1}-c(t)\right)-\alpha\left(c_{2}-c(t)\right)$, we have

$$
\begin{aligned}
& \left\|\mathrm{D}_{c(t)} A[v]\right\|_{E}=|\alpha|\left\|\mathrm{D}_{c(t)} A\left[c_{1}-c(t)\right]-\mathrm{D}_{c(t)} A\left[c_{2}-c(t)\right]\right\|_{E} \\
& \stackrel{3.3}{=}|\alpha|\left\|\lim _{h \rightarrow 0^{+}} \frac{A\left(t, c(t)+h\left(c_{1}-c(t)\right), u\right)-A\left(t, c(t)+h\left(c_{2}-c(t)\right), u\right)}{h}\right\|_{E} \\
& =|\alpha|\left\|\lim _{h \rightarrow 0^{+}} \frac{A\left(t,(1-h) c(t)+h c_{1}, u\right)-A\left(t,(1-h) c(t)+h c_{2}, u\right)}{h}\right\|_{E}
\end{aligned}
$$

$$
\stackrel{(i)}{\leq} l|\alpha|\left\|c_{1}-c_{2}\right\|_{E}=l\|v\|_{E} .
$$

Hence, using again a density argument and recalling that $E_{C}=\overline{\mathbb{R}(C-C)}$, we deduce that $\left\|\mathrm{D}_{c(t)} A\right\|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \leq l$ in $[0, T]$.

From now on we use $\mathrm{D}_{c} A$ instead of $B_{(t, c, u)}$. We define the adjoint operator $\mathrm{D}_{c}^{*} A \in$ $\mathcal{L}\left(E_{C}^{*}, E_{C}^{*}\right)$ of $\mathrm{D}_{c} A$ as

$$
\left\langle\mathrm{D}_{c}^{*} A[p], v\right\rangle:=\left\langle p, \mathrm{D}_{c} A[v]\right\rangle \quad \forall p \in E_{C}^{*}, v \in E_{C} .
$$

Remark 3.2. By its definition

$$
\left\|\mathrm{D}_{c}^{*} A\right\|_{\mathcal{L}\left(E_{C}^{*}, E_{C}^{*}\right)}=\left\|\mathrm{D}_{c} A\right\|_{\mathcal{L}\left(E_{C} ; E_{C}\right)}
$$

Then it follows from the second statement of Lemma 3.1 that

$$
\begin{equation*}
\mathrm{D}_{c(t)}^{*} A \in L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C}^{*}, E_{C}^{*}\right)\right), \tag{3.4}
\end{equation*}
$$

for every $c \in A C([0, T] ; C)$ and $u \in \mathcal{U}$.
(H $\varphi$ ode) Let $\varphi: C \rightarrow \mathbb{R}$ be $C$-differentiable for every $c \in C$ with $C$-differential $\mathrm{D}_{c} \varphi \in E_{C}^{*}$ defined as in (vi), i.e.

$$
\lim _{C \ni c^{\prime} \rightarrow c} \frac{\left|\varphi\left(c^{\prime}\right)-\varphi(c)-\left\langle\mathrm{D}_{c} \varphi, c^{\prime}-c\right\rangle\right|}{\left\|c^{\prime}-c\right\|_{E}}=0 .
$$

Under the sets of assumptions (HAode) and (H $\varphi$ ode), we aim at finding necessary conditions for optimal solutions $(\mathfrak{c}(t), \mathfrak{u}(t)) \in A C([0, T] ; C) \times \mathcal{U}$ of the following Mayer problem with free terminal point:

$$
\begin{equation*}
\min _{u \in \mathcal{U}} \varphi(c(T, u)) \tag{3.5}
\end{equation*}
$$

subject to

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} c(t)=A(t, c(t), u(t)) \quad \text { in }(0, T]  \tag{3.6}\\
c(0)=c_{0} \in C
\end{array}\right.
$$

Notice that from the assumptions on $A$ and from Theorem I. 4 of [19] we deduce that (3.6) admits a unique weak solution $c \in A C([0, T ; C)$ for every control $u$ and every initial condition $c_{0} \in C$.

We have the following result:
Theorem 3.3. Consider the optimal control problem (3.5)-(3.6), under the assumptions (HAode) and (H $\varphi$ ode). Let $\mathfrak{u} \in \mathcal{U}$ be an admissible control whose corresponding trajectory $\mathfrak{c} \in A C([0, T] ; C)$ is optimal. Let $p:[0, T] \rightarrow E_{C}^{*}$ the weak solution of the adjoint equation

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} p(t)=-\mathrm{D}_{\mathbf{c}(t)}^{*} A[p(t)] \quad \text { in }[0, T) \\
p(T)=-\mathrm{D}_{\mathbf{c}(T)} \varphi
\end{array}\right.
$$

Then the maximality condition

$$
\langle p(t), A(t, \mathfrak{c}(t), \mathfrak{u}(t))\rangle=\max _{\omega \in U}\{\langle p(t), A(t, \mathfrak{c}(t), \omega)\rangle\}
$$

holds for almost everywhere $t$ in $[0, T]$.

Proof. Fix any time $\tau \in(0, T]$ and any admissible control value $\omega \in U$. For $\varepsilon>0$ sufficiently small, we consider the following function called needle variation:

$$
u_{\varepsilon}(t)= \begin{cases}\omega & \text { if } t \in[\tau-\varepsilon, \tau] \\ \mathfrak{u}(t) & \text { otherwise }\end{cases}
$$

It follows from the assumptions on $A$ and from Theorem I. 4 of [19 that

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} c(t)=A\left(t, c(t), u_{\varepsilon}(t)\right) \quad \text { in }(0, T]  \tag{3.7}\\
c(0)=c_{0} \in C
\end{array}\right.
$$

has a unique weak solution $c_{\varepsilon}(t) \in C$ for a.e. $t \in[0, T]$. By its definition $c_{\varepsilon}(t)=\mathfrak{c}(t)$ for a.e. $t \in[0, \tau-\varepsilon]$. Now we prove that $c_{\varepsilon}(t)$ converges to $\mathfrak{c}(t)$ uniformly in $[0, T]$. By assumptions on $A$ we deduce that $A(t, \mathfrak{c}(t), \mathfrak{u}(t)), A\left(t, \mathfrak{c}(t), u_{\varepsilon}(t)\right)$ and $A\left(t, c_{\varepsilon}(t), u_{\varepsilon}(t)\right)$ are integrable on $[0, T]$. Moreover $R>0$ depending only on $c_{0}$ and $T$ exists such that $\|\mathfrak{c}(t)\|_{E}+\left\|c_{\varepsilon}(t)\right\|_{E} \leq R$ in $[0, T]$. Hence, we have

$$
\begin{array}{rll}
\left\|\mathfrak{c}(t)-c_{\varepsilon}(t)\right\|_{E} \stackrel{\sqrt[3.61]{ }, \sqrt{3.7}}{\leq} & \int_{0}^{T}\left\|A(s, \mathfrak{c}(s), \mathfrak{u}(s))-A\left(s, \mathfrak{c}(s), u_{\varepsilon}(s)\right)\right\|_{E} \mathrm{~d} s \\
& +\int_{0}^{T}\left\|A\left(s, \mathfrak{c}(s), u_{\varepsilon}(s)\right)-A\left(s, c_{\varepsilon}(s), u_{\varepsilon}(s)\right)\right\|_{E} \mathrm{~d} s \\
\stackrel{\sqrt{3.11}}{\leq} & \int_{0}^{T}\left\|A(s, \mathfrak{c}(s), \mathfrak{u}(s))-A\left(s, \mathfrak{c}(s), u_{\varepsilon}(s)\right)\right\|_{E} \mathrm{~d} s \\
& & L_{R} \int_{0}^{T}\left\|\mathfrak{c}(s)-c_{\varepsilon}(s)\right\|_{E} \mathrm{~d} s . \tag{3.8}
\end{array}
$$

Since $u_{\varepsilon}$ converges to $\mathfrak{u}$ in $L^{1}([0, T] ; U)$ as $\varepsilon$ tends to $0^{+}$, by (iv) we obtain that $A\left(t, \mathfrak{c}(t), u_{\varepsilon}(t)\right)$ converges to $A(t, \mathfrak{c}(t), \mathfrak{u}(t))$ almost everywhere in $[0, T]$ as $\varepsilon$ tends to $0^{+}$. Using (ii) we have that $\left\|A(s, \mathfrak{c}(s), \mathfrak{u}(s))-A\left(s, \mathfrak{c}(s), u_{\varepsilon}(s)\right)\right\|_{E} \leq 2 M(1+R)$ for some $M>0$. Then applying the Lebesgue theorem we obtain that $A\left(t, \mathfrak{c}(t), u_{\varepsilon}(t)\right)$ converges to $A(t, \mathfrak{c}(t), \mathfrak{u}(t))$ in $L^{1}([0, T] ; E)$ as $\varepsilon$ tends to $0^{+}$. This, together with (3.8), implies

$$
\left\|\mathfrak{c}(t)-c_{\varepsilon}(t)\right\|_{E} \leq \delta_{\varepsilon}+L_{R} \int_{0}^{T}\left\|\mathfrak{c}(s)-c_{\varepsilon}(s)\right\|_{E} \mathrm{~d} s
$$

where $\delta_{\varepsilon} \rightarrow 0$ as $\varepsilon \rightarrow 0^{+}$. Applying the integral formulation of the Grönwall inequality we get

$$
\left\|\mathfrak{c}(t)-c_{\varepsilon}(t)\right\|_{E} \leq \delta_{\varepsilon} \mathrm{e}^{L_{R} T}
$$

whence

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0^{+}} c_{\varepsilon}(t)=\mathfrak{c}(t) \quad \text { uniformly in }[0, T] . \tag{3.9}
\end{equation*}
$$

Note that, since $A(t, \mathfrak{c}(t), \omega)$ and $A(t, \mathfrak{c}(t), \mathfrak{u}(t))$ are integrable in $[0, T]$, then a.e. time $\tau$ in $[0, T]$ is one of their Lebesgue points. From now on $\tau$ is a time with such property. Notice that such set of Lebesgue points can be made independent of $\omega \in U$, since $U$ is a compact metric space. It follows that a positive $\delta_{\varepsilon}$ tending to 0 as $\varepsilon$ tends to 0 exists such that

$$
\begin{aligned}
& \left\|\frac{c_{\varepsilon}(\tau)-\mathfrak{c}(\tau)}{\varepsilon}-(A(\tau, \mathfrak{c}(\tau), \omega)-A(\tau, \mathfrak{c}(\tau), \mathfrak{u}(\tau)))\right\|_{E} \\
& \leq \frac{1}{\varepsilon} \int_{\tau-\varepsilon}^{\tau}\left\|A\left(s, c_{\varepsilon}(s), \omega\right)-A(s, \mathfrak{c}(s), \omega)\right\|_{E} \mathrm{~d} s+\left\|\frac{1}{\varepsilon} \int_{\tau-\varepsilon}^{\tau} A(s, \mathfrak{c}(s), \omega) \mathrm{d} s-A(\tau, \mathfrak{c}(\tau), \omega)\right\|_{E}
\end{aligned}
$$

$$
\begin{aligned}
& +\left\|A(\tau, \mathfrak{c}(\tau), \mathfrak{u}(\tau))-\frac{1}{\varepsilon} \int_{\tau-\varepsilon}^{\tau} A(s, \mathfrak{c}(s), \mathfrak{u}(s)) \mathrm{d} s\right\|_{E} \\
& \stackrel{\sqrt{3.11}}{\leq} L_{R}\left\|\mathfrak{c}-c_{\varepsilon}\right\|_{L^{\infty}([0, T] ; E)}+\delta_{\varepsilon} .
\end{aligned}
$$

Therefore, by (3.9), we obtain

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0^{+}} \frac{c_{\varepsilon}(\tau)-\mathfrak{c}(\tau)}{\varepsilon}=A(\tau, \mathfrak{c}(\tau), \omega)-A(\tau, \mathfrak{c}(\tau), \mathfrak{u}(\tau)) \tag{3.10}
\end{equation*}
$$

whence we get $A(\tau, \mathfrak{c}(\tau), \omega)-A(\tau, \mathfrak{c}(\tau), \mathfrak{u}(\tau)) \in \overline{E_{\mathfrak{c}(\tau)}}$.
Now we consider the following equation a priori defined in $E_{C}$

$$
\begin{cases}\frac{\mathrm{d}}{\mathrm{~d} t} v(t)=\mathrm{D}_{\mathfrak{c}(t)} A[v(t)] & \text { in }(\tau, T]  \tag{3.11}\\ v(\tau)=A(\tau, \mathfrak{c}(\tau), \omega)-A(\tau, \mathfrak{c}(\tau), \mathfrak{u}(\tau))\end{cases}
$$

By assumption (vi) and Lemma 3.1, using Theorem I. 4 of [19] we infer the existence of a unique weak solution $v \in A C\left([\tau, T] ; E_{C}\right)$ such that $\frac{\mathrm{d}}{\mathrm{d} t} v \in L^{1}\left([\tau, T] ; E_{C}\right)$. We want to prove that

$$
\begin{equation*}
v(t)=\lim _{\varepsilon \rightarrow 0^{+}} \frac{c_{\varepsilon}(t)-\mathfrak{c}(t)}{\varepsilon} \quad \text { uniformly in }[\tau, T], \tag{3.12}
\end{equation*}
$$

which in particular implies that $v(t) \in \overline{E_{\mathfrak{c}(t)}}$ a.e. $t \in[\tau, T]$. Let us start by noting that $c_{\varepsilon}(t)$ and $\mathfrak{c}(t)$ in $[\tau, T]$ solve the same equation (i.e. (3.6) with $u(t)$ replaced by $\mathfrak{u}(t)$ ) with $c_{\varepsilon}(\tau)$ and $\mathfrak{c}(\tau)$ respectively as initial datum. We define

$$
v_{\varepsilon}(t):=\frac{c_{\varepsilon}(t)-\mathfrak{c}(t)}{\varepsilon} \quad \text { for } t \in[\tau, T] .
$$

Hence by its definition $v_{\varepsilon}(t) \in E_{\mathbf{c}(t)}$ for every $t \in[\tau, T]$. Thus (3.12) is equivalent to

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0^{+}}\left\|v_{\varepsilon}(t)-v(t)\right\|_{L^{\infty}([\tau, T] ; E)}=0 \tag{3.13}
\end{equation*}
$$

Moreover by (3.10) $v_{\varepsilon}(\tau)$ converges to $v(\tau)$ in $E$. Now, recalling that $\left\|c_{\varepsilon}(t)\right\|_{E}+\|\mathfrak{c}(t)\|_{E} \leq$ $R$ in $[0, T]$ and denoting by $L$ the constant given by (3.1) (we omit the subscript R ), by continuity with respect to the initial data (see Theorem 1 of [1]) we obtain for $t \in[\tau, T]$

$$
\begin{equation*}
\left\|v_{\varepsilon}(t)\right\|_{E} \leq \mathrm{e}^{L T}\left\|v_{\varepsilon}(\tau)\right\|_{E} \stackrel{\sqrt{3.100}}{\leq} 2 \mathrm{e}^{L T}\|v(\tau)\|_{E} . \tag{3.14}
\end{equation*}
$$

Now we note that by its definition $v_{\varepsilon}(t)$ satisfies

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} v_{\varepsilon}(t)=\frac{1}{\varepsilon}\left(A\left(t, c_{\varepsilon}(t), \mathfrak{u}(t)\right)-A(t, \mathfrak{c}(t), \mathfrak{u}(t))\right)=\mathrm{D}_{\mathfrak{c}(t)} A\left[v_{\varepsilon}(t)\right]+r_{\varepsilon}(t) \tag{3.15}
\end{equation*}
$$

where

$$
r_{\varepsilon}(t):=\frac{1}{\varepsilon}\left(A\left(t, c_{\varepsilon}(t), \mathfrak{u}(t)\right)-A(t, \mathfrak{c}(t), \mathfrak{u}(t))\right)-\mathrm{D}_{\mathfrak{c}(t)} A\left[v_{\varepsilon}(t)\right] .
$$

It follows from (3.9), (3.2) and (3.14) that for every $t \in(\tau, T]$
$\lim _{\varepsilon \rightarrow 0}\left\|r_{\varepsilon}(t)\right\|_{E}=\lim _{\varepsilon \rightarrow 0} \frac{\left\|A\left(t, c_{\varepsilon}(t), \mathfrak{u}(t)\right)-A(t, \mathfrak{c}(t), \mathfrak{u}(t))-\mathrm{D}_{\mathfrak{c}(t)} A\left[c_{\varepsilon}(t)-\mathfrak{c}(t)\right]\right\|_{E}}{\left\|c_{\varepsilon}(t)-\mathfrak{c}(t)\right\|_{E}}\left\|v_{\varepsilon}(t)\right\|_{E}=0$.
Moreover by Lemma 3.1 we have

$$
\begin{aligned}
\left\|r_{\varepsilon}(t)\right\|_{E} & \stackrel{\sqrt{3.15]}}{\leq} \\
& \frac{1}{\varepsilon}\left\|A\left(t, c_{\varepsilon}(t), \mathfrak{u}(t)\right)-A(t, \mathfrak{c}(t), \mathfrak{u}(t))\right\|_{E}+\left\|\mathrm{D}_{\mathfrak{c}(t)} A\left[v_{\varepsilon}(t)\right]\right\|_{E} \\
& \stackrel{\sqrt{3.1]}}{\leq} \\
& \left(L+\left\|\mathrm{D}_{\mathfrak{c}(t)} A\right\|_{L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)}\right)\left\|v_{\varepsilon}(t)\right\|_{E}
\end{aligned}
$$

Hence using (3.14) and applying the Lebesgue theorem we obtain that $r_{\varepsilon} \rightarrow 0$ in $L^{1}([\tau, T] ; E)$ as $\varepsilon \rightarrow 0^{+}$. Therefore, recalling that $v$ solves (3.11) and using again Lemma 3.1, we obtain

$$
\begin{align*}
\left\|v_{\varepsilon}(t)-v(t)\right\|_{E} & \leq\left\|v_{\varepsilon}(\tau)-v(\tau)\right\|_{E}+\left\|\int_{\tau}^{t} \mathrm{D}_{\mathbf{c}(s)} A\left[v_{\varepsilon}(s)-v(s)\right] \mathrm{d} s\right\|_{E}+\left\|\int_{\tau}^{t} r_{\varepsilon}(s) \mathrm{d} s\right\|_{E} \\
& \stackrel{\text { 3.10) }}{\leq} \delta_{\varepsilon}+\left\|\mathrm{D}_{\mathbf{c}(t)} A\right\|_{L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)} \int_{\tau}^{t}\left\|v_{\varepsilon}(s)-v(s)\right\|_{E} \mathrm{~d} s \tag{3.16}
\end{align*}
$$

where $\delta_{\varepsilon}:=\left\|v_{\varepsilon}(\tau)-v(\tau)\right\|_{E}+\left\|r_{\varepsilon}\right\|_{L^{1}([\tau, T] ; E)} \rightarrow 0$ as $\varepsilon \rightarrow 0^{+}$. Since $\left\|v_{\varepsilon}(t)-v(t)\right\|_{E}$ is a continuous function from $[\tau, T]$ to $\mathbb{R}$, by (3.16), applying the integral form of the Grönwall inequality and recalling that $\left\|\mathrm{D}_{\mathfrak{c}(t)} A\right\|_{L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)} \leq L$ it follows that

$$
\left\|v_{\varepsilon}(t)-v(t)\right\|_{E} \leq \delta_{\varepsilon} \mathrm{e}^{L T} .
$$

This implies (3.13) and, consequentially, (3.12).
Since the control $\mathfrak{u}(t)$ is optimal in $\mathcal{U}$, we deduce for every $\varepsilon>0$

$$
\begin{equation*}
\varphi(\mathfrak{c}(T)) \leq \varphi\left(c_{\varepsilon}(T)\right) \tag{3.17}
\end{equation*}
$$

Moreover by (3.12) we have $c_{\varepsilon}(T)=\mathfrak{c}(T)+\varepsilon v(T)+o(\varepsilon)$. Hence, recalling that $\varphi$ is $C$-differentiable and that

$$
c_{\varepsilon}(T)-\mathfrak{c}(T)=\varepsilon v(T)+o(\varepsilon) \in E_{\mathfrak{c}(t)}
$$

we have

$$
\begin{align*}
\varphi\left(c_{\varepsilon}(T)\right)-\varphi(\mathfrak{c}(T)) & =\varphi(\mathfrak{c}(T)+\varepsilon v(T)+o(\varepsilon))-\varphi(\mathfrak{c}(T)) \\
& =\left\langle\mathrm{D}_{\mathfrak{c}(T)} \varphi, \varepsilon v(T)+o(\varepsilon)\right\rangle+o(\varepsilon) . \tag{3.18}
\end{align*}
$$

Therefore

$$
\begin{equation*}
0 \stackrel{\sqrt{3.17}}{\leq} \lim _{\varepsilon \rightarrow 0^{+}} \frac{\varphi\left(c_{\varepsilon}(T)\right)-\varphi(\mathfrak{c}(T))}{\varepsilon} \stackrel{\sqrt{3.18}}{=}\left\langle\mathrm{D}_{\mathfrak{c}(T)} \varphi, v(T)\right\rangle \tag{3.19}
\end{equation*}
$$

Now we consider the adjoint equation of (3.11) which transports $p(T)=\mathrm{D}_{\mathfrak{c}(T)} \varphi \in E_{C}^{*}$ backward in time, i.e.

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} p(t)=-\mathrm{D}_{\mathbf{c}(t)}^{*} A[p(t)] \quad \text { in }[\tau, T),  \tag{3.20}\\
p(T)=-\mathrm{D}_{\mathbf{c}(T)} \varphi \in E_{C}^{*}
\end{array}\right.
$$

It follows from (3.4) and Theorem I. 4 of [19] that $p(t) \in A C\left([\tau, T] ; E_{C}^{*}\right)$ is the unique weak solution of $(3.20)$. Then $\frac{\mathrm{d}}{\mathrm{d} t} p \in L^{1}\left([\tau, T] ; E_{C}^{*}\right)$. The density of $C^{1}\left([\tau, T] ; E_{C}\right)$ and $C^{1}\left([\tau, T] ; E_{C}^{*}\right)$ in $W^{1,1}\left([\tau, T] ; E_{C}\right)$ and $W^{1,1}\left([\tau, T] ; E_{C}^{*}\right)$ respectively implies that the function $\langle p(\cdot), v(\cdot)\rangle:[\tau, T] \rightarrow \mathbb{R}$ is weak differentiable and its weak derivative is

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\langle p(t), v(t)\rangle=\left\langle\frac{\mathrm{d}}{\mathrm{~d} t} p(t), v(t)\right\rangle+\left\langle p(t), \frac{\mathrm{d}}{\mathrm{~d} t} v(t)\right\rangle .
$$

Therefore, using equations (3.11) and (3.20) and the definition of adjoint operator, we obtain

$$
\frac{\mathrm{d}}{\mathrm{~d} t}\langle p(t), v(t)\rangle=0 \quad \text { a.e. } t \in[\tau, T] .
$$

This implies that $\langle p(t), v(t)\rangle$ is constant in $[\tau, T]$, more precisely
for a.e. $\tau \in[0, T]$. Since $\omega \in U$ is arbitrary, we have that

$$
\langle p(\tau), A(\tau, \mathfrak{c}(\tau), \mathfrak{u}(\tau))\rangle=\max _{\omega \in U}\{\langle p(\tau), A(\tau, \mathfrak{c}(\tau), \omega)\rangle\} \quad \text { for a.e. } \tau \in[0, T] .
$$

Now we focus on the Bolza problem with running cost $L$, i.e.

$$
\begin{equation*}
\min _{u \in \mathcal{U}}\left\{\int_{0}^{T} L(t, c(t), u(t)) \mathrm{d} t\right\} \tag{3.21}
\end{equation*}
$$

subject to (3.6). We assume:
(HLode) Let $L:[0, T] \times C \times U \rightarrow \mathbb{R}$ be a map satisfying:
(a) for every $R>0$ there exists a constant $L_{R} \geq 0$ such that for every $c_{1}, c_{2} \in$ $B_{R, C}, t \in[0, T]$ and $u \in U$

$$
\left|L\left(t, c_{1}, u\right)-L\left(t, c_{2}, u\right)\right| \leq L_{R}\left\|c_{1}-c_{2}\right\|_{E}
$$

(b) for every $c \in C$ and $u \in U$ the map $t \mapsto L(t, c, u)$ belongs to $L^{1}([0, T] ; \mathbb{R})$;
(c) for every $t \in[0, T]$ and $c \in C$ the map $u \mapsto L(t, c, u)$ belongs to $C^{0}(U ; \mathbb{R})$;
(d) for every $c \in C, t \in[0, T]$ and $u \in U, L$ is $C$-differentiable at $c \in C$ with $C$-differential $\mathrm{D}_{c} L \in \mathcal{L}\left(E_{C} ; \mathbb{R}\right)$.

Under these assumptions on $L$ with the same argument contained in Lemma 3.1 we have that $\mathrm{D}_{c} L$ is single valued in $E_{C}$ and $\mathrm{D}_{c(t)} L \in L^{\infty}\left([0, T] ; E_{C}^{*}\right)$ for every $c \in A C([0, T] ; C)$ and $u \in \mathcal{U}$.

We can rewrite the Bolza problem in Mayer form introducing the auxiliary variable

$$
c_{a u}(t)=\int_{0}^{t} L(s, c(s), u(s)) \mathrm{d} s
$$

with $c_{a u}(0)=0$. We introduce the following notations:

$$
\begin{align*}
& \underline{c}:=\binom{c}{c_{a u}} \in C \times \mathbb{R}, \quad \underline{v}:=\binom{v}{v_{a u}} \in E_{C} \times \mathbb{R}, \quad \underline{p}:=\binom{p}{p_{a u}} \in E_{C}^{*} \times \mathbb{R} \\
& \underline{A}(t, \underline{c}, u):=\binom{A(t, c, u)}{L(t, c, u)} . \tag{3.22}
\end{align*}
$$

Since $A$ and $L$ are $C$-differentiable, it follows that $\underline{A}:[0, T] \times C \times \mathbb{R} \times U \rightarrow E_{C}$ is $C$-differentiable and

$$
\mathrm{D}_{\underline{\underline{A}}} \underline{A}:=\left(\begin{array}{cc}
\mathrm{D}_{c} A & 0  \tag{3.23}\\
\mathrm{D}_{c} L & 0
\end{array}\right) \in \mathcal{L}\left(E_{C} \times \mathbb{R} ; E_{C} \times \mathbb{R}\right) .
$$

Finally, we denote by $\mathrm{D}_{\underline{\underline{c}}}^{*} \underline{A}$ the adjoint operator of $\mathrm{D}_{\underline{\underline{c}}} \underline{A}$. Note that for every $c \in$ $A C([0, T] ; C)$ and $u \in \mathcal{U}$, since $\mathrm{D}_{\underline{c}(t)} \underline{A} \in L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C} \times \mathbb{R} ; E_{C} \times \mathbb{R}\right)\right)$, with the same argument of Remark 3.2 we obtain $\mathrm{D}_{\underline{c}(t)}^{*} \underline{A} \in L^{\infty}\left([0, T] ; \mathcal{L}\left(E_{C}^{*} \times \mathbb{R} ; E_{C}^{*} \times \mathbb{R}\right)\right.$. Recalling that $A(t, c, u) \in E_{c} \subset E_{C}$, we define $\mathcal{H}:[0, T] \times C \times E_{C}^{*} \times U \rightarrow \mathbb{R}$ as

$$
\begin{equation*}
\mathcal{H}(t, c, p, u)=\langle p, A(t, c, u)\rangle-L(t, c, u) . \tag{3.24}
\end{equation*}
$$

By its definition and since $A$ and $L$ are $C$-differentiable, it follows that $\mathcal{H}$ is $C$-differentiable at $c \in C$ and its $C$-differential $\mathrm{D}_{c} \mathcal{H} \in E_{C}^{*}$ has the form

$$
\mathrm{D}_{c} \mathcal{H}[v]=\left\langle p, \mathrm{D}_{c} A[v]\right\rangle-D_{c} L[v]=\left\langle\mathrm{D}_{c}^{*} A[p], v\right\rangle-D_{c} L[v] \quad \forall v \in E_{C} .
$$

Moreover $\mathcal{H}$ is Fréchet-differentiable with respect to $p$ and $\mathrm{D}_{p} \mathcal{H}=A(t, c, u) \in E_{C} \subset E_{C}^{* *}$. It follows that $\mathcal{H}$ is differentiable in $C \times E_{C}^{*}$ and

$$
\mathrm{D}_{(c, p)} \mathcal{H}(t, c, p, u)=\left(\begin{array}{c}
\mathrm{D}_{c}^{*} A[p]-D_{c} L  \tag{3.25}\\
A(t, c, u) \\
11
\end{array}\right) \in E_{C}^{*} \times E_{C} .
$$

In this way the Bolza problem (3.21)-( $\sqrt{3.6)}$ is equivalent to the Mayer problem

$$
\begin{equation*}
\min _{u \in \mathcal{U}} c_{a u}(T, u) \tag{3.26}
\end{equation*}
$$

subject to

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \underline{c}(t)=\underline{A}(t, \underline{c}(t), u(t)) \quad \text { in }(0, T]  \tag{3.27}\\
\underline{c}(0)=\binom{c_{0}}{0} \in C \times \mathbb{R}
\end{array}\right.
$$

Thanks to the assumptions on $A$ and $L$ we can apply Theorem 3.3 to problem (3.26)-(3.27) obtaining the following result.
Theorem 3.4. Under the assumptions (HAode) and (HLode), let $\mathfrak{u} \in \mathcal{U}$ be an admissible control whose corresponding trajectory $\mathfrak{c} \in A C([0, T] ; C)$ is optimal for problem (3.21)-(3.6). Then there exists $\mathfrak{p} \in A C\left([0, T] ; E_{C}^{*}\right)$ such that $(\mathfrak{c}, \mathfrak{p})$ solves in distributional sense

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t}\binom{\mathfrak{c}(t)}{\mathfrak{p}(t)}=J \mathrm{D}_{(\mathfrak{c}(t), \mathfrak{p}(t))} \mathcal{H}(t, \mathfrak{c}(t), \mathfrak{p}(t), \mathfrak{u}(t)) \quad \text { in }[0, T), \\
\mathfrak{c}(0)=c_{0} \in C \\
\mathfrak{p}(T)=\mathbf{0} \in E_{C}^{*}
\end{array}\right.
$$

where $\mathcal{H}$ is defined by (3.24) and $J: E_{C}^{*} \times E_{C} \ni(T, e) \mapsto(e,-T) \in E_{C} \times E_{C}^{*}$. Moreover the maximality condition

$$
\mathcal{H}(t, \mathfrak{c}(t), \mathfrak{p}(t), \mathfrak{u}(t))=\max _{\omega \in U} \mathcal{H}(t, \mathfrak{c}(t), \mathfrak{p}(t), \omega)
$$

holds for almost everywhere $t$ in $[0, T]$.
Proof. Let us start by noting that, if $\mathfrak{u} \in \mathcal{U}$ be an admissible control whose corresponding trajectory $\mathfrak{c} \in A C([0, T] ; C)$ is optimal for problem (3.21)-(3.6), then by (3.22) and (3.21) we have that $\underline{\mathfrak{c}}=\left(\mathfrak{c}, \mathfrak{c}_{a u}\right) \in A C([0, T] ; C \times \mathbb{R})$ is optimal for problem (3.26)-(3.27). Since $\varphi(\underline{c}(T, u))=c_{a u}(T, u)$, we deduce that $\mathrm{D}_{\underline{\mathfrak{c}}(T)} \varphi=\binom{\mathbf{0}}{1} \in E_{C}^{*} \times \mathbb{R}$. Hence it follows from Theorem 3.3 that

$$
\begin{align*}
& \langle\underline{\mathfrak{p}}(t), \underline{A}(t, \underline{\mathfrak{c}}(t), \mathfrak{u}(t))\rangle \stackrel{\sqrt{3.22]}}{=}\langle\mathfrak{p}(t), A(t, \mathfrak{c}(t), \mathfrak{u}(t))\rangle+\mathfrak{p}_{a u}(t) L(t, \mathfrak{c}(t), \mathfrak{u}(t)) \\
& =\max _{\omega \in U}\left\{\langle\mathfrak{p}(t), A(t, \mathfrak{c}(t), \omega)\rangle+\mathfrak{p}_{a u}(t) L(t, \mathfrak{c}(t), \omega)\right\} \tag{3.28}
\end{align*}
$$

holds for almost everywhere $t$ in $[0, T]$, where $\mathfrak{p}$ is the weak solution of

$$
\begin{cases}\left.\frac{\mathrm{d}}{\mathrm{~d} t}(t)=-\mathrm{D}_{\underline{\mathfrak{c}}(t)}^{*} \underline{A} \underline{p}(t)\right] & \text { in }[0, T)  \tag{3.29}\\ \underline{p}(T)=-\mathrm{D}_{\underline{\mathfrak{c}}(T)} \varphi=\binom{0}{-1} \in E_{C}^{*} \times \mathbb{R} & \end{cases}
$$

Then for $(v, 0) \in E_{C} \times \mathbb{R}$ it holds

$$
\begin{align*}
\left\langle\frac{\mathrm{d}}{\mathrm{~d} t} \mathfrak{p}(t), v\right\rangle & \left.=\left\langle\frac{\mathrm{d}}{\mathrm{~d} t} \mathfrak{p}(t),\binom{v}{0}\right\rangle \stackrel{\sqrt{3.29 \mathfrak{p}}}{=}\left\langle-\mathrm{D}_{\underline{\mathfrak{c}}(t)}^{*} \underline{A} \underline{\mathfrak{p}}(t)\right],\binom{v}{0}\right\rangle \\
& =-\left\langle\underline{\mathfrak{p}}(t), \mathrm{D}_{\underline{\mathfrak{c}}(t)} \underline{A}\left[\begin{array}{c}
v \\
0
\end{array}\right]\right\rangle \stackrel{(3.22), \sqrt{3.23}}{=}-\left\langle\binom{\mathfrak{p}(t)}{\mathfrak{p}_{a u}(t)},\binom{\mathrm{D}_{\mathfrak{c}(t)} A[v]}{\mathrm{D}_{\mathfrak{c}(t)} L[v]}\right\rangle \\
& =-\left\langle\mathfrak{p}(t), \mathrm{D}_{\mathfrak{c}(t)} A[v]\right\rangle-\mathfrak{p}_{a u}(t) \mathrm{D}_{\mathfrak{c}(t)} L[v] \\
& =-\left\langle\mathrm{D}_{\mathfrak{c}(t)}^{*} A[\mathfrak{p}(t)], v\right\rangle-\mathfrak{p}_{a u}(t) \mathrm{D}_{\mathfrak{c}(t)} L[v] . \tag{3.30}
\end{align*}
$$

Therefore, recalling (3.24) and (3.25), in order to have the result it remains to prove that $\mathfrak{p}_{a u}(t)=-1$ almost everywhere $t$ in $[0, T]$ or, equivalently, since $\mathfrak{p}_{a u}(T)=-1$, to prove that $\frac{\mathrm{d}}{\mathrm{d} t} \mathfrak{p}_{a u}(t)=0$ almost everywhere $t$ in $[0, T]$. Using (3.29) and the definition of adjoint operator we obtain that

$$
\begin{aligned}
\frac{\mathrm{d}}{\mathrm{~d} t} \mathfrak{p}_{a u}(t) & \left.=\left\langle\frac{\mathrm{d}}{\mathrm{~d} t} \mathfrak{p}(t),\binom{\mathbf{0}}{1}\right\rangle=\left\langle-\mathrm{D}_{\underline{\mathfrak{c}}(t)}^{*} \underline{A} \underline{\mathfrak{p}}(t)\right],\binom{\mathbf{0}}{1}\right\rangle \\
& =-\left\langle\underline{\mathfrak{p}}(t), \mathrm{D}_{\underline{\mathfrak{c}}(t)} \underline{A}\left[\begin{array}{l}
\mathbf{0} \\
1
\end{array}\right]\right\rangle \stackrel{\sqrt{3.23}}{=} 0
\end{aligned}
$$

Hence, by substituting $\mathfrak{p}_{a u}(t)=-1$ in (3.28) and (3.30), the result follows from (3.24) and (3.25).
3.1. A generalization to finite particle control problems. In this subsection we rewrite the results obtained in Theorems 3.3 and 3.4 for a control problem involving $N$ particles, $N \in \mathbb{N}$. Before stating the assumptions on the velocity field and the cost functions, we introduce some notation. We denote by $\mathbf{c}$ a generic element $\left(c_{1}, \ldots, c_{N}\right) \in$ $C^{N}$ and by $\mathbf{c}(t)$ a generic element $\left(c_{1}(t), \ldots, c_{N}(t)\right) \in L^{1}\left([0, T] ; C^{N}\right)$, by $\mathbf{p}$ an element $\left(p_{1}, \ldots, p_{N}\right) \in\left(E_{C}^{*}\right)^{N}$ and by $\mathbf{p}(t)$ an element $\left(p_{1}(t), \ldots, p_{N}(t)\right) \in L^{1}\left([0, T] ;\left(E_{C}^{*}\right)^{N}\right)$, by $\mathbf{u}$ an element $\left(u_{1}, \ldots, u_{N}\right) \in U^{N}$ and by $\mathbf{u}(t)$ an element $\left(u_{1}(t), \ldots, u_{N}(t)\right) \in \mathcal{U}^{N}$ respectively.

We make the following assumptions on the velocity field $\mathbf{A}$ and on the cost functional $\mathbf{L}$, which are the finite-particle counterparts of (HAode) and of (HLode), respectively.
(HAsym) Let A: $[0, T] \times C^{N} \times C \times U \rightarrow E$ be an operator satisfying the following properties:
(i) for every $c \in C, \mathbf{c} \in C^{N}, t \in[0, T]$ and $u \in U$ it holds
$\mathbf{A}(t, \mathbf{c}, c, u)=\mathbf{A}(t, \sigma(\mathbf{c}), c, u) \quad$ for every permutation $\sigma: C^{N} \rightarrow C^{N} ;$
(ii) for every $R>0$ there exists a constant $L_{R} \geq 0$ such that for every $c, \tilde{c} \in$ $B_{R, C}, t \in[0, T], \mathbf{c}, \tilde{\mathbf{c}} \in B_{R, C}^{N}$ and $u \in U$

$$
\|\mathbf{A}(t, \mathbf{c}, c, u)-\mathbf{A}(t, \tilde{\mathbf{c}}, \tilde{c}, u)\|_{E} \leq L_{R}\left(\|c-\tilde{c}\|_{E}+W_{1}\left(\frac{1}{N} \sum_{i=1}^{N} \delta_{c_{i}}, \frac{1}{N} \sum_{i=1}^{N} \delta_{\tilde{c}_{i}}\right)\right)
$$

(iii) there exists $M>0$ such that for every $c \in C, \mathbf{c} \in C^{N}, t \in[0, T]$ and $u \in U$ it holds

$$
\|\mathbf{A}(t, \mathbf{c}, c, u)\|_{E} \leq M\left(1+\|c\|_{E}+m_{1}\left(\frac{1}{N} \sum_{i=1}^{N} \delta_{c_{i}}\right)\right)
$$

(iv) for every $c \in C, \mathbf{c} \in C^{N}$ and $u \in U$ the map $t \mapsto \mathbf{A}(t, \mathbf{c}, c, u)$ belongs to $L^{1}([0, T] ; E)$;
$(v)$ for every $t \in[0, T], c \in C$ and $\mathbf{c} \in C^{N}$ the map $u \mapsto \mathbf{A}(t, \mathbf{c}, c, u)$ belongs to $C^{0}(U ; E) ;$
(vi) for every $R>0$ there exists $\theta>0$ such that for every $t \in[0, T], \mathbf{c} \in B_{R, C}^{N}$ and $u \in U$

$$
c \in C,\|c\|_{E} \leq R \Rightarrow c+\theta \mathbf{A}(t, \mathbf{c}, c, u) \in C
$$

(vii) for every $c \in C, \mathbf{c}=\left(c_{1}, \ldots, c, \ldots, c_{N}\right) \in C^{N}, t \in[0, T]$ and $u \in U, \mathbf{A}$ is $C$-differentiable at $c \in C$, i.e. there exists $\mathrm{D}_{c} \mathbf{A} \in \mathcal{L}\left(E_{C} ; E_{C}\right)$

$$
\lim _{C \ni c^{\prime} \rightarrow c} \frac{\left\|\mathbf{A}\left(t, \mathbf{c}^{\prime}, c^{\prime}, u\right)-\mathbf{A}(t, \mathbf{c}, c, u)-\mathbf{D}_{c} \mathbf{A}\left[c^{\prime}-c\right]\right\|_{E}}{\left\|c^{\prime}-c\right\|_{E}}=0 .
$$

(HLsym) Let $\mathbf{L}:[0, T] \times C^{N} \times U^{N} \rightarrow \mathbb{R}$ be a continuous function such that
(a) for every $\mathbf{c} \in C^{N}, t \in[0, T]$ and $\mathbf{u} \in U^{N}$ it holds
$\mathbf{L}(t, \mathbf{c}, \mathbf{u})=\mathbf{L}(t, \sigma(\mathbf{c}), \iota(\mathbf{u}))$ for every permutations $\sigma: C^{N} \rightarrow C^{N}, \iota: U^{N} \rightarrow U^{N} ;$
(b) for every $R>0$ there exists a constant $L_{R} \geq 0$ such that for every $\mathbf{c}, \tilde{\mathbf{c}} \in$ $\left(B_{R}^{C}\right)^{N}, t \in[0, T]$ and $\mathbf{u} \in U^{N}$

$$
|\mathbf{L}(t, \mathbf{c}, \mathbf{u})-\mathbf{L}(t, \tilde{\mathbf{c}}, \mathbf{u})| \leq L_{R} W_{1}\left(\frac{1}{N} \sum_{i=1}^{N} \delta_{c_{i},}, \frac{1}{N} \sum_{i=1}^{N} \delta_{\tilde{c}_{i}}\right) ;
$$

(c) for every $c \in C, \mathbf{c}=\left(c_{1}, \ldots, c, \ldots, c_{N}\right) \in C^{N}, t \in[0, T]$ and $u \in U, \mathbf{L}$ is $C$-differentiable at $c \in C$ with $C$-differential $\mathrm{D}_{c} \mathbf{L} \in \mathcal{L}\left(E_{C} ; \mathbb{R}\right)$.

We consider the following Bolza problem

$$
\begin{equation*}
\min _{\mathbf{u}(t) \in \mathcal{U}^{N}}\left\{\int_{0}^{T} \mathbf{L}(t, \mathbf{c}(t), \mathbf{u}(t)) \mathrm{d} t\right\} \tag{3.31}
\end{equation*}
$$

subject to the system

$$
\begin{cases}\frac{\mathrm{d}}{\mathrm{~d} t} c_{i}(t)=\mathbf{A}\left(t, \mathbf{c}(t), c_{i}(t), u_{i}(t)\right) & \text { in }(0, T],  \tag{3.32}\\ c_{i}(0)=c_{0, i} \in C, & \text { for } i=1, \ldots, N .\end{cases}
$$

Thanks to the assumptions on $\mathbf{A}$ and by Theorem I. 4 of 19$]$ for every $\mathbf{u} \in \mathcal{U}^{N}$ there exists a unique weak solution $\mathbf{c}(t) \in A C\left([0, T] ; C^{N}\right)$ of system (3.32).

Assume that the minimum control problem (3.31)-(3.32) admits a solution $(\mathfrak{c}, \mathfrak{u}) \in$ $A C\left([0, T] ; C^{N}\right) \times \mathcal{U}^{N}$. Then, after an application of Theorem 3.4, the following result follows.

Theorem 3.5. Let $\mathbf{A}$ and $\mathbf{L}$ satisfy (HAsym) and (HLsym) respectively. Let $\mathfrak{u} \in \mathcal{U}^{N}$ be an admissible control whose corresponding trajectory $\mathfrak{c} \in A C\left([0, T] ; C^{N}\right)$ is optimal for problem (3.31)-(3.32). Then there exists a family of co-state curves $\mathfrak{p} \in A C\left([0, T] ;\left(E_{C}^{*}\right)^{N}\right)$ such that $(\mathfrak{c}, \mathfrak{p}, \mathfrak{u})$ solve in distributional sense the following system

$$
\begin{cases}\frac{\mathrm{d}}{\mathrm{~d} t}\binom{\mathfrak{c}_{i}(t)}{\mathfrak{p}_{i}(t)}=J \mathrm{D}_{\left(\mathfrak{c}_{i}(t), \mathfrak{p}_{i}(t)\right)} \mathcal{H}(t, \mathfrak{c}(t), \mathfrak{p}(t), \mathfrak{u}(t)) & \text { in }[0, T), \\ \mathfrak{c}_{i}(0)=c_{0, i} \in C, & \text { for every } i=1, \ldots, N, \\ \mathfrak{p}_{i}(T)=\mathbf{0} \in E_{C}^{*}, & \\ \mathfrak{u}(t) \in \underset{\mathbf{u} \in U^{N}}{\operatorname{argmax}} \mathcal{H}(t, \mathfrak{c}(t), \mathfrak{p}(t), \mathbf{u}), & \end{cases}
$$

where the Hamiltonian $\mathcal{H}:[0, T] \times C^{N} \times\left(E_{C}^{*}\right)^{N} \times U^{N} \rightarrow \mathbb{R}$ is defined by

$$
\mathcal{H}(t, \mathbf{c}, \mathbf{p}, \mathbf{u})=\sum_{i=1}^{N}\left\langle p_{i}, \mathbf{A}\left(t, \mathbf{c}, c_{i}, u_{i}\right)\right\rangle-\mathbf{L}(t, \mathbf{c}, \mathbf{u})
$$

## 4. Pontryagin maximum principle in the Wasserstein space of convex SPACES

We will henceforth need stronger assumptions than those in Section 3. Let the following assumptions hold:

- $\left(E,\|\cdot\|_{E}\right)$ is a separable Banach space;
- $C \subset E$ is a closed and convex subset of $E$;
- $E_{C}$ is reflexive;
- $\left(Z,\|\cdot\|_{Z}\right)$ is a separable Banach space;
- $\mathcal{U}=L^{1}([0, T] ; U)$ where $U$ is a not empty compact subset of the normed space
$\left(C_{b}^{1}(C ; Z),\|\cdot\|_{C_{b}^{1}}\right):=\left\{u \in C^{1}(C ; Z):\|u\|_{C_{b}^{1}}:=\sup _{c \in C}\|u(c)\|_{z}+\sup _{c \in C}\left\|\mathrm{D}_{c} u\right\|_{\mathcal{L}\left(E_{C} ; Z\right)}<+\infty\right\}$,
where $C^{1}(C ; Z)$ means continuous $C$-differentiability.
From now on we will often use for any $u \in \mathcal{U}$ the identification $u(t)(c)=u(t, c)$.

Remark 4.1. Since $C$ is closed in $E$ (which is Banach and separable) we deduce that $C$ with the induced metric is a complete separable metric space. This implies that $C$ is a separable Radon space (it is actually a Polish space). Similarly, since $E_{C}$ is reflexive and separable, then $E_{C}^{*}$ is a separable Banach space, and so a separable Radon space. It follows that $C \times E_{C}^{*}$ is a separable Radon space and that $\mathcal{P}_{p}(C), \mathcal{P}_{p}\left(E_{C}^{*}\right)$ and $\mathcal{P}_{p}\left(C \times E_{C}^{*}\right)$ are complete separable metric space for every $p \geq 1$.

We define

$$
\mathcal{P}_{c}(C):=\{\mu \in \mathcal{P}(C): \operatorname{supp}(\mu) \text { is compact in } C\}
$$

and

$$
\mathcal{P}_{b}(C):=\{\mu \in \mathcal{P}(C): \operatorname{supp}(\mu) \text { is bounded in } C\}
$$

From now on, unless otherwise specified, when we write $\mathcal{P}_{c}(C)$ and $\mathcal{P}_{b}(C)$ we mean the separable metric spaces $\left(\mathcal{P}_{c}(C), W_{1}\right)$ and $\left(\mathcal{P}_{b}(C), W_{1}\right)$ respectively, both seen as subset of the separable complete metric space $\left(\mathcal{P}_{1}(C), W_{1}\right)$.

Note that $E_{C \times E_{C}^{*}}=\overline{\mathbb{R}\left(C \times E_{C}^{*}-C \times E_{C}^{*}\right)}=E_{C} \times E_{C}^{*}$ and $\left(E_{C} \times E_{C}^{*}\right)^{*}=E_{C}^{*} \times E_{C}^{* *}=$ $E_{C}^{*} \times E_{C}$ (recall that $E_{C}$ is reflexive). We define the linear and continuous operator

$$
\begin{aligned}
J: E_{C}^{*} \times E_{C} & \rightarrow E_{C} \times E_{C}^{*} \\
(T, e) & \mapsto(e,-T)
\end{aligned}
$$

We define the $R$-fattening of the support of a measure $\mu \in \mathcal{P}_{c}(C)$ as

$$
B_{\mu}(R):=\bigcup_{c \in \operatorname{supp}(\mu)} B_{R, C}(c),
$$

where $B_{R, C}(c)$ is the closed subset $\left\{e \in E:\|e-c\|_{E} \leq R\right\} \cap C$. Note that, since $\mu$ has compact support, if $\eta \in \mathcal{P}\left(B_{\mu}(R)\right)$ then $\eta \in \mathcal{P}_{b}(C)$.

From now on let $\phi: \mathcal{P}_{1}(C) \rightarrow Y$ be such that $\mathcal{P}_{b}(C) \subset D(\phi):=\left\{\mu \in \mathcal{P}_{1}(C)\right.$ : $\left.\|\phi(\mu)\|_{Y}<+\infty\right\}$, with $\left(Y,\|\cdot\|_{Y}\right)$ a Banach space. We introduce a definition of local differentiability of $\phi$ at $\mu \in \mathcal{P}_{c}(C)$.

Definition 4.2. A functional $\phi: \mathcal{P}_{1}(C) \rightarrow Y$ is locally differentiable at $\mu \in \mathcal{P}_{c}(C)$ if there exists a map $\nabla_{\mu} \phi(\mu) \in L_{\mu}^{2}\left(C ; \mathcal{L}\left(E_{C} ; Y\right)\right)$ in Bochner sense such that for every $R>0$ and for every $\nu \in \mathcal{P}\left(B_{\mu}(R)\right)$ it holds

$$
\phi(\nu)-\phi(\mu)=\int_{C \times C} \nabla_{\mu} \phi(\mu)\left(c_{1}\right)\left[c_{2}-c_{1}\right] \mathrm{d} \gamma\left(c_{1}, c_{2}\right)+o_{R}\left(W_{2, \gamma}(\mu, \nu)\right)
$$

for any $\gamma \in \Gamma(\mu, \nu)$, where

$$
W_{2, \gamma}^{2}(\mu, \nu):=\int_{C \times C}\left\|c_{1}-c_{2}\right\|_{E}^{2} \mathrm{~d} \boldsymbol{\gamma}\left(c_{1}, c_{2}\right) .
$$

This new notion of local differentiability for functional with values in Banach spaces enjoys the following chain rule.
Proposition 4.3. Let $\mu \in \mathcal{P}_{c}(C)$ and $V$ be a closed and bounded subset of $C$ such that $\operatorname{supp}(\mu) \subset V$. Suppose that $\phi: \mathcal{P}(V) \rightarrow Y$ is locally differentiable over $\mathcal{P}_{c}(V)$. Let $\psi: C \rightarrow C$ be a $C$-differentiable map with $C$-differential $D_{c} \psi$ such that
(i) $D \psi \in L_{\mu}^{\infty}\left(C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)$ where $D \psi(c):=D_{c} \psi$;
(ii) $\operatorname{supp}\left(\psi_{\#} \mu\right) \subset V$.

Let $\bar{\varepsilon}>0$ and $\mathcal{G}:(-\bar{\varepsilon}, \bar{\varepsilon}) \times V \rightarrow C$ be a map such that:
(iii) $\mathcal{G}(0, \cdot)=I d_{V}$ and the map $\varepsilon \mapsto \mathcal{G}(\varepsilon, c) \in C \subset E$ is Fréchet-differentiable at $\varepsilon=0$ uniformly in $V$;
(iv) the map $\mathcal{F}: V \rightarrow E_{C}$ which is defined by $c \mapsto \mathcal{F}(c):=\left.\frac{\mathrm{d}}{\mathrm{d} \varepsilon} \mathcal{G}(\varepsilon, c)\right|_{\varepsilon=0}$ belongs to $L_{\mu}^{\infty}\left(V ; E_{C}\right)$;
(v) $\operatorname{supp}\left((\psi \circ \mathcal{G}(\varepsilon, \cdot))_{\# \mu} \mu\right) \subset V$ for every $\varepsilon \in(-\bar{\varepsilon}, \bar{\varepsilon})$.

Then the map $\varepsilon \mapsto \phi\left((\psi \circ \mathcal{G}(\varepsilon, \cdot))_{\#} \mu\right)$ is Fréchet-differentiable at $\varepsilon=0$ and

$$
\begin{equation*}
\left.\frac{\mathrm{d}}{\mathrm{~d} \varepsilon} \phi\left((\psi \circ \mathcal{G}(\varepsilon, \cdot))_{\#} \mu\right)\right|_{\varepsilon=0}=\int_{C} \nabla_{\psi_{\#} \mu} \phi\left(\psi_{\#} \mu\right)(\psi(c))\left[D_{c} \psi[\mathcal{F}(c)]\right] \mathrm{d} \mu(c) \tag{4.1}
\end{equation*}
$$

Proof. For notational convenience set $\psi_{\varepsilon}:=(\psi \circ \mathcal{G}(\varepsilon, \cdot)): C \rightarrow C$. First we observe that, by $(i i)$, since $\psi: C \rightarrow C$ is continuous and $\mu \in \mathcal{P}_{c}(C)$, then $\operatorname{supp}\left(\psi_{\#} \mu\right)=\psi(\operatorname{supp}(\mu))$ is compact in $V$ (see formula (5.2.6) of $[7]$ ). Hence $\psi_{\#} \mu \in \mathcal{P}_{c}(V)$ and, by $(v),\left(\psi_{\varepsilon}\right)_{\#} \mu \in$ $\mathcal{P}\left(B_{\mu}(R)\right)$ for $R>0$ large enough. Then, using that $\boldsymbol{\mu}_{\varepsilon}:=\left(\psi, \psi_{\varepsilon}\right)_{\#} \mu \in \Gamma\left(\psi_{\#} \mu,\left(\psi_{\varepsilon}\right)_{\#} \mu\right)$ and that $\phi$ is locally differentiable over $\mathcal{P}_{c}(V)$, from Definition 4.2 we obtain

$$
\begin{align*}
& \frac{\phi\left(\left(\psi_{\varepsilon}\right)_{\#} \mu\right)-\phi\left(\psi_{\#} \mu\right)}{\varepsilon} \\
& =\int_{C \times C} \nabla_{\psi_{\#} \mu} \phi\left(\psi_{\#} \mu\right)(c)\left[\frac{\tilde{c}-c}{\varepsilon}\right] \mathrm{d} \boldsymbol{\mu}_{\varepsilon}(c, \tilde{c})+\frac{1}{\varepsilon} o_{R}\left(W_{2, \mu_{\varepsilon}}\left(\psi_{\#} \mu,\left(\psi_{\varepsilon}\right)_{\#} \mu\right)\right) \\
& =\int_{C} \nabla_{\psi_{\# \mu} \mu} \phi\left(\psi_{\#} \mu\right)(\psi(c))\left[\frac{\psi_{\varepsilon}(c)-\psi(c)}{\varepsilon}\right] \mathrm{d} \mu(c) \\
& \quad+\frac{1}{\varepsilon} o_{R}\left(\left(\int_{C}\left\|\psi_{\varepsilon}(c)-\psi(c)\right\|_{E}^{2} \mathrm{~d} \mu(c)\right)^{\frac{1}{2}}\right) \tag{4.2}
\end{align*}
$$

Using that $\mathcal{G}$ is Frechet differentiable at $\varepsilon=0$ uniformly in $V$ and since $\psi$ is $C$ differentiable, we have

$$
\begin{equation*}
\psi_{\varepsilon}(c)=\psi(\mathcal{G}(\varepsilon, c))=\psi\left(c+\varepsilon \mathcal{F}(c)+r_{\mathcal{F}}(\varepsilon)\right)=\psi(c)+\varepsilon D_{c} \psi[\mathcal{F}(c)]+r(\varepsilon, c), \tag{4.3}
\end{equation*}
$$

where

$$
r(\varepsilon, c)=D_{c} \psi\left[r_{\mathcal{F}}(\varepsilon)\right]+o\left(\left\|\varepsilon \mathcal{F}(c)+r_{\mathcal{F}}(\varepsilon)\right\|\right)
$$

with $r_{\mathcal{F}}(\varepsilon)=o(\varepsilon)$ not depending on $c$. It follows from (i) and (iv) that

$$
\begin{equation*}
\|r(\varepsilon, c)\|_{L_{\mu}^{\infty}\left(C ; E_{C}\right)} \leq 2\|D \psi\|_{L_{\mu}^{\infty}\left(C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)}\left\|r_{\mathcal{F}}(\varepsilon)\right\|_{E}=o(\varepsilon) . \tag{4.4}
\end{equation*}
$$

Now we can pass to the limit in (4.2) as $\varepsilon \rightarrow 0$. Indeed, by (4.3)-(4.4), we have

$$
\lim _{\varepsilon \rightarrow 0} \frac{\psi_{\varepsilon}(c)-\psi(c)}{\varepsilon}=D_{c} \psi[\mathcal{F}(c)] \quad \text { a.e. } c \in C
$$

and

$$
\left\|\frac{\psi_{\varepsilon}(c)-\psi(c)}{\varepsilon}\right\|_{L_{\mu}^{\infty}\left(C ; E_{C}\right)} \leq 2\|D \psi\|_{L_{\mu}^{\infty}\left(C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)}\|\mathcal{F}\|_{L_{\mu}^{\infty}\left(C ; E_{C}\right)} .
$$

Hence we deduce that

$$
o_{R}\left(\left(\int_{C}\left\|\psi_{\varepsilon}(c)-\psi(c)\right\|_{E}^{2} \mathrm{~d} \mu(c)\right)^{\frac{1}{2}}\right)=o_{R}(\varepsilon) \quad \text { as } \varepsilon \rightarrow 0
$$

Therefore, recalling that $\nabla_{\psi_{\#} \mu} \phi\left(\psi_{\#} \mu\right) \in L_{\mu}^{2}\left(C ; \mathcal{L}\left(E_{C} ; Y\right)\right.$ ), after an application of the Lebesgue theorem we obtain (4.1).

We start by considering the following Mayer optimal control problem:

$$
\begin{equation*}
\min _{u \in \mathcal{U}} \varphi(\mu(T)) \tag{4.5}
\end{equation*}
$$

subject to

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \mu(t)+\operatorname{div}(A(t, \mu(t), \cdot, u(t, \cdot)) \mu(t))=0 \quad \text { in }(0, T]  \tag{4.6}\\
\mu(0)=\mu^{0} \in \mathcal{P}_{c}(C)
\end{array}\right.
$$

We assume that the non-local velocity field $A:[0, T] \times \mathcal{P}_{b}(C) \times C \times C_{b}^{0}(C ; Z) \rightarrow E_{C}$ satisfies the following.
(HA1): Assumptions on $A$ for the well-posedness of 4.6)
(i) there exists a constant $L>0$ such that for every $t \in[0, T], \mu, \tilde{\mu} \in \mathcal{P}_{c}(C)$, $c, \tilde{c} \in C$ and $u, \tilde{u} \in C_{b}^{0}(C ; Z)$

$$
\|A(t, \mu, c, u)-A(t, \tilde{\mu}, \tilde{c}, \tilde{u})\|_{E} \leq L\left(\|c-\tilde{c}\|_{E}+W_{1}(\mu, \tilde{\mu})+\|u-\tilde{u}\|_{C_{b}^{0}}\right)
$$

(ii) there exists $M>0$ such that for every $t \in[0, T], \mu \in \mathcal{P}_{c}(C), c \in C$ and $u \in U$ there holds

$$
\|A(t, \mu, c, u)\|_{E} \leq M\left(1+\|c\|_{E}+m_{1}(\mu)\right)
$$

(iii) for every $\mu \in \mathcal{P}_{c}(C), c \in C$ and $u \in U$ the map $t \mapsto A(t, \mu, c, u)$ belongs to $L^{1}([0, T] ; E) ;$
(iv) for every $R>0$ there exists $\theta>0$ such that for every $t \in[0, T], \mu \in \mathcal{P}_{c}(C)$ and $u \in U$

$$
c \in C,\|c\|_{E} \leq R \Rightarrow c+\theta A(t, \mu, c, u) \in C .
$$

Under (HA1), it follows from Theorem 3.3 of [6] that, for every $u \in \mathcal{U}$, there exists a unique solution $\mu \in A C\left([0, T] ; \mathcal{P}_{c}(C)\right)$ of 4.6). The curve $\mu(t)$ can be represented as

$$
\begin{equation*}
\mu(t)=\left(\Phi_{(0, t)}^{\mu^{0}}\right)_{\#}^{\mu^{0}} \quad \text { for } t \in[0, T] \tag{4.7}
\end{equation*}
$$

where, for $0 \leq s \leq t \leq T, \Phi_{(s, t)}^{\mu(s)}$ : $C \rightarrow C$ denotes the family of non-local flows defined by

$$
\begin{equation*}
\Phi_{(s, t)}^{\mu^{s}}(c):=c+\int_{s}^{t} A\left(t, \mu(\sigma), \Phi_{(s, \sigma)}^{\mu^{s}}(c), u\left(\sigma, \Phi_{(s, \sigma)}^{\mu^{s}}(c)\right)\right) \mathrm{d} \sigma, \quad \text { with } \mu^{s}:=\mu(s) \tag{4.8}
\end{equation*}
$$

Moreover, by Theorem 1 of [1], for every $0 \leq s \leq t \leq T$ we have

$$
\begin{equation*}
\Phi_{(s, t)}^{\mu^{s}} \in C_{b}^{0}(C ; C) \tag{4.9}
\end{equation*}
$$

In order to find necessary conditions for the optimal control problem (4.5)-(4.6) we need further assumptions.
(HA2): Assumptions on $A$ for the well-posedness of 4.14)
$(v)$ for every $t \in[0, T], \mu \in \mathcal{P}_{c}(C)$ and $c \in C$ the map $u \mapsto A(t, \mu, c, u)$ belongs to $C^{1}\left(C_{b}^{1}(C ; Z) ; E\right)$;
(vi) for every $t \in[0, T], \mu \in \mathcal{P}_{c}(C)$ and $u \in U$ the map $c \mapsto A(t, \mu, c, u(t, c))$ is $C$-differentiable with $C$-differential $\mathrm{D}_{c} A$ and the application $\mathrm{D} A: \mathcal{P}_{c}(C) \times$ $C \times U \rightarrow \mathcal{L}\left(E_{C} ; E_{C}\right)$ defined by $(\mu, c, u) \mapsto \mathrm{D}_{c} A:=\mathrm{D}_{c} A(t, \mu, c, u(t, c))$ is continuous for every $t \in[0, T]$;
(vii) for every $t \in[0, T], c \in C$ and $u \in U$ the map $\mathcal{P}_{c}(C) \ni \mu \mapsto A(t, \mu, c, u) \in$ $E_{C}$ is locally differentiable at any $\mu$ in the sense of the Definition 4.2 with differential $\nabla_{\mu} A:=\nabla_{\mu} A(t, \mu, c, u)$ and the application $\mathcal{P}_{c}(C) \times C \times U \times C \ni$ $(\mu, c, u, \tilde{c}) \mapsto \nabla_{\mu} A(\tilde{c}) \in \mathcal{L}\left(E_{C} ; E_{C}\right)$ is continuous for every $t \in[0, T]$.

With an extended but similar argument to the one used to prove Lemma 3.1, we can get a boundedness result for $\mathrm{D}_{\Phi_{(s, t)}^{\mu^{s}}(c)} A$. We state such result in the following lemma which will be proved in Appendix A.
Lemma 4.4. Under the assumptions (HA1)-(i) and (HA2)-(vi), it holds that

$$
\begin{equation*}
\left\|\mathrm{D}_{\Phi_{(s, t)}^{\mu^{s}}(c)} A\right\|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \leq L \quad \text { for every } 0 \leq s \leq t \leq T \text { and } c \in C \tag{4.10}
\end{equation*}
$$

where $L$ is a positive constant which only depends on the Lipschitz constant of $A$ and on $U$. In particular,

$$
\text { the map } \quad(t, c) \mapsto \mathrm{D}_{\Phi_{(s, t)}^{\mu^{s}}(c)} A \quad \text { belongs to } L_{\mathcal{L} \times \mu^{s}}^{\infty}\left([s, T] \times C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)
$$

Thanks to Lemma 4.4 we can deduce a regularity result for the family of non-local flows defined in (4.8). In the next lemma we show that this family is $C$-differentiable and characterize its $C$-differential. Since the techniques used in the proof are similar to those seen in Section 3, we postpone the proof in Appendix A.
Lemma 4.5. Under the assumptions (HA1) and (HA2)-(vi), it holds that for every $s \in[0, T]$ the $\operatorname{map} \Phi_{(s, t)}^{\mu^{s}}: C \rightarrow C$ is $C$-differentiable with $C$-differential $\mathrm{D}_{c} \Phi_{(s, t)}^{\mu^{s}}$ at $c \in C$ and

$$
\text { the map } \quad(t, c) \mapsto \mathrm{D}_{c} \Phi_{(s, t)}^{\mu^{s}} \quad \text { belongs to } L_{\mathcal{L} \times \mu^{s}}^{\infty}\left([s, T] \times C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)
$$

Moreover, for every $c \in C$ and $f \in E_{C}, D_{c} \Phi_{(s, t)}^{\mu^{s}}[f]$ is the unique solution to the linear differential equation

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} z(t, c)=\mathrm{D}_{\Phi_{(s, t)}^{\mu^{s}}(c)} A[z(t, c)] \quad \text { in }(s, T]  \tag{4.11}\\
z(s, c)=f
\end{array}\right.
$$

To adapt the needle variation technique to the infinite-dimensional case we need another lemma describing how the flow (4.8) and the solution of (4.6) behave when we do an infinitesimal variation of the control. Also in this case we postpone the proof to Appendix A. For simplicity of notation we give the result for initial time $s=0$.

Lemma 4.6. Assume (HA1) and (HA2). Let $\varepsilon>0$. Let $u, u_{\varepsilon} \in \mathcal{U}$ be such that $u_{\varepsilon} \rightarrow u$ in $L^{1}\left([0, T] ;\left(U,\|\cdot\|_{C_{b}^{1}}\right)\right)$ as $\varepsilon \rightarrow 0$. Let $\mu^{0}, \mu_{\varepsilon}^{0} \in \mathcal{P}_{c}(C)$ be such that $W_{1}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right) \rightarrow 0$ as $\varepsilon \rightarrow 0$. Let $\mu$ be the solution of (4.6) with initial datum $\mu^{0}$ and corresponding family of flows $\Phi_{(0, t)}^{\mu^{0}}$. Let $\mu_{\varepsilon}$ be the solution of (4.6) in which $u$ is replaced by $u_{\varepsilon}$ with initial datum $\mu_{\varepsilon}^{0}$ and corresponding family of flows $\Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}$. Then the following hold as $\varepsilon \rightarrow 0$ :
(a) $\Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}} \rightarrow \Phi_{(0, t)}^{\mu^{0}}$ in $C_{b}^{0}([0, T] \times C ; C)$;
(b) $W_{1}\left(\mu_{\varepsilon}(t), \mu(t)\right) \rightarrow 0$ uniformly in $[0, T]$;
(c) $\mathrm{D}_{c} \Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}} \rightarrow \mathrm{D}_{c} \Phi_{(0, t)}^{\mu^{0}}$ in $\mathcal{L}\left(E_{C} ; E_{C}\right)$ for every $c \in C$ and $t \in[0, T]$.

As for the final cost $\varphi: \mathcal{P}_{c}(C) \rightarrow \mathbb{R}$ we assume:
$(\mathbf{H} \varphi): \varphi$ is Lipschitz continuous w.r.t. the $W_{1}$-metric and locally differentiable over $\mathcal{P}_{c}(C)$ with differential $\nabla_{\mu} \varphi$. Moreover, the map $C \ni c \mapsto \nabla_{\mu} \varphi(c) \in E_{C}^{*}$ is continuous.

We define the Hamiltonian $\mathcal{H}:[0, T] \times \mathcal{P}_{b}\left(C \times E_{C}^{*}\right) \times U \rightarrow \mathbb{R}$ of the Mayer problem (4.5)-(4.6) as

$$
\begin{equation*}
\mathcal{H}(t, \nu, \omega)=\int_{C \times E_{C}^{*}}\left\langle p, A\left(t, \pi_{\#}^{1} \nu, c, \omega(c)\right)\right\rangle \mathrm{d} \nu(c, p) . \tag{4.12}
\end{equation*}
$$

Note that, by (HA1)-(ii), $\mathcal{H}(t, \nu, \omega)$ is finite for every $\nu \in \mathcal{P}_{b}\left(C \times E_{C}^{*}\right)$. Moreover, thanks to the assumptions (HA1)-(HA2) we can apply Lemma B. 2 to obtain that $\mathcal{H}$ is locally differentiable at any $\nu \in \mathcal{P}_{c}\left(C \times E_{C}^{*}\right)$ in the sense of Definition 4.2. The following explicit formula for its differential $\nabla_{\nu} \mathcal{H}(t, \nu, u): C \times E_{C}^{*} \rightarrow E_{C}^{*} \times E_{C}$ holds:

$$
\begin{equation*}
\nabla_{\nu} \mathcal{H}(t, \nu, u)(c, p)=\binom{D_{c}^{*} A\left(t, \pi_{\#}^{1} \nu, c, u\right)[p]+\int_{C \times E_{C}^{*}} \nabla_{\pi_{\#}^{1}}^{*} A\left(t, \pi_{\#}^{1} \nu, \tilde{c}, u\right)(c)[\tilde{p}] \mathrm{d} \nu(\tilde{c}, \tilde{p})}{A\left(t, \pi_{\#}^{1} \nu, c, u\right)} \tag{4.13}
\end{equation*}
$$

where $D_{c}^{*} A\left(t, \pi_{\#}^{1} \nu, c, u\right)$ is the adjoint operator of $D_{c} A\left(t, \pi_{\#}^{1} \nu, c, u\right)$ and $\nabla_{\pi_{\#}^{1}}^{*} A\left(t, \pi_{\#}^{1} \nu, \tilde{c}, u\right)(c)$ is the adjoint operator of $\nabla_{\pi_{\#}^{1}} \nu\left(t, \pi_{\#}^{1} \nu, \tilde{c}, u\right)(c)$.

We state the Pontryagin maximum principle for the infinite-dimensional Mayer problem (4.5)-4.6).

Theorem 4.7. Let $(\boldsymbol{\mu}, \mathfrak{u}) \in A C\left([0, T] ; \mathcal{P}_{c}(C)\right) \times \mathcal{U}$ be an optimal pair control trajectory for (4.5)-(4.6). Then there exists $\boldsymbol{\nu} \in A C\left([0, T] ; \mathcal{P}_{c}\left(C \times E_{C}^{*}\right)\right)$ which solves in distributional sense

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \boldsymbol{\nu}(t)=-\operatorname{div}_{(c, p)}\left(\left(J \nabla_{\boldsymbol{\nu}(t)} \mathcal{H}(t, \boldsymbol{\nu}(t), \mathfrak{u}(t))(\cdot, \cdot)\right) \boldsymbol{\nu}(t)\right) \quad \text { in }[0, T),  \tag{4.14}\\
\pi_{\#}^{1} \boldsymbol{\nu}(t)=\boldsymbol{\mu}(t), \\
\boldsymbol{\nu}(T)=\left(I d,-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))\right)_{\#} \boldsymbol{\mu}(T),
\end{array}\right.
$$

where $\mathcal{H}$ and $\nabla_{\nu} \mathcal{H}$ are defined by (4.12) and (4.13) respectively. Moreover the following maximality condition holds

$$
\begin{equation*}
\mathcal{H}(t, \boldsymbol{\nu}(t), \mathfrak{u}(t))=\max _{\omega \in U}\{\mathcal{H}(t, \boldsymbol{\nu}(t), \omega)\} \quad \text { for a.e. } t \in[0, T] . \tag{4.15}
\end{equation*}
$$

In what follows we are concerned with the proof of Theorem 4.7.
First we give a boundedness result for $\nabla_{\mu} A$ which will be useful to derive the linearised non-local equation.
Lemma 4.8. Let A satisfy (HA1)-(i) and (HA2)-(vi). Then

$$
\left\|\nabla_{\mu} A(\tilde{c})\right\|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \leq L
$$

for every $(\mu, c, u) \in \mathcal{P}_{c}(C) \times C \times C_{b}^{1}(C ; Z)$ and a.e. $t \in[0, T]$ and $\mu$-a.e. $\tilde{c} \in \operatorname{supp}(\mu)$, where $L$ is the Lipschitz constant of $A$. In particular,

$$
\nabla_{\mu} A \in L_{\mathcal{L} \times \mu \times \mu}^{\infty}\left([0, T] \times \underset{19}{C} \times \operatorname{supp}(\mu) ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right) .
$$

Proof. Let $\mu \in \mathcal{P}_{c}(C), c \in C$ and $u \in C_{b}^{1}(C ; Z)$. Fix $c_{0} \in \operatorname{supp}(\mu)$ and $c_{1}, c_{2}, \bar{c} \in C$. We define for $r>0,0<\varepsilon<1$ and $i=1,2$

$$
\varphi_{\varepsilon}^{i}(\tilde{c}):= \begin{cases}c_{i} & \text { for } \tilde{c} \in B_{r}\left(c_{0}\right) \cap C  \tag{4.16}\\ \left(1-\frac{\|\tilde{c}\|-r}{\varepsilon}\right) c_{i}+\left(\frac{\|\tilde{\|}\|-r}{\varepsilon}\right) \bar{c} & \text { for } \tilde{c} \in\left(B_{r+\varepsilon}\left(c_{0}\right) \backslash B_{r}\left(c_{0}\right)\right) \cap C \\ \bar{c} & \text { otherwise }\end{cases}
$$

Then $\varphi_{\varepsilon}^{i}$ are continuous functions from $C$ to $C$ for $i=1,2$. Since $C$ is convex, we deduce that $(1-\varepsilon) I d+\varepsilon \varphi_{\varepsilon}^{i}: C \rightarrow C$, and therefore $\mu_{\varepsilon}^{i}:=\left((1-\varepsilon) I d+\varepsilon \varphi_{\varepsilon}^{i}\right)_{\#} \mu$ are in $\mathcal{P}_{c}(C)$. It follows choosing in Definition 4.2

$$
\boldsymbol{\gamma}_{i}=\left((1-\varepsilon) I d+\varepsilon \varphi_{\varepsilon}^{i}, I d\right)_{\#} \mu \in \Gamma\left(\mu_{\varepsilon}^{i}, \mu\right)
$$

that

$$
\begin{aligned}
& A\left(t, \mu_{\varepsilon}^{2}, c, u\right)-A\left(t, \mu_{\varepsilon}^{1}, c, u\right)=A\left(t, \mu_{\varepsilon}^{2}, c, u\right)-A(t, \mu, c, u)-\left(A\left(t, \mu_{\varepsilon}^{1}, c, u\right)-A(t, \mu, c, u)\right) \\
& =\varepsilon \int_{C} \nabla_{\mu} A(t, \mu, c, u)(\tilde{c})\left[\varphi_{\varepsilon}^{2}(\tilde{c})-\varphi_{\varepsilon}^{1}(\tilde{c})\right] \mathrm{d} \mu(\tilde{c}) \\
& +o_{R}\left(\varepsilon\left(\int_{C}\left\|\varphi_{\varepsilon}^{2}(\tilde{c})-\tilde{c}\right\|_{E}^{2} \mathrm{~d} \mu(\tilde{c})\right)^{\frac{1}{2}}\right)+o_{R}\left(\varepsilon\left(\int_{C}\left\|\varphi_{\varepsilon}^{1}(\tilde{c})-\tilde{c}\right\|_{E}^{2} \mathrm{~d} \mu(\tilde{c})\right)^{\frac{1}{2}}\right) .
\end{aligned}
$$

Hence, by (4.16), using (HA1)-(i) and recalling that $\operatorname{supp}(\mu)$ is compact in $C$, we have

$$
\begin{align*}
& \left\|\int_{B_{r+\varepsilon}\left(c_{0}\right) \cap C} \nabla_{\mu} A(t, \mu, c, u)(\tilde{c})\left[\varphi_{\varepsilon}^{2}(\tilde{c})-\varphi_{\varepsilon}^{1}(\tilde{c})\right] \mathrm{d} \mu(\tilde{c})\right\|_{E} \\
& \leq \frac{1}{\varepsilon}\left\|A\left(t, \mu_{\varepsilon}^{2}, c, u\right)-A\left(t, \mu_{\varepsilon}^{1}, c, u\right)\right\|_{E}+\left\|\frac{1}{\varepsilon} o_{R}(M \varepsilon)\right\|_{E} \\
& \leq \frac{L}{\varepsilon} W_{1}\left(\mu_{\varepsilon}^{1}, \mu_{\varepsilon}^{2}\right)+\left\|\frac{1}{\varepsilon} o_{R}(M \varepsilon)\right\|_{E} \\
& \leq L\left\|c_{1}-c_{2}\right\|_{E} \mu\left(B_{r+\varepsilon}\left(c_{0}\right) \cap C\right)+\left\|\frac{1}{\varepsilon} o_{R}(M \varepsilon)\right\|_{E} \tag{4.17}
\end{align*}
$$

where $M$ is a positive constant dependent on $\bar{c}, c_{1}, c_{2}$ and $\operatorname{supp}(\mu)$. Now, by Definition 4.2, $\nabla_{\mu} A(t, \mu, c, u) \in L_{\mu}^{2}\left(C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)$. Moreover,

$$
\left\|\varphi_{\varepsilon}^{2}(\tilde{c})-\varphi_{\varepsilon}^{1}(\tilde{c})\right\|_{E} \leq\left(\left\|c_{1}\right\|_{E}+\left\|c_{2}\right\|_{E}+\|\bar{c}\|_{E}\right) \in L_{\mu}^{2}\left(C ; E_{C}\right)
$$

and

$$
\varphi_{\varepsilon}^{2}(\tilde{c})-\varphi_{\varepsilon}^{1}(\tilde{c}) \rightarrow\left(c_{2}-c_{1}\right) \chi_{\left\{B_{r}\left(c_{0}\right) \cap C\right\}} \quad \text { pointwise in } E .
$$

Therefore, after an application of the Lebesgue theorem for general measures, we infer that as $\varepsilon \rightarrow 0^{+}$

$$
\begin{equation*}
\left\|\frac{1}{\mu\left(B_{r}\left(c_{0}\right) \cap C\right)} \int_{B_{r}\left(c_{0}\right) \cap C} \nabla_{\mu} A(t, \mu, c, u)(\tilde{c})\left[c_{2}-c_{1}\right] \mathrm{d} \mu(\tilde{c})\right\|_{E} \leq L\left\|c_{2}-c_{1}\right\|_{E} . \tag{4.18}
\end{equation*}
$$

Finally, by (HA2)-(vi), $\nabla_{\mu} A(t, \mu, c, u) \in C^{0}\left(C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right)$ for a.e. $t \in[0, T]$. Then, applying Lebesgue differentiation theorem to (4.18), we deduce as $r \rightarrow 0^{+}$that

$$
\left\|\nabla_{\mu} A(t, \mu, c, u)\left(c_{0}\right)\left[c_{2}-c_{1}\right]\right\|_{E} \leq L\left\|c_{2}-c_{1}\right\|_{E}
$$

for any $c, c_{1}, c_{2} \in C, \mu \in \mathcal{P}_{c}(C), u \in C_{b}^{1}(C ; Z)$ and $c_{0} \in \operatorname{supp}(\mu)$. Thanks to the arbitrariness of $c_{2}-c_{1} \in E_{C}$, the result follows.

Now we derive the expression of the linearised Cauchy problem associated to the nonlocal continuity equation 4.6).

Proposition 4.9. Let $V$ be a closed and bounded subset of $C$ and $\mu^{0} \in \mathcal{P}_{c}(V)$. Assume (HA1)-(HA2). Fix $u \in \mathcal{U}$ and let $\mu(t)$ and $\Phi_{(0, t)}^{\mu^{0}}$ be the solution of 4.6 and the associated non-local flow defined in (4.8) respectively.
Let $\bar{\varepsilon}>0$ and $\mathcal{G}:(-\bar{\varepsilon}, \bar{\varepsilon}) \times V \rightarrow C$ be a continuous map such that:
(i) $\mathcal{G}(0, \cdot)=I d_{V}$ and the map $\varepsilon \mapsto \mathcal{G}(\varepsilon, c) \in C \subset E$ is Fréchet-differentiable at $\varepsilon=0$ uniformly in $V$;
(ii) the map $\mathcal{F}: V \rightarrow E_{C}$ which is defined by $c \mapsto \mathcal{F}(c):=\left.\frac{\mathrm{d}}{\mathrm{d} \varepsilon} \mathcal{G}(\varepsilon, c)\right|_{\varepsilon=0}$ belongs to $L_{\mu^{0}}^{\infty}\left(V ; E_{C}\right)$;
(iii) $\operatorname{supp}\left(\left(\Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{G}(\varepsilon, \cdot)\right)_{\#} \mu^{0}\right) \subset V$ for every $\varepsilon \in(-\bar{\varepsilon}, \bar{\varepsilon})$.

Set $\mu_{\varepsilon}^{0}:=\mathcal{G}(\varepsilon, \cdot)_{\#} \mu^{0}$. Then, the map $(-\bar{\varepsilon}, \bar{\varepsilon}) \ni \varepsilon \mapsto \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c) \in C$ is Fréchet-differentiable at $\varepsilon=0$ for all $(t, c) \in[0, T] \times V$. Moreover, its differential $v(t, c)$ is continuous with respect to $c$ and it is the unique solution of the non-local Cauchy problem

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} v(t, c)=\mathrm{D}_{\Phi_{(0, t)}^{\mu^{0}}(c)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)[v(t, c)]  \tag{4.19}\\
\quad+\int_{C} \nabla_{\mu(t)} A\left(t, \mu(t) \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\left[\mathrm{D}_{\tilde{c}} \Phi_{(0, t)}^{\mu^{0}}(\tilde{c}) \mathcal{F}(\tilde{c})+v(t, \tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c}), \\
v(0, c)=0 .
\end{array}\right.
$$

Proof. For notational convenience we define $\mathcal{G}^{\varepsilon}(c):=\mathcal{G}(\varepsilon, c)$ for every $(\varepsilon, c) \in(-\bar{\varepsilon}, \bar{\varepsilon}) \times V$ and $\mu_{\varepsilon}(t):=\left(\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}\right)_{\#} \mu_{\varepsilon}^{0}$. Note that $\mu_{\varepsilon}(t)$ is the unique solution of 4.6) with initial datum $\mu_{\varepsilon}^{0}$. Moreover, by (HA2), we can apply Lemma B. 1 obtaining that there exists a unique distributional solution $v(t, c)$ of (4.19) which is continuous with respect to $c$.
Let $v_{\varepsilon}:[0, T] \times V \rightarrow E_{C}$ be

$$
\begin{equation*}
v_{\varepsilon}(t, c):=\frac{\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)}{\varepsilon} \quad \text { for } \varepsilon \in(-\bar{\varepsilon}, \bar{\varepsilon}) \tag{4.20}
\end{equation*}
$$

Then the thesis is equivalent to

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} v_{\varepsilon}(t, c)=v(t, c) \quad \text { for any }(t, c) \in[0, T] \times V \tag{4.21}
\end{equation*}
$$

To prove (4.21) we divide the proof into steps.
Step 1. In this step we prove that $v_{\varepsilon}(t, c)$ is uniformly bounded in $[0, T] \times V$ for $|\varepsilon|$ small enough.
By (4.8) we have for a positive constant $L_{V}$ only dependent on $V$ that

$$
\left.\begin{array}{rll}
\left\|v_{\varepsilon}(t, c)\right\|_{E} & \leq & \frac{1}{|\varepsilon|} \int_{0}^{t} \| A\left(s, \mu_{\varepsilon}(s), \Phi_{(0, s)}^{\mu_{0}^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu_{\omega^{0}}^{0}}(c)\right)\right)-A\left(s, \mu(s), \Phi_{(0, s)}^{\mu^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu_{0}^{0}}(c)\right)\right)
\end{array} \|_{E} \mathrm{~d} s\right)
$$

where in the last inequality we have applied Theorem 4.1 of [6]. Therefore we deduce that

$$
\begin{equation*}
\left\|v_{\varepsilon}(t, c)\right\|_{E} \leq L \int_{0}^{t}\left\|v_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} s+\left(\mathrm{e}^{2 L T}-1\right) \frac{W_{1}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right)}{2|\varepsilon|} . \tag{4.22}
\end{equation*}
$$

To estimate the second term on the right-hand side of (4.22), we use the definition of 1-Wasserstain distance recalling that $\left(\mathcal{G}^{\varepsilon}, I d\right)_{\#} \mu^{0} \in \Gamma\left(\mu_{\varepsilon}^{0}, \mu^{0}\right)$ to obtain

$$
\begin{aligned}
W_{1}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right) & =W_{1}\left(\mathcal{G}_{\#}^{\varepsilon} \mu^{0}, \mu^{0}\right) \leq \int_{V}\|\mathcal{G}(\varepsilon, \tilde{c})-\tilde{c}\|_{E} \mathrm{~d} \mu^{0}(\tilde{c}) \\
& \stackrel{(i)-(i i)}{=} \int_{V}\|\varepsilon \mathcal{F}(\tilde{c})+o(\varepsilon)\|_{E} \mathrm{~d} \mu^{0}(\tilde{c}) \stackrel{(i i)}{\leq} 2|\varepsilon|\|\mathcal{F}\|_{L_{\mu^{0}}^{\infty}\left(V ; E_{C}\right)} .
\end{aligned}
$$

Hence, in view of (4.22), we infer that

$$
\left\|v_{\varepsilon}(t, c)\right\|_{E} \leq L \int_{0}^{t}\left\|v_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} s+\left(\mathrm{e}^{2 L T}-1\right)\|\mathcal{F}\|_{L_{\mu}^{\infty}\left(V ; E_{C}\right)}
$$

which implies, applying the Grönwall inequality, that for $|\varepsilon|$ small enough

$$
\begin{equation*}
\left\|v_{\varepsilon}(t, c)\right\|_{E} \leq \mathrm{e}^{3 L T}\|\mathcal{F}\|_{L_{\mu}{ }^{\infty}\left(V ; E_{C}\right)} \quad \text { uniformly in }[0, T] \times V . \tag{4.23}
\end{equation*}
$$

Step 2. In this step we explicitly write $v_{\varepsilon}(t, c)-v(t, c)$.
By (4.20) and 4.19, we have, noting that $v_{\varepsilon}(0, c)=0$, that for every $(t, c) \in[0, T] \times V$

$$
\begin{align*}
& v_{\varepsilon}(t, c)-v(t, c)=\int_{0}^{t} \mathrm{D}_{\Phi_{(0, s)}^{\mu^{0}}(c)} A\left(s, \mu(s), \Phi_{(0, s)}^{\mu^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right)\left[v_{\varepsilon}(s, c)-v(s, c)\right] \mathrm{d} s \\
& \quad+\int_{0}^{t} \int_{C} \nabla_{\mu(s)} A\left(s, \mu(s) \Phi_{(0, s)}^{\mu^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right)\left(\Phi_{(0, s)}^{\mu^{0}}(\tilde{c})\right)\left[v_{\varepsilon}(s, \tilde{c})-v(s, \tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c}) \mathrm{d} s \\
& \quad+\int_{0}^{t} r_{\varepsilon}(s, c) \mathrm{d} s \tag{4.24}
\end{align*}
$$

where

$$
\begin{align*}
& r_{\varepsilon}(t, c)=\frac{1}{\varepsilon}\left\{A\left(t, \mu_{\varepsilon}^{(t), \Phi_{(0, t)}^{\mu_{( }^{0}}(c), u}\left(t, \Phi_{(0, t)}^{\mu_{( }^{0}}(c)\right)\right)-A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{0}^{0}}(c)\right)\right)\right\} \\
& -\mathrm{D}_{\Phi_{(0, t)}^{\mu^{0}}(c)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left[v_{\varepsilon}(t, c)\right]  \tag{4.25}\\
& -\int_{C} \nabla_{\mu(t)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu_{0}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\left[D_{\tilde{c}} \Phi_{(0, t)}^{\mu^{0}}(\tilde{c}) \mathcal{F}(\tilde{c})+v_{\varepsilon}(s, \tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c}) .
\end{align*}
$$

By adding and subtracting to 4.25) the term

$$
A\left(t, \mu(t), \Phi_{(0, t)}^{\mu_{E}^{0}}(c), u\left(t, \Phi_{(0, t)}^{0_{E}^{0}}(c)\right)\right)+A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right) \mu_{\#}^{\left.\mu^{0} \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)}\right.
$$

we obtain

$$
\begin{equation*}
\left\|r_{\varepsilon}(t, c)\right\|_{E} \leq \mathrm{I}+\mathrm{II}+\mathrm{III}, \tag{4.26}
\end{equation*}
$$

where

$$
\begin{align*}
& \mathrm{I}:=\| \frac{1}{\varepsilon}\left\{A\left(t, \mu(t), \Phi_{(0, t)}^{\Phi_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)-A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\right\} \\
& -\mathrm{D}_{\Phi_{(0, t)}^{\mu^{0}}(c)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left[v_{\varepsilon}(t, c)\right] \|_{E},  \tag{4.27}\\
& \mathrm{II}:=\| \frac{1}{\varepsilon}\left\{A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{\mu^{0}, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)}\right)-A\left(t, \mu(t), \Phi_{(0, t)}^{\mu_{E}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)\right\}
\end{align*}
$$

$$
-\int_{C} \nabla_{\mu(t)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\left[\mathrm{D}_{\tilde{c}} \Phi_{(0, t)}^{\mu^{0}}(\tilde{c}) \mathcal{F}(\tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c}) \|_{E}(4.28)
$$

and

$$
\begin{align*}
\text { III }:= & \| \frac{1}{\varepsilon}\left\{A\left(t, \mu_{\varepsilon}(t), \Phi_{(0, t)}^{\mu_{(, t)}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)-A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{O}^{\varepsilon}\right)\right)_{\#}^{\left.\left.\mu^{0}, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{(, t)}^{0}}(c)\right)\right)\right\}}\right. \\
& -\int_{C} \nabla_{\mu(t)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu_{0}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\left[v_{\varepsilon}(t, \tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c}) \|_{E} . \quad(4 \tag{4.29}
\end{align*}
$$

Step 3. In this step we prove that $r_{\varepsilon}(t, c) \rightarrow 0$ as $\varepsilon \rightarrow 0$ in $E$ for any $(t, c) \in[0, T] \times V$. To estimate I, we rewrite the right-hand side of (4.27) in the following way

$$
\begin{aligned}
I= & \frac{1}{\|} \begin{aligned}
\mu_{(0, t)}^{0}(c)-\Phi_{(0, t)}^{\mu^{0}}(c) \|_{E}
\end{aligned} \| A\left(t, \mu(t), \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)-A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{0}^{0}}(c)\right)\right) \\
& \left.-D_{\Phi_{(0, t)}^{\mu^{0}}(c)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu_{0}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)\left[\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right] \|_{E}\right\}\left\|v_{\varepsilon}(t, c)\right\|_{E},
\end{aligned}
$$

then, it follows from (HA2)-(vi) and using (4.23) that $I \rightarrow 0$ as $\varepsilon \rightarrow 0$ for any $(t, c) \in$ $[0, T] \times V$.

The fact of II $\rightarrow 0$ as $\varepsilon \rightarrow 0$ follows from Lemma 4.5 and assumption (HA2)-(vii) which allow us to apply Proposition 4.3 in (4.28) with $\phi(\mu)=A(t, \mu, c, u)$ and $\psi=\Phi_{(0, t)}^{\mu^{0}}$. Thus we have

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0}\{I+I I\}=0 \tag{4.30}
\end{equation*}
$$

The estimate of term III is trickier. By (HA2)-(vii), recalling Definition 4.2 and that $\mu_{\varepsilon}(t)=\left(\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{0}$, we have

$$
\stackrel{4.20]}{=} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right)_{\#}^{\left.\mu^{0}, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)}\right.
$$

$$
\begin{equation*}
+o\left(|\varepsilon|\left(\int_{C}\left\|v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)\right\|_{E}^{2} \mathrm{~d} \mu^{0}(\tilde{c})\right)^{\frac{1}{2}}\right) \tag{4.31}
\end{equation*}
$$

where in the last equality we have chosen $\gamma=\left(\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}} \circ \mathcal{G}^{\varepsilon}, \Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{0}$. Combining (4.29) and (4.31), we deduce

$$
\begin{aligned}
& -\int_{C} \nabla_{\left.\left(\Phi_{(0, t)}^{\mu^{0}}\right)\right)_{\#}^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right){ }_{\#}^{\left.\mu^{0}, \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right)}\right)_{\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)}\left[v_{\varepsilon}(t, \tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c})
\end{aligned}
$$

$$
\begin{aligned}
& A\left(t, \mu_{\varepsilon}(t), \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)=A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right) \#_{\#}^{\left.\mu^{0}, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)}\right. \\
& +\int_{C \times C} \nabla_{\left.\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right)\right)_{\#}^{\mu^{0}}} A\left(t, \Phi_{(0, t)}^{{\mu^{0}}^{\mu^{\varepsilon}}}\right)_{\#} \mu^{\left.\mu^{0}, \Phi_{(0, t)}^{\mu_{\tilde{\varepsilon}}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c)\right)\right)\left(c_{1}\right)\left[c_{2}-c_{1}\right] \mathrm{d} \gamma\left(c_{1}, c_{2}\right) .} \\
& +o\left(W_{2, \gamma}\left(\left(\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{0},\left(\Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{0}\right)\right)
\end{aligned}
$$

$$
\begin{equation*}
+\frac{1}{\varepsilon} o\left(|\varepsilon|\left(\int_{C}\left\|v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)\right\|_{E}^{2} \mathrm{~d} \mu^{0}(\tilde{c})\right)^{\frac{1}{2}}\right) \|_{E} \tag{4.32}
\end{equation*}
$$

By adding and subtracting to the terms within the norm in 4.32) the term

$$
\left.\left.\int_{C} \nabla_{\left(\Phi_{(0, t)}^{\mu^{0}}\right){ }_{\#}{ }^{\mu^{0^{2}}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right)\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\left[v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)\right] \mathrm{d} \mu^{0}(\tilde{c})
$$

we obtain

$$
\begin{equation*}
\mathrm{III} \leq \mathrm{III}^{\prime}+\mathrm{III}^{\prime \prime}+\left\|\frac{1}{\varepsilon} o\left(|\varepsilon|\left(\int_{C}\left\|v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)\right\|_{E}^{2} \mathrm{~d} \mu^{0}(\tilde{c})\right)^{\frac{1}{2}}\right)\right\|_{E} \tag{4.33}
\end{equation*}
$$

where

$$
\begin{aligned}
& \left.\left.-\nabla_{\left(\Phi_{(0, t)}^{\mu^{0}}\right)_{\#}^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right)\right)_{\#}^{\mu^{0} \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)}\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\right)\left[v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)\right] \mathrm{d} \mu^{0}(\tilde{c}) \|_{E},
\end{aligned}
$$

and
III" :=

$$
\left.\| \int_{C} \nabla_{\left(\Phi_{(0, t)}^{\mu^{0}}\right)_{\neq}^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right)\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)}\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)\left[v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)-v_{\varepsilon}(t, \tilde{c})\right] \mathrm{d} \mu^{0}(\tilde{c}) \|_{E}
$$

First we focus on $\mathrm{III}^{\prime}$. Using 4.23) and applying Hölder inequality, we obtain for a positive constant $M$

$$
\begin{align*}
\mathrm{III}^{\prime} & \leq M\left(\int_{C} \| \nabla_{\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu_{\tilde{c}}^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)\right) \\
& \left.\left.\left.-\nabla_{\left(\Phi_{(0, t)}^{\mu^{0}}\right)_{\#} \mu^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right)\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right) \|_{\mathcal{L}\left(E_{C} ; E_{C}\right)}^{2} \mathrm{~d} \mu^{0}(\tilde{c})\right)^{\frac{1}{2}} \tag{4.34}
\end{align*}
$$

By the continuity assumption (HA2)-(vii) and recalling that in Step 1 we have seen that $W_{1}\left(\left(\Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#} \mu^{0},\left(\Phi_{(0, t)}^{\mu^{0}}\right)_{\#} \mu^{0}\right) \rightarrow 0$ as $\varepsilon \rightarrow 0$, we deduce

$$
\begin{aligned}
& \left.\left.\lim _{\varepsilon \rightarrow 0}\left\{\nabla_{\left(\Phi_{(0, t)}^{\mu^{0}} \circ \mathcal{G}^{\varepsilon}\right)_{\#}^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right)\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)\right)\right\} \\
& \left.\stackrel{\mathcal{L}\left(E_{C} ; E_{C}\right)}{=} \nabla_{\left(\Phi_{(0, t)}^{\mu^{0}}\right) \#_{\#}^{\mu^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right) \#_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right)
\end{aligned}
$$

uniformly for $\tilde{c} \in \operatorname{supp}\left(\mu^{0}\right)$ (recall that $\operatorname{supp}\left(\mu^{0}\right)$ is compact in $C$ ). Moreover, by Lemma 4.8, it follows that

$$
\begin{aligned}
& \left.\left.\| \nabla_{\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right) \#_{\#}^{\mu^{0^{0}}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}} \mathcal{G}^{\varepsilon}\right)\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu^{(0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)\right) \\
& \left.\left.-\nabla_{\left(\Phi_{(0, t)}^{\mu^{0}}\right){ }_{\#}{ }^{0^{0}}} A\left(t,\left(\Phi_{(0, t)}^{\mu^{0}}\right)\right)_{\#}^{\mu^{0}, \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right.}\right)\right)\left(\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right) \|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \leq 2 L .
\end{aligned}
$$

Thus, after an application of the Lebesgue theorem for general measures to the right-hand side of (4.34), we infer that

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \mathrm{III}^{\prime}=0 \tag{4.35}
\end{equation*}
$$

As for III", using again Lemma 4.8, we infer

$$
\begin{equation*}
\mathrm{III}^{\prime \prime} \leq M \int_{C}\left\|v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)-v_{\varepsilon}(t, \tilde{c})\right\|_{E} \mathrm{~d} \mu^{0}(\tilde{c}) \tag{4.36}
\end{equation*}
$$

for a positive constant $M$. Moreover, by assumption $(i)-(i i)$ and thanks to the property of uniform $C$-differentiability of $\Phi_{(0, t)}^{\mu^{0}}$ and $\Phi_{(0, t)}^{\mu_{\mathrm{e}}^{0}}$ on $\operatorname{supp}\left(\mu^{0}\right)$ given by Lemma 4.5, we have as $\varepsilon \rightarrow 0$

$$
\begin{aligned}
& v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)-v_{\varepsilon}(t, \tilde{c}) \stackrel{\stackrel{4.20}{=}}{ } \frac{\Phi_{(0, t)}^{\mu^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)-\Phi_{(0, t)}^{\mu^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)}{\varepsilon}-\frac{\Phi_{(0, t)}^{\mu_{0}^{0}}(\tilde{c})-\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})}{\varepsilon} \\
&=\frac{1}{\varepsilon}\left\{\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)-\Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}(\tilde{c})-\Phi_{(0, t)}^{\mu^{0}}\left(\mathcal{G}^{\varepsilon}(\tilde{c})\right)+\Phi_{(0, t)}^{\mu^{0}}(\tilde{c})\right\} \\
&=\frac{1}{\varepsilon}\left\{\varepsilon \mathrm{D}_{\tilde{c}} \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}[\mathcal{F}(\tilde{c})]-\varepsilon \mathrm{D}_{\tilde{c}} \Phi_{(0, t)}^{\mu^{0}}[\mathcal{F}(\tilde{c})]+o(\varepsilon)\right\} \\
&=\left(\mathrm{D}_{\tilde{c}} \Phi_{(0, t)}^{\mu_{\varepsilon}^{0}}-\mathrm{D}_{\tilde{c}} \Phi_{(0, t)}^{\mu^{0}}\right)[\mathcal{F}(\tilde{c})]+\frac{o(\varepsilon)}{\varepsilon} \xrightarrow{E} 0,
\end{aligned}
$$

where in the limit we have applied Lemma 4.6. Thanks to 4.23) and applying the Lebesgue theorem to the right-hand side of (4.36), we obtain

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} \mathrm{III}^{\prime \prime}=0 \tag{4.37}
\end{equation*}
$$

Finally, using again (4.23), we infer that

$$
\lim _{\varepsilon \rightarrow 0}\left\{\frac{1}{\varepsilon} o\left(|\varepsilon|\left(\int_{C}\left\|v_{\varepsilon}\left(t, \mathcal{G}^{\varepsilon}(\tilde{c})\right)\right\|_{E}^{2} \mathrm{~d} \mu^{0}(\tilde{c})\right)^{\frac{1}{2}}\right)\right\}=0
$$

Hence, combining the last limit with (4.33), 4.35) and 4.37), we conclude that

$$
\lim _{\varepsilon \rightarrow 0} \mathrm{III}=0
$$

This, together with 4.26) and 4.30, implies that

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0} r_{\varepsilon}(t, c)=0 \quad \forall(t, c) \in[0, T] \times V . \tag{4.38}
\end{equation*}
$$

Step 4. In this last step we prove the result.
Using (HA1)- $(i)$, the fact that $W_{1}\left(\mu_{\varepsilon}(t), \mu(t)\right) \leq M|\varepsilon|$ (as seen in Step 1) and 4.20) in the first term of the right-hand side of (4.25), inequalities (4.10) and (4.23) in the second term and Lemma 4.8, Lemma 4.5, assumption (ii) and again 4.23) in the third term respectively, we deduce for a positive constant $M$ that

$$
\left\|r_{\varepsilon}(t, c)\right\|_{E} \leq M \quad \forall(t, c) \in[0, T] \times V .
$$

Therefore, by 4.38) and applying the Lebesgue theorem, we infer that

$$
\begin{equation*}
r_{\varepsilon}(t, c) \rightarrow 0 \quad \text { in } L_{\mathcal{L} \times \mu^{0}}^{1}([0, T] \times C ; E) \quad \text { as } \varepsilon \rightarrow 0 . \tag{4.39}
\end{equation*}
$$

Finally, we estimate 4.24). By (4.10) and Lemma 4.8, we deduce

$$
\begin{align*}
\left\|v_{\varepsilon}(t, c)-v(t, c)\right\|_{E} \leq & L \int_{0}^{t}\left\|v_{\varepsilon}(s, c)-v(s, c)\right\|_{E} \mathrm{~d} s \\
& +L \int_{0}^{t} \int_{C}\left\|v_{\varepsilon}(s, \tilde{c})-v(s, \tilde{c})\right\|_{E} \mathrm{~d} \mu^{0}(\tilde{c}) \mathrm{d} s \\
& +\int_{0}^{t}\left\|r_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} s . \tag{4.40}
\end{align*}
$$

Integrating over $C$, recalling that $\mu^{0}(C)=1$ and applying Fubini theorem, we have

$$
\begin{aligned}
\int_{C}\left\|v_{\varepsilon}(t, c)-v(t, c)\right\|_{E} \mathrm{~d} \mu^{0}(c) \leq & 2 L \int_{0}^{t} \int_{C}\left\|v_{\varepsilon}(s, c)-v(s, c)\right\|_{E} \mathrm{~d} \mu^{0}(c) \mathrm{d} s \\
& +\int_{0}^{t} \int_{C}\left\|r_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} \mu^{0}(c) \mathrm{d} s
\end{aligned}
$$

which implies, applying Grönwall inequality and using (4.39), that
$\int_{C}\left\|v_{\varepsilon}(t, c)-v(t, c)\right\|_{E} \mathrm{~d} \mu^{0}(c) \leq \delta_{\varepsilon} \mathrm{e}^{2 L t}, \quad \delta_{\varepsilon}:=\int_{0}^{T} \int_{C}\left\|r_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} \mu^{0}(c) \mathrm{d} s \rightarrow 0$ as $\varepsilon \rightarrow 0$.
Inserting this last inequality in (4.40), we obtain for every $(t, c) \in[0, T] \times V$

$$
\left\|v_{\varepsilon}(t, c)-v(t, c)\right\|_{E} \leq L \int_{0}^{t}\left\|v_{\varepsilon}(s, c)-v(s, c)\right\|_{E} \mathrm{~d} s+\delta_{\varepsilon} \mathrm{e}^{2 L T}+\int_{0}^{t}\left\|r_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} s
$$

Hence, setting

$$
\hat{\delta}_{\varepsilon}:=\delta_{\varepsilon} \mathrm{e}^{2 L T}+\int_{0}^{T}\left\|r_{\varepsilon}(s, c)\right\|_{E} \mathrm{~d} s
$$

recalling that, by (4.39), $\hat{\delta}_{\varepsilon} \rightarrow 0$ uniformly in $V$ as $\varepsilon \rightarrow 0$, and applying again Grönwall inequality, we infer

$$
\left\|v_{\varepsilon}(t, c)-v(t, c)\right\|_{E} \leq \hat{\delta}_{\varepsilon} \mathrm{e}^{L t}
$$

which gives 4.21.
Proof of Theorem 4.7. We divide the proof into steps.
Step 1: Needle variations. Fix any time $\tau \in(0, T]$ and any admissible control value $\omega \in U$. Up to a null set we can assume that $\tau$ is a Lebesgue point of $t \mapsto A(t, \boldsymbol{\mu}(t), c, z)$ for all $z \in Z$. Notice that this is possible by separability of $Z$ and the uniform Lipschitz continuity of $A$ in the last variable. For $\varepsilon \in[0, \bar{\varepsilon})$ with $\bar{\varepsilon}>0$ sufficiently small, we consider the following needle variation:

$$
u_{\varepsilon}(t)= \begin{cases}\omega & \text { if } t \in[\tau-\varepsilon, \tau]  \tag{4.41}\\ \mathfrak{u}(t) & \text { otherwise }\end{cases}
$$

Thanks to (HA1), by Theorem 3.3 of [6], there exists a unique solution $\mu_{\varepsilon} \in A C\left([0, T] ; \mathcal{P}_{c}(C)\right)$ of

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \mu(t)+\operatorname{div}\left(A\left(t, \mu(t), \cdot, u_{\varepsilon}(t, \cdot)\right) \mu(t)\right)=0 \quad \text { in }(0, T] \\
\mu(0)=\mu^{0} \in \mathcal{P}_{c}(C)
\end{array}\right.
$$

The curve $\mu_{\varepsilon}(t)$ can be represented using the associated family of non-local flows as

$$
\mu_{\varepsilon}(t)=\left(\Phi_{(0, t)}^{\varepsilon, \mu^{0}}\right)_{\#} \mu^{0} \quad \text { for } t \in[0, T] .
$$

Moreover, by Proposition 4 of [1], $R>0$ exists such that $\operatorname{supp}(\boldsymbol{\mu}(t)) \cup \operatorname{supp}\left(\mu_{\varepsilon}(t)\right) \subset B_{R, C}$ and $\left\|\Phi_{(0, t)}^{\varepsilon, \mu^{0}}\right\|_{C_{b}^{0}(C ; C)} \leq R$ for every $(\varepsilon, t) \in[0, \bar{\varepsilon}) \times[0, T]$.

Let $t \in(\tau, T]$. We want to compute for every $(t, c) \in(\tau, T] \times B_{R, C}$

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0^{+}} \frac{\Phi_{(0, t)}^{\varepsilon, \mu^{0}} \circ \Phi_{(t, 0)}^{\boldsymbol{\mu}(t)}(c)-\operatorname{Id}(c)}{2^{\varepsilon}} \tag{4.42}
\end{equation*}
$$

where Id: $C \rightarrow C$ is the identity function. It follows from (4.41) and from the definition of $\Phi_{(s, t)}^{\varepsilon, \mu_{\varepsilon}(s)}$ and $\Phi_{(s, t)}^{\mu(s)}(0 \leq s<t \leq T)$ that

$$
\begin{equation*}
\Phi_{(0, t)}^{\varepsilon, \mu^{0}} \circ \Phi_{(t, 0)}^{\mu(t)}(c)=\Phi_{(\tau, t)}^{\mu_{\varepsilon}(\tau)} \circ \Phi_{(\tau-\varepsilon, \tau)}^{\varepsilon, \mu(\tau-\varepsilon)} \circ \Phi_{(\tau, \tau-\varepsilon)}^{\mu(\tau)}\left(c_{\tau}\right) \tag{4.43}
\end{equation*}
$$

where $c_{\tau}:=\Phi_{(t, \tau)}^{\mu(t)}(c)$. Using again (HA1), we deduce that the map

$$
\sigma \mapsto A\left(\sigma, \boldsymbol{\mu}(\sigma), \Phi_{(\sigma, \tau-\varepsilon)}^{\mu(\tau)}\left(c_{\tau}\right), \mathfrak{u}\left(\sigma, \Phi_{(\sigma, \tau-\varepsilon)}^{\mu(\tau)}\left(c_{\tau}\right)\right)\right) \in L^{\infty}([0, T] ; E) .
$$

Then, after an application of the Lebesgue differentiation theorem for vector-valued functions, we obtain

$$
\begin{aligned}
& \Phi_{(\tau, \tau-\varepsilon)}^{\mu(\tau)}\left(c_{\tau}\right) \stackrel{4.88}{=} c_{\tau}-\int_{\tau-\varepsilon}^{\tau} A\left(\sigma, \boldsymbol{\mu}(\sigma), \Phi_{(\sigma, \tau-\varepsilon)}^{\mu(\tau)}\left(c_{\tau}\right), \mathfrak{u}\left(\sigma, \Phi_{(\sigma, \tau-\varepsilon)}^{\mu(\tau)}\left(c_{\tau}\right)\right)\right) \mathrm{d} \sigma \\
&=c_{\tau}-\varepsilon A\left(\tau, \boldsymbol{\mu}(\tau), c_{\tau}, \mathfrak{u}\left(\tau, c_{\tau}\right)\right)+\delta_{\varepsilon}^{1},
\end{aligned}
$$

with $\frac{\delta_{\varepsilon}^{1}}{\varepsilon} \rightarrow 0$ in $E$ as $\varepsilon \rightarrow 0^{+}$. Now, choose $d=c_{\tau}-\varepsilon A\left(\tau, \boldsymbol{\mu}(\tau), c_{\tau}, \mathfrak{u}\left(\tau, c_{\tau}\right)\right)+\delta_{\varepsilon}^{1}$. Observe that

$$
\frac{1}{\varepsilon} \int_{\tau-\varepsilon}^{\tau}\left[W_{1}\left(\mu_{\varepsilon}(\sigma), \boldsymbol{\mu}(\tau)\right)+\left\|\Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \boldsymbol{\mu}(\tau-\varepsilon)}(d)-c_{\tau}\right\|_{E}+\left\|\omega\left(\Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \boldsymbol{\mu}(\tau-\varepsilon)}(d)\right)-\omega\left(c_{\tau}\right)\right\|_{Z}\right] \mathrm{d} \sigma \rightarrow 0
$$

Using (HA1)-(i) we deduce that

$$
\frac{1}{\varepsilon} \int_{\tau-\varepsilon}^{\tau} A\left(\sigma, \mu_{\varepsilon}(\sigma) \Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \mu(\tau-\varepsilon)}(d), \omega\left(\Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \mu(\tau-\varepsilon)}(d)\right)\right)-A\left(\sigma, \mu_{\varepsilon}(\sigma), \Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \mu(\tau-\varepsilon)}(d), \omega\left(c_{\tau}\right)\right) \mathrm{d} \sigma \rightarrow 0
$$

With the above equalities and the Lebesgue differentiation theorem we get

$$
\begin{aligned}
\Phi_{(\tau-\varepsilon, \tau)}^{\varepsilon, \mu(\tau-\varepsilon)}(d) \stackrel{\boxed{4.8}}{=} & d+\int_{\tau-\varepsilon}^{\tau} A\left(\sigma, \mu_{\varepsilon}(\sigma), \Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \mu(\tau) \varepsilon}(d), \omega\left(\Phi_{(\tau-\varepsilon, \sigma)}^{\varepsilon, \mu(\tau) \varepsilon}(d)\right)\right) \mathrm{d} \sigma \\
& =c_{\tau}-\varepsilon A\left(\tau, \boldsymbol{\mu}(\tau), c_{\tau}, \mathfrak{u}\left(\tau, c_{\tau}\right)\right)+\varepsilon A\left(\tau, \boldsymbol{\mu}(\tau), c_{\tau}, \omega\left(c_{\tau}\right)\right)+\delta_{\varepsilon}^{2},
\end{aligned}
$$

with $\frac{\delta_{\varepsilon}^{2}}{\varepsilon} \rightarrow 0$ in $E$ as $\varepsilon \rightarrow 0^{+}$. Recalling the definition of $d$, this amounts to

$$
\begin{equation*}
\Phi_{(\tau-\varepsilon, \tau)}^{\varepsilon, \boldsymbol{\mu}(\tau-\varepsilon)} \circ \Phi_{(\tau, \tau-\varepsilon)}^{\boldsymbol{\mu}(\tau)}\left(c_{\tau}\right)=c_{\tau}+\varepsilon\left(A\left(\tau, \boldsymbol{\mu}(\tau), c_{\tau}, \omega\left(c_{\tau}\right)\right)-A\left(\tau, \boldsymbol{\mu}(\tau), c_{\tau}, \mathfrak{u}\left(\tau, c_{\tau}\right)\right)\right)+\delta_{\varepsilon} \tag{4.44}
\end{equation*}
$$

with $\frac{\delta_{\varepsilon}}{\varepsilon} \rightarrow 0$ in $E$ as $\varepsilon \rightarrow 0^{+}$. Recalling that $\boldsymbol{\mu}(t)$ and $\mu_{\varepsilon}(t)$ can be represented using the associated non-local flows, we have the following expression which links $\mu_{\varepsilon}(\tau)$ and $\boldsymbol{\mu}(\tau)$ :

$$
\begin{align*}
& \mu_{\varepsilon}(\tau)=\left(\Phi_{(\tau-\varepsilon, \tau)}^{\varepsilon, \boldsymbol{\mu}(\tau-\varepsilon)}\right)_{\#} \boldsymbol{\mu}(\tau-\varepsilon)=\left(\Phi_{(\tau-\varepsilon, \tau)}^{\varepsilon, \boldsymbol{\mu}(\tau-\varepsilon)} \circ \Phi_{(\tau, \tau-\varepsilon)}^{\boldsymbol{\mu}(\tau)}\right)_{\#} \boldsymbol{\mu}(\tau) \\
& \stackrel{\text { 4.44] }}{=}\left(\operatorname{Id}+\varepsilon(A(\tau, \boldsymbol{\mu}(\tau), \cdot, \omega(\cdot))-A(\tau, \boldsymbol{\mu}(\tau), \cdot \cdot, \mathfrak{u}(\tau, \cdot)))+\delta_{\varepsilon}\right)_{\#} \boldsymbol{\mu}(\tau) . \tag{4.45}
\end{align*}
$$

We then define for fixed $c \in B_{R, C}$ the function $\mathcal{G}:[0, \bar{\varepsilon}) \rightarrow E$ as

$$
\begin{equation*}
\mathcal{G}(\varepsilon, \tau, c):=\operatorname{Id}+\varepsilon(A(\tau, \boldsymbol{\mu}(\tau), c, \omega(c))-A(\tau, \boldsymbol{\mu}(\tau), c, \mathfrak{u}(\tau, c)))+\delta_{\varepsilon}, \tag{4.46}
\end{equation*}
$$

and, since $\mathcal{G}$ is Fréchet-differentiable from the right with differential

$$
\begin{equation*}
\mathcal{F}(\tau, c):=A(\tau, \boldsymbol{\mu}(\tau), c, \omega(c))-A(\tau, \boldsymbol{\mu}(\tau), c, \mathfrak{u}(\tau, c)), \tag{4.47}
\end{equation*}
$$

applying Lemma 2.11 of 35 we can extend $\mathcal{G}$ from $[0, \bar{\varepsilon})$ to $(-\bar{\varepsilon}, \bar{\varepsilon})$ preserving the Fréchetdifferentiability at $\varepsilon=0$. Consequently, by (HA1)-(ii), we have for $\varepsilon \in(-\bar{\varepsilon}, \bar{\varepsilon})$

$$
\left\{\begin{array}{l}
\mathcal{G}(0, \tau, \cdot)=\mathrm{Id}  \tag{4.48}\\
\left.\frac{\mathrm{~d}}{\mathrm{~d} \varepsilon} \mathcal{G}(\varepsilon, \tau, c)\right|_{\varepsilon=0}=\mathcal{F}(\tau, c) \in L_{\boldsymbol{\mu}(\tau)}^{\infty}\left(B_{R, C} ; E\right), \\
\mu_{\varepsilon}(\tau) \stackrel{\boxed{4.45 \mid}}{=} \mathcal{G}(\varepsilon, \tau, \cdot)_{\#}^{\boldsymbol{\mu}(\tau)}
\end{array}\right.
$$

Therefore, combining (4.42), (4.43), (4.44) and (4.46), we infer that

$$
\begin{equation*}
\frac{\Phi_{(0, t)}^{\varepsilon, \mu^{0}} \circ \Phi_{(t, 0)}^{\mu(t)}(c)-\operatorname{Id}(c)}{\varepsilon}=\frac{\Phi_{(\tau, t)}^{\mu_{\varepsilon}(\tau)}\left(\mathcal{G}\left(\varepsilon, \tau, c_{\tau}\right)\right)-\Phi_{(\tau, t)}^{\mu(\tau)}\left(c_{\tau}\right)}{\varepsilon} \tag{4.49}
\end{equation*}
$$

for every $(\varepsilon, t, c) \in(-\bar{\varepsilon}, \bar{\varepsilon}) \times(\tau, T] \times B_{R, C}$. Thanks to 4.48) we can apply Proposition 4.9 and Lemma 4.5 obtaining

$$
\begin{align*}
\Phi_{(\tau, t)}^{\mu_{\varepsilon}(\tau)}\left(\mathcal{G}\left(\varepsilon, \tau, c_{\tau}\right)\right) & =\Phi_{(\tau, t)}^{\mu(\tau)}\left(\mathcal{G}\left(\varepsilon, \tau, c_{\tau}\right)\right)+\varepsilon v\left(t, \mathcal{G}\left(\varepsilon, \tau, c_{\tau}\right)\right)+\hat{\delta}_{\varepsilon} \\
& =\Phi_{(\tau, t)}^{\mu(\tau)}\left(c_{\tau}\right)+\varepsilon\left(\mathrm{D}_{c_{\tau}} \Phi_{(\tau, t)}^{\mu(\tau)}\left[\mathcal{F}\left(\tau, c_{\tau}\right)\right]+v\left(t, \mathcal{G}\left(\varepsilon, \tau, c_{\tau}\right)\right)\right)+\tilde{\delta}_{\varepsilon}, 4 \tag{4.50}
\end{align*}
$$

where $v(t, \cdot)$ is the Fréchet derivative at $\varepsilon=0$ of $\Phi_{(\tau, t)}^{\mu_{\varepsilon}(\tau)}(\cdot)$ and $\frac{\hat{\delta}_{\varepsilon}}{\varepsilon}, \frac{\tilde{\delta}_{\varepsilon}}{\varepsilon} \rightarrow 0$ as $\varepsilon \rightarrow 0^{+}$. Hence, by the continuity of $\mathcal{G}$ with respect to $\varepsilon$ and $v$ with respect to $c$ given by Lemma B.1, we conclude that for every $(t, c) \in(\tau, T] \times B_{R, C}$

$$
\lim _{\varepsilon \rightarrow 0^{+}} \frac{\Phi_{(0, t)}^{\varepsilon, \mu^{0}} \circ \Phi_{(t, 0)}^{\mu(t)}(c)-\operatorname{Id}(c)}{\varepsilon} \stackrel{4.49), \sqrt[4.50]{=}}{\mathrm{D}_{c_{\tau}} \Phi_{(\tau, t)}^{\mu(\tau)}\left[\mathcal{F}\left(\tau, c_{\tau}\right)\right]+v\left(t, c_{\tau}\right) .}
$$

We define $\mathcal{F}:(\tau, T] \times B_{R, C} \rightarrow E_{C}$ as

$$
\begin{array}{rlll}
\mathcal{F}(t, c) & := & \mathrm{D}_{c} \Phi_{(\tau, t)}^{\boldsymbol{\mu}(\tau)}[\mathcal{F}(\tau, c)]+v(t, c) \\
& \stackrel{4.47}{=} & \mathrm{D}_{c} \Phi_{(\tau, t)}^{\boldsymbol{\mu}(\tau)}[A(\tau, \boldsymbol{\mu}(\tau), c, \omega(c))-A(\tau, \boldsymbol{\mu}(\tau), c, \mathfrak{u}(\tau, c))]+v(t, c), \tag{4.51}
\end{array}
$$

thus obtaining

$$
\begin{equation*}
\lim _{\varepsilon \rightarrow 0^{+}} \frac{\Phi_{(0, t)}^{\varepsilon, \mu^{0}} \circ \Phi_{(t, 0)}^{\mu(t)}(c)-\operatorname{Id}(c)}{\varepsilon} \stackrel{4.51)}{=} \mathcal{F}\left(t, \Phi_{(t, \tau)}^{\boldsymbol{\mu}(t)}(c)\right) . \tag{4.52}
\end{equation*}
$$

Finally, again applying Lemma 4.5 and Proposition 4.9, we characterize $\mathcal{F}(t, c)$ as the unique solution of the following linear ODE defined in $E_{C}$ for $t \in(\tau, T]$ :

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \mathcal{F}(t, c)=\mathrm{D}_{\Phi_{(\tau, t)}^{\mu(\tau)}(c)} A\left(t, \boldsymbol{\mu}(t), \Phi_{(\tau, t)}^{\mu(\tau)}(c), \mathfrak{u}\left(t, \Phi_{(\tau, t)}^{\boldsymbol{\mu}(\tau)}(c)\right)\right)[\mathcal{F}(t, c)]  \tag{4.53}\\
\quad+\int_{C} \nabla_{\boldsymbol{\mu}(t)} A\left(t, \boldsymbol{\mu}(t), \Phi_{(\tau, t)}^{\left.\boldsymbol{\mu}(\tau)(c), \mathfrak{u}\left(t, \Phi_{(\tau, t)}^{\mu(\tau)}(c)\right)\right)\left(\Phi_{(\tau, t)}^{\mu(\tau)}(\tilde{c})\right)[\mathcal{F}(t, \tilde{c})] \mathrm{d} \boldsymbol{\mu}(\tau)(\tilde{c}),}\right. \\
\mathcal{F}(\tau, c) \stackrel{(4,47]}{=} A(\tau, \boldsymbol{\mu}(\tau), c, \omega(c))-A(\tau, \boldsymbol{\mu}(\tau), c, \mathfrak{u}(\tau, c)) .
\end{array}\right.
$$

Step 2: Optimality condition. Note that, by the definitions of $\boldsymbol{\mu}(t)$ and $\mu_{\varepsilon}(t)$ and using (4.52), we deduce

$$
\begin{equation*}
\mu_{\varepsilon}(t)=\left(\Phi_{(0, t)}^{\varepsilon, \mu^{0}} \circ \Phi_{(t, 0)}^{\boldsymbol{\mu}(t)}\right)_{\#} \boldsymbol{\mu}(t)=\left(\operatorname{Id}+\varepsilon \mathcal{F}\left(t, \Phi_{(t, \tau)}^{\boldsymbol{\mu}(t)}(\cdot)\right)+\delta_{\varepsilon}\right)_{\#} \boldsymbol{\mu}(t), \tag{4.54}
\end{equation*}
$$

for every $t \in(\tau, T]$, where $\frac{\delta_{\varepsilon}}{\varepsilon} \rightarrow 0$ as $\varepsilon \rightarrow 0^{+}$. Thus, it follows from the assumption ( $\mathbf{H} \varphi$ ) and using (4.54) with $t=\stackrel{\varepsilon}{T}$ that we can apply Proposition 4.3 obtaining

$$
\begin{equation*}
\varphi\left(\mu_{\varepsilon}(T)\right)=\varphi(\boldsymbol{\mu}(T))+\varepsilon \int_{C}\left\langle\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c), \mathcal{F}\left(T, \Phi_{(T, \tau)}^{\boldsymbol{\mu}(T)}(c)\right)\right\rangle \mathrm{d} \boldsymbol{\mu}(T)(c)+\delta_{\varepsilon}^{1}, \tag{4.55}
\end{equation*}
$$

where $\frac{\delta_{\varepsilon}^{1}}{\varepsilon} \rightarrow 0$ as $\varepsilon \rightarrow 0^{+}$. By the optimality condition on $\mathfrak{u} \in \mathcal{U}$, we know that

$$
\frac{\varphi\left(\mu_{\varepsilon}(T)\right)-\varphi(\boldsymbol{\mu}(T))}{\varepsilon} \geq 0 \quad \text { for every } \varepsilon>0
$$

hence, using (4.55) and letting $\varepsilon$ go to 0 , we infer

$$
\begin{equation*}
\int_{C}\left\langle\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c), \mathcal{F}\left(T, \Phi_{(T, \tau)}^{\boldsymbol{\mu}(T)}(c)\right)\right\rangle \mathrm{d} \boldsymbol{\mu}(T)(c) \geq 0 \tag{4.56}
\end{equation*}
$$

Step 3: Adjoint problem. In this Step we follow Step 3 of Section 3.2 in [18]. We consider the unique distributional solution $\psi_{t}(c, p):[0, T] \times C \times E_{C}^{*} \rightarrow E_{C}^{*}$ of the adjoint equation of (4.53) (which is a linear ODE defined on $E_{C}^{*}$ ):

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} w(t, c, p)=-\mathrm{D}_{\Phi_{(T, t)}^{\mu(T)}(c)}^{*} A\left(t, \boldsymbol{\mu}(t), \Phi_{(T, t)}^{\boldsymbol{\Phi}^{\mu}(T)}(c), \mathfrak{u}\left(t, \Phi_{(T, t)}^{\boldsymbol{\mu}(T)}(c)\right)\right)[w(t, c, p)]  \tag{4.57}\\
-\int_{C \times E_{C}^{*}} \nabla_{\boldsymbol{\mu}(t)}^{*} A\left(t, \boldsymbol{\mu}(t), \Phi_{(T, t)}^{\boldsymbol{\mu}(T)}(\tilde{c}), \mathfrak{u}\left(t, \Phi_{(T, t)}^{\left.\boldsymbol{\mu}_{(T)}^{(T)}(\tilde{c})\right)}\right)\left(\Phi_{(T, t)}^{\boldsymbol{\mu}(T)}(c)\right)[w(t, \tilde{c}, \tilde{p})] \mathrm{d}\left(\left(\mathrm{Id},-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))\right)_{\neq} \boldsymbol{\mu ( T )}\right)(\tilde{c}, \tilde{p}),\right. \\
w(T, c, p)=p \in E_{C}^{*} .
\end{array}\right.
$$

Note that, since $D_{c} A(t, \mu, c, u)$ and $\nabla_{\mu} A(t, \mu, \tilde{c}, u)(c)$ are uniformly bounded operators in $\mathcal{L}\left(E_{C} ; E_{C}\right)$, then their adjoint operators are uniformly bounded in $\mathcal{L}\left(E_{C}^{*} ; E_{C}^{*}\right)$. Hence, applying Lemma B.1, there exists a unique distributional solution $\psi_{t}(c, p) \in A C\left([0, T] ; E_{C}^{*}\right)$ of 4.57). Moreover, $\psi_{t}(c, p)$ is continuous with respect to $c$ and $p$. Thus, we define a curve of measure in $E_{C}^{*}$ as

$$
\begin{equation*}
\sigma_{c}: t \mapsto \psi_{t}(c, \cdot)_{\#} \delta_{-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c)} . \tag{4.58}
\end{equation*}
$$

Since $\delta_{-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c)} \in \mathcal{P}_{c}\left(E_{C}^{*}\right)$, we have $\sigma_{c} \in A C\left([0, T] ; \mathcal{P}_{c}\left(E_{C}^{*}\right)\right)$. Now we extend this curve of measures on $C \times E_{C}^{*}$. First we define

$$
\begin{equation*}
\nu_{T}: t \mapsto \int_{C} \sigma_{c}(t) \mathrm{d} \boldsymbol{\mu}(T)(c) \tag{4.59}
\end{equation*}
$$

which belongs to $A C\left([0, T] ; \mathcal{P}_{c}\left(C \times E_{C}^{*}\right)\right)$ since $\mu(T)$ and $\sigma_{c}(t)$ have compact support in $C$ and $E_{C}^{*}$ respectively. Then we set

$$
\begin{equation*}
\boldsymbol{\nu}: t \mapsto\left(\Phi_{(T, t)}^{\boldsymbol{\mu}(T)} \circ \pi^{1}, \pi^{2}\right)_{\#} \nu_{T}(t), \tag{4.60}
\end{equation*}
$$

thus obtaining $\boldsymbol{\nu} \in A C\left([0, T] ; \mathcal{P}_{c}\left(C \times E_{C}^{*}\right)\right)$. By its definition we deduce that

$$
\begin{array}{rcl}
\mathrm{d} \boldsymbol{\nu}(T)(c, p) & \stackrel{4.60}{=} & \mathrm{d} \nu_{T}(T)(c, p) \stackrel{\boxed{4.59}}{=} \mathrm{d} \sigma_{c}(T)(p) \mathrm{d} \boldsymbol{\mu}(T)(c) \\
& \begin{array}{l}
4.58), \sqrt[4.57]{=} \\
\\
=
\end{array} & \mathrm{d} \delta_{-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c)}(p) \mathrm{d} \boldsymbol{\mu}(T)(c) \\
& \mathrm{d}\left(\left(I d,-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))\right)_{\#} \boldsymbol{\mu}(T)\right)(c, p) . \tag{4.61}
\end{array}
$$

Moreover, we have

$$
\begin{equation*}
\pi_{\#}^{1} \boldsymbol{\nu}(t) \stackrel{\sqrt{4.60}, \sqrt{4.59}}{=}\left(\Phi_{(T, t)}^{\boldsymbol{\mu}(T)}\right){ }_{\#} \boldsymbol{\mu}(T)=\boldsymbol{\mu}(t) \quad \text { for every } t \in[0, T] \tag{4.62}
\end{equation*}
$$

Now we prove that $\boldsymbol{\nu}$ is a distributional solution of

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \nu(t)=-\operatorname{div}_{(c, p)}(\Gamma(t, \nu(t), c, p, \mathfrak{u}(t, c)) \nu(t)) \quad \text { in }[0, T),  \tag{4.63}\\
\pi_{\#}^{1} \nu(t)=\boldsymbol{\mu}(t), \\
\nu(T)=\left(I d,-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))\right)_{\#} \boldsymbol{\mu}(T),
\end{array}\right.
$$

where $\Gamma=\left(\Gamma_{1}, \Gamma_{2}\right):[0, T] \times \mathcal{P}_{b}\left(C \times E_{C}^{*}\right) \times C \times E_{C}^{*} \times U \rightarrow E_{C} \times E_{C}^{*}$ is defined componentwise by

$$
\begin{align*}
\Gamma_{1}(t, \nu, c, p, u) & :=A\left(t, \pi_{\#}^{1} \nu, c, u\right)  \tag{4.64}\\
\Gamma_{2}(t, \nu, c, p, u) & :=-D_{c}^{*} A\left(t, \pi_{\#}^{1} \nu, c, u\right)[p]-\int_{C \times E_{C}^{*}} \nabla_{\pi_{\#}^{1} \nu}^{*} A\left(t, \pi_{\#}^{1} \nu, \tilde{c}, u\right)(c)[\tilde{p}] \mathrm{d} \nu(\tilde{c}, \tilde{p}) .
\end{align*}
$$

Let $\xi \in C_{c}^{\infty}\left(E \times E_{C}^{*} ; \mathbb{R}\right)$ with differential $\left(\mathrm{D}_{E} \xi, \mathrm{D}_{E_{C}^{*}} \xi\right)$. Then

$$
\sqrt{4.611, \sqrt[4.58)]{ }, \sqrt{4.599}, \sqrt{4.60}} \int_{C \times E_{C^{*}}}\left\langle\mathrm{D}_{E} \xi(c, p), A(t, \boldsymbol{\mu}(t), c, \boldsymbol{u}(t, c))\right\rangle \mathrm{d} \boldsymbol{\nu}(t)(c, p)
$$

$$
-\int_{C \times E_{C^{*}}}\left\langle\mathrm{D}_{E_{C}^{*}} \xi(c, p), \mathrm{D}_{c}^{*} A(t, \boldsymbol{\mu}(t), c, \boldsymbol{u}(t, c))[p]+\int_{C \times E_{C}^{*}} \nabla_{\boldsymbol{\mu}(t)}^{*} A(t, \boldsymbol{\mu}(t), \tilde{c}, \boldsymbol{u}(t, \tilde{c}))(c)[\tilde{p}] \mathrm{d} \boldsymbol{\nu}(t)(\tilde{c}, \tilde{p})\right\rangle \mathrm{d} \boldsymbol{\nu}(t)(c, p)
$$

$$
\xlongequal{4.62), \sqrt{4.64}} \int_{C \times E_{C^{*}}}\left\langle\mathrm{D}_{E} \xi(c, p), \Gamma_{1}(t, \boldsymbol{\nu}(t), c, p, \mathfrak{u}(t, c))\right\rangle \mathrm{d} \boldsymbol{\nu}(t)(c, p)
$$

$$
+\int_{C \times E_{C}^{*}}\left\langle\mathrm{D}_{E_{C}^{*}} \xi(c, p), \Gamma_{2}(t, \boldsymbol{\nu}(t), c, p, \mathfrak{u}(t, c))\right\rangle \mathrm{d} \boldsymbol{\nu}(t)(c, p) .
$$

It follows, recalling (4.62) and (4.61), that $\boldsymbol{\nu}$ is a distributional solution of (4.63).
Finally we define $\mathcal{K}_{\omega, \tau}:[\tau, T] \rightarrow \mathbb{R}$ as

$$
\begin{array}{rcc}
\mathcal{K}_{\omega, \tau}(t) & := & \int_{C \times E_{C}^{*}}\left\langle p, \mathcal{F}\left(t, \Phi_{(t, \tau)}^{\boldsymbol{\mu}(t)}(c)\right)\right\rangle \mathrm{d} \boldsymbol{\nu}(t)(c, p) \\
& \int_{C \times E_{C}^{*}}\left\langle\psi_{t}(c, p), \mathcal{F}\left(t, \Phi_{(T, \tau)}^{\mu(T)}(c)\right)\right\rangle \mathrm{d} \boldsymbol{\nu}(T)(c, p), \tag{4.65}
\end{array}
$$

where $\mathcal{F}$ is defined by 4.51). Using that $\psi_{t}(c, p)$ and $\mathcal{F}(t, c)$ are weak solutions of 4.57) and (4.53) respectively, by the same density argument used in the proof of Theorem 3.3 (cf. the comments after (3.20)), recalling the definition of adjoint operator and applying Fubini's Theorem, we obtain

$$
\begin{equation*}
\frac{\mathrm{d}}{\mathrm{~d} t} \mathcal{K}_{\omega, \tau}(t)=0 \quad \text { a.e } t \in[\tau, T] . \tag{4.66}
\end{equation*}
$$

Step 4: Conclusion of the proof. Note that, by 4.66) and since

$$
\mathcal{K}_{\omega, \tau}(T) \stackrel{44.65],(4.61]}{=} \int_{C}\left\langle-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c), \mathcal{F}\left(T, \Phi_{(T, \tau)}^{\boldsymbol{\mu}(T)}(c)\right)\right\rangle \mathrm{d} \boldsymbol{\mu}(T)(c) \stackrel{\sqrt[4.56]]{\leq}}{\leq} 0,
$$

$$
\begin{aligned}
& \frac{\mathrm{d}}{\mathrm{~d} t}\left[\int_{C \times E_{C}^{*}} \xi(c, p) \mathrm{d} \boldsymbol{\nu}(t)(c, p)\right] \stackrel{4.60}{=} \frac{\mathrm{d}}{\mathrm{~d} t}\left[\int_{C \times E_{C}^{*}} \xi\left(\Phi_{(T, t)}^{\boldsymbol{\mu}^{T}}(c), p\right) \mathrm{d} \nu_{T}(t)(c, p)\right] \\
& \stackrel{(4.59)}{=} \frac{\mathrm{d}}{\mathrm{~d} t}\left[\int_{C \times E_{C}^{*}} \xi\left(\Phi_{(T, t)}^{\boldsymbol{\mu}^{T}}(c), p\right) \mathrm{d} \sigma_{c}(t)(p) \mathrm{d} \boldsymbol{\mu}(T)(c)\right] \\
& \stackrel{4.58]}{=} \frac{\mathrm{d}}{\mathrm{~d} t}\left[\int_{C \times E_{C}^{*}} \xi\left(\Phi_{(T, t)}^{\boldsymbol{\mu}^{T}}(c), \psi_{t}(c, p)\right) \mathrm{d} \delta_{-\nabla_{\boldsymbol{\mu}(T)} \varphi(\boldsymbol{\mu}(T))(c)}(p) \mathrm{d} \boldsymbol{\mu}(T)(c)\right] \\
& \stackrel{(4.611}{=} \frac{\mathrm{d}}{\mathrm{~d} t}\left[\int_{C \times E_{C}^{*}} \xi\left(\Phi_{(T, t)}^{\mu^{T}}(c), \psi_{t}(c, p)\right) \mathrm{d} \boldsymbol{\nu}(T)(c, p)\right] \\
& \sqrt[{4.87, \sqrt[4.57]{ }}]{\int} \int_{C \times E_{C^{*}}}\left\langle\mathrm{D}_{E} \xi\left(\Phi_{(T, t)}^{\boldsymbol{\mu}^{T}}(c), \psi_{t}(c, p)\right), A\left(t, \boldsymbol{\mu}(t), \Phi_{(T, t)}^{\mu^{T}}(c), u\left(t, \Phi_{(T, t)}^{\mu^{T}}(c)\right)\right)\right\rangle \mathrm{d} \boldsymbol{\nu}(T)(c, p) \\
& -\int_{C \times E_{C^{*}}}\left\langle\mathrm{D}_{E_{C}^{*}} \xi\left(\Phi_{(T, t)}^{\mu^{T}}(c), \psi_{t}(c, p)\right), \mathrm{D}_{\Phi_{(T, t)}^{*}(c)}^{*} A\left(t, \boldsymbol{\mu}(t), \Phi_{(T, t)}^{\mu^{T}}(c), u\left(t, \Phi_{(T, t)}^{\mu^{T}}(c)\right)\right)\left[\psi_{t}(c, p)\right]\right. \\
& \left.+\int_{C \times E_{C}^{*}} \nabla_{\boldsymbol{\mu}(t)}^{*} A\left(t, \boldsymbol{\mu}(t), \Phi_{(T, t)}^{\mu^{T}}(\tilde{c}), u\left(t, \Phi_{(T, t)}^{\boldsymbol{\mu}^{T}}(\tilde{c})\right)\right)\left(\Phi_{(T, t)}^{\boldsymbol{\mu}^{T}}(c)\right)\left[\psi_{t}(\tilde{c}, \tilde{p})\right] \mathrm{d} \boldsymbol{\nu}(T)(\tilde{c}, \tilde{p})\right\rangle \mathrm{d} \boldsymbol{\nu}(T)(c, p)
\end{aligned}
$$

we deduce that $\mathcal{K}_{\omega, \tau}(t) \leq 0$ for $t \in[\tau, T]$. In particular, it holds for every $\omega \in U$ and a.e. $\tau \in[0, T]$ that

$$
\begin{aligned}
& \mathcal{K}_{\omega, \tau}(\tau) \stackrel{\boxed{4.65)}}{=} \int_{C \times E_{C}^{*}}\langle p, \mathcal{F}(\tau, c)\rangle \mathrm{d} \boldsymbol{\nu}(\tau)(c, p) \\
& \stackrel{44.47}{=} \int_{C \times E_{C}^{*}}\langle p, A(\tau, \boldsymbol{\mu}(\tau), c, \omega(c))-A(\tau, \boldsymbol{\mu}(\tau), c, \mathfrak{u}(\tau, c))\rangle \mathrm{d} \boldsymbol{\nu}(\tau)(c, p) \leq 0 .
\end{aligned}
$$

Thus, recalling (4.12), we get 4.15). Finally, it follows from (4.13) and (4.64) that

$$
\Gamma(t, \nu, c, p, u)=J \nabla_{\nu} \mathcal{H}(t, \nu, u)(c, p) .
$$

Hence, by (4.63), $\boldsymbol{\nu}$ is a distributional solution of (4.14).
Now we focus on the Bolza problem and we give the infinite-dimensional version of Theorem 3.4. Let $L:[0, T] \times \mathcal{P}_{b}(C) \times C_{b}^{1}(C ; Z) \rightarrow \mathbb{R}$. We consider

$$
\begin{equation*}
\min _{u \in \mathcal{U}}\left\{\int_{0}^{T} L(t, \mu(t), u(t)) \mathrm{d} t\right\}, \tag{4.67}
\end{equation*}
$$

subject to (4.6). We assume for the running cost $L$ that a function $l:[0, T] \times \mathcal{P}_{b}(C) \times$ $C \times C_{b}^{1}(C ; Z) \rightarrow \mathbb{R}$ exists such that

$$
L(t, \mu, u)=\int_{C} l(t, \mu, c, u(c)) \mathrm{d} \mu(c) .
$$

Moreover, the following hold.
(HL):
(a) there exists a constant $L>0$ such that for every $t \in[0, T], \mu, \tilde{\mu} \in \mathcal{P}_{c}(C)$ and $u \in U$

$$
|L(t, \mu, u)-L(t, \tilde{\mu}, u)| \leq L W_{1}(\mu, \tilde{\mu}) ;
$$

(b) there exists $M>0$ such that for every $t \in[0, T], \mu \in \mathcal{P}_{c}(C)$ and $u \in U$ there holds

$$
|L(t, \mu, u)| \leq M\left(1+m_{1}(\mu)\right) ;
$$

(c) for every $\mu \in \mathcal{P}_{c}(C)$ and $u \in U$ the map $t \mapsto L(t, \mu, u)$ belongs to $L^{1}([0, T] ; \mathbb{R})$;
(d) for every $t \in[0, T]$ and $\mu \in \mathcal{P}_{c}(C)$ the map $u \mapsto L(t, \mu, u)$ belongs to $C^{0}\left(C_{b}^{1}(C ; Z) ; \mathbb{R}\right) ;$
(e) for every $t \in[0, T]$ and $u \in U$ the map $\mathcal{P}_{c}(C) \ni \mu \mapsto L(t, \mu, u) \in \mathbb{R}$ is locally differentiable at any $\mu$ in the sense of the Definition 4.2 with differential $\nabla_{\mu} L:=\nabla_{\mu} L(t, \mu, u)$ and the application $\mathcal{P}_{c}(C) \times U \times C \ni(\mu, u, \tilde{c}) \mapsto$ $\nabla_{\mu} L(\tilde{c}) \in E_{C}^{*}$ is continuous for every $t \in[0, T]$.

We define the Hamiltonian $\mathcal{H}:[0, T] \times \mathcal{P}_{b}\left(C \times E_{C}^{*}\right) \times U \rightarrow \mathbb{R}$ for the problem (4.67)-(4.6) as

$$
\begin{equation*}
\mathcal{H}(t, \nu, \omega)=\int_{C \times E_{C}^{*}}\left\langle p, A\left(t, \pi_{\#}^{1} \nu, c, \omega(c)\right)\right\rangle \mathrm{d} \nu(c, p)-L\left(t, \pi_{\#}^{1} \nu, \omega\right) . \tag{4.68}
\end{equation*}
$$

As seen for the Mayer problem we can compute explicitly $\nabla_{\nu} \mathcal{H}(t, \nu, u): C \times E_{C}^{*} \rightarrow E_{C}^{*} \times$ $E_{C}$ :

$$
\begin{equation*}
\nabla_{\nu} \mathcal{H}(t, \nu, u)(c, p)= \tag{4.69}
\end{equation*}
$$

$$
\binom{D_{c}^{*} A\left(t, \pi_{\#}^{1} \nu, c, u\right)[p]+\int_{C \times E_{C}^{*}} \nabla_{\pi_{\#}^{1}}^{*} A\left(t, \pi_{\#}^{1} \nu, \tilde{c}, u\right)(c)[\tilde{p}] \mathrm{d} \nu(\tilde{c}, \tilde{p})-\nabla_{\pi_{\# \nu}^{1}} L\left(t, \pi_{\#}^{1} \nu, u\right)(c)}{A\left(t, \pi_{\#}^{1} \nu, c, u\right)}
$$

Finally we state the Pontryagin maximum principle for the Bolza problem.
Theorem 4.10. Let $(\boldsymbol{\mu}, \mathfrak{u}) \in A C\left([0, T] ; \mathcal{P}_{c}(C)\right) \times \mathcal{U}$ be an optimal pair control trajectory for 4.67)-(4.6). Then there exists $\boldsymbol{\nu} \in A C\left([0, T] ; \mathcal{P}_{c}\left(C \times E_{C}^{*}\right)\right)$ which solves in distributional sense

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} \boldsymbol{\nu}(t)=-\operatorname{div}_{(c, p)}\left(\left(J \nabla_{\boldsymbol{\nu}(t)} \mathcal{H}(t, \boldsymbol{\nu}(t), \mathfrak{u}(t))(\cdot, \cdot)\right) \boldsymbol{\nu}(t)\right) \quad \text { in }[0, T), \\
\pi_{\#}^{1} \boldsymbol{\nu}(t)=\boldsymbol{\mu}(t), \\
\boldsymbol{\nu}(T)=\mu(T) \times \delta_{\mathbf{0}} \in \mathcal{P}_{c}\left(C \times E_{C}^{*}\right)
\end{array}\right.
$$

where $\mathcal{H}$ and $\nabla_{\nu} \mathcal{H}$ are defined by (4.68) and (4.69) respectively. Moreover the following maximality condition holds

$$
\mathcal{H}(t, \boldsymbol{\nu}(t), \mathfrak{u}(t))=\max _{\omega \in U}\{\mathcal{H}(t, \boldsymbol{\nu}(t), \omega)\} \quad \text { for a.e. } t \in[0, T] .
$$

Proof. The proof is an infinite-dimensional adaptation of the proof of Theorem 3.4. We do not report the details here as they are very similar of [15, Section 4.2].

## 5. A model example

In this section we briefly discuss an example that fits into the theoretical framework presented in Section 4. The general setting we present here is that of entropy regularised mean-field particle systems presented in [4, 10] for inhomogeneous evolutionary games [6].

We fix a compact metric space $(V, d)$ of pure strategies, a probability measure $\eta \in \mathcal{P}(\bar{V})$ with $\operatorname{supp}(\eta)=V$, and $p \in(1,+\infty)$. We define as ambient space $E:=\mathbb{R}^{d} \times L^{p}(V, \eta)$, where

$$
L^{p}(V, \eta):=\left\{\ell: V \rightarrow \mathbb{R}: \int_{V}|\ell(v)|^{p} \mathrm{~d} \eta(v)<+\infty\right\}
$$

Denoting by $|\cdot|$ the Euclidian norm in $\mathbb{R}^{d}$ and by $\|\cdot\|_{p}$ the $L^{p}$-norm in $L^{p}(V, \eta)$, we endow $E$ with the norm $\|\cdot\|_{E}=|\cdot|+\|\cdot\|_{p}$. With such choice, $E$ is a separable and reflexive Banach space. For $0<r<R<+\infty$ we set

$$
C_{r, R}:=\mathbb{R}^{d} \times\left\{\ell \in L^{p}(V, \eta): r \leq \ell(v) \leq R \text { for } \eta \text {-a.e. } v \in V\right\} .
$$

In particular, $C_{r, R}$ is a convex and closed subset of $E$. To shorten the notation, we denote by $c=(x, \ell)$ the generic element of $E$.

Let us fix $0<r<R<+\infty$. As in Section 4, we consider $\mathcal{U}=L^{1}([0, T] ; U)$ for $U$ a not empty compact subset of $\left(C_{b}^{1}\left(C_{r, R} ; \mathbb{R}^{d}\right),\|\cdot\|_{C_{b}^{1}}\right)$. As velocity field, for $\Psi \in \mathcal{P}_{1}\left(C_{r, R}\right)$ and $u \in U$ we consider

$$
\begin{equation*}
A(\Psi, c, u):=\left(v_{\Psi}(c)+u, \mathcal{T}_{\Psi}(c)+\varepsilon \mathcal{S}(\ell)\right) \tag{5.1}
\end{equation*}
$$

for some $\varepsilon>0$. Here, $v_{\Psi}: E \mapsto \mathbb{R}^{d}$ takes the form

$$
v_{\Psi}(c):=\int_{C_{r, R}} K\left(x^{\prime}-x\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right)
$$

for a suitable interaction kernel $K: \mathbb{R}^{d} \rightarrow \mathbb{R}^{d}$ such that
(K) $K \in C_{b}^{1}\left(\mathbb{R}^{d} ; \mathbb{R}^{d}\right)$.

As for the operators $\mathcal{T}_{\Psi}: C_{r, R} \rightarrow L^{p}(V, \eta)$ and $\mathcal{S}: C_{r, R} \rightarrow L^{0}(V, \eta)$ we consider

$$
\begin{aligned}
& \mathcal{T}_{\Psi}(c):=\left(\int_{C_{r, R}} J\left(x, \cdot, x^{\prime}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right)-\int_{V} \int_{C_{r, R}} J\left(x, v^{\prime}, x^{\prime}\right) \ell\left(v^{\prime}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right) \mathrm{d} \eta\left(v^{\prime}\right)\right) \ell \\
& \mathcal{S}(\ell):=\left(\int_{V} \ell(v) \log (\ell(v)) \mathrm{d} \eta(v)-\log (\ell)\right) \ell
\end{aligned}
$$

for a payoff function $J: \mathbb{R}^{d} \times V \times \mathbb{R}^{d} \rightarrow \mathbb{R}$ satisfying
(J.1) $J$ is bounded on $\mathbb{R}^{d} \times V \times \mathbb{R}^{d}$ and there exists a constant $L>0$ such that for every $x_{1}, x_{1}^{\prime}, x_{2}, x_{2}^{\prime} \in \mathbb{R}^{d}$ and every $v \in V$

$$
\left|J\left(x_{1}, v, x_{1}^{\prime}\right)-J\left(x_{2}, v, x_{2}^{\prime}\right)\right| \leq L\left(\left|x_{1}-x_{2}\right|+\left|x_{1}^{\prime}-x_{2}^{\prime}\right|\right)
$$

(J.2) for every $v \in V$, the map $\left(x, x^{\prime}\right) \mapsto J\left(x, v, x^{\prime}\right)$ is differentiable with $\left(x, x^{\prime}\right) \mapsto$ $\left(\nabla_{x} J\left(x, v, x^{\prime}\right), \nabla_{x^{\prime}} J\left(x, v, x^{\prime}\right)\right)$ continuous in $\mathbb{R}^{d} \times \mathbb{R}^{d}$.
Under the assumptions (K) and (J.1)-(J.2), it has been shown in [4, Proposition 3.2] and in [31, Proposition 5.8] that for every $\varepsilon>0$ there exists $0<r_{\varepsilon}<R_{\varepsilon}<+\infty$ such that, setting $C:=C_{r_{\varepsilon}, R_{\varepsilon}}$, the velocity field $A$ defined in (5.1) complies with assumptions (i)-(iv) in (HA1). Moreover, we have that

$$
E_{C}=\overline{\mathbb{R}(C-C)}=\left\{\ell \in L^{p}(V, \eta): \int_{V} \ell(v) \mathrm{d} \eta(v)=0\right\}
$$

Hence, $E_{C}$ is reflexive.
We now verify that $A$ also satisfies the set of assumptions (HA2). As before, condition (v) of (HA2) is trivial in view of the linear dependence of $A$ on the control variable $u$. As for the $C$-differentiability of $A$, we write the $C$-differential piece-by-piece as

$$
\begin{aligned}
& D_{c} v_{\Psi}(c)\left[c_{1}-c_{2}\right]=-\int_{C} \nabla K\left(x^{\prime}-x\right)\left(x_{1}-x_{2}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right) \\
& \begin{aligned}
& D_{c} \mathcal{T}_{\Psi}(c) {\left[c_{1}-c_{2}\right]=\left(\int_{C} \nabla_{x} J\left(x, \cdot, x^{\prime}\right) \cdot\left(x_{1}-x_{2}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right)\right.} \\
&\left.-\int_{V} \int_{C} \nabla_{x} J\left(x, v^{\prime}, x^{\prime}\right) \cdot\left(x_{1}-x_{2}\right) \ell\left(v^{\prime}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right) \mathrm{d} \eta\left(v^{\prime}\right)\right) \ell \\
& \quad+\left(\int_{C} J\left(x, \cdot, x^{\prime}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right)-\int_{V} \int_{C} J\left(x, v^{\prime}, x^{\prime}\right) \ell\left(v^{\prime}\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right) \mathrm{d} \eta\left(v^{\prime}\right)\right)\left(\ell_{1}-\ell_{2}\right) \\
& \quad-\ell \int_{V} \int_{C} J\left(x, v^{\prime}, x^{\prime}\right)\left(\ell_{1}\left(v^{\prime}\right)-\ell_{2}\left(v^{\prime}\right)\right) \mathrm{d} \Psi\left(x^{\prime}, \ell^{\prime}\right) \mathrm{d} \eta\left(v^{\prime}\right) \\
& D_{c} \mathcal{S}(c) {\left[c_{1}-c_{2}\right]=\left(\int_{V} \ell(v) \log (\ell(v)) \mathrm{d} \eta(v)-\log (\ell)-1\right)\left(\ell_{1}-\ell_{2}\right) } \\
& \quad+\ell \int_{V}\left(\ell_{1}(v)-\ell_{2}(v)\right)(1+\log (\ell(v))) \mathrm{d} \eta(v),
\end{aligned}
\end{aligned}
$$

for every $\Psi \in \mathcal{P}(C)$ and every $c_{1}, c_{2} \in C$ (recall that $c_{i}=\left(x_{i}, \ell_{i}\right)$ for $\left.i=1,2\right)$. In particular, (K) and (J.1)-(J.2) imply that $A$ satisfies (vi) of (HA2).

Finally, the map $\mathcal{P}_{c}(C) \ni \Psi \mapsto A(\Psi, c, u) \in E_{C}$ is locally differentiable at any $\Psi$ in the sense of the Definition 4.2 with differential

$$
\begin{aligned}
& \left(\nabla_{\Psi} A(\Psi, c, u)\right)\left(c_{1}\right)\left[c_{2}-c_{1}\right] \\
& \quad=\binom{\nabla K\left(x_{1}-x\right)\left[x_{2}-x_{1}\right]}{\left[\nabla_{x^{\prime}} J\left(x, \cdot, x_{1}\right) \cdot\left(x_{2}-x_{1}\right)-\int_{V} \nabla_{x^{\prime}} J\left(x, v, x_{1}\right) \cdot\left(x_{2}-x_{1}\right) \ell(v) \mathrm{d} \eta(v)\right] \ell(\cdot)},
\end{aligned}
$$

for every $c, c_{1}, c_{2} \in C$. In particular, (K) and (J.1)-(J.2) yield (vii) in (HA2).

## Appendix

We first prove the results whose proof we have postponed, then we give two technical results that we used in Section 4.

## A: Proofs of Lemma 4.4, Lemma 4.5 and Lemma 4.6.

Proof of Lemma 4.4. Fix $c \in C$. We recall that $\Phi_{(s, t)}^{\mu^{s}} \in C_{b}^{0}(C ; C)$. For brevity of notation we set $\tilde{c}:=\Phi_{(s, t)}^{\mu^{s}}(c)$. Let $v \in \mathbb{R}(C-C)$. Then $v=\alpha\left(c_{1}-c_{2}\right)=\alpha\left(c_{1}-\tilde{c}\right)-\alpha\left(c_{2}-\tilde{c}\right)$. It follows by definition of $C$-differential (see (3.2)) that

$$
\begin{aligned}
& \left\|\mathrm{D}_{\tilde{c}} A[v]\right\|_{E}=|\alpha|\left\|\mathrm{D}_{\tilde{c}} A\left[c_{1}-\tilde{c}\right]-\mathrm{D}_{\tilde{c}} A\left[c_{2}-\tilde{c}\right]\right\|_{E} \\
& \stackrel{(3.3)}{=}|\alpha| \lim _{h \rightarrow 0^{+}}\left\|\frac{1}{h}\left(A\left(t, \mu(t), \tilde{c}+h\left(c_{1}-\tilde{c}\right), u\left(t, \tilde{c}+h\left(c_{1}-\tilde{c}\right)\right)\right)-A\left(t, \mu(t), \tilde{c}+h\left(c_{2}-\tilde{c}\right), u\left(t, \tilde{c}+h\left(c_{2}-\tilde{c}\right)\right)\right)\right)\right\|_{E} \\
& \stackrel{4.9),(\mathbf{H A 1})-(i)}{\leq}|\alpha| L \lim _{h \rightarrow 0^{+}}\left\{h\left\|c_{1}-c_{2}\right\|_{E}+\left\|u\left(t, \tilde{c}+h\left(c_{1}-\tilde{c}\right)\right)-u\left(t, \tilde{c}+h\left(c_{2}-\tilde{c}\right)\right)\right\|_{C_{b}^{0}}\right\} .
\end{aligned}
$$

Since $u \in U$ which is compact in $C_{b}^{1}(C ; Z)$, by Remark A. 5 of [6], there exists a positive constant $L_{U}$ such that, starting from the last inequality, we deduce

$$
\left\|\mathrm{D}_{\tilde{c}} A[v]\right\|_{E} \leq L\left(1+L_{U}\right)\|v\|_{E} .
$$

By an abuse of notation we use the symbol $L$ to denote $L\left(1+L_{U}\right)$, hence, the result follows.

Proof of Lemma 4.5. Let us first note that, by 4.10), the Cauchy problem (4.11) is welldefined and admits a unique solution $z_{f}(t, c) \in A C\left([s, T] ; E_{C}\right)$. In particular, we can define the family of diffeomorphisms $L_{(s, t)}^{c}: E_{C} \rightarrow E_{C}$ as $L_{(s, t)}^{c}[f]:=z_{f}(t, c)$. Since the differential equation in (4.11) is linear, we have that $L_{(s, t)}^{c}$ is a linear operator for every $(t, c) \in[s, T] \times C$. Moreover, by its definition

$$
\begin{equation*}
z_{f}(t, c)=f+\int_{s}^{t} \mathrm{D}_{\Phi_{(s, \sigma)}^{\mu s}(c)} A\left[z_{f}(\sigma, c)\right] \mathrm{d} \sigma \tag{A.1}
\end{equation*}
$$

which, using 4.10, implies

$$
\left\|z_{f}(t, c)\right\|_{E} \leq\|f\|_{E}+L \int_{s}^{t}\left\|z_{f}(\sigma, c)\right\|_{E} \mathrm{~d} \sigma .
$$

Then, after an application of the Grönwall inequality, we obtain

$$
\left\|L_{(s, t)}^{c}[f]\right\|_{E}=\left\|z_{f}(t, c)\right\|_{E} \leq\|f\|_{E} \mathrm{e}^{L T} .
$$

Hence $L_{(s, t)}^{c} \in \mathcal{L}\left(E_{C} ; E_{C}\right)$ and $\left\|L_{(s, t)}^{c}\right\|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \leq \mathrm{e}^{L T}$, consequently we deduce that

$$
\text { the map } \quad(t, c) \mapsto L_{(s, t)}^{c} \quad \text { belongs to } L_{\mathcal{L} \times \mu^{s}}^{\infty}\left([s, T] \times C ; \mathcal{L}\left(E_{C} ; E_{C}\right)\right) .
$$

Let us now show that $L_{(s, t)}^{c}$ is exactly the $C$-differential of $\Phi_{(s, t)}^{\mu^{s}}$ at $c \in C$, thus proving the statement. By definition of $C$-differential, defining

$$
h:=\left\|c^{\prime}-c\right\|_{E}, \quad z_{h}(t):=\frac{\Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right)-\Phi_{(s, t)}^{\mu^{s}}(c)}{h} \quad \text { and } \quad f:=\frac{c^{\prime}-c}{h},
$$

we just have to prove that

$$
\begin{equation*}
\lim _{c^{\prime} \rightarrow c}\left\|z_{h}(t)-z_{f}(t, c)\right\|_{E}=0 \tag{A.2}
\end{equation*}
$$

Note that, by its definition, $z_{h}(t)$ solves
$\left\{\begin{array}{l}\frac{\mathrm{d}}{\mathrm{d} t} z_{h}(t)=\frac{1}{h}\left\{A\left(t, \mu(t), \Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right), u\left(t, \Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right)\right)\right)-A\left(t, \mu(t), \Phi_{(s, t)}^{\mu^{s}}(c), u\left(t, \Phi_{(s, t)}^{\mu^{s}}(c)\right)\right)\right\} \quad \text { in }(s, T], \\ z_{h}(s)=\frac{c^{\prime}-c}{h} .\end{array}\right.$
We can therefore rewrite the equation for $z_{h}(t)$ as

$$
\frac{\mathrm{d}}{\mathrm{~d} t} z_{h}(t)=\mathrm{D}_{\Phi_{(s, t)}^{\mu^{s}}(c)} A\left[z_{h}(t)\right]+r_{h}(t),
$$

where

$$
\begin{aligned}
& r_{h}(t):=\frac{\left\|\Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right)-\Phi_{(s, t)}^{\mu^{s}}(c)\right\|_{E}}{h} \times \\
& \frac{A\left(t, \mu(t), \Phi_{(s, t)}^{s^{s}}\left(c^{\prime}\right), u\left(t, \Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right)\right)\right)-A\left(t, \mu(t), \Phi_{(s, t)}^{\mu^{s}}(c), u\left(t, \Phi_{(s, t)}^{\mu^{s}}(c)\right)\right)-\mathrm{D}_{(s, t)}^{\mu^{s}}(c)}{} A\left[\Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right)-\Phi_{(s, t)}^{\mu^{s}}(c)\right] \\
& \left\|\Phi_{t}^{\mu^{s}}\left(c^{\prime}\right)-\Phi_{t}^{\mu^{s}}(c)\right\|_{E}
\end{aligned}
$$

leading to

$$
\begin{equation*}
z_{h}(t)=f+\int_{s}^{t} \mathrm{D}_{\Phi_{(s, \sigma)}^{\mu^{s}}(c)} A\left[z_{h}(\sigma)\right] \mathrm{d} \sigma+\int_{s}^{t} r_{h}(\sigma) \mathrm{d} \sigma . \tag{A.3}
\end{equation*}
$$

By (4.8), (HA1)-(i), and again applying the Grönwall inequality, we obtain

$$
\begin{equation*}
\left\|\Phi_{(s, t)}^{\mu^{s}}\left(c^{\prime}\right)-\Phi_{(s, t)}^{\mu^{s}}(c)\right\|_{E} \leq h \mathrm{e}^{L T} \tag{A.4}
\end{equation*}
$$

which, together to (HA2)-(vi), implies that

$$
\lim _{h \rightarrow 0}\left\|r_{h}(t)\right\|_{E}=0 \quad \text { for } t \in[s, T] .
$$

Moreover, using again (HA1)-(i), 4.10) and (A.4), it easy to check that

$$
\left\|r_{h}(t)\right\|_{E} \leq 2 L \mathrm{e}^{L T} \quad \text { for } t \in[s, T],
$$

so that, from an application of the Lebesgue dominated convergence theorem, we have

$$
\begin{equation*}
r_{h}(t) \rightarrow 0 \quad \text { in } \quad L^{1}\left([s, T] ; E_{C}\right) \quad \text { as } h \rightarrow 0 . \tag{A.5}
\end{equation*}
$$

Finally, by A.1 and A.3), we deduce that

$$
\begin{aligned}
\left\|z_{h}(t)-z_{f}(t, c)\right\|_{E} & \leq \int_{s}^{t}\left\|\mathrm{D}_{\Phi_{(s, \sigma)}^{\mu^{s}}(c)} A\left[z_{h}(\sigma)-z_{f}(\sigma, c)\right]\right\|_{E} \mathrm{~d} \sigma+\int_{s}^{t}\left\|r_{h}(\sigma)\right\|_{E} \mathrm{~d} \sigma \\
& \stackrel{4.10)}{\leq} L \int_{s}^{t}\left\|z_{h}(\sigma)-z_{f}(\sigma, c)\right\|_{E} \mathrm{~d} \sigma+\delta_{h},
\end{aligned}
$$

where

$$
\begin{equation*}
\delta_{h}:=\int_{s}^{T}\left\|r_{h}(\sigma)\right\|_{E} \mathrm{~d} \sigma \stackrel{\text { A.55 }}{\Rightarrow} 0 \quad \text { as } h \rightarrow 0 \tag{A.6}
\end{equation*}
$$

Applying again the Grönwall inequality, we conclude that

$$
\left\|z_{h}(t)-z_{f}(t, c)\right\|_{E} \leq \delta_{h} \mathrm{e}^{L T} .
$$

Therefore, recalling that $h=\left\|c^{\prime}-c\right\|_{E}$ and by A.6), we get A.2).
Proof of Lemma 4.6. We divide the proof in agreement with the three statements.
Proof of (a) and (b). Let $\hat{\mu}_{\varepsilon}$ be the solution of (4.6) with $u$ replaced by $u_{\varepsilon}$ and with initial datum $\mu^{0}$. Let $\Phi_{(0, t)}^{\varepsilon, \mu^{0}}$ be the family of non-local flows associated to $\hat{\mu}_{\varepsilon}$. We know by (4.9) that $\Phi_{(0, t)}^{\mu^{0}}, \Phi_{(0, t)}^{\varepsilon, \mu^{0}}, \Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}$ belong to $C_{b}^{0}([0, T] \times C ; C)$. Fix $c \in C$. It follows from triangle inequality that

$$
\begin{equation*}
\left\|\Phi_{(0, t)}^{\varepsilon, \mu_{\epsilon}^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right\|_{E} \leq\left\|\Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c)-\underset{(0, t)}{\Phi^{\varepsilon, \mu^{0}}}(c)\right\|_{E}+\left\|\Phi_{(0, t)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right\|_{E} . \tag{A.7}
\end{equation*}
$$

We focus on the second term of the right-hand side of A.7). We define

$$
F_{\varepsilon}(t):=\sup _{C}\left\|\Phi_{(0, t)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right\|_{E} \quad \text { for } t \in[0, T] .
$$

It follows from (4.8) and (HA1)-(i) that

$$
\begin{align*}
& \left\|\Phi_{(0, t)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right\|_{E}  \tag{A.8}\\
& \leq \int_{0}^{t}\left\|A\left(s, \hat{\mu}_{\varepsilon}(s), \Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c), u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)\right)\right)-A\left(s, \mu(s), \Phi_{(0, s)}^{\mu^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right)\right\|_{E} \mathrm{~d} s \\
& \leq L \int_{0}^{t}\left\{W_{1}\left(\hat{\mu}_{\varepsilon}(s), \mu(s)\right)+\left\|\Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, s)}^{\mu^{0}}(c)\right\|_{E}+\left\|u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)\right)-u\left(s, \Phi_{(0, s)}^{\Phi^{0}}(c)\right)\right\|_{C_{b}^{0}}\right\} \mathrm{d} s \\
& \leq L \int_{0}^{t} W_{1}\left(\hat{\mu}_{\varepsilon}(s), \mu(s)\right) \mathrm{d} s+L \int_{0}^{t}\left\|\Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, s)}^{\mu^{0}}(c)\right\|_{E} \mathrm{~d} s \\
& +L \int_{0}^{t}\left\|u_{\varepsilon\left(s, \Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)\right)}-u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right\|_{C_{b}^{0} \mathrm{~d} s}+L \int_{0}^{t}\left\|u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)-u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right\|_{C_{b}^{0}} \mathrm{~d} s .
\end{align*}
$$

Now we estimate the terms on the right-hand side of (A.8). As for the first term, by definition of Wasserstein distance, we deduce that

$$
\begin{equation*}
\int_{0}^{t} W_{1}\left(\hat{\mu}_{\varepsilon}(s), \mu(s)\right) \mathrm{d} s \leq \int_{0}^{t} \int_{C}\left\|\Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, s)}^{\mu^{0}}(c)\right\|_{E} \mathrm{~d} \mu_{0}(c) \mathrm{d} s \leq \int_{0}^{t} F_{\varepsilon}(s) \mathrm{d} s \tag{A.9}
\end{equation*}
$$

Regarding the third term, since $u_{\varepsilon}, u \in U$ which is compact in $\left(C_{b}^{1}(C ; Z),\|\cdot\|_{C_{b}^{1}}\right)$, there exists a positive constant $L_{U}$ such that

$$
\begin{equation*}
\int_{0}^{t}\left\|u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\varepsilon, \mu^{0}}(c)\right)-u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right\|_{C_{b}^{0}} \mathrm{~d} s \leq L_{U} \int_{0}^{t} F_{\varepsilon}(s) \mathrm{d} s . \tag{A.10}
\end{equation*}
$$

Finally, by assumption, $u_{\varepsilon} \rightarrow u$ in $L^{1}\left([0, T] ;\left(U,\|\cdot\|_{C_{b}^{1}}\right)\right)$, then

$$
\begin{equation*}
\int_{0}^{t}\left\|u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)-u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right\|_{C_{b}^{0}} \mathrm{~d} s \rightarrow 0 \quad \text { as } \varepsilon \rightarrow 0 . \tag{A.11}
\end{equation*}
$$

Combining (A.9), A.10) and A.11) with A.8), we obtain that

$$
F_{\varepsilon}(t) \leq L\left(2+L_{U}\right) \int_{0}^{t} F_{\varepsilon}(s) \mathrm{d} s+\delta_{\varepsilon}
$$

where $\delta_{\varepsilon}$ is a positive constant not depending on $c$ and $t$ which goes to 0 as $\varepsilon \rightarrow 0$. Hence, applying Grönwall inequality, we conclude that

$$
\begin{equation*}
F_{\varepsilon}(t) \leq \delta_{\varepsilon} \mathrm{e}^{L\left(2+L_{U}\right) T}, \tag{A.12}
\end{equation*}
$$

whence

$$
\begin{equation*}
\sup _{C}\left\|\Phi_{(0, t)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right\|_{E} \rightarrow 0 \quad \text { as } \varepsilon \rightarrow 0 \quad \text { uniformly in } t \in[0, T] . \tag{A.13}
\end{equation*}
$$

As regards the first term on the right-hand side of (A.7), since by Theorem 3.3 of [6] we have for a positive constant $M$

$$
\begin{equation*}
W_{1}\left(\mu_{\varepsilon}(t), \hat{\mu}_{\varepsilon}(t)\right) \leq \mathrm{e}^{M t} W_{1}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right) \quad \text { for } t \in[0, T] \tag{A.14}
\end{equation*}
$$

and recalling that, by assumption, $W_{1}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right) \rightarrow 0$ as $\varepsilon \rightarrow 0$, we can proceed in the same way used for the second term to obtain that

$$
\begin{equation*}
\sup _{C}\left\|\Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c)-\Phi_{(0, t)}^{\varepsilon, \mu^{0}}(c)\right\|_{E} \rightarrow 0 \quad \text { as } \varepsilon \rightarrow 0 \quad \text { uniformly in } t \in[0, T] . \tag{A.15}
\end{equation*}
$$

Hence, combining (A.13) and (A.15) with (A.7), we deduce (a). The proof of (b) is an easy consequence of what we have just seen. Indeed, by triangle inequality,

$$
\begin{aligned}
W_{1}\left(\mu_{\varepsilon}(t), \mu(t)\right) & \leq \\
& \stackrel{\text { A.14] }}{\leq} \mathrm{W}_{1}\left(\mu_{\varepsilon}(t), \hat{\mu}_{\varepsilon}(t)\right)+W_{1}\left(\hat{\mu}_{\varepsilon}(t), \mu(t)\right) \\
& \stackrel{A T}{ }{ }^{M T}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right)+\int_{C}\left\|\Phi_{(0, t)}^{\varepsilon, \mu^{0}}(c)-\Phi_{(0, t)}^{\mu^{0}}(c)\right\|_{E} \mathrm{~d} \mu^{0}(c) \\
& \mathrm{e}^{M T} W_{1}\left(\mu_{\varepsilon}^{0}, \mu^{0}\right)+\delta_{\varepsilon} \mathrm{e}^{L\left(2+L_{U}\right) T},
\end{aligned}
$$

so (b) follows.
Proof of ( $c$ ). Fix $t \in[0, T]$ and $c \in C$. We define for a fixed $f \in E_{C}$

$$
z_{\varepsilon}(t, c):=\mathrm{D}_{c} \Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}[f] \quad \text { and } \quad z(t, c):=\mathrm{D}_{c} \Phi_{(0, t)}^{\mu^{0}}[f] .
$$

By Lemma 4.5 we know that $\|z(t, c)\|_{E} \leq M\|f\|_{E}$ for a positive constant $M$ and that

$$
\begin{align*}
& \left\|z_{\varepsilon}(t, c)-z(t, c)\right\|_{E} \tag{A.16}
\end{align*}
$$

$$
\begin{aligned}
& +\int_{0}^{t}\left\|\left(\mathrm{D}_{\Phi_{(0, s)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c)} A\left(s, \mu_{\varepsilon}(s), \Phi_{(0, s)}^{\left.\varepsilon_{\left(0, \mu_{\varepsilon}^{0}\right.}^{0}(c), u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c)\right)\right)}\right)_{\Phi_{(0, s)}^{\mu^{0}}(c)} A\left(s, \mu(s), \Phi_{(0, s)}^{\mu^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right)\right)[z(s, c)]\right\|_{E} \mathrm{~d} s \\
& \stackrel{\sqrt[4.100]{\leq}}{\leq} L \int_{0}^{t}\left\|z_{\varepsilon}(s, c)-z(s, c)\right\|_{E} \mathrm{~d} s \\
& +M\|f\|_{E} \int_{0}^{t}\left\|\mathrm{D}_{\Phi_{(0, s)}^{\varepsilon, 4}(c)} A\left(s, \mu_{\varepsilon}^{0}(s), \Phi_{(0, s)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c), u_{\varepsilon}\left(s, \Phi_{(0, s)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c)\right)\right)-\mathrm{D}_{\Phi_{(0, s)}^{\mu^{0}}(c)} A\left(s, \mu(s), \Phi_{(0, s)}^{\mu_{0}^{0}}(c), u\left(s, \Phi_{(0, s)}^{\mu^{0}}(c)\right)\right)\right\|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \mathrm{d} s .
\end{aligned}
$$

Furthermore, by assumption (HA2)-(vi) and using (a) and the convergence of $u_{\varepsilon}$ to $u$ in $L^{1}\left([0, T] ;\left(U,\|\cdot\|_{C_{b}^{1}}\right)\right)$, we have for a.e. $t \in[0, T]$ that as $\varepsilon \rightarrow 0$

$$
\| \mathrm{D}_{\Phi_{(0, t)}^{\varepsilon} \mu_{( }^{0}(c)} A\left(t, \mu_{\varepsilon}^{\left.(t), \Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c), u_{\varepsilon}^{( }\left(t, \Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}(c)\right)\right)} \mathrm{D}_{\Phi_{(0, t)}^{\mu^{0}}(c)} A\left(t, \mu(t), \Phi_{(0, t)}^{\mu^{0}}(c), u\left(t, \Phi_{(0, t)}^{\mu^{0}}(c)\right)\right) \|_{\mathcal{L}\left(E_{C} ; E_{C}\right)} \rightarrow 0 .\right.
$$

Since, by Lemma 4.4 , both $\mathrm{D}_{\Phi_{\left(,, \omega_{t}\right)}^{0}(c)} A$ and $\mathrm{D}_{\Phi_{(0, t)}^{\mu^{0}}(c)} A$ are uniformly bounded in $[0, T]$, we can apply the Lebesgue theorem to the last term on the right-hand side of A.16) obtaining that

$$
\left\|z_{\varepsilon}(t, c)-z(t, c)\right\|_{E} \leq L \int_{0}^{t}\left\|z_{\varepsilon}(s, c)-z(s, c)\right\|_{E} \mathrm{~d} s+M\|f\|_{E} \delta_{\varepsilon}(c)
$$

where $\delta_{\varepsilon}(c)$ is a positive constant depending on $c$ which goes to 0 as $\varepsilon \rightarrow 0$. It follows by an application of the Grönwall inequality that

$$
\left\|z_{\varepsilon}(t, c)-z(t, c)\right\|_{E} \leq M\|f\|_{E} \delta_{\varepsilon}(c) \mathrm{e}^{L T} .
$$

Hence, recalling the definitions of $z_{\varepsilon}$ and $z$, we conclude that

$$
\sup _{E_{C} \ni f \neq 0} \frac{\left\|\left(\mathrm{D}_{c} \Phi_{(0, t)}^{\varepsilon, \mu_{\varepsilon}^{0}}-\mathrm{D}_{c} \Phi_{(0, t)}^{\mu^{0}}\right)[f]\right\|_{E}}{\|f\|_{E}} \leq M \delta_{\varepsilon}(c) \mathrm{e}^{L T} \rightarrow 0 \quad \text { as } \varepsilon \rightarrow 0,
$$

thus finding (c).

B: Auxiliary lemmas. In what follows let $X$ be a closed and convex subset of a separable Banach space $Y$ and let $E_{X}$ be the closure of the vectorial subspace $\mathbb{R}(X-X)$ in $Y$.

Lemma B.1. Let $\mu$ belong to $\mathcal{P}_{c}(X)$. Let $B_{1}:[0, T] \times X \rightarrow \mathcal{L}\left(E_{X} ; E_{X}\right)$ and $B_{2}:[0, T] \times$ $X \times X \rightarrow \mathcal{L}\left(E_{X} ; E_{X}\right)$ be operators satisfying
(i) the map $x \mapsto B_{1}(t, x)$ belongs to $C\left(X ; \mathcal{L}\left(E_{X} ; E_{X}\right)\right)$ for a.e. $t \in[0, T]$ and $B_{1} \in$ $L_{\mathcal{L} \times \mu}^{\infty}\left([0, T] \times X ; \mathcal{L}\left(E_{X} ; E_{X}\right)\right)$;
(ii) the map $(x, \tilde{x}) \mapsto B_{2}(t, x, \tilde{x})$ belongs to $C\left(X \times X ; \mathcal{L}\left(E_{X} ; E_{X}\right)\right)$ for a.e. $t \in[0, T]$ and $B_{1} \in L_{\mathcal{L} \times \mu \times \mu}^{\infty}\left([0, T] \times X \times X ; \mathcal{L}\left(E_{X} ; E_{X}\right)\right)$.
Then there exists a unique continuous weak solution $w:[0, T] \times X \rightarrow E_{X}$ of

$$
\left\{\begin{array}{l}
\frac{\mathrm{d}}{\mathrm{~d} t} w(t, x)=B_{1}(t, x)[w(t, x)]+\int_{X} B_{2}(t, x, \tilde{x})[w(t, \tilde{x})] \mathrm{d} \mu(\tilde{x}) \quad \text { in }(0, T], \\
w(0, x)=y \in E_{X} .
\end{array}\right.
$$

Proof. We consider the metric space $C_{b}^{0}\left([0, T] \times X ; E_{X}\right)$ equipped with the norm

$$
\|w\|_{\alpha}:=\sup _{[0, T] \times X} \mathrm{e}^{-2 \alpha t}\|w(t, x)\|_{Y} .
$$

Set $F:=\left(C_{b}^{0}\left([0, T] \times X ; E_{X}\right),\|\cdot\|_{\alpha}\right)$. Since $E_{X}$ is a closed subset of $Y$ which is a Banach space, it follows that $E_{X}$ is a complete metric space. Then $F$ is a complete metric space. We define the operator $S: F \rightarrow F$ as

$$
S(w)(t, x):=y+\int_{0}^{t} B_{1}(s, x)[w(s, x)] \mathrm{d} s+\int_{0}^{t} \int_{X} B_{2}(s, x, \tilde{x})[w(s, \tilde{x})] \mathrm{d} \mu(\tilde{x}) \mathrm{d} s .
$$

Thanks to the assumptions $(i)-(i i), S$ is well-defined. Following the same reasoning as in the first part of the proof of Proposition 5 in [18], it is easy to prove that, with a suitable choice of $\alpha, S$ is a contraction mapping. Then, applying the Banach-Caccioppoli theorem, the result follows.

Lemma B.2. Let $\mathcal{H}: \mathcal{P}_{b}(X) \rightarrow \mathbb{R}$ be

$$
\mathcal{H}(\mu)=\int_{X} H(\mu, x) \mathrm{d} \mu(x),
$$

where $H: \mathcal{P}_{b}(X) \times X \rightarrow \mathbb{R}$. Assume that $H$ is $X$-differentiable in the sense of (3.2) with $X$-differential satisfying
(i) the map $(\mu, x) \mapsto D_{x} H(\mu, x)$ belongs to $C\left(\mathcal{P}_{b}(X) \times X ; E_{X}^{*}\right)$;
(ii) there exists a positive constant $L$ not depending on $\mu \in \mathcal{P}_{b}(X)$ and $x \in X$ such that $\left\|D_{x} H(\mu, x)\right\|_{E_{C}^{*}} \leq L$ for every $\mu \in \mathcal{P}_{b}(X)$ and $x \in X$.
Moreover, assume that $H$ is differentiable in the sense of Definition 4.2 with $\mu$-differential satisfying
(iii) the map $(\mu, \tilde{x}, x) \mapsto \nabla_{\mu} H(\mu, \tilde{x})(x)$ belongs to $C\left(\mathcal{P}_{c}(X) \times X \times X ; E_{X}^{*}\right)$;
(iv) the $\operatorname{map}(\tilde{x}, x) \mapsto \nabla_{\mu} H(\mu, \tilde{x})(x)$ belongs to $L_{\mu \times \mu}^{\infty}\left(X \times X ; E_{X}^{*}\right)$ for every $\mu \in \mathcal{P}_{c}(X)$.

Then, $\mathcal{H}$ is differentiable at any $\mu \in \mathcal{P}_{c}(X)$ in the sense of Definition 4.2 and

$$
\nabla_{\mu} \mathcal{H}(\mu)(x)=D_{x} H(\mu, x)+\int_{X} \nabla_{\mu} H(\mu, \tilde{x})(x) \mathrm{d} \mu(\tilde{x})
$$

Proof. Let $R>0$ and $\eta \in \mathcal{P}\left(B_{\mu}(R)\right)$. It follows from the disintegration theorem that for $\gamma \in \Gamma(\mu, \eta)$

$$
\begin{align*}
& \mathcal{H}(\eta)-\mathcal{H}(\mu)=\int_{X} H(\eta, x) \mathrm{d} \eta(x)-\int_{X} H(\mu, x) \mathrm{d} \mu(x)  \tag{B.1}\\
& =\int_{X \times X}\left(H\left(\eta, x_{2}\right)-H\left(\mu, x_{1}\right)\right) \mathrm{d} \gamma\left(x_{1}, x_{2}\right) \\
& =\underbrace{\int_{X \times X}\left(H\left(\eta, x_{2}\right)-H\left(\eta, x_{1}\right)\right) \mathrm{d} \gamma\left(x_{1}, x_{2}\right)}_{=: I_{1}}+\underbrace{\int_{X}\left(H\left(\eta, x_{1}\right)-H\left(\mu, x_{1}\right)\right) \mathrm{d} \mu\left(x_{1}\right)}_{=: I_{2}} .
\end{align*}
$$

Now we focus on $I_{1}$. Recalling that $X$ is convex and using (i) and (ii), we have for $x(s)=x_{1}+s\left(x_{2}-x_{1}\right) \in X$ with $s \in[0,1]$ and applying Fubini theorem for Bochner integral

$$
\begin{align*}
I_{1}= & \int_{X \times X} \int_{0}^{1}\left\langle\mathrm{D}_{x(s)} H(\eta, x(s)), x_{2}-x_{1}\right\rangle \mathrm{d} s \mathrm{~d} \gamma\left(x_{1}, x_{2}\right)  \tag{B.2}\\
= & \int_{0}^{1} \int_{X \times X}\left\langle\mathrm{D}_{x(s)} H(\eta, x(s)), x_{2}-x_{1}\right\rangle \mathrm{d} \gamma\left(x_{1}, x_{2}\right) \mathrm{d} s \\
= & \int_{X \times X}\left\langle\mathrm{D}_{x_{1}} H\left(\mu, x_{1}\right), x_{2}-x_{1}\right\rangle \mathrm{d} \gamma\left(x_{1}, x_{2}\right) \\
& +\underbrace{\int_{X \times X}\left\langle\mathrm{D}_{x_{1}} H\left(\eta, x_{1}\right)-\mathrm{D}_{x_{1}} H\left(\mu, x_{1}\right), x_{2}-x_{1}\right\rangle \mathrm{d} \gamma\left(x_{1}, x_{2}\right)}_{=: I_{1,1}} \\
& +\underbrace{\int_{0}^{1} \int_{X \times X}\left\langle\mathrm{D}_{x(s)} H(\eta, x(s))-\mathrm{D}_{x_{1}} H\left(\eta, x_{1}\right), x_{2}-x_{1}\right\rangle \mathrm{d} \gamma\left(x_{1}, x_{2}\right) \mathrm{d} s}_{=: I_{1,2}} .
\end{align*}
$$

As for $I_{1,1}$, using Hölder inequality, we deduce that

$$
\left|I_{1,1}\right| \leq\left(\int_{X}\left\|\mathrm{D}_{x_{1}} H\left(\eta, x_{1}\right)-\mathrm{D}_{x_{1}} H\left(\mu, x_{1}\right)\right\|_{E_{X}^{*}}^{2} \mathrm{~d} \mu\left(x_{1}\right)\right)^{\frac{1}{2}} W_{2, \gamma}(\mu, \eta)
$$

Since, $\gamma \in \mathcal{P}(X \times X)$, we have that if $W_{2, \gamma}(\mu, \eta) \rightarrow 0$ then $W_{1}(\mu, \eta) \rightarrow 0$. Hence, using (i) and (ii), we have

$$
\left\{\begin{array}{l}
\mathrm{D}_{x_{1}} H\left(\eta, x_{1}\right) \rightarrow \mathrm{D}_{x_{1}} H\left(\mu, x_{1}\right) \quad \text { as } W_{2, \gamma}(\mu, \eta) \rightarrow 0, \\
\left\|\mathrm{D}_{x_{1}} H\left(\eta, x_{1}\right)-\mathrm{D}_{x_{1}} H\left(\mu, x_{1}\right)\right\|_{E_{X}^{*}} \leq 2 L .
\end{array}\right.
$$

Therefore, by an application of the Lebesgue theorem, we get

$$
\begin{equation*}
I_{1,1}=o_{R}\left(W_{2, \gamma}(\mu, \eta)\right) \tag{B.3}
\end{equation*}
$$

As regards $I_{1,2}$, using again $(i)$ and (ii), we know that for every $\varepsilon>0$ there exists $\delta>0$ depending on $\eta$ such that

$$
\begin{aligned}
\left|I_{1,2}\right| & \leq \varepsilon \int_{\left\{\left\|x_{1}-x_{2}\right\|_{Y}<\delta\right\}}\left\|x_{1}-x_{2}\right\|_{Y} \mathrm{~d} \gamma\left(x_{1}, x_{2}\right)+2 L \int_{\left\{\left\|x_{1}-x_{2}\right\|_{Y} \geq \delta\right\}}\left\|x_{1}-x_{2}\right\|_{Y} \mathrm{~d} \gamma\left(x_{1}, x_{2}\right) \\
& \leq \varepsilon W_{2, \gamma}(\mu, \eta)+\frac{2 L}{\delta} W_{2, \gamma}^{2}(\mu, \eta) .
\end{aligned}
$$

It follows that, if $W_{2, \gamma}(\mu, \eta) \rightarrow 0$, then we can choose $W_{2, \gamma}(\mu, \eta) \leq \frac{\varepsilon \delta}{2 L}$ so that from the last inequality we deduce for every $\varepsilon>0$ that $\left|I_{1,2}\right| \leq 2 \varepsilon W_{2, \gamma}(\mu, \eta)$. Hence,

$$
\begin{equation*}
I_{1,2}=o_{R}\left(W_{2, \gamma}(\mu, \eta)\right) . \tag{B.4}
\end{equation*}
$$

Finally we focus on $I_{2}$. We have using (iii) and (iv) and applying Fubini theorem for Bochner integral that

$$
\begin{align*}
I_{2} & =\int_{X} \int_{X \times X}\left\langle\nabla_{\mu} H\left(\mu, x_{1}\right)\left(x_{3}\right), x_{4}-x_{3}\right\rangle \mathrm{d} \gamma\left(x_{3}, x_{4}\right) \mathrm{d} \mu\left(x_{1}\right)+o_{R}\left(W_{2, \gamma}(\mu, \eta)\right) \\
& =\int_{X \times X}\left\langle\int_{X} \nabla_{\mu} H\left(\mu, x_{1}\right)\left(x_{3}\right) \mathrm{d} \mu\left(x_{1}\right), x_{4}-x_{3}\right\rangle \mathrm{d} \gamma\left(x_{3}, x_{4}\right)+o_{R}\left(W_{2, \gamma}(\mu, \eta)\right)( \tag{B.5}
\end{align*}
$$

Combining ( $\overline{\mathrm{B} .2}$ ), ( $\overline{\mathrm{B} .3}$ ), ( $\overline{\mathrm{B} .4})$ and (B.5) with ( B .1$)$ we obtain the result.
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