# PATHWISE CONVERGENCE OF THE HARD SPHERES KAC PROCESS 

By Daniel Heydecker ${ }^{1}$<br>University of Cambridge


#### Abstract

We derive two estimates for the deviation of the $N$-particle, hard-spheres Kac process from the corresponding Boltzmann equation, measured in expected Wasserstein distance. Particular care is paid to the long-time properties of our estimates, exploiting the stability properties of the limiting Boltzmann equation at the level of realisations of the interacting particle system. As a consequence, we obtain an estimate for the propagation of chaos, uniformly in time and with polynomial rates, as soon as the initial data has a $k$ th moment, $k>2$. Our approach is similar to Kac's proposal of relating the long-time behaviour of the particle system to that of the limit equation. Along the way, we prove a new estimate for the continuity of the Boltzmann flow measured in Wasserstein distance.


## CONTENTS

1. Introduction and main results ..... 3063
1.1. Main results ..... 3067
1.2. Plan of the paper ..... 3071
1.3. Literature review ..... 3072
1.4. Discussion of our results ..... 3074
2. Moment estimates ..... 3078
3. Regularity and stability estimates ..... 3080
3.1. Stability estimates ..... 3080
3.2. Regularity estimates ..... 3082
4. Proof of Theorem 1.6 ..... 3083
5. The interpolation decomposition for Kac's process ..... 3087
5.1. Proof of Lemma 5.2 ..... 3090
6. Proof of Theorem 1.1 ..... 3096
7. Proof of Theorem 1.2 ..... 3103
8. Proof of Theorem 1.3 ..... 3113
9. Proof of Theorem 1.4 ..... 3116
10. Proof of Theorem 1.7 ..... 3121
Appendix: Calculus of martingales ..... 3124
Acknowledgements ..... 3125
References ..... 3125
[^0]1. Introduction and main results. Kac [23] introduced a Markov model for the behaviour of a dilute gas, corresponding to the spatially homogeneous Boltzmann equation. We consider an ensemble of $N$ indistinguishable particles, with velocities $v_{1}(t), \ldots, v_{N}(t) \in \mathbb{R}^{d}$ at time $t \geq 0$, which are are encoded in the empirical velocity distribution

$$
\begin{equation*}
\mu_{t}^{N}=N^{-1} \sum_{i=1}^{N} \delta_{v_{i}(t)} \tag{1.1}
\end{equation*}
$$

Throughout, unless specified otherwise, we consider only the following example, known as the hard spheres kernel, of Kac processes, which is one of two main examples of physical interest. The dynamics are as follows:

1. For every (unordered) pair of particles with velocities $v, v_{\star} \in \operatorname{supp}\left(\mu_{t}^{N}\right)$, the particles collide at a rate $2\left|v-v_{\star}\right| / N$.
2. When two particles collide, take an independent random variable $\Sigma$, distributed uniformly on $S^{d-1}$. The particles then separate in direction $\Sigma$.
3. The velocities change to $v^{\prime}\left(v, v_{\star}, \Sigma\right)$ and $v_{\star}^{\prime}\left(v, v_{\star}, \Sigma\right)$, given by conservation of energy and momentum as

$$
\begin{align*}
v^{\prime}\left(v, v_{\star}, \Sigma\right) & =\frac{v+v_{\star}+\Sigma\left|v-v_{\star}\right|}{2} \\
v_{\star}^{\prime}\left(v, v_{\star}, \Sigma\right) & =\frac{v+v_{\star}-\Sigma\left|v-v_{\star}\right|}{2} \tag{1.2}
\end{align*}
$$

The measure changes to

$$
\begin{equation*}
\mu \mapsto \mu^{N, v, v_{\star}, \Sigma}=\mu+\frac{1}{N}\left(\delta_{v^{\prime}}+\delta_{v_{\star}^{\prime}}-\delta_{v}-\delta_{v_{\star}}\right) . \tag{1.3}
\end{equation*}
$$

More formally, we consider the space $\mathcal{S}$ of Borel measures on $\mathbb{R}^{d}$, satisfying

$$
\begin{equation*}
\left.\langle 1, \mu\rangle=1 ; \quad\langle v, \mu\rangle=0 ;\left.\quad\langle | v\right|^{2}, \mu\right\rangle=1 \tag{1.4}
\end{equation*}
$$

the notational conventions that angle brackets $\langle\cdot, \cdot\rangle$ denote integration against a measure, and $v$ denotes the identity function on $\mathbb{R}^{d} . \mathcal{S}$ is called the Boltzmann sphere, and consists of those measures with normalised mass, momentum, and energy. We write $\mathcal{S}^{k}$ for the subspace of $\mathcal{S}$ where the $k$ th moment $\left.\left.\langle | v\right|^{k}, \mu\right\rangle$ is finite, and define the following family of weights:

$$
\begin{equation*}
\Lambda_{k}(\mu):=\left\langle\left(1+|v|^{2}\right)^{\frac{k}{2}}, \mu\right\rangle \tag{1.5}
\end{equation*}
$$

This leads to a natural family of subspaces:

$$
\begin{equation*}
\mathcal{S}_{a}^{k}:=\left\{\mu \in \mathcal{S}: \Lambda_{k}(\mu) \leq a\right\} \tag{1.6}
\end{equation*}
$$

For shorthand, we will often write $\Lambda_{k}(\mu, \nu):=\max \left(\Lambda_{k}(\mu), \Lambda_{k}(\nu)\right)$.

Let $\mathcal{S}_{N}$ be the subset of $\mathcal{S}$ consisting of normalised empirical measures on $N$ points; we will typically write $\mu^{N}$ for a generic element of $\mathcal{S}_{N}$. Formally, the Kac process is the Markov process on $\mathcal{S}_{N}$ with kernel

$$
\begin{align*}
& \mathcal{Q}_{N}\left(\mu^{N}\right)(A) \\
& \quad=N \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} 1\left(\mu^{N, v, v_{\star}, \sigma} \in A\right)\left|v-v_{\star}\right| \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right) d \sigma . \tag{1.7}
\end{align*}
$$

Note that, since the map $\mu^{N} \mapsto \mu^{N, v, v_{\star}, \sigma}$ preserves particle number, momentum, and kinetic energy, $\mathcal{Q}_{N}\left(\mu^{N}\right)$ is supported on $\mathcal{S}_{N}$ whenever $\mu^{N} \in \mathcal{S}_{N}$. We write $\left(\mu_{t}^{N}\right)_{t \geq 0}$ for a Kac process on $N$ particles. Observe that the rates are bounded by $2 N$, and so for any initial datum $\mu_{0}^{N}$, the law of a Kac process started from $\mu_{0}^{N}$ exists, and is unique, and the process is almost surely nonexplosive.

Measure solutions to the Boltzmann equation. Following many previous works, [24, 29, 33], we study measure-valued solutions to the Boltzmann equation. We define the Boltzmann collision operator $Q(\mu, \nu)$ for measures $\mu, \nu \in \mathcal{S}$ as

$$
\begin{equation*}
Q(\mu, v)=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left\{\delta_{v^{\prime}}+\delta_{v_{\star}^{\prime}}-\delta_{v}-\delta_{v_{\star}}\right\}\left|v-v_{\star}\right| d \sigma \mu(d v) v\left(d v_{\star}\right) \tag{1.8}
\end{equation*}
$$

For brevity, we will denote $Q(\mu, \mu)$ by $Q(\mu)$. We say that a family $\left(\mu_{t}\right)_{t \geq 0}$ of measures in $\mathcal{S}$ satisfies the Boltzmann equation if, for any bounded measurable $f$ of compact support,

$$
\begin{equation*}
\forall t \geq 0 \quad\left\langle f, \mu_{t}\right\rangle=\left\langle f, \mu_{0}\right\rangle+\int_{0}^{t}\left\langle f, Q\left(\mu_{s}\right)\right\rangle d s \tag{BE}
\end{equation*}
$$

The Boltzmann equation is known to have a unique fixed point $\gamma \in \mathcal{S}$, which is given by the Maxwellian, or Gaussian, density:

$$
\begin{equation*}
\gamma(d v)=\frac{e^{-\frac{d}{2}|v|^{2}}}{\left(2 \pi d^{-1}\right)^{d / 2}} d v \tag{1.9}
\end{equation*}
$$

Measuring convergence to the Boltzmann equation. To discuss the convergence of Kac's process to the Boltzmann equation, we will work with the following Wasserstein metric on $\mathcal{S}$. Consider the Sobolev space of test functions

$$
\begin{align*}
X & =W^{1, \infty}\left(\mathbb{R}^{d}\right)=\left\{\text { Bounded, Lipschitz functions } f: \mathbb{R}^{d} \rightarrow \mathbb{R}\right\} ;  \tag{1.10}\\
\|f\|_{X} & :=\max \left(\sup _{v}|f|(v), \sup _{v \neq w} \frac{|f(v)-f(w)|}{|v-w|}\right) . \tag{1.11}
\end{align*}
$$

We write $B_{X}$ for the unit ball of $X$; that is, those functions which are 1-bounded and 1-Lipschitz. Given a function $f$ on $\mathbb{R}^{d}$, we write $\hat{f}$ for the function

$$
\begin{equation*}
\hat{f}(v)=\frac{f(v)}{1+|v|^{2}} \tag{1.12}
\end{equation*}
$$

We write $\mathcal{A}$ for the space of weighted-Lipschitz functions:

$$
\begin{equation*}
\mathcal{A}:=\left\{f: \mathbb{R}^{d} \rightarrow \mathbb{R}: \hat{f} \in X,\|\hat{f}\|_{X} \leq 1\right\} \tag{1.13}
\end{equation*}
$$

We will also write

$$
\begin{equation*}
\mathcal{A}_{0}=\left\{f: \mathbb{R}^{d} \rightarrow \mathbb{R}: \hat{f} \in L^{\infty}\left(\mathbb{R}^{d}\right),\|\hat{f}\|_{\infty} \leq 1\right\} \tag{1.14}
\end{equation*}
$$

The weighted Wasserstein metric $W$ is given by the duality:

$$
\begin{equation*}
W(\mu, v):=\sup _{f \in \mathcal{A}}|\langle f, \mu-v\rangle| . \tag{1.15}
\end{equation*}
$$

We make the following remark on alternative possible choices of metric. Our metric $W$ is closely related to the $p$-Wasserstein metrics $W_{p}$ on the subspaces $\mathcal{S}^{p}$, given by

$$
\begin{equation*}
W_{p}(\mu, v)=\inf \left\{\int_{\mathbb{R}^{d}}|v-w|^{p} \pi(d v, d w): \pi \text { is a coupling of } \mu \text { and } v\right\} \tag{1.16}
\end{equation*}
$$

In the special case $p=1$, the metric $W_{1}$ is known as the Monge-KantorovichWasserstein (MKW) metric, and can alternatively be given by

$$
\begin{equation*}
W_{1}(\mu, v)=W\left(\frac{\mu}{1+|v|^{2}}, \frac{v}{1+|v|^{2}}\right) . \tag{1.17}
\end{equation*}
$$

It is straightforward to check that, on the space $\mathcal{S}$, the metrics $W, W_{1}, W_{2}$ all induce the same topology, and that for some absolute constant $C$, we have the bound $W_{1} \leq C W$ on $\mathcal{S}$. Moreover, on the subspaces $\mathcal{S}_{a}^{k}$ defined in (1.6), with $k>2$, we can find explicit bounds $W \leq C W_{1}^{\alpha}$, with $\alpha \in(0,1)$.

We now state the motivating result of [33] on the convergence of the Kac process to the Boltzmann equation.

Proposition 1 ([33], Theorem 10.1). Let $k>2$. We say that a family $\left(\mu_{t}\right)_{t \geq 0}$ is locally $\mathcal{S}^{k}$-bounded if $\sup _{s \leq t} \Lambda_{k}\left(\mu_{s}\right)<\infty$ for any $t \geq 0$.

For any $\mu_{0} \in \mathcal{S}^{k}$, there is a unique locally $\mathcal{S}^{k}$-bounded solution to the Boltzmann equation (BE), starting from $\mu_{0}$; we write this solution as $\left(\phi_{t}\left(\mu_{0}\right)\right)_{t \geq 0}$.

Moreover, for any $\epsilon>0, t_{\mathrm{fin}}<\infty, \lambda<\infty$, there exist constants $C\left(\epsilon, \lambda, k, t_{\mathrm{fin}}\right)<$ $\infty$ and $\alpha(d, k)>0$ such that, whenever $\left(\mu_{t}^{N}\right)_{t \geq 0}$ is a Kac process on $N \geq 1$ particles, with $\Lambda_{k}\left(\mu_{0}^{N}\right) \leq \lambda, \Lambda_{k}\left(\mu_{0}\right) \leq \lambda$, we have

$$
\begin{equation*}
\mathbb{P}\left(\sup _{t \leq t_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)>C\left(W\left(\mu_{0}^{N}, \mu_{0}\right)+N^{-\alpha}\right)\right)<\epsilon \tag{1.18}
\end{equation*}
$$

For $d \geq 3$ and $k>8$, we can take $\alpha=\frac{1}{d}$.
While the study of the convergence of the Kac process to the Boltzmann equation is a well-known and extensively studied topic, this is most usually studied through the propagation of chaos, discussed below, by contrast to the pathwise
style of estimate here which we seek to emulate. We note that the existence of solutions is known [24] for the case $k=2$, but that no quantitative results are known for the convergence of the Kac process in this case.

From existence and uniqueness, we can consider the Boltzmann equation as describing a nonlinear semigroup of flow operators on $\left(\phi_{t}\right)_{t \geq 0}$ on $\bigcup_{k>2} \mathcal{S}^{k}$. To prove Proposition 1, Norris [33] introduces a family of random linear operators $E_{\mathrm{st}}$, and develops a representation formula in terms of these operators, which will be reviewed in Sections 4, 8. Crucial to the proof are estimates for the operator norms of $E_{\text {st }}$, which are obtained by Grönwall-style estimates, and as a result, the constant $C$ depends badly on the terminal time $t_{\text {fin }}$, with a priori exponential growth. Our work was inspired by the observation that strong stability estimates for the nonlinear semigroup $\left(\phi_{t}\right)$, proven by Mischler and Mouhot [29], allow us to avoid using Grönwall-style estimates, and hence obtain estimates with better long-time properties.

Chaoticity We will also discuss the notion of chaoticity, which is the usual framework used to analyse the convergence of the Kac process to the Boltzmann equation. In this context, it is natural to preserve the labels on the particles, and to consider the labelled Kac process $\mathcal{V}_{t}^{N}=\left(v_{1}(t), \ldots, v_{N}(t)\right)$, taking values in the labelled Boltzmann sphere

$$
\begin{equation*}
\mathbb{S}^{N}=\left\{\left(v_{1}, \ldots, v_{N}\right) \in\left(\mathbb{R}^{d}\right)^{N}: \sum_{i=1}^{N} v_{i}=0, \sum_{i=1}^{N}\left|v_{i}\right|^{2}=N\right\} . \tag{1.19}
\end{equation*}
$$

We may recover recover $\mathcal{S}_{N}$ by taking empirical measures:

$$
\begin{equation*}
\theta_{N}: \mathbb{S}^{N} \rightarrow \mathcal{S}_{N} ; \quad\left(v_{1}, \ldots, v_{N}\right) \mapsto \frac{1}{N} \sum_{i=1}^{N} \delta_{v_{i}} \tag{1.20}
\end{equation*}
$$

Moreover, if $\mathcal{V}_{t}^{N}$ is a labelled Kac process, then $\mu_{t}^{N}=\theta_{N}\left(\mathcal{V}_{t}^{N}\right)$ is an unlabelled Kac process. We write $\mathcal{L} \mathcal{V}_{t}^{N}$ for the law of $\left(v_{1}(t), \ldots, v_{N}(t)\right)$ on $\mathbb{S}^{N}$. We will measure chaoticity using the following (unweighted) Wasserstein metrics on probability measures on $\left(\mathbb{R}^{d}\right)^{l}$ for all $l \geq 1$, defined in a similar way to (1.15):

$$
\begin{equation*}
\mathcal{W}_{1, l}\left(\mathcal{L}, \mathcal{L}^{\prime}\right)=\sup \left\{\int_{\left(\mathbb{R}^{d}\right)^{l}} f(V)\left(\mathcal{L}(d V)-\mathcal{L}^{\prime}(d V)\right)\right\} \tag{1.21}
\end{equation*}
$$

where the supremum is over all functions $f$ of the form $f=f_{1} \otimes f_{2} \otimes \cdots \otimes f_{l}$, with each $f_{i}$ a bounded and Lipschitz test function, $f_{i} \in B_{X}$, and the subscript $l$ recalls the relevant dimension. We now recall the following definition from [23].

DEFINITION (Finite dimensional chaos). For each $N$, let $\mathcal{L}^{N}$ be a law on $\mathbb{S}^{N}$, which is symmetric under permutations of the indexes. We say that $\left(\mathcal{L}^{N}\right)_{N \geq 2}$ is $\mu$-chaotic, if, for all $l \geq 1$, we have

$$
\begin{equation*}
\mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{L}^{N}\right], \mu^{\otimes l}\right) \rightarrow 0 \tag{1.22}
\end{equation*}
$$

where $\Pi_{l}$ denotes the marginal distribution on the first $l$ factors.

A stronger notion, put forward by Mischler and Mouhot [29], is that of infinitedimensional chaos, which allows the number of marginals $l$ to vary with $N$ :

$$
\begin{equation*}
\max _{1 \leq l \leq N}\left[\frac{1}{l} \mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{L}^{N}\right], \mu^{\otimes l}\right)\right] \rightarrow 0 \tag{1.23}
\end{equation*}
$$

Kac proposed the following propagation of chaos property. Let $\left(\mathcal{V}_{t}^{N}\right)_{t \geq 0}$ be a labelled Kac process, such that the initial distribution $\mathcal{L} \mathcal{V}_{t}^{N}$ is $\mu_{0}$-chaotic. Then, for all times $t \geq 0$, the law $\mathcal{L} \mathcal{V}_{t}^{N}$ will be $\phi_{t}\left(\mu_{0}\right)$-chaotic, where $\phi_{t}\left(\mu_{0}\right)$ is the solution to the Boltzmann equation starting at $\mu_{0}$. This is the original sense in which Kac proposed to study the convergence of his model to the Boltzmann equation, and has been extensively studied; key previous results in this direction will be discussed in our literature review.
1.1. Main results. We now state the main results of the paper, concerning the long-time nature of the convergence to the Boltzmann flow. Our first theorem controls the deviation from the Boltzmann flow at a single, deterministic time $t \geq 0$, which we refer to as a pointwise estimate. We highlight that this estimate is uniform in time.

THEOREM 1.1. Let $0<\epsilon<\frac{1}{d}$ and let $a \geq 1$. For sufficiently large $k$, depending on $\epsilon$, d, let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process in dimension $d \geq 3$, and let $\mu_{0} \in \mathcal{S}^{k}$, satisfying the moment bounds

$$
\begin{equation*}
\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a ; \quad \Lambda_{k}\left(\mu_{0}\right) \leq a \tag{1.24}
\end{equation*}
$$

Then for some $C=C(\epsilon, d, k)<\infty$ and $\zeta=\zeta(d)>0$, we have the uniform bound

$$
\begin{equation*}
\sup _{t \geq 0}\left\|W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)\right\|_{L^{2}(\mathbb{P})} \leq C a\left(N^{\epsilon-1 / d}+W\left(\mu_{0}^{N}, \mu_{0}\right)^{\zeta}\right) \tag{1.25}
\end{equation*}
$$

This generalises, by conditioning, to the case where the initial data $\mu_{0}^{N}$ is random, provided that $\mathbb{E} \Lambda_{k}\left(\mu_{0}^{N}\right) \leq a$.

This result is, to the best of our knowledge, new, although an equivalent result is known for Maxwell molecules [8]. We will see, in Theorem 1.7, that estimates of this form imply the propagation of chaos for hard spheres, in the sense of (1.22)(1.23), with better rates than found in [29] for the hard spheres process.

Our second main theorem controls, in $L^{p}(\mathbb{P})$, the maximum deviation from the Boltzmann flow up to a time $t_{\text {fin }}$, in analogy with Proposition 1. We refer to this as a pathwise, local uniform in time estimate.

THEOREM 1.2. Let $0<\epsilon<\frac{1}{2 d}, a \geq 1$ and $p \geq 2$. For sufficiently large $k \geq 0$, depending on $\epsilon$, det $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process on $N \geq 2$ particles and let $\mu_{0} \in$ $\mathcal{S}^{k}$, with initial moments

$$
\begin{equation*}
\Lambda_{k p}\left(\mu_{0}^{N}\right) \leq a^{p} ; \quad \Lambda_{k}\left(\mu_{0}\right) \leq a \tag{1.26}
\end{equation*}
$$

For some $\alpha=\alpha(\epsilon, d, p)>0$ and $C=C(\epsilon, d, p, k)<\infty$ and $\zeta=\zeta(d)>0$, we can estimate, for all $t_{\mathrm{fin}} \geq 0$,

$$
\begin{equation*}
\left.\left\|\sup _{t \leq t_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)\right\|_{L^{p}(\mathbb{P})} \leq C a\left(\left(1+t_{\mathrm{fin}}\right)^{1 / p} N^{-\alpha}+W\left(\mu_{0}^{N}, \mu_{0}\right)^{\zeta}\right)\right) \tag{1.27}
\end{equation*}
$$

$\alpha$ is given explicitly by

$$
\begin{equation*}
\alpha=\frac{p^{\prime}}{2 d}-\epsilon, \tag{1.28}
\end{equation*}
$$

where $1<p^{\prime} \leq 2$ is the Hölder conjugate to $p$.

At the end of this section, we will discuss related results, and how they may be compared to this estimate.

An unfortunate feature of these approximation theorems is the dependence on the unknown, and potentially large, moment index $k$; a trivial reformulation which avoids this is to ask instead for an exponential moment bound $\left\langle e^{z|v|}, \mu_{0}^{N}\right\rangle \leq b$, for some $z>0$. We will also prove the following variant of the theorems above which allows us to use any moment estimate higher than second.

THEOREM 1.3 (Convergence with few moment estimates). Let $k>2$ and $a \geq 1$. Let $\left(\mu_{t}^{N}\right)$ be an $N$-particle Kac process, and $\mu_{0}$ in $\mathcal{S}$ with initial moment estimates

$$
\begin{equation*}
\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a ; \quad \Lambda_{k}\left(\mu_{0}\right) \leq a \tag{1.29}
\end{equation*}
$$

There exists $\epsilon=\epsilon(d, k)>0$ and a constant $C=C(d, k)$ such that

$$
\begin{equation*}
\sup _{t \geq 0}\left\|W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)\right\|_{L^{1}(\mathbb{P})} \leq C a\left(N^{-\epsilon}+W\left(\mu_{0}^{N}, \mu_{0}\right)^{\epsilon}\right) \tag{1.30}
\end{equation*}
$$

For a local uniform estimate, if $p \geq 2$, then there exists a constant $C=C(d, k, p)$ and $\epsilon=\epsilon(d, k, p)>0$ such that, for all $t_{\mathrm{fin}}<\infty$,

$$
\begin{equation*}
\left\|\sup _{t \leq f_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)\right\|_{L^{1}(\mathbb{P})} \leq C a\left(\left(1+t_{\mathrm{fin}}\right)^{1 / p} N^{-\epsilon}+W\left(\mu_{0}^{N}, \mu_{0}\right)^{\epsilon}\right) . \tag{1.31}
\end{equation*}
$$

In the course of proving this result, we will see that the higher moment conditions are only required to obtain the optimal rates on a very short time interval $\left[0, u_{N}\right]$ and, in particular, we can obtain very good time-dependence without higher moment estimates.

We also study the long-time behaviour of the Kac process. We cannot extend Theorem 1.2 to control the maximum deviations over all times $t \geq 0$, due to the following recurrence features of the Kac process.

THEOREM 1.4. There exists a universal constant $C>0$ such that, for every $N$, for every $k>2$ and $a>1$, there exists a Kac process $\left(\mu_{t}^{N}\right)_{t \geq 0}$ with initial moment $\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a$ but, almost surely,

$$
\begin{equation*}
\limsup _{t \rightarrow \infty} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right) \geq 1-\frac{C}{\sqrt{N}} \tag{1.32}
\end{equation*}
$$

Hence we cannot omit the factor of $\left(1+t_{\mathrm{fin}}\right)^{1 / p}$ in Theorem 1.2.
In keeping with the terminology above, we say that there is no pathwise, uniform in time estimate. In the course of proving Theorem 1.4, we will show that the longtime deviation (1.32) is typical for the Kac process. We will show that the Kac process returns, infinitely often, to 'highly ordered' subsets of $\mathcal{S}_{N}$, which are far from the Boltzmann flow. However, we make the following remark on the times necessary for such deviations to occur.

## Corollary 1.5. Define

$$
\begin{equation*}
T_{N, \epsilon}=\inf \left\{t \geq 0: W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)>\epsilon\right\} . \tag{1.33}
\end{equation*}
$$

Let $\left(\mu_{t}^{N}\right)$ be a family of Kac processes with an initial exponential moment bound: $\left\langle e^{z|v|}, \mu_{0}^{N}\right\rangle \leq b$, for some $z>0$ and $b>0$. Let $\mu_{0} \in \mathcal{S}$ satisfy $\left\langle e^{z|v|}, \mu_{0}\right\rangle \leq b$, and suppose that $W\left(\mu_{0}^{N}, \mu_{0}\right) \rightarrow 0$ in probability.

Let $t_{N, \epsilon, \delta}$ be the quantile constants of $T_{N, \epsilon}$ under $\mathbb{P}$; that is,

$$
\begin{equation*}
\mathbb{P}\left(T_{N, \epsilon} \leq t_{N, \epsilon, \delta}\right) \geq \delta \tag{1.34}
\end{equation*}
$$

Then, for fixed $\epsilon, \delta>0, t_{N, \epsilon, \delta} \rightarrow \infty$, faster than any power of $N$.
This follows as an immediate consequence of Theorem 1.2. Taken together with Theorem 1.4, we see that macroscopic deviations occur, but typically at times growing faster than any power of $N$.

In the course of proving Theorems 1.1, 1.2, we will establish the following continuity estimate for the Boltzmann flow $\phi_{t}$ measured in the Wasserstein distance $W$, which may be of independent interest.

THEOREM 1.6. There exist constants $k, C, w$ depending only on $d$ such that, whenever $a \geq 1$ and $\mu, v \in \mathcal{S}_{a}^{k}$, we have the estimate

$$
\begin{equation*}
W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) \leq C e^{w t} a W(\mu, \nu) \tag{1.35}
\end{equation*}
$$

Moreover, for all $k>2$, there exist constants $C=C(k, d)$ and $\zeta=\zeta(k, d)>0$ such that, whenever $\mu, \nu \in \mathcal{S}_{a}^{k}$, we have the estimate

$$
\begin{equation*}
\sup _{t \geq 0} W\left(\phi_{t}(\mu), \phi_{t}(v)\right) \leq C a W(\mu, v)^{\zeta} \tag{1.36}
\end{equation*}
$$

In the second part of the theorem, and in Theorems 1.1, 1.2 above, the exponent $\zeta$ can be taken to be $\lambda_{0} /\left(\lambda_{0}+2 w\right)$ by making $k$ large enough, where $w$ is as in the first part of the theorem, and $\lambda_{0}=\lambda_{0}(d)>0$ is the spectral gap of the linearised Boltzmann operator. While it may be possible to obtain better continuity results, with $\zeta$ close to 1 , we will not explore this here.

Due to a result of Sznitman [37], the property of chaoticity is equivalent to convergence of the empirical measures in expected Wasserstein distance $W$. Therefore, as mentioned before, the theorems displayed above are closely related to the propagation of chaos for the hard-spheres Kac process, proven in [29]. We now give a chaoticity result which may be derived from the previous theorems.

THEOREM 1.7 (Theorems 1.1, 1.3 as a chaos estimate). We can view Theorems 1.1, 1.3 as propagation of chaos and conditional propagation of chaos, as follows.

We denote $\mathcal{P}_{t}^{N}\left(\mathcal{V}^{N}, \cdot\right)$ the transition probabilities of the $N$-particle labelled Kac process, started at $\mathcal{V}^{N} \in \mathbb{S}^{N}$. We form the symmetrised version, which we denote $\mathcal{P}_{t}^{N}\left(\mu^{N}, \cdot\right)$ by

$$
\begin{equation*}
\mathcal{P}_{t}^{N}\left(\mu^{N}, A\right)=\frac{1}{\# \theta_{N}^{-1}\left(\mu^{N}\right)} \sum_{\mathcal{V}^{N} \in \theta_{N}^{-1}\left(\mu^{N}\right)} \mathcal{P}_{t}^{N}\left(\mathcal{V}^{N}, A\right) \tag{1.37}
\end{equation*}
$$

Let $k>2$ and $a \geq 1$, and suppose $\mu_{0}^{N} \in \mathcal{S}_{N}$ satisfies a moment bound $\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a$. Then we can estimate

$$
\begin{equation*}
\sup _{t \geq 0} \max _{1 \leq l \leq N} \frac{\mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{P}_{t}^{N}\left(\mu_{0}^{N}, \cdot\right)\right], \phi_{t}\left(\mu_{0}^{N}\right)^{\otimes l}\right)}{l} \leq C a N^{-\beta} \tag{1.38}
\end{equation*}
$$

for some constants $C=C(d, k)<\infty ; \beta=\beta(d, k)>0$. This has the following consequences:
(i) (Chaotic case) Let $k$, a be as above, and suppose $\mu_{0} \in \mathcal{S}$ satisfies $\Lambda_{k}\left(\mu_{0}\right) \leq$ $a$.

Construct initial data $\mathcal{V}_{0}^{N}=\left(v_{1}(0), \ldots, v_{N}(0)\right)$ as follows. Let $u_{1}, \ldots, u_{N}$ be an independent, and identically distributed sample from $\mu_{0}$. Define

$$
\begin{equation*}
\bar{u}_{N}=\frac{1}{N} \sum_{i=1}^{N} u_{i} ; \quad s_{N}=\frac{1}{N} \sum_{i=1}^{N}\left|u_{i}-\bar{u}_{N}\right|^{2} \tag{1.39}
\end{equation*}
$$

and set

$$
\begin{align*}
v_{i}(0) & =s_{N}^{-1 / 2}\left(u_{i}-\bar{u}_{N}\right), \quad i=1,2, \ldots, N  \tag{1.40}\\
\mathcal{V}_{0}^{N} & =\left(v_{1}(0), \ldots, v_{N}(0)\right) .
\end{align*}
$$

Let $\mathcal{V}_{t}^{N}$ be a labelled Kac process starting from $\mathcal{V}_{0}^{N}$. Then there exist constants $C=C(d, k)<\infty ; \beta=\beta(d, k)>0$ such that

$$
\begin{equation*}
\sup _{t \geq 0} \max _{1 \leq l \leq N} \frac{\mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{L} \mathcal{V}_{t}^{N}\right], \phi_{t}\left(\mu_{0}\right)^{\otimes l}\right)}{l} \leq C N^{-\beta} \tag{1.41}
\end{equation*}
$$

(ii) (General case) Let a, $k$ be as above, and suppose that $\left(\mathcal{V}_{t}^{N}\right)_{t \geq 0}$ are labelled Kac processes such that the empirical measures $\mu_{0}^{N}$ satisfy

$$
\begin{equation*}
\mathbb{E} \Lambda_{k}\left(\mu_{0}^{N}\right) \leq a \tag{1.42}
\end{equation*}
$$

Then we have the estimate

$$
\begin{equation*}
\sup _{t \geq 0} \max _{1 \leq l \leq N} \frac{\mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{L} \mathcal{V}_{t}^{N}\right], \mathcal{L}_{t}^{l}\right)}{l} \leq C a N^{-\beta} \tag{1.43}
\end{equation*}
$$

for $C$ and $\beta$ as in the main statement, and where $\mathcal{L}_{t}^{l}$ is the probability measure given by

$$
\begin{equation*}
\mathcal{L}_{t}^{l}=\mathbb{E}\left[\phi_{t}\left(\mu_{0}^{N}\right)^{\otimes l}\right] \tag{1.44}
\end{equation*}
$$

## REMARK 1.8.

(i) Roughly, (1.38) says that, conditional on the observation of the empirical data $\mu_{0}^{N}$ at time 0 , the law $\mathcal{L} \mathcal{V}_{t}^{N}$ is quantitatively $\phi_{t}\left(\mu_{0}^{N}\right)$-chaotic. This may be viewed as propagation of chaos, with the heuristic that 'conditional on $\mu_{0}^{N}, \mathcal{V}_{0}^{N}$ is $\mu_{0}^{N}$-chaotic'. We term this conditional propagation of chaos. In this spirit, we may view the main estimate (1.38) and point (ii) as a quenched and annealed pair.
(ii) The polynomial result obtained here improves on the previously known result [29], Theorem 6.2, for the hard spheres chaos. This improvement is due to the continuity estimate (1.36), which improves on the corresponding estimate in [29], equations $6.39,6.42$; we could derive the chaoticity estimate (1.38) by using the estimate (1.36) in the arguments of [29], Section 6, at the cost of potentially requiring a stronger initial moment control. We will recall the relevant arguments for completeness, and this will be discussed in the literature review.
(iii) This construction of chaotic initial data in point (i) is due to [33], Proposition 9.2, which may be thought of as 'as close to perfect independence as possible'.
(iv) We will show that the main point can be deduced from Theorems 1.1 or 1.3. However, we will see in Section 10 that deriving either of these from this result appears to be no less technical than the main proof presented in Section 6.

In our arguments, we will frequently encounter numerical constants which are ultimately absorbed into the constants $C$ whose dependence is specified in the relevant theorem. To ease notation, we will denote inequality, up to such a constant, by $\lesssim$.

### 1.2. Plan of the paper. Our programme will be as follows:

i. In the remainder of this section, we will present a review of known results in the study of the Kac process and similar models. We will then discuss several aspects of our results, and how they may be interpreted.
ii. For later convenience, we discuss some classical moment estimates for the Kac process and the Boltzmann equation. These allow us to stochastically control the weights $\Lambda_{k}$ in appropriate $L^{p}$ spaces.
iii. We cite the analytical regularity and stability estimates from Mischler and Mouhot [29]. The stability estimates, in particular, are crucial to obtaining the good time-dependence in Theorems 1.1, 1.2.
iv. As a first application of the stability estimates, we analyse the continuity of the Boltzmann flow $\phi_{t}$ on subsets $\mathcal{S}_{a}^{k}$, with respect to the metric $W$, and uniformly in time. This is the content of Theorem 1.6, and allows us to reduce Theorems 1.1, 1.2 to the special case $\mu_{0}=\mu_{0}^{N}$.
v. We use ideas of infinite-dimensional differential calculus, developed by [29], to prove an interpolation decomposition of the difference $\mu_{t}^{N}-\phi_{t}\left(\mu_{0}^{N}\right)$. This is the key identity used for the proofs of Theorems 1.1, 1.2, as all of the terms appearing in our formula can be controlled by the stability estimates.
vi. We then turn to the proof of Theorem 1.1. The main technical aspect is the control of a family of martingales $\left(M_{t}^{N, f}\right)_{f \in \mathcal{A}}$, uniformly in $f$. This is obtained using a quantitative compactness argument similar to that in [33].
vii. For a local uniform analysis, we first adopt the ideas of Theorem 1.1 to a local uniform setting, with suitable adaptations, to state a local uniform martingale estimate, and deduce a preliminary, weak version of Theorem 1.2 with worse dependence in $t_{\mathrm{fin}}$. We then use the stability estimates to 'bootstrap' to the improved estimate Theorem 1.2, and finally return to prove the local martingale estimate.
viii. We next prove Theorem 1.3. The strategy here is to use a localised form of the main argument from [33] to control behaviour on a very short time interval [ $0, u_{N}$ ], and use the previous results, together with the moment production property recalled in Section 2, to control behaviour at times larger than $u_{N}$.
ix. We prove Theorem 1.4, based on relaxation to equilibrium.
x. Finally, we prove the chaoticity result Theorem 1.7. This proof follows a similar pattern to the proof in [29], using our estimates.
1.3. Literature review. We will now briefly discuss related works, to which our results may be compared.

1. Probabilistic techniques for the Kac process and Boltzmann equation. The probabilistic, pathwise approach to the Kac process was pioneered by Tanaka [40, 41], who constructed a Markov process describing the velocity of a 'typical' particle in the Kac process with Maxwell molecules, and whose law at time $t$ is the solution to the associated Boltzmann equation. This was generalised by Fournier and Méléard [15] to include the cases without cutoff, and for non-Maxwellian molecules. A similar idea was used by Rousset [35] to prove convergence to equilibrium as $t \rightarrow \infty$.

Our main convergence results may be compared to the motivating work of Norris [33], of which the main result is recalled in Proposition 1 above. Theorem 1.2
improves on Proposition 1 in two notable ways. First, we have much better asymptotic behaviour in the time-horizon $t_{\text {fin }}$, which was the original motivation for our work. Second, we control the deviation in the stronger sense of $L^{p}$, rather than in probability; this arises as a result of using moment estimates within the framework of a 'growth control', rather than excluding events of small probability where the moments are large. We also remark that the analysis of the martingale term in Sections 6, 7 is simplified from the equivalent analysis in [33], Theorem 1.1, by our 'interpolation decomposition', Formula 5.1, which removes anticipating behaviour.
2. Propagation of chaos for the Kac process. The problem of propagation of chaos for the Kac process and Boltzmann equation has been extensively studied. The earliest results in this direction are due to McKean [27], Grünbaum [18], and Sznitman [36], and prove the qualitative statement (1.22) for the cases of the hard spheres kernel considered here, or for the related case of Maxwell molecules. Recent work has produced quantitative estimates: Mischler and Mouhout [29] showed propagation of infinite-dimensional chaos (1.23) for both hard spheres and Maxwell molecules. The estimates are uniform in time, with a quantitative estimate going as $(\log N)^{-r}$ for the hard spheres case. As remarked above, our estimates (Theorem 1.1, 1.3, 1.7) improve this rate; this improvement is due to the improvement of Theorem 1.6 over the corresponding estimate in [29], and this will be discussed further below. More recently, [8] proved a chaoticity estimate for Maxwell molecules in $d=3$, measured in the $L^{2}(\mathbb{P})$ norm of Wasserstein ${ }_{2}$ distance (1.16), and with an almost optimal rate $N^{\epsilon-1 / 3}$, which is almost completely analogous to Theorem 1.1.
3. Propagation of chaos for related models. We also mention the study of other models in kinetic theory where chaoticity has been studied. Malrieu [25] studied a McKean-Vlasov model related to granular media equations, and deduced chaoticity for a related system. The main estimate here is a uniform in time estimate, similar in nature to Theorem 1.1. Similarly, Bolley, Guillin and Malrieu [2] have also proven propagation of chaos for a particle system associated to a Vlasov-Focker-Plank equation, through a pointwise convergence result. Most recently, Durmus et al. [12] have proved a uniform in time chaoticity estimate based on a coupling approach, for the case with a confinement potential. Both of these models are amenable to the general framework of [29], and propagation of chaos for these models has been proven using the same techniques in a companion paper [30].

We may also compare Theorem 1.2 to a result of Bolley, Guillin and Villani [3], Theorem 2.9, which proves exponential concentration of the maximum $\sup _{t \leq \mathrm{f}_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}(\mu)\right)$ about 0 , for McKean-Vlasov dynamics. This improves upon the rates $\mathcal{O}\left(N^{-\infty}\right)$ which would be obtained using Theorem 1.2, but does not produce an explicit $L^{p}(\mathbb{P})$ bound. More recently, Holding [21] proved a result similar to Theorem 1.2 for McKean-Vlasov systems interacting through a Hölder
continuous force, in order to deduce propagation of chaos. However, neither of these results track the dependence in the terminal time $t_{\mathrm{fin}}$, and so may have much weaker time dependence than our result. To the best of our knowledge, no local uniform estimate for the McKean-Vlasov system exists which seeks to optimise time dependence in the spirit of Theorem 1.2; the applicability of our methods to this system will be considered in the discussion section below.

The notion of chaoticity has also been studied in more abstract settings. Sznitman [37] has studied equivalent conditions for a family of measures to be chaotic, and Gottlieb [16] has produced a necessary and sufficient condition for families of Markov chains to propagate chaoticity.
4. Relaxtion to equilibrium of the Kac process. Kac [23] proposed to relate the asymptotic behaviour of the Boltzmann flow $\phi_{t}\left(\mu_{0}\right)$ to the asymptotic relaxation to equilibrium of the particle system, and conjectured the existence of a spectral gap for the master equation. This has been extensively studied, and Kac's conjecture on the spectral gap positively answered [5, 7, 22, 26]. However, this is not an entirely satisfactory answer for Kac's question on convergence to equilibrium; for chaotic initial data, this still requires times order $\mathcal{O}(N)$ to show relaxation to equilibrium. Carlen et al. also considered in a later paper [6] the more intricate notion of convergence in relative entropy, which somewhat avoids this problem. Mischler and Mouhot [29] answered Kac's question, proving relaxation to equilibrium in Wasserstein distance, uniformly in $N$, for the cases of hard spheres and Maxwell molecules.

We remark that our philosophy is similar to Kac's proposal. Rather than investigating the long-time behaviour of the law $\mathcal{L} \mathcal{V}_{t}^{N}$ of the Kac process, our results use the asymptotics of the Boltzmann equation to partially understand the asymptotics of realisations of the Kac process. Moreover, Theorem 1.4 shows that this cannot be extended to completely understand the full, long-time asymptotics in this sense.
1.4. Discussion of our results. In this subsection, we will discuss the interpretation of our results, especially in view of the framework of chaoticity set out above.

1. Theorems 1.1, 1.2 as a pathwise interpretation of the Boltzmann equation. The main philosophy of our approach follows [33], in considering the Kac process as a Markov chain, and adapting techniques [9, 32] from the general scaling limits of Markov processes.

It is instructive to compare this to the case of a particle system evolving under Vlasov dynamics. In this case, we write $\mu_{t}^{N, \mathrm{~V} 1}$ for the $N$-particle empirical measure, evolving under (nonrandom) Hamiltonian dynamics; Dobrushin [11] showed that $\mu_{t}^{N, \mathrm{Vl}}$ is a weak measure solution to the associated mean field PDE, the Vlasov equation. For the case of Kac dynamics, we may interpret Theorems 1.1, 1.2 as saying that

$$
\begin{equation*}
\forall t \geq 0 \quad \mu_{t}^{N}=\phi_{t}\left(\mu_{0}^{N}\right)+\mathcal{N}_{t}^{N} \tag{1.45}
\end{equation*}
$$

where $\mathcal{N}_{t}^{N}$ is a stochastic noise term, which is small in an appropriate sense. This is a general phenomenon in the 'fluid limit' scaling of Markov processes [9, 32, 33]. In this sense, we may interpret the Boltzmann equation in a pathwise sense; we stress that this interpretation of the Boltzmann equation does not require any chaoticity assumptions on the initial data.
2. Theorem 1.1 as propagation of chaos. It is natural, and instructive, to compare our chaoticity result Theorem 1.7 and our techniques to those of [29], on whose work we build.

In Theorem 1.7, we have improved the rate of chaoticity, from $(\log N)^{-r}$ to a polynomial estimate $N^{-\alpha}$. In proving this result, we will compare our estimates to the estimates of the three error terms $\mathcal{T}_{1}, \mathcal{T}_{2}, \mathcal{T}_{3}$ in the abstract result [29], Theorem 3.1:
(i) The first term $\mathcal{T}_{1}$ is a purely combinatorial term which may be controlled by general, elementary arguments.
(ii) The second error term $\mathcal{T}_{2}$ may be controlled by $\mathbb{E} W\left(\phi_{t}\left(\mu_{0}^{N}\right), \mu_{t}^{N}\right)$, which is a special case of Theorems 1.1, 1.3 with $\mu_{0}=\mu_{0}^{N}$.
(iii) The third error $\mathcal{T}_{3}$ depends on the continuity of the Boltzmann flow $\phi_{t}$ in Wasserstein distance, which is controlled by the Hölder estimates Theorem 1.6.

As mentioned above, the improvement over [29], Theorem 6.2, is due to the improved control on $\mathcal{T}_{3}$, using the estimate (1.36). The controls on $\mathcal{T}_{1}, \mathcal{T}_{2}$ are similar to those in [29], and the claimed result (1.38) follows by using our estimates (6.25), (1.36) in the arguments of [29], Section 6. In order to give a self-contained proof, we will recall the relevant arguments in Section 10.

We also remark that we use each of the assumptions (A1)-(A5) from [29] in our analysis:
(i) Assumption (A1) corresponds to the moment bounds, which follow from the discussion of moment bounds in Proposition 2.
(ii) Assumption (A2(i)) and (A5) concern the continuity of the Boltzmann flow $\phi_{t}$, which is addressed in Theorem 1.6. Assumption (A2(ii)) concerns the continuity of the collision operator $Q$, which is discussed in Section 3.
(iii) Assumption (A3) is the convergence of the generators. A special case of this is the content of Lemma 5.2, which is used to prove our 'interpolation decomposition' Formula 5.1.
(iv) Assumption (A4) is the differential stability of the Boltzmann flow $\phi_{t}$, recalled in Proposition 3, which is crucial to obtaining estimates with good long-time properties.

We will also see that, in order to recover Theorem 1.1 theorem from either of the chaoticity results (Theorem 1.7 or [29], Theorem 6.2), we would need to move a supremum over test functions $f$ inside an expectation, which corresponds to one
of the most technical steps in our proof (Lemmas 6.1, 7.1). Moreover, this technique cannot generalise to produce a pathwise, local uniform convergence result analogous to Theorem 1.2 or Proposition 1.
3. Theorems 1.1, 1.2 without chaoticity. We also remark that neither of the approximation results Theorems 1.1, 1.2 require special preparation of the initial data, beyond a moment estimate; in particular, both are valid even if the initial data $\mathcal{V}_{0}^{N}$ are not chaotic. We will now give an explicit example of such a distribution where this chaoticity property fails.

Example 1.9 (Nonchaotic initial data). Assume that $N$ is a multiple of $2^{d}$. Choose $\Sigma \in S^{d-1}$ uniformly at random, and let $P_{1}, P_{2}, \ldots, P_{2^{d}}$ be the $2^{d}$ points obtained from $\Sigma$ by all reflections in coordinate axes. Let $\mathcal{V}_{0}^{N}$ be given by giving $\frac{N}{2^{d}}$ particles velocity $P_{i}$, for each $i=1, \ldots, 2^{d}$, such that the resulting law $\mathcal{L} \mathcal{V}_{0}^{N}$ is symmetric. Then each marginal distribution is the uniform distribution Uniform $\left(S^{d-1}\right) \in \mathcal{S}$, but there exists a constant $\delta>0$, uniform in $N$, such that

$$
\begin{equation*}
W\left(\mu_{0}^{N}, \operatorname{Uniform}\left(S^{d-1}\right)\right) \geq \delta>0 \tag{1.46}
\end{equation*}
$$

almost surely, where $\mu_{0}^{N}$ is the empirical measure of $\mathcal{V}_{0}^{N}$. In particular, by Sznitman's characterisation, $\mathcal{V}_{0}^{N}$ is not $\operatorname{Uniform}\left(S^{d-1}\right)$-chaotic.

In cases such as this, we may still understand the Boltzmann equation as 'nearly' holding pathwise, in the sense of point 1 . Alternatively, we may view the result Theorem 1.1, and its consequence in Theorem 1.7, as a chaoticity estimate for $\mathcal{V}_{t}^{N}$ about $\phi_{t}\left(\mu_{0}^{N}\right)$, conditional on the initial measure $\mu_{0}^{N}$.
4. Theorem 1.4 in view of the $H$-theorem. As commented after the statement of Theorem 1.4, the key idea of the proof of Theorem 1.4 is that the Kac process $\mu_{t}^{N}$ will, infinitely often, return to 'highly ordered' subsets of the state space $\mathcal{S}_{N}$. However, this appears to contradict a naïve statement of Boltzmann's celebrated $H$-Theorem [4], that "entropy increases". Indeed, this is highly reminiscent of Zermelo's objection, based on Poincaré recurrence of deterministic dynamical systems [43].

However, our results are compatible with the $H$-Theorem, which is rigorously established in [29]. This apparent paradox arises because the $H$-functional, representing the negative of entropy, is a statistical, and not pathwise, concept; that is, $H_{t}$ depends on the data $\mathcal{V}_{t}^{N}$ through the law $\mathcal{L} \mathcal{V}_{t}^{N}$, rather than being a random variable depending directly on a particular observation $\mathcal{V}_{t}^{N}(\omega)$. In particular, for our case, the time $T_{N}$ of reaching the 'ordered state' is a large, random time, and observing a particular realisation $T_{N}(\omega)=t$ tells us very little about the general behaviour $\mathcal{L} \mathcal{V}_{t}^{N}$, and so about the entropy at time $t$.
5. Sharpness of our results. We will now discuss how sharp the main results (Theorems 1.1, 1.2) are, with regards to dependencies in $N$, and the terminal time $t_{\text {fin }}$ in the case of Theorem 1.2.

5a. $N$-dependence. It is instructive to first consider the 'optimal' case of independent particles, for which the empirical measure converges in Wasserstein distance at rate $N^{-1 / d}$. More precisely, for $d \geq 3$, let $\mu \in \mathcal{S}_{a}^{k}$ for $k \geq \frac{3 d}{d-1}$, and let $\mu^{N}$ be an empirical measure for $N$ independent draws from $\mathcal{S}$. Then, for some $C=C(a, k, d)$, we have

$$
\begin{equation*}
\left\|W\left(\mu^{N}, \mu\right)\right\|_{L^{2}(\mathbb{P})} \leq C N^{-1 / d} \tag{1.47}
\end{equation*}
$$

This is shown in [33], Proposition 9.3. Moreover, this rate is optimal: if $\mu$ is absolutely continuous with respect to the underlying Lebesgue measure, then the optimal approximation in $W$ metric is of the order $N^{-1 / d}$, for $d \geq 3$. Results of Talagrand ([38, 39], and discussion in [14]) suggest that this may also be true for higher $L^{p}$ norms, at least for the simple case of the uniform distribution on $(-1,1]^{d}$.

In view of this, we see that the exponent for the pointwise bound is almost sharp, in the sense that we obtain exponents $\epsilon-\frac{1}{d}$ which are arbitrarily close to the optimal exponent $-\frac{1}{d}$, but cannot obtain the optimal exponent itself. This appears to be a consequence of using a particular estimate (3.7) from [29], which is 'almost Lipschitz' in a similar sense. For the local uniform estimate Theorem 1.2, we obtain exponent $-\alpha$, where $\alpha$ is given by

$$
\begin{equation*}
\alpha=-\epsilon+\frac{p^{\prime}}{2 d} ; \quad \frac{1}{p}+\frac{1}{p^{\prime}}=1 \tag{1.48}
\end{equation*}
$$

In the special case $p=2$, this produces the almost sharp exponent as discussed above. However, for $p>2$, the exponents are bounded away from $-\frac{1}{d}$, and so do not appear to be sharp.

5b. Time dependence. In light of Theorem 1.4, we see that we cannot exclude the factor $\left(1+t_{\mathrm{fin}}\right)^{1 / p}$ in Theorem 1.2. Hence, this time dependence is sharp among power laws. However, we do not know what the true sharpest time-dependence is. Similar techniques to those of Graversen and Peskir [17] may be able to provide a sharper bound; we do not explore this here.

We remark that Theorem 1.2 interpolates between almost optimal $N$ dependence at $p=2$, and almost optimal $t_{\mathrm{fin}}$ dependence as $p \rightarrow \infty$. Moreover, by taking $p \rightarrow \infty$, we sacrifice optimal dependence in $N$, but the exponent $\alpha(d, p)$ is bounded away from 0 , and so we have good convergence, on any polynomial time scale. This is the content of Corollary 1.5.
6. Further applicability of our methods in kinetic theory. Finally, we will mention other models in kinetic theory which may be amenable to our techniques.
(a) Sharp $N$ dependence for hard spheres. We believe that our techniques could be modified to prove an estimate for Theorem 1.1, and Theorem 1.2 in the case $p=2$, in order to obtain the optimal rate $N^{-1 / d}$ discussed above; however, this would likely come at the cost of poor dependence in time. Since a similar result (Proposition 1) is already known, and since this is not the spirit of this work in seeking to optimise time dependence, we will not consider this further.
(b) The Kac process on Maxwell molecules. In addition to the hard spheres case analysed here, the main collision kernel of physical interest is the case of Maxwell molecules with or without cutoff. Many of the estimates used in our argument for the hard spheres kernel have an analagous version for Maxwell molecules, including the stability estimates proven in [29]. For this case, a result similar to Theorem 1.1 is already known [8], Theorem 2.
(c) McKean-Vlasov dynamics, and inelastic collisions. Other kinetic system which may be analysed in the framework of [29] include cases of McKean-Vlasov dynamics, and inelastic collisions, coupled to a heat bath, which have been studied in the functional framework of [29] by Mischler, Mouhot and Wennburg in a companion paper [30]. In these cases, the analagous estimates for stability and differentiability, computed in [30], have potentially poor dependence in time. As a result, our methods would still apply, but with correspondingly poor time dependence.

For the case of McKean-Vlasov dynamics without confinement potential, this is a fundamental limitation; Malrieu [25] showed that the propagation of chaos is not uniform in time. Instead, he proposed to study a projected particle system, which satisfies uniform propagation of chaos, and whose limiting flow has exponential convergence to equilibrium [25], Theorem 6.2. This suggests that it may be possible to use our bootstrap method, used in the proof of Theorem 1.2, to obtain a pathwise estimates with good long-time properties, analogous to Theorem 1.2.

We remark that, in the case of McKean-Vlasov dynamics, the presence of Brownian noise may complicate the derivation of the interpolation decomposition (Formula 5.1), which is the key identity required for our argument.
2. Moment estimates. In order to deal with the appearance of the momentbased weights $\Lambda_{k}$ in future calculations, we discuss the moment structure of Kac's process and the Boltzmann equation. That is, we seek bounds on $\Lambda_{k}\left(\mu_{t}\right)$ where $\mu_{t}$ is, correspondingly, either a Kac process, or a solution to the Boltzmann equation.

The results presented here are mostly classical, and the arguments are well known for the Boltzmann equation. Central to the proof is an inequality due to Povzner [34], from which Elmroth [13] deduced global moment bounds for the (function-valued) Boltzmann equation in terms of the moments of the initial data. This conclusion was strengthened to moment production by Desvillettes [10] provided control of an initial moment $\Lambda_{s}\left(\mu_{0}\right)$ for any $s>2$. Wennberg [31, 42] demonstrated an optimal version of this result, only requiring finite initial energy $\left.\left.\langle | v\right|^{2}, \mu_{0}\right\rangle$. Bobylev [1] proved propagation of exponential moments, which may also be applied here as a simplification. These results have been proven for measure-valued solutions of the Boltzmann equation by Lu and Mouhot [24], and the techniques have been applied to the Kac process by Mischler and Mouhot [29] and Norris [33]. We collect below the precise results which we will use.

Proposition 2 (Moment inequalities for the Kac process and Boltzmann equation). We have the following moment bounds for polynomial velocity moments:
(i) Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process on $N \geq 1$ particles, and let $q>2, p \geq 2$ with $q \geq p$. Then there exists a constant $C(p, q)<\infty$ such that, for all $t \geq 0$,

$$
\begin{equation*}
\mathbb{E}\left[\Lambda_{q}\left(\mu_{t}^{N}\right)\right] \leq C\left(1+t^{p-q}\right) \Lambda_{p}\left(\mu_{0}^{N}\right) \tag{2.1}
\end{equation*}
$$

and, for another constant $C=C(q)$,

$$
\begin{equation*}
\mathbb{E}\left(\sup _{0 \leq t \leq t_{\mathrm{fin}}} \Lambda_{q}\left(\mu_{t}^{N}\right)\right) \leq\left(1+C(q) t_{\mathrm{fin}}\right) \Lambda_{q}\left(\mu_{0}^{N}\right) \tag{2.2}
\end{equation*}
$$

(ii) Let $p, q$ be as above, and let, and $\mu_{0} \in \bigcup_{k>2} \mathcal{S}^{k}$. Then there exists a constant $C=C(p, q)$ such that the solution $\phi_{t}\left(\mu_{0}\right)$ to (BE) satisfies

$$
\begin{equation*}
\Lambda_{q}\left(\phi_{t}\left(\mu_{0}\right)\right) \leq C\left(1+t^{p-q}\right) \Lambda_{p}\left(\mu_{0}\right) \tag{2.3}
\end{equation*}
$$

(iii) There exist constants $C_{1}, C_{2}<\infty$ such that, whenever $\mu_{0} \in \bigcup_{k>2} \mathcal{S}^{k}$, we have the bound for all $t \geq 0$,

$$
\begin{equation*}
\int_{0}^{t} \Lambda_{3}\left(\phi_{s}\left(\mu_{0}\right)\right) d s \leq C_{1} t+C_{2}\left\langle\left(1+|v|^{2}\right) \log \left(1+|v|^{2}\right), \mu_{0}\right\rangle \tag{2.4}
\end{equation*}
$$

As a consequence, if $c \geq 0$, then there exists $w<\infty, k<\infty$ such that, for all $t \geq 0$,

$$
\begin{equation*}
\exp \left(c \int_{0}^{t} \Lambda_{3}\left(\phi_{s}\left(\mu_{0}\right)\right) d s\right) \leq e^{w t} \Lambda_{k}\left(\mu_{0}\right) \tag{2.5}
\end{equation*}
$$

The first item is exactly [33], Proposition 3.1. For the second item, if $\phi_{t}\left(\mu_{0}\right)$ is locally $\mathcal{S}^{k}$ bounded for all $k$, then we can apply the same reasoning as the cited proposition to the Boltzmann equation. To remove this condition, we consider the Boltzmann equation started from $\mu_{\delta}=\phi_{\delta}\left(\mu_{0}\right)$ : thanks to the qualitative moment creation property $[10,31]$, the Boltzmann flow started at $\mu_{\delta}$ is locally $\mathcal{S}^{k}$ bounded for all $k$, and so the claimed result holds with $\mu_{\delta}$ in place of $\mu_{0}$. The claimed result may then be obtained by carefully taking the limit $\delta \downarrow 0$.

The first conclusion of item (iii) is proven in [29], equation (6.20), and the final point follows, using the interpolation, for all $\mu \in \mathcal{S}$,

$$
\begin{equation*}
\left\langle\left(1+|v|^{2}\right) \log \left(1+|v|^{2}\right), \mu\right\rangle \leq 8\left(1+\log \Lambda_{5}(\mu)\right) \tag{2.6}
\end{equation*}
$$

In our estimates for the various terms of the interpolation decomposition, we will frequently encounter the weightings $\Lambda_{k}\left(\mu_{t}^{N}\right)$ appearing in the integrand. We refer to points (i)-(ii) of Proposition 2, along with the following lemma, as growth control of the weightings, which allows us to control these factors in suitable $L^{p}$ norms.

Lemma 2.1. Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process on $N \geq 1$ particles, and fix an exponent $k \geq 2$. Then for any time $t \geq 0$, and any measure $\mu^{N}$ which can be obtained from $\mu_{t}^{N}$ by a collision,

$$
\begin{equation*}
\Lambda_{k}\left(\mu^{N}\right) \leq 2^{\frac{k}{2}+1} \Lambda_{k}\left(\mu_{t}^{N}\right) \tag{2.7}
\end{equation*}
$$

Proof. This is immediate, by noting that if $v, v_{\star}$ are pre-collision velocities leading to post-collision $v^{\prime}, v_{\star}^{\prime}$, we have the bound

$$
\begin{align*}
\left(1+\left|v^{\prime}\right|^{2}\right)^{k} & \leq\left(\left(1+|v|^{2}\right)+\left(1+\left|v_{\star}\right|^{2}\right)\right)^{\frac{k}{2}} \\
& \leq 2^{k / 2}\left(\left(1+|v|^{2}\right)^{\frac{k}{2}}+\left(1+\left|v_{\star}\right|^{2}\right)^{\frac{k}{2}}\right) . \tag{2.8}
\end{align*}
$$

Using the same bound for $v_{\star}^{\prime}$ leads to the claimed result.
A final property of the weighting estimates which will prove useful is the following correlation inequality.

Lemma 2.2. Let $k_{1}, k_{2} \geq 2$, and let $\mu \in \mathcal{S}^{k_{1}+k_{2}}$. Then we have

$$
\begin{equation*}
\Lambda_{k_{1}}(\mu) \Lambda_{k_{2}}(\mu) \leq \Lambda_{k_{1}+k_{2}}(\mu) \tag{2.9}
\end{equation*}
$$

Proof. Since the maps $x \mapsto\left(1+|x|^{2}\right)^{k_{i} / 2}$, for $i=1,2$, are both monotonically increasing on $[0, \infty)$, for any $v, v_{\star}$ we have the bound

$$
\begin{equation*}
\left\{\left(1+|v|^{2}\right)^{k_{1} / 2}-\left(1+\left|v_{\star}\right|^{2}\right)^{k_{1} / 2}\right\}\left\{\left(1+|v|^{2}\right)^{k_{2} / 2}-\left(1+\left|v_{\star}\right|^{2}\right)^{k_{2} / 2}\right\} \geq 0 \tag{2.10}
\end{equation*}
$$

Integrating both variables with respect to $\mu$ produces the result.
3. Regularity and stability estimates. In this section, we give precise statements of analytical results concerning the flow maps $\left(\phi_{t}\right)_{t \geq 0}$, and the drift operator $Q$, which will be used in our convergence theorems. We need a combination of regularity for the drift map $Q$, which appears in the proof of Lemma 7.1, and differentiability and stability results for the flow maps $\left(\phi_{t}\right)_{t \geq 0}$.
3.1. Stability estimates. The key component to our analysis of the Kac process is the stability of the limiting Boltzmann equation, that is, that the limit flow suppresses errors, rather than allowing exponential amplification. We begin by defining appropriate linear structures.

DEFINITION 3.1. Consider the space $Y$ of signed measures, given by

$$
\begin{equation*}
Y=\left\{\xi:\|\xi\|_{\mathrm{TV}}<\infty ;\langle 1, \xi\rangle=0\right\} . \tag{3.1}
\end{equation*}
$$

We equip $Y$ with the total variation norm $\|\cdot\|_{\mathrm{Tv}}$. For real $q \geq 0$, we define the subspace $Y_{q}$ of measures with finite $q$ th moments:

$$
\begin{equation*}
\left.Y_{q}=\left\{\xi \in Y:\left.\langle 1+| v\right|^{q},|\xi|\right\rangle<\infty\right\} \tag{3.2}
\end{equation*}
$$

We define the norm with $q$-weighting on $Y_{q}$ by

$$
\begin{equation*}
\left.\|\xi\|_{\mathrm{TV}+q}=\left.\langle 1+| v\right|^{q},|\xi|\right\rangle \tag{3.3}
\end{equation*}
$$

The notation $\|\cdot\|_{\mathrm{TV}+q}$ is chosen to emphasise that this is a total variation norm, with additional polynomial weighting of order $q$, while avoiding potential ambiguity with the $L^{q}$ norms of random variables.

REMARK 3.1. The total variation norms $\|\cdot\|_{\mathrm{TV}+q}$ appearing in the following analysis are much stronger than the Wasserstein distance appearing in Theorems $1.1,1.2,1.3$. We can understand this as follows. Recalling the definitions of $\mathcal{A}, \mathcal{A}_{0}$ in (1.13), (1.14), we note that the TV +2 distance is given by a duality

$$
\begin{equation*}
\|\mu-v\|_{\mathrm{TV}+2}=\sup _{f \in \mathcal{A}_{0}}|\langle f, \mu-v\rangle| \tag{3.4}
\end{equation*}
$$

and, if we write $\left.\mathcal{A}\right|_{r},\left.\mathcal{A}\right|_{r}$ for the restriction of functions to $[-r, r]^{d}$, then the inclusion

$$
\begin{equation*}
\left.\left.\mathcal{A}\right|_{r} \subset \mathcal{A}_{0}\right|_{r} \tag{3.5}
\end{equation*}
$$

is compact in the norm of $\left.\mathcal{A}_{0}\right|_{r}$, by the classical theorem of Arzelá-Ascoli. This is at the heart of a quantitative compactness argument in Lemmas 6.1, 7.1, which allows us to to take the supremum over $f \in \mathcal{A}$ inside the expectation.

We can now state the precise results as they appear in [29], Lemma 6.6.

Proposition 3. Let $\eta \in(0,1)$. Then there are absolute constants $C \in(0, \infty)$ and $\lambda_{0}>0$ such that, for $k$ large enough (depending only on $\eta$ ), and all $\mu, \nu \in \mathcal{S}^{k}$, there is a unique solution $\left(\xi_{t}\right)_{t \geq 0} \subset Y_{2}$ to the linearised differential equation

$$
\begin{equation*}
\xi_{0}=v-\mu ; \quad \partial_{t} \xi_{t}=2 Q\left(\phi_{t}(\mu), \xi_{t}\right) \tag{3.6}
\end{equation*}
$$

This solution satisfies the bounds

$$
\begin{align*}
\left\|\phi_{t}(v)-\phi_{t}(\mu)\right\|_{\mathrm{TV}+2} & \leq C e^{-\lambda_{0} t / 2} \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\eta}  \tag{3.7}\\
\left\|\xi_{t}\right\|_{\mathrm{TV}+2} & \leq C e^{-\lambda_{0} t / 2} \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\eta}  \tag{3.8}\\
\left\|\phi_{t}(v)-\phi_{t}(\mu)-\xi_{t}\right\|_{\mathrm{TV}+2} & \leq C e^{-\lambda_{0} t / 2} \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{1+\eta} \tag{3.9}
\end{align*}
$$

This allows us to define a linear map $\mathcal{D} \phi_{t}(\mu)$ by

$$
\begin{equation*}
\mathcal{D} \phi_{t}(\mu)[v-\mu]:=\xi_{t} \tag{3.10}
\end{equation*}
$$

This linear map will play the role of a functional derivative for the Boltzmann flow $\phi_{t}$ in the calculus developed by [29].

To obtain estimates with the weighted metric $W$, we will use a version of Proposition 3 with the difference $\phi_{t}(\mu)-\phi_{t}(\nu)$ measured in stronger norms $\|\cdot\|_{\mathrm{TV}+q}$. The following estimate may be obtained by a simple interpolation between Propositions 2, 3.

Corollary 3.2. Let $q \geq 2, \eta \in(0,1)$ and $\lambda<\lambda_{0}$. Then for all $k$ large enough, depending on $\eta, \lambda$ and $q$, there exists a constant $C$ such that

$$
\begin{equation*}
\forall \mu, v \in \mathcal{S}^{k}, \quad\left\|\phi_{t}(\mu)-\phi_{t}(v)\right\|_{\mathrm{TV}+q} \leq C e^{-\lambda t / 2} \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\eta} \tag{3.11}
\end{equation*}
$$

We emphasise that the rapid decay is the key property that allows us to obtain good long-time behaviour for our estimates. The pointwise estimate Theorem 1.1 and the initial estimate for pathwise local uniform convergence Lemma 7.2 would hold for estimates

$$
\begin{align*}
\left\|\phi_{t}(v)-\phi_{t}(\mu)\right\|_{\mathrm{TV}+5} & \leq F(t) \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\eta} ;  \tag{3.12}\\
\left\|\phi_{t}(v)-\phi_{t}(\mu)-\xi_{t}\right\|_{\mathrm{TV}+2} & \leq G(t) \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{1+\eta} \tag{3.13}
\end{align*}
$$

for functions $F, G$ such that

$$
\begin{equation*}
\left(\int_{0}^{\infty} F^{2} d t\right)^{1 / 2}<\infty ; \quad \int_{0}^{\infty} G d t<\infty \tag{3.14}
\end{equation*}
$$

The full strength of exponential decay is used to 'bootstrap' to the pathwise local uniform estimate Theorem 1.2, which provides better behaviour in the time horizon $t_{\mathrm{fin}}$, with only a logarithmic loss in the number of particles $N$. Provided that $F \rightarrow 0$ as $t \rightarrow \infty$, we could use the same 'bootstrap', but with a potentially much larger loss in $N$.
3.2. Regularity estimates. For the proof of the local uniform estimate Lemma 7.1, it will be important to control the continuity of $Q$ after application of the flow maps $\phi_{t}$; for brevity, we will write the composition as $Q_{t}=Q \circ \phi_{t}$. We can exploit the use of the stronger $\|\cdot\|_{\mathrm{TV}+2}-$ norm in the stability estimates Proposition 3, to prove a strong notion of continuity for $Q_{t}$, including the dependence on $t$.

It is well known that, for $q \geq 1$, and $\mu, v \in \mathcal{S}^{q+1}$, we have the bilinear estimate

$$
\begin{equation*}
\|Q(\mu)-Q(v)\|_{\mathrm{TV}+q} \lesssim \Lambda_{q+1}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}+(q+1)} \tag{3.15}
\end{equation*}
$$

and, by interpolating, this leads to

$$
\begin{equation*}
\|Q(\mu)-Q(v)\|_{\mathrm{TV}+q} \lesssim \Lambda_{3(q+1)}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\frac{1}{2}} \tag{3.16}
\end{equation*}
$$

Combining this the stability estimate in Corollary 3.2, we deduce the following. For $q \geq 1, \eta \in(0,1)$ and $\lambda<\lambda_{0}$, then there exists $k$ such that, for $\mu, v \in \mathcal{S}^{k}$, we have the estimate

$$
\begin{equation*}
\left\|Q_{t}(\mu)-Q_{t}(v)\right\|_{\mathrm{TV}+q} \lesssim e^{-\lambda t} \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\eta} \tag{3.17}
\end{equation*}
$$

4. Proof of Theorem 1.6. As a first application of the stability estimates, we will now prove Theorem 1.6, which establishes a continuity result for the Boltzmann flow $\left(\phi_{t}\right)$ with respect to our weighted Wasserstein metric $W$. For Theorems 1.1, 1.2, we wish to approximate a given starting point $\mu_{0}$ by an empirical measure $\mu_{0}^{N} \in \mathcal{S}_{N}$ on $N$ points; in this context, the total variation distance is too strong, as there is no discrete approximation to any continuous measure $\mu_{0}$. We therefore seek a continuity estimate for the Boltzmann flow $\phi_{t}$, measured in the Wasserstein distance $W$ defined in (1.15), and which is uniform in time.

The proof combines a representation formula, and associated estimates, from [33], which establishes the first claim; the second claim will then follow using a long-time estimate recalled in Proposition 3. We will first review the definition, and claimed representation formula for the Boltzmann flow.

Definition 4.1 (Linearised Kac process). Write $V=\mathbb{R}^{d}$ and $V^{*}$ for the signed space $V^{*}=V \times\{ \pm 1\}=V^{+} \sqcup V^{-}$. We write $\pi: V^{*} \rightarrow V$ as the projection onto the first factor, and $\pi_{ \pm}: V^{ \pm} \rightarrow V$ for the obvious bijections.

Let $\left(\rho_{t}\right)_{t \geq 0}$ be family of measures on $V=\mathbb{R}^{d}$ such that

$$
\begin{align*}
\left\langle 1, \rho_{t}\right\rangle & \left.=1 ;\left.\quad\langle | v\right|^{2}, \rho_{t}\right\rangle=1 ;  \tag{4.1}\\
\int_{0}^{t} \Lambda_{3}\left(\rho_{s}\right) d s<\infty & \text { for all } t<\infty \tag{4.2}
\end{align*}
$$

The linearised Kac process in environment $\left(\rho_{t}\right)_{t \geq 0}$ is the branching process on $V^{*}$ where each particle of type $(v, 1)$, at rate $2\left|v-v_{\star}\right| \rho\left(d v_{\star}\right) d \sigma$, dies, and is replaced by three particles, of types

$$
\begin{equation*}
\left(v^{\prime}\left(v, v_{\star}, \sigma\right), 1\right) ; \quad\left(v_{\star}^{\prime}\left(v, v_{\star}, \sigma\right), 1\right) ; \quad\left(v_{\star},-1\right) \tag{4.3}
\end{equation*}
$$

where $v^{\prime}, v_{\star}^{\prime}$ are the post-collisional velocities given by (1.2). The dynamics are identical for particles of type $(v,-1)$, with the signs exchanged.

We write $\Xi_{t}^{*}$ for the associated process of unnormalised empirical measures on $V^{*}$, and define a signed measure $\Xi_{t}$ on $V$ by including the sign at each particle:

$$
\begin{equation*}
\Xi_{t}=\Xi_{t}^{+}-\Xi_{t}^{-} ; \quad \Xi_{t}^{ \pm}=\Xi_{t}^{\star} \circ \pi_{ \pm}^{-1} \tag{4.4}
\end{equation*}
$$

We can also consider the same branching process, started from a time $s \geq 0$ instead. We write $E$ for the expectation over the branching process, which is not the full expectation in the case where $\rho$ is itself random. When we wish to emphasise the initial velocity $v$ and starting time $s$, we will write $E_{(s, v)}$ when the process is started from $\Lambda_{0}^{*}=\delta_{(v, 1)}$ at time $s$, and $E_{v}$ in the case $s=0$.

Provided that the initial data $\Xi_{0}$ is finitely supported, one can show that the branching process is almost surely nonexplosive, and that

$$
\begin{equation*}
\left.\left.E_{v_{0}}\langle 1+| v\right|^{2},\left|\Xi_{t}\right|\right\rangle \leq\left(1+|v|^{2}\right) \exp \left[8 \int_{0}^{t} \Lambda_{3}\left(\rho_{s}\right) d s\right] \tag{4.5}
\end{equation*}
$$

REMARK 4.1. We can connect this branching process with a different proof of existence and uniqueness for the difference $\xi_{t}$ in Theorem 3. For existence, consider the linearised Kac process $\left(\Xi_{t}\right)_{t \geq 0}$ in environment $\rho_{t}=\phi_{t}(\mu)$, where particles are initialised at $t=0$ according to a Poisson random measure of intensity

$$
\theta(d v)= \begin{cases}\xi_{0}^{+}(d v)=v(d v) & \text { on } V^{+}  \tag{4.6}\\ \xi_{0}^{-}(d v)=\mu(d v) & \text { on } V^{-}\end{cases}
$$

Let $\xi_{t}=\mathbb{E}\left(\Xi_{t}\right)$, which may be formalised in the sense of a Bochner integral in the weighted space $\left(Y_{2},\|\cdot\|_{\mathrm{TV}+2}\right)$ defined in (3.1). Then the same proof of the representation formula [33], Proposition 4.2, shows that $\partial_{t} \xi_{t}=2 Q\left(\phi_{t}(\mu), \xi_{t}\right)$, and that this solution is unique.

Recall from the Introduction that $\mathcal{A}$ is the set of all functions $f$ on $\mathbb{R}^{d}$, such that $\widehat{f}(v)=\left(1+|v|^{2}\right)^{-1} f(v)$ satisfies

$$
\begin{equation*}
|\widehat{f}(v)| \leq 1 ; \quad \frac{|f(v)-f(w)|}{|v-w|} \leq 1 \quad \text { for all } v \neq w \tag{4.7}
\end{equation*}
$$

From the bound (4.5), we can now define, for functions of quadratic growth,

$$
\begin{equation*}
f_{\mathrm{st}}\left(v_{0}\right)=E_{\left(s, v_{0}\right)}\left[\left\langle f, \Xi_{t}\right\rangle\right] . \tag{4.8}
\end{equation*}
$$

When we wish to emphasise the environment, we will write $f_{\text {st }}[\rho]\left(v_{0}\right)$. We now recall the following estimates from [33].

Proposition 4 (Continuity estimates for $f_{\mathrm{st}}$ ). Fix $t \geq 0$, and let $z_{t}$ be given by

$$
\begin{equation*}
z_{t}=3 \exp \left[8 \int_{0}^{t} \Lambda_{3}\left(\rho_{u}\right) d u\right] \tag{4.9}
\end{equation*}
$$

Then, for $f \in \mathcal{A}$ and $s \leq t$, we have $f_{\mathrm{st}} \in z_{t} \mathcal{A}$. This is, in our notation, a reformulation of [33], Propositions 4.3.

The other result which we will use is the representation formula [33], Proposition 4.2, which expresses the difference of two Boltzmann flows $\phi_{t}(\mu)-\phi_{t}(\nu)$ in terms of the functions $f_{0 t}$. This may be obtained from the proof of [33], Proposition 4.2, without essential modification, as in the proof of [33], Theorem 10.1.

Proposition 5 (Representation formula). Let $\mu, \nu \in \mathcal{S}^{k}$ for some $k>2$, and let $\left(\rho_{t}\right)_{t \geq 0}$ be given by

$$
\begin{equation*}
\rho_{t}=\frac{1}{2}\left(\phi_{t}(\mu)+\phi_{t}(\nu)\right), \tag{4.10}
\end{equation*}
$$

where $\phi_{t}(\mu)$ is the unique, locally $\mathcal{S}^{k}$-bounded solution to the Boltzmann equation, starting at $\mu$, and similarly for $\nu$. Then, for all $f \in \mathcal{A}$, we have

$$
\begin{equation*}
\left\langle f, \phi_{t}(\mu)-\phi_{t}(v)\right\rangle=\left\langle f_{0 t}[\rho], \mu-v\right\rangle . \tag{4.11}
\end{equation*}
$$

Note that the moment production property in Proposition 2 guarantees that (4.2) holds for this environment. This will allow us to find an estimate for the Boltzmann flow $\phi_{t}$ which behaves well in short time. We now give the proof of Theorem 1.6.

Proof of Theorem 1.6. From the representation formula (4.11) and continuity estimate Proposition 4 , for any $f \in \mathcal{A}$,

$$
\begin{equation*}
\left\langle f, \phi_{t}(\mu)-\phi_{t}(v)\right\rangle=\left\langle f_{0 t}[\rho], \mu-v\right\rangle \leq z_{t} W(\mu, v) \tag{4.12}
\end{equation*}
$$

where $\rho_{t}=\left(\phi_{t}(\mu)+\phi_{t}(\nu)\right) / 2$. It therefore suffices to bound

$$
\begin{equation*}
z_{t}:=3 \exp \left(4 \int_{0}^{t}\left[\Lambda_{3}\left(\phi_{s}(\mu)\right)+\Lambda_{3}\left(\phi_{s}(\nu)\right)\right] d s\right) \tag{4.13}
\end{equation*}
$$

Using the logarithmic moment production for the Boltzmann equation recalled in Proposition 2, there exist constants $k, w$ such that

$$
\begin{align*}
z_{t} & \lesssim e^{w t} \Lambda_{k / 2}(\mu) \Lambda_{k / 2}(v) \\
& \lesssim e^{w t} \Lambda_{k / 2}(\mu, v)^{2} \lesssim e^{w t} \Lambda_{k}(\mu, v) \tag{4.14}
\end{align*}
$$

This proves the first claim. For the second claim, we first deal with the case where $k \geq 3$ is large enough that the above holds, and such that the stability estimate Proposition 3 holds with Hölder exponent $\eta=\frac{1}{2}$. Fix $\mu, \nu \in \mathcal{S}_{a}^{k}$, and assume without loss of generality that $0<W(\mu, \nu)<1$. From the stability estimate (3.7), we have

$$
\begin{equation*}
\left\|\phi_{t}(\mu)-\phi_{t}(\nu)\right\|_{\mathrm{TV}+2} \lesssim a^{\frac{1}{2}} e^{-\lambda_{0} t / 2} \tag{4.15}
\end{equation*}
$$

for some constants $\lambda_{0}>0$. It is immediate from the definitions that

$$
\begin{equation*}
W(\mu, v) \leq\|\mu-v\|_{\mathrm{TV}+2} \tag{4.16}
\end{equation*}
$$

and so combining with the previous result, we have

$$
\begin{equation*}
W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) \lesssim a \min \left(e^{-\lambda_{0} t / 2}, W(\mu, \nu) e^{w t}\right) \tag{4.17}
\end{equation*}
$$

The right-hand side is maximised when $e^{-\lambda_{0} t / 2}=W(\mu, v) e^{w t}$, which occurs when

$$
\begin{equation*}
t=-\frac{2}{\lambda_{0}+2 w} \log W(\mu, v) \tag{4.18}
\end{equation*}
$$

Therefore, the maximum value of the right-hand side is

$$
\begin{align*}
\sup _{t \geq 0} W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) & \lesssim a \exp \left(\frac{\lambda_{0}}{\lambda_{0}+2 w} \log W(\mu, v)\right)  \tag{4.19}\\
& =a W(\mu, v)^{\zeta}
\end{align*}
$$

with

$$
\begin{equation*}
\zeta(d)=\frac{\lambda_{0}}{\lambda_{0}+2 w} \tag{4.20}
\end{equation*}
$$

which is the claimed Hölder continuity, for $k$ sufficiently large.

Finally, we deal with the second point for arbitrary $k>2$. This argument uses a localisation principle to control the moments on a very short initial interval $[0, u]$, and may be read as a warm-up to the more involved arguments in the proof of Theorem 1.3.

Let $k_{0}$ be large enough such that the estimate (4.19) holds, and let $\zeta_{0}$ be the resulting exponent. Let $\beta=\frac{k-2}{2}$, let $\mu, \nu$ be as in the statement of the result, and let $u \in(0,1]$ be chosen later. Define

$$
\begin{equation*}
T=\inf \left\{t \geq 0: \Lambda_{3}\left(\rho_{t}\right)>\frac{\beta t^{\beta-1}+1}{2}\right\} \tag{4.21}
\end{equation*}
$$

where $\rho_{t}$ is as above. We now deal with the two cases $T>u, T \leq u$ separately.
If $T>u$, then we have the estimate

$$
\begin{align*}
z_{u} & :=3 \exp \left(4 \int_{0}^{u} \Lambda_{3}\left(\rho_{s}\right) d s\right) \\
& \leq 3 \exp \left(4 \int_{0}^{1} \frac{\beta s^{\beta-1}+1}{2} d s\right) \lesssim 1 . \tag{4.22}
\end{align*}
$$

Using the representation formula in Proposition 5 as in (4.12), we therefore obtain

$$
\begin{equation*}
\sup _{t \leq u} W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) \lesssim W(\mu, \nu) . \tag{4.23}
\end{equation*}
$$

Using (4.19) on $\phi_{u}(\mu), \phi_{u}(\nu)$, and using the moment production property recalled in Proposition 2, we have the estimate

$$
\begin{equation*}
\sup _{t \geq u} W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) \lesssim u^{2-k_{0}} W(\mu, v)^{\zeta_{0}} \tag{4.24}
\end{equation*}
$$

We next deal with the case $T \leq u$. In this case, comparing the moment production property to the definition of $T$ shows that

$$
\begin{equation*}
T^{\beta-1} \lesssim \Lambda_{3}\left(\phi_{T}(\mu)\right)+\Lambda_{3}\left(\phi_{T}(\nu)\right) \lesssim a T^{k-3} ; \quad T \leq u \tag{4.25}
\end{equation*}
$$

which rearranges to produce the bound $1 \lesssim a u^{k / 2-1}$. In particular, in this case, we have

$$
\begin{equation*}
\sup _{t \geq 0} W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) \leq 4 \lesssim a u^{k / 2-1} \tag{4.26}
\end{equation*}
$$

Combining estimates (4.23), (4.24, 4.26), we see that in all cases,

$$
\begin{equation*}
\sup _{t \geq 0} W\left(\phi_{t}(\mu), \phi_{t}(\nu)\right) \lesssim u^{2-k_{0}} W(\mu, v)^{\zeta_{0}}+a u^{k / 2-1} \tag{4.27}
\end{equation*}
$$

Now, if we choose $u=\min \left(1, W(\mu, v)^{\delta}\right)$ for sufficiently small $\delta>0$, we obtain

$$
\begin{equation*}
\sup _{t \geq 0} W\left(\phi_{t}(\mu), \phi_{t}(v)\right) \lesssim a W(\mu, v)^{\zeta} \tag{4.28}
\end{equation*}
$$

for a new exponent $\zeta=\zeta(d, k)>0$.
5. The interpolation decomposition for Kac's process. We introduce a pair of random measures associated to the Markov process $\left(\mu_{t}^{N}\right)_{t \geq 0}$. The jump measure $m^{N}$ is the un-normalised empirical measure on $(0, \infty) \times \mathcal{S}_{N}$, of all pairs $\left(t, \mu^{N}\right)$, such that the system collides at time $t$, with new measure $\mu^{N}$. Its compensator $\bar{m}^{N}$ is the random measure on $(0, \infty) \times \mathcal{S}_{N}$ given by

$$
\begin{equation*}
\bar{m}^{N}\left(d t, d \mu^{N}\right)=\mathcal{Q}_{N}\left(\mu_{t-}^{N}, d \mu^{N}\right) d t \tag{5.1}
\end{equation*}
$$

where $\mathcal{Q}_{N}(\cdot, \cdot)$ is the transition kernel of the Kac process, given by (1.7). The goal of this section is to prove the following 'interpolation decomposition' for the difference between Kac's process and the Boltzmann flow, which is the key identity required for the proofs of Theorems 1.1, 1.2. This is based on an idea of Norris [32], which was inspired by [29], Section 3.3.

Formula 5.1. Let $\mu_{t}^{N}$ be a Kac process on $N \geq 2$ particles, and suppose $f \in \mathcal{A}_{0}$ is a test function. To ease notation, we write

$$
\Delta\left(s, t, \mu^{N}\right)=\phi_{t-s}\left(\mu^{N}\right)-\phi_{t-s}\left(\mu_{s-}^{N}\right) ; \quad 0 \leq s \leq t, \mu^{N} \in \mathcal{S}_{N}
$$

$$
\begin{equation*}
\psi(u, \mu, v)=\phi_{u}(v)-\phi_{u}(\mu)-\mathcal{D} \phi_{u}(\mu)[v-\mu] ; \quad u \geq 0, \mu, v \in \bigcap_{k>2} \mathcal{S}^{k}, \tag{5.3}
\end{equation*}
$$

where $\mathcal{D} \phi_{t}$ is the derivative of the Boltzmann flow $\phi_{t}$, defined in Proposition 3; this makes sense, provided that all moments of $\mu, v$ are finite. Then we can decompose

$$
\begin{equation*}
\left\langle f, \mu_{t}^{N}-\phi_{t}\left(\mu_{0}^{N}\right)\right\rangle=M_{t}^{N, f}+\int_{0}^{t}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s \tag{5.4}
\end{equation*}
$$

where

$$
\begin{equation*}
M_{t}^{N, f}=\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \Delta\left(s, t, \mu^{N}\right)\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu_{s}^{N}\right) \tag{5.5}
\end{equation*}
$$

and where $\rho^{N}$ is given in terms of the transition kernel $\mathcal{Q}_{N}(1.7)$ by

$$
\begin{equation*}
\left\langle f, \rho^{N}\left(u, \mu^{N}\right)\right\rangle=\int_{\mathcal{S}_{N}}\left\langle f, \psi\left(u, \mu^{N}, \nu\right)\right\rangle \mathcal{Q}_{N}\left(\mu^{N}, d \nu\right) \tag{5.6}
\end{equation*}
$$

REMARK 5.1.
(i) This is the key identity needed for Theorems $1.1,1.2$; the remainder of the proofs are to establish suitable controls over each of the two terms.
(ii) This representation formula offers two major advantages over the equivalent representation formula in [33], which will be recalled in Proposition 6.

- First, all the quantities appearing in our formula are adapted to the natural filtration of $\left(\mu_{t}^{N}\right)_{t \geq 0}$, and so we can use martingale estimates directly; by contrast, [33], Proposition 4.2, contains anticipating terms. This allows us to prove convergence in $L^{p}$ spaces, rather than simply in probability.
- Second, all terms appearing in our formula may be controlled by the stability estimates (3.7), (3.9). This allows us to exploit the stability of the limit equation, at the level of individual realisations of the empirical particle system $\mu_{0}^{N}$.

The main technicality in the proof of this is to derive a Chapman-Kolmogorovstyle equation, which allows us to manipulate the functional derivatives $\mathcal{D} \phi_{t}$. This is the content of the following lemma.

Lemma 5.2 (Exchange lemma). Let $\mu^{N} \in \mathcal{S}_{N}$ and $f \in \mathcal{A}$. Then for all times $t \geq 0$, we have the equalities

$$
\begin{align*}
\frac{d}{d t} & \left\langle f, \phi_{t}\left(\mu^{N}\right)\right\rangle \\
& =\left\langle f, \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right]\right\rangle  \tag{5.7}\\
& =\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left\langle f, \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right]\right\rangle \\
& \times\left|v-v_{\star}\right| N d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right),
\end{align*}
$$

where $\mu^{N, v, v_{\star}, \sigma}$ is the post-collision measure given by (1.3), $\mathcal{Q}_{N}$ is the generator of the Kac process (1.7) and where $\mathcal{D} \phi_{t}$ is the functional derivative given by Proposition 3.

The first equality is familiar from semigroup theory, but is complicated by the nonlinearity of the flow maps; we resolve this by using ideas of the infinite dimensional differential calculus developed in [29]. The second equality can be thought of as a continuity property for the linear map $\mathcal{D} \phi_{t}\left(\mu^{N}\right)[\cdot]$, and is justified in Lemma 5.2 by the explicit construction of the derivative in Proposition 3.

Assuming this for the moment, we now prove the interpolation decomposition Formula 5.1.

Proof of Formula 5.1. Fix $t \geq 0$ and $f \in \mathcal{A}$, and consider the process $\Gamma_{s}^{N, f, t}=\left\langle f, \phi_{t-s}\left(\mu_{s}^{N}\right)\right\rangle$, for $0 \leq s \leq t$. Then $\Gamma^{N, f, t}$ is càdlàg, and is differentiable
on intervals where $\mu_{s}^{N}$ is constant. On such intervals, Lemma 5.2 tells us that

$$
\begin{align*}
\frac{d}{d s} & \left\langle f, \phi_{t-s}\left(\mu_{s}^{N}\right)\right\rangle \\
= & -\left.\frac{d}{d u}\right|_{u=t-s}\left\langle f, \phi_{u}\left(\mu_{s}^{N}\right)\right\rangle \\
= & -\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left\langle f, \mathcal{D} \phi_{t-s}\left(\mu_{s}^{N}\right)\left[\mu^{N, v, v_{\star}, \sigma}-\mu_{s}^{N}\right]\right\rangle  \tag{5.8}\\
& \times\left|v-v_{\star}\right| N \mu_{s}^{N}(d v) \mu_{s}^{N}\left(d v_{\star}\right) d \sigma \\
= & -\int_{\mathcal{S}_{N}}\left\langle f, \mathcal{D} \phi_{t-s}\left(\mu_{s}^{N}\right)\left[\mu^{N}-\mu_{s}^{N}\right]\right\rangle \mathcal{Q}_{N}\left(\mu_{s}^{N}, d \mu^{N}\right)
\end{align*}
$$

where the final equality is to rewrite integral in terms of the transition kernel $\mathcal{Q}_{N}$ of the Kac process, defined in (1.7). Writing $\mathcal{I}_{t}$ for the (finite) set of jumps $\mathcal{I}_{t}=$ $\left\{s \leq t: \mu_{s}^{N} \neq \mu_{s-}^{N}\right\}$, the contribution to $\Gamma_{t}^{N, f, t}-\Gamma_{0}^{N, f, t}$ from drift between jumps is

$$
\begin{align*}
& \int_{(0, t] \backslash \mathcal{I}_{t}} \frac{d}{d s}\left\langle f, \phi_{t-s}\left(\mu_{s}^{N}\right)\right\rangle d s  \tag{5.9}\\
& \quad=-\int_{\left((0, t] \backslash \mathcal{I}_{t}\right) \times \mathcal{S}_{N}}\left\langle f, \mathcal{D} \phi_{t-s}\left(\mu_{s}^{N}\right)\left[\mu^{N}-\mu_{s}^{N}\right]\right\rangle \mathcal{Q}_{N}\left(\mu_{s}^{N}, d \mu^{N}\right) d s
\end{align*}
$$

Using the definitions (5.2), (5.3) of $\psi$ and $\Delta$, the integrand can be expressed as

$$
\begin{equation*}
\left\langle f, \mathcal{D} \phi_{t-s}\left(\mu_{s}^{N}\right)\left[\mu^{N}-\mu_{s}^{N}\right]\right\rangle=\left\langle f, \Delta\left(s, t, \mu^{N}\right)-\psi\left(t-s, \mu_{s}^{N}, \mu^{N}\right)\right\rangle \tag{5.10}
\end{equation*}
$$

for any $s \notin \mathcal{I}_{t}$. Since the set $\mathcal{I}_{t}$ has 0 Lebesgue measure, the set $\mathcal{I}_{t} \times \mathcal{S}_{N}$ has 0 measure with respect to $\mathcal{Q}_{N}\left(\mu_{s}^{N}, d \mu^{N}\right) d s$, and so the inclusion of this set does not change the integral. Using the definitions (5.1), (5.6) of $\bar{m}^{N}$ and $\rho^{N}$, we can rewrite the integral as

$$
\begin{align*}
& \int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \psi\left(t-s, \mu_{s}^{N}, \mu^{N}\right)-\Delta\left(s, t, \mu^{N}\right)\right\rangle \mathcal{Q}_{N}\left(\mu_{s}^{N}, d \mu^{N}\right) d s \\
& \quad=\int_{0}^{t}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s-\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \Delta\left(s, t, \mu^{N}\right)\right| \bar{m}^{N}\left(d s, d \mu^{N}\right) . \tag{5.11}
\end{align*}
$$

On the other hand, at the times when $\mu_{s}^{N}$ jumps, we have

$$
\begin{equation*}
\Gamma_{s}^{N, f, t}-\Gamma_{s-}^{N, f, t}=\left\langle f, \phi_{t-s}\left(\mu_{s}^{N}\right)-\phi_{t-s}\left(\mu_{s-}^{N}\right)\right\rangle=\left\langle f, \Delta\left(s, t, \mu_{s}^{N}\right)\right\rangle \tag{5.12}
\end{equation*}
$$

Therefore, the contribution to $\Gamma_{t}^{N, f, t}-\Gamma_{0}^{N, f, t}$ from jumps is

$$
\begin{align*}
\sum_{s \in \mathcal{I}_{t}} \Gamma_{s}^{N, f, t}-\Gamma_{s-}^{N, f, t} & =\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \Delta\left(s, t, \mu^{N}\right)\right| m^{N}\left(d s, d \mu^{N}\right)  \tag{5.13}\\
& =M_{t}^{N, f}+\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \Delta\left(s, t, \mu^{N}\right)\right| \bar{m}^{N}\left(d s, d \mu^{N}\right)
\end{align*}
$$

Combining the contributions (5.11), (5.13), we see that

$$
\begin{align*}
\left\langle f, \mu_{t}^{N}-\phi_{t}\left(\mu_{0}^{N}\right)\right\rangle & =\Gamma_{t}^{N, f, t}-\Gamma_{0}^{N, f, t} \\
& =\int_{(0, t\rfloor \backslash \mathcal{I}_{t}} \frac{d}{d s}\left\langle f, \phi_{t-s}\left(\mu_{s}^{N}\right)\right\rangle d s+\sum_{s \in \mathcal{I}_{t}} \Gamma_{s}^{N, f, t}-\Gamma_{s-}^{N, f, t}  \tag{5.14}\\
& =M_{t}^{N, f}+\int_{0}^{t}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right) d s\right.
\end{align*}
$$

as desired.
5.1. Proof of Lemma 5.2. In this subsection, we will prove the ChapmanKolmogorov property Lemma 5.2, which is crucial to the interpolation decomposition. We prove the two claimed equalities separately.

Lemma 5.3. Let $N \geq 2$ and let $\mu^{N} \in \mathcal{S}_{N}$. Then, for all $t>0$ and $f \in \mathcal{A}$, we have the differentiability

$$
\begin{equation*}
\frac{d}{d t}\left\langle f, \phi_{t}\left(\mu^{N}\right)\right\rangle=\left\langle f, \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right]\right\rangle \tag{5.15}
\end{equation*}
$$

At $t=0$, this is a one-sided, right differentiability.
The following proof uses ideas of [29], notably the concept of the infinitedimensional differential calculus and building on ideas of [29], Lemma 2.11.

Proof. Throughout, fix $\mu^{N} \in \mathcal{S}_{N}$ and $f \in \mathcal{A}$. Recall, for clarity, the notation $Q_{t}(\mu)=Q\left(\phi_{t}(\mu)\right)$. Using the boundedness of appropriate moments of $\mu^{N} \in \mathcal{S}_{N}$, together with the continuity estimate (3.16), it is straightforward to see that the map $t \mapsto Q_{t}\left(\mu^{N}\right)$ is Hölder continuous in time, with respect to the weighted norm $\|\cdot\|_{\mathrm{TV}+2}$ : for some constant $C_{1}=C_{1}(N)$, we have the estimate

$$
\begin{equation*}
\left\|Q_{t}\left(\mu^{N}\right)-Q_{s}\left(\mu^{N}\right)\right\|_{\mathrm{TV}+2} \leq C_{1}|t-s|^{\frac{1}{2}} \tag{5.16}
\end{equation*}
$$

From the definition (BE) of the Boltzmann dynamics, together with dominated convergence, we have that

$$
\begin{equation*}
\left\langle f, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle=\left\langle f, \mu^{N}\right\rangle+\int_{0}^{t}\left\langle f, Q_{s}\left(\mu^{N}\right)\right\rangle d s \tag{5.17}
\end{equation*}
$$

Therefore, the map $t \mapsto\left\langle f, \phi_{t}\left(\mu^{N}\right)\right\rangle$ is continuously differentiable in time, with derivative

$$
\begin{equation*}
\frac{d}{d t}\left\langle f, \phi_{t}\left(\mu^{N}\right)\right\rangle=\left\langle f, Q_{t}\left(\mu^{N}\right)\right\rangle \tag{5.18}
\end{equation*}
$$

where, at $t=0$, this is a one-sided, right derivative. It therefore suffices to show that (5.15) holds as a right derivative.

Fix $t \geq 0$, and observe that, for $s>0$ small enough, $v_{s}^{N}=\mu^{N}+s Q\left(\mu^{N}\right)$ defines a measure $\nu_{s}^{N} \in \mathcal{S}$. From the semigroup property, it follows that $\phi_{t}\left(\phi_{s}\left(\mu^{N}\right)\right)=$ $\phi_{t+s}\left(\mu^{N}\right)$, and we can therefore expand

$$
\begin{align*}
& \left\langle f, \phi_{t+s}\left(\mu^{N}\right)-\phi_{t}\left(\mu^{N}\right)-s \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right]\right\rangle \\
& \quad=\underbrace{\left\langle f, \phi_{t}\left(\phi_{s}\left(\mu^{N}\right)\right)-\phi_{t}\left(v_{s}^{N}\right)\right\rangle}_{:=\mathcal{T}_{1}(s)}  \tag{5.19}\\
& \quad+\underbrace{\left\langle f, \phi_{t}\left(v_{s}^{N}\right)-\phi_{t}\left(\mu^{N}\right)-s \mathcal{D} \phi_{t}(\mu)\left[Q\left(\mu^{N}\right)\right]\right\rangle}_{:=\mathcal{T}_{2}(s)} .
\end{align*}
$$

We will now show that each of the two terms $\mathcal{T}_{1}, \mathcal{T}_{2}$ are $o(s)$, which implies the result.

Estimate on $\mathcal{T}_{1}(s)$. Let $\eta \in\left(\frac{2}{3}, 1\right)$, and choose $k$ large enough that the stability estimates (3.7), (3.9) hold with exponent $\eta$. As $s \downarrow 0$, the probability measures $v_{s}^{N}=\mu^{N}+s Q\left(\mu^{N}\right)$ and $\phi_{s}\left(\mu^{N}\right)$ are bounded in $\mathcal{S}^{k}$; therefore, from (3.7), there exists a constant $C_{2}=C_{2}(N)<\infty$ such that, for all $s>0$ small enough,

$$
\begin{equation*}
\left\|\phi_{t}\left(\phi_{s}(\mu)\right)-\phi_{t}\left(v_{s}\right)\right\|_{\mathrm{TV}+2} \leq C_{2}\left\|\phi_{s}(\mu)-v_{s}\right\|_{\mathrm{TV}+2}^{\eta} . \tag{5.20}
\end{equation*}
$$

The left-hand side is a bound for $\mathcal{T}_{1}(s)$. Using the estimate (5.16) above, we estimate the right-hand side, following [29], Lemma 2.11:

$$
\begin{align*}
\left\|\phi_{s}\left(\mu^{N}\right)-v_{s}^{N}\right\|_{\mathrm{TV}+2} & =\left\|\int_{0}^{s}\left(Q_{u}\left(\mu^{N}\right)-Q_{0}\left(\mu^{N}\right)\right) d u\right\|_{\mathrm{TV}+2} \\
& \leq \int_{0}^{s}\left\|Q_{u}\left(\mu^{N}\right)-Q_{0}\left(\mu^{N}\right)\right\|_{\mathrm{TV}+2} d u  \tag{5.21}\\
& \leq C_{1}(N) \int_{0}^{s} u^{\frac{1}{2}} d u=\frac{2}{3} C_{1}(N) s^{\frac{3}{2}}
\end{align*}
$$

Combining the estimates (5.20), (5.21), we see that

$$
\begin{equation*}
\mathcal{T}_{1}(s) \leq C_{2}\left(\frac{2}{3} C_{1}\right)^{\eta} s^{\frac{3 \eta}{2}} \tag{5.22}
\end{equation*}
$$

Since we chose $\eta>\frac{2}{3}$, this shows that $\mathcal{T}_{1}$ is $o(s)$ as $s \downarrow 0$.
Estimate on $\mathcal{T}_{2}(s)$. Let $\eta$ and $k$ be as above, and recall that in (3.9), $\xi_{t}$ is the definition of $\mathcal{D} \phi_{t}(\mu)[v-\mu]$. We now apply this estimate to $\mu^{N}$ and $v_{s}^{N}$, noting that $v_{s}^{N}=\mu^{N}+s Q\left(\mu^{N}\right)$ and $\phi_{s}\left(\mu^{N}\right)$ are bounded in $\mathcal{S}^{k}$ as $s \downarrow 0$, and that $v_{s}^{N}-\mu^{N}=$ $s Q\left(\mu^{N}\right)$. The bound (3.9) now shows that, for some constants $C_{3}, C_{4}<\infty$,

$$
\begin{align*}
& \| \phi_{t}\left(v_{s}^{N}\right)-\phi_{t}\left(\mu^{N}\right)-s \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right] \|_{\mathrm{TV}+2} \\
& \quad \leq C_{3}\left\|v_{s}^{N}-\mu^{N}\right\|_{\mathrm{TV}}^{1+\eta}  \tag{5.23}\\
& \quad=C_{3}\left\|s Q\left(\mu^{N}\right)\right\|_{\mathrm{TV}}^{1+\eta} \\
& \leq C_{4} s^{1+\eta} .
\end{align*}
$$

The left-hand side is a bound for $\mathcal{T}_{2}(s)$, which implies that $\mathcal{T}_{2}$ is $o(s)$, as desired. Together with the previous estimate on $\mathcal{T}_{1}$, this concludes the proof.

We now turn to the proof of the second equality in (5.7), that is,

$$
\begin{align*}
& \left\langle f, \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right]\right\rangle \\
& =\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left\langle f, \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right]\right\rangle  \tag{5.24}\\
& \quad \times N\left|v-v_{\star}\right| d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right) .
\end{align*}
$$

Using the definition (1.3), we see that the integral on the right-hand side is equivalent to that defining $Q\left(\mu^{N}\right)$ in (1.8). However, we cannot simply exchange the integration with the linear map $\mathcal{D} \phi_{t}$, as the construction in Proposition 3 does not guarantee that $\mathcal{D} \phi_{t}\left(\mu^{N}\right)$ is bounded as a linear map. We will instead prove (5.24) from the explicit way in which $\mathcal{D} \phi_{t}\left(\mu^{N}\right)$ is constructed in Proposition 3, and show that this construction implies 'enough' continuity.

This is closely related to, and may be derived from, condition (A3), convergence of the generators, in [29]. We present here a more direct proof, to avoid introducing additional spaces and notation. The crucial observation of our argument is that 'enough' small perturbations of a discrete measure $\mu^{N} \in \mathcal{S}_{N}$ will remain in $\mathcal{S}$; this is made precise in equation (5.39). The same idea is present in the corresponding argument [29], Section 5.5, but not made explicit.

Before turning to the proof of (5.24), we will prove the following auxiliary lemma. In order to justify the exchange of various integrals, we wish to improve the moments of the derivative $\xi_{t}=\mathcal{D} \phi_{t}(\mu)[v-\mu]$ in Proposition 3. The following argument combines ideas of [33], Proposition 4.2, and [29], Lemma 6.3.

Lemma 5.4. Suppose $\mu, v \in \bigcap_{k \geq 2} \mathcal{S}^{k}$, and let $\left(\xi_{t}\right)_{t \geq 0}$ be the solution to the differential equation (3.6). Then, for all $k \geq 2$, there exists a constant $c=c(k)$ such that, for all $t \geq 0$,

$$
\begin{equation*}
\left\|\xi_{t}\right\|_{\mathrm{TV}+k} \leq 2 \Lambda_{k}(\mu, v) \exp \left(c t \Lambda_{k+1}(\mu)\right) \tag{5.25}
\end{equation*}
$$

Moreover, if $k^{\prime}>2$ is large enough, then we have the continuity estimate, for all $0 \leq s \leq t$, and for some absolute constants $C_{1}, C_{2}$,

$$
\begin{equation*}
\left\|\xi_{t}-\xi_{s}\right\|_{\mathrm{TV}+k} \leq C_{1} \Lambda_{k+k^{\prime}}(\mu, v)^{\frac{1}{2}} \exp \left(\frac{1}{2} C_{2} \Lambda_{2(k+1)}(\mu) t\right)(t-s)^{\frac{1}{2}} \tag{5.26}
\end{equation*}
$$

Proof. First, we observe that, by hypothesis, the map $t \mapsto \xi_{t}$ is continuous in the norm $\|\cdot\|_{\mathrm{TV}+2}$, and is therefore locally bounded. We have the estimate on total variation

$$
\begin{equation*}
\left\|Q\left(\phi_{t}(\mu), \xi_{t}\right)\right\|_{\mathrm{TV}} \leq 4 \int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}\left|v-v_{\star}\right| \phi_{t}(\mu)(d v)\left|\xi_{t}\right|\left(d v_{\star}\right) \leq 8\left\|\xi_{t}\right\|_{\mathrm{TV}+2} \tag{5.27}
\end{equation*}
$$

where we have used the bound $\left|v-v_{\star}\right| \leq\left(1+|v|^{2}\right)\left(1+\left|v_{\star}\right|^{2}\right)$. Similarly, we estimate

$$
\begin{align*}
& \left\|Q\left(\phi_{t}(\mu), \xi_{t}\right)-Q\left(\phi_{s}(\mu), \xi_{s}\right)\right\|_{\mathrm{TV}} \\
& \quad \leq\left\|Q\left(\phi_{t}(\mu)-\phi_{s}(\mu), \xi_{t}\right)\right\|_{\mathrm{TV}}+\left\|Q\left(\phi_{s}(\mu), \xi_{t}-\xi_{s}\right)\right\|_{\mathrm{TV}}  \tag{5.28}\\
& \quad \leq 4\left(\left\|\xi_{t}\right\|_{\mathrm{TV}+2}\left\|\phi_{t}(\mu)-\phi_{s}(\mu)\right\|_{\mathrm{TV}+2}+2\left\|\xi_{t}-\xi_{s}\right\|_{\mathrm{TV}+2}\right)
\end{align*}
$$

Since $t \mapsto \phi_{t}(\mu)$ is continuous in $\|\cdot\|_{\mathrm{TV}+2}$, it follows that the map

$$
\begin{equation*}
t \mapsto \partial_{t} \xi_{t}=2 Q\left(\phi_{t}(\mu), \xi_{t}\right) \tag{5.29}
\end{equation*}
$$

is continuous and locally bounded in $\|\cdot\|_{\mathrm{TV}}$. Therefore, for all $t \geq 0$, the measure $\pi_{t}=\int_{0}^{t}\left|\partial_{s} \xi_{s}\right| d s$ is a finite measure, and $\partial_{s} \xi_{s}$ is absolutely continuous with respect to $\pi_{t}$ for all $0 \leq s \leq t$. Therefore, by a result of Norris [33], Lemma 11.1, on the time variation of signed measures, there exists a measurable map $f:[0, \infty) \times$ $\mathbb{R}^{d} \rightarrow\{-1,0,1\}$ such that

$$
\begin{equation*}
\xi_{t}=f_{t}\left|\xi_{t}\right| ; \quad\left|\xi_{t}\right|=\left|\xi_{0}\right|+\int_{0}^{t} f_{s} \partial_{s} \xi_{s} d s \tag{5.30}
\end{equation*}
$$

Writing $\check{f}_{s}(v)=\left(1+|v|^{k}\right) f_{s}$, we have the bound

$$
\begin{aligned}
\langle 1+ & \left.|v|^{k},\left|\xi_{t}\right|-\left|\xi_{0}\right|\right\rangle \\
= & \int_{0}^{t} d s \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left(\check{f}\left(v^{\prime}\right)+\check{f}\left(v_{\star}^{\prime}\right)-\check{f}\left(v_{\star}\right)-\check{f}(v)\right) \\
& \times\left|v-v_{\star}\right| \phi_{s}(\mu)(d v) \xi_{s}\left(d v_{\star}\right) d \sigma \\
\leq & \int_{0}^{t} d s \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left(2+\left|v^{\prime}\right|^{k}+\left|v_{\star}^{\prime}\right|^{k}+\left|v_{\star}\right|^{k}-|v|^{k}\right) \\
& \times\left|v-v_{\star}\right| \phi_{s}(\mu)\left(d v_{\star}\right)\left|\xi_{s}\right|(d v) d \sigma .
\end{aligned}
$$

Now, there exists a constant $C_{1}=C_{1}(k)$ such that, for all $v, v_{\star}, \sigma$, we have the bound

$$
\begin{equation*}
\left|v^{\prime}\right|^{k}+\left|v_{\star}^{\prime}\right|^{k}+\left|v_{\star}\right|^{k}-|v|^{k} \leq C_{1}(k)\left(|v|^{k-2}\left|v_{\star}\right|^{2}+\left|v_{\star}\right|^{k}\right) \tag{5.32}
\end{equation*}
$$

Therefore, for a different constant $C_{2}=C_{2}(k)$,
(5.33) $2\left|v-v_{\star}\right|\left(2+\left|v^{\prime}\right|^{k}+\left|v_{\star}^{\prime}\right|^{k}+\left|v_{\star}\right|^{k}-|v|^{k}\right) \leq C_{2}(k)\left(1+|v|^{k}\right)\left(1+\left|v_{\star}\right|^{k+1}\right)$.

Using the moment bounds in Proposition 2, we obtain for some $c=c(k)$,

$$
\begin{align*}
\langle 1+ & \left.|v|^{k},\left|\xi_{t}\right|\right\rangle \\
\leq & \left.\left.\langle 1+| v\right|^{k},\left|\xi_{0}\right|\right\rangle \\
& +C_{2} \int_{0}^{t} \int_{\mathbb{R}^{d} \times \mathbb{R}^{d}}\left(1+|v|^{k}\right)\left(1+\left|v_{\star}\right|^{k+1}\right)\left|\xi_{s}\right|(d v) \phi_{s}(\mu)\left(d v_{\star}\right)  \tag{5.34}\\
\leq & \left.\left.\left.\langle 1+| v\right|^{k},\left|\xi_{0}\right|\right\rangle+\left.c \Lambda_{k+1}(\mu) \int_{0}^{t}\langle 1+| v\right|^{k},\left|\xi_{s}\right|\right\rangle d s
\end{align*}
$$

Grönwall's lemma now gives the claimed moment bound. For the continuity statement, if $k^{\prime}$ is chosen large enough that (3.8) holds for some $\eta<1$, then (5.27) gives the bound

$$
\begin{equation*}
\left\|Q\left(\phi_{t}(\mu), \xi_{t}\right)\right\|_{\mathrm{TV}} \leq C_{3} \Lambda_{k^{\prime}}(\mu, \nu) \tag{5.35}
\end{equation*}
$$

and, therefore, for all $0 \leq s \leq t$,

$$
\begin{equation*}
\left\|\xi_{t}-\xi_{s}\right\|_{\mathrm{TV}} \leq C_{3} \Lambda_{k^{\prime}}(\mu, v)(t-s) \tag{5.36}
\end{equation*}
$$

The continuity statement follows by combining (5.36) with the moment bound for $2 k$, with the interpolation

$$
\begin{equation*}
\left\|\xi_{t}-\xi_{s}\right\|_{\mathrm{TV}+k} \leq\left\|\xi_{t}-\xi_{s}\right\|_{\mathrm{TV}}^{1 / 2}\left\|\xi_{t}+\xi_{s}\right\|_{\mathrm{TV}+2 k}^{1 / 2} \tag{5.37}
\end{equation*}
$$

and using the correlation property (Lemma 2.2) to absorb both moment terms.
We can now prove the second claimed equality in Lemma 5.2.
LEmma 5.5. Let $\mu^{N} \in \mathcal{S}_{N}$, for $N \geq 2$. Then we have the equality

$$
\begin{align*}
& \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right] \\
& =\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right]  \tag{5.38}\\
& \quad \times N\left|v-v_{\star}\right| d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right)
\end{align*}
$$

where the right-hand side is a Bochner integral in the space $\left(Y_{2},\|\cdot\|_{\mathrm{TV}+2}\right)$. In particular, the equality (5.24) holds.

Proof. We exploit the fact that, for $\delta>0$ small enough, we have

$$
\begin{equation*}
\mu^{N}+\delta Q\left(\mu^{N}\right) \in \mathcal{S}, \quad \forall v, v_{\star}, \sigma, \mu^{N}+\delta\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right] \in \mathcal{S} . \tag{5.39}
\end{equation*}
$$

We will assume that $\delta>0$ is chosen so that this holds. For $v, v_{\star} \in \operatorname{Supp}\left(\mu^{N}\right)$ and $\sigma \in S^{d-1}$, we define $\xi_{t}^{N, v, v_{\star}, \sigma}$ by the differential equation

$$
\begin{align*}
\xi_{0}^{N, v, v_{\star}, \sigma} & =\delta\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right] \\
\partial_{t} \xi_{t}^{N, v, v_{\star}, \sigma} & =2 Q\left(\phi_{t}\left(\mu^{N}\right), \xi_{t}^{N, v, v_{\star}, \sigma}\right) . \tag{5.40}
\end{align*}
$$

From Proposition 3, the solution to this equation exists, and is unique. By the characterisation of the derivative $\mathcal{D} \phi_{t}\left(\mu^{N}\right)$, we also have

$$
\begin{equation*}
\xi_{t}^{N, v, v_{\star}, \sigma}=\delta \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right] . \tag{5.41}
\end{equation*}
$$

From Lemma 5.4, we also have a bound that $\left\|\xi_{s}^{N, v, v_{\star}, \sigma}\right\|_{\mathrm{TV}+4} \leq C$ for all $s \leq t$, and for some constant $C=C\left(\mu^{N}, N, t\right)$ independent of $v, v_{\star}$ and $\sigma$. In this notation, we wish to establish the equality

$$
\begin{equation*}
\mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right] \stackrel{?}{=} \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} \xi^{N, v, v_{\star}, \sigma}\left|v-v_{\star}\right| \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right) d \sigma \tag{5.42}
\end{equation*}
$$

From the bound above, the right-hand side is well-defined as a Bochner integral in $\left(Y_{2},\|\cdot\|_{\mathrm{TV}+2}\right)$.

Firstly, arguing as in (5.27), for all $t \geq 0$, there is a constant $C=C\left(\mu^{N}, N, t\right)$ such that, for all $v, v_{\star}, \sigma$ and $s \leq t$, we have

$$
\begin{equation*}
\left\|Q\left(\phi_{t}\left(\mu^{N}\right), \xi_{s}^{N, v, v_{\star}, \sigma}\right)\right\|_{\mathrm{TV}+3} \leq C . \tag{5.43}
\end{equation*}
$$

We now define

$$
\begin{equation*}
\xi_{t}=\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} \xi_{t}^{N, v, v_{\star}, \sigma}\left|v-v_{\star}\right| d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right) \tag{5.44}
\end{equation*}
$$

where the right-hand side is a Bochner integral in $\left(Y_{3},\|\cdot\|_{\mathrm{TV}+3}\right)$. From the definition (1.8) of $Q$, we have

$$
\begin{equation*}
\xi_{0}=\delta N^{-1} Q\left(\mu^{N}\right) \tag{5.45}
\end{equation*}
$$

Moreover, using Fubini, we can express

$$
\begin{align*}
\xi_{t}- & \xi_{0} \\
= & \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}}\left\{\int_{0}^{t} 2 Q\left(\phi_{s}\left(\mu^{N}\right), \xi_{s}^{N, v, v_{\star}, \sigma}\right) d s\right\} \\
& \times\left|v-v_{\star}\right| d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right)  \tag{5.46}\\
= & \int_{0}^{t}\left\{\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} 2 Q\left(\phi_{s}\left(\mu^{N}\right), \xi_{s}^{N, v, v_{\star}, \sigma}\right)\right. \\
& \left.\times\left|v-v_{\star}\right| d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right)\right\} d s
\end{align*}
$$

The same argument as in (5.27) shows that, for fixed $\mu \in \mathcal{S}^{3}$, the map

$$
\begin{equation*}
Q(\mu, \cdot):\left(Y_{3},\|\cdot\|_{\mathrm{TV}+3}\right) \rightarrow\left(Y_{2},\|\cdot\|_{\mathrm{TV}+2}\right) ; \quad \xi \mapsto Q(\mu, \xi) \tag{5.47}
\end{equation*}
$$

is a bounded linear map. It follows that, for all $s \geq 0$,

$$
\begin{align*}
Q\left(\phi_{s}\left(\mu^{N}\right), \xi_{s}\right)= & \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} Q\left(\phi_{s}\left(\mu^{N}\right), \xi_{s}^{N, v, v_{\star}, \sigma}\right)  \tag{5.48}\\
& \times\left|v-v_{\star}\right| d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right)
\end{align*}
$$

as an equality of Bochner integrals in $\left(Y_{2},\|\cdot\|_{\mathrm{TV}+2}\right)$. Therefore, (5.46) shows that, for all $t \geq 0$,

$$
\begin{equation*}
\xi_{t}=\xi_{0}+\int_{0}^{t} 2 Q\left(\phi_{s}\left(\mu^{N}\right), \xi_{s}\right) d s \tag{5.49}
\end{equation*}
$$

From Lemma 5.4, there exists a constant $C=C\left(\mu^{N}, N, t\right)$ such that, for all $v, v_{\star}, \sigma$ and $0 \leq s \leq t$,

$$
\begin{equation*}
\left\|\xi_{t}^{N, v, v_{\star}, \sigma}-\xi_{s}^{N, v, v_{\star}, \sigma}\right\|_{\mathrm{TV}+2} \leq C(t-s)^{\frac{1}{2}} \tag{5.50}
\end{equation*}
$$

and, therefore, for a different constant $C^{\prime}$,

$$
\begin{equation*}
\left\|\xi_{t}-\xi_{s}\right\|_{\mathrm{TV}+2} \leq C^{\prime}(t-s)^{\frac{1}{2}} \tag{5.51}
\end{equation*}
$$

By the same reasoning as (5.28), we see that the map $t \mapsto 2 Q\left(\phi_{t}\left(\mu^{N}\right), \xi_{t}\right)$ is continuous with respect to the norm $\|\cdot\|_{\mathrm{TV}+2}$, and so we may differentiate (5.49) to obtain $\partial_{t} \xi_{t}=2 Q\left(\phi_{t}\left(\mu^{N}\right), \xi_{t}\right)$. From Proposition 3, this uniquely characterises the derivative $\mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[\delta N^{-1} Q\left(\mu^{N}\right)\right]$, and so we have proven the claimed equality

$$
\begin{align*}
& \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[Q\left(\mu^{N}\right)\right]=\delta^{-1} N \xi_{t} \\
& =\delta^{-1} \int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} \xi_{t}^{N, v, v_{\star}, \sigma}\left|v-v_{\star}\right| N d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right)  \tag{5.52}\\
& =\int_{\mathbb{R}^{d} \times \mathbb{R}^{d} \times S^{d-1}} \mathcal{D} \phi_{t}\left(\mu^{N}\right)\left[\mu^{N, v, v_{\star}, \sigma}-\mu^{N}\right] \\
& \quad \times\left|v-v_{\star}\right| N d \sigma \mu^{N}(d v) \mu^{N}\left(d v_{\star}\right) .
\end{align*}
$$

6. Proof of Theorem 1.1. The main difficulty in obtaining a pathwise statement is the martingale term $M_{t}^{N, f}$ in Formula 5.1, which we defined above as

$$
\begin{equation*}
M_{t}^{N, f}=\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \phi_{t-s}\left(\mu^{N}\right)-\phi_{t-s}\left(\mu_{s-}^{N}\right)\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu^{N}\right) \tag{6.1}
\end{equation*}
$$

Recall the definition of $\mathcal{A}$ as those functions $f: \mathbb{R}^{d} \rightarrow \mathbb{R}$ satisfying

$$
\begin{equation*}
\forall v, v^{\prime} \in \mathbb{R}^{d}, \quad|\hat{f}(v)| \leq 1 ; \quad\left|\hat{f}(v)-\hat{f}\left(v^{\prime}\right)\right| \leq\left|v-v^{\prime}\right| \tag{6.2}
\end{equation*}
$$

where $\widehat{f}(v)=\frac{f(v)}{1+|v|^{2}}$. We will be interested in controlling an expression of the form $\sup _{f \in \mathcal{A}}\left|M_{t}^{N, f}\right|$, either pointwise in time, or (pathwise) locally uniformly in time. However, unlike in the finite dimensional cases in [9], we cannot directly apply estimates from the elementary theory of martingales, as such estimates degrade in large dimensions. Instead, we will use the relative compactness discussed in Remark 3.1 to argue that this is an effectively finite dimensional problem. More precisely, we show that it can be approximated by a discretised, finite dimensional martingale approximation problem, with the following trade off: that making the truncation error small requires taking a large (finite) dimensional martingale. As in $[9,33]$, the martingale term is 'small', as a function of $N$, but will increase as a function of the dimension of the approximation. By optimising over the discretisation, we will be able to balance the two terms to find a useful estimate on the family of processes. This is the same approach as used for an equivalent problem in [33], Theorem 1.1.

Finding the best exponents of $N$, we have been able to obtain uses a 'hierarchical decomposition'. This approach was inspired by an equivalent technique used in [33], Proposition 7.1.

Lemma 6.1. Let $\epsilon>0, a \geq 1$ and $0<\lambda<\lambda_{0}$. Let $k$ be large enough that Corollary 3.2 holds with $q=4$, exponent $\lambda$ and Hölder exponent $1-\epsilon$.

Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process in dimension $d \geq 3$, with initial moment $\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a$. Let $M_{t}^{N, f}$ be the processes given by (5.5). Then we have, uniformly in $t \geq 0$,

$$
\begin{equation*}
\left\|\sup _{f \in \mathcal{A}}\left|M_{t}^{N, f}\right|\right\|_{L^{2}(\mathbb{P})} \lesssim a^{1 / 2} N^{\epsilon-1 / d} \tag{6.3}
\end{equation*}
$$

Once we have obtained the control of the martingale term, the remaining proof of Theorem 1.1 is straightforward.

Proof of Theorem 1.1. Take $k=k(\epsilon)$ as in Lemma 6.1, and such that Proposition 3 holds with exponent $\max \left(1-\epsilon, \frac{1}{2}\right)$.

We first note that it is sufficient to prove the case $\mu_{0}=\mu_{0}^{N}$. Given this case, we use the continuity established in Theorem 1.6 to estimate the difference

$$
\begin{equation*}
W\left(\phi_{t}\left(\mu_{0}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right) \lesssim a^{1 / 2} W\left(\mu_{0}^{N}, \mu_{0}\right)^{\zeta} \tag{6.4}
\end{equation*}
$$

for some $\zeta=\zeta(d, k)$, which implies the claimed result.
From now on, we assume that $\mu_{0}=\mu_{0}^{N}$. From the interpolation decomposition Formula 5.1, we majorise

$$
\begin{equation*}
W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right) \leq \sup _{f \in \mathcal{A}}\left|M_{t}^{N, f}\right|+\int_{0}^{t} \sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s, \tag{6.5}
\end{equation*}
$$

where, as in (5.3), (5.6), the integrand is given by

$$
\begin{align*}
\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle & =\int_{\mathcal{S}_{N}}\left\langle f, \psi\left(t-s, \mu_{s}^{N}, v\right)\right\rangle \mathcal{Q}_{N}\left(\mu^{N}, d v\right)  \tag{6.6}\\
\psi(u, \mu, v) & =\phi_{u}(v)-\phi_{u}(\mu)-\mathcal{D} \phi_{u}(\mu)[v-\mu] \tag{6.7}
\end{align*}
$$

and $\mathcal{Q}_{N}$ is the transition kernel (1.7) of the Kac process.
The first term of (6.5) is controlled in $L^{2}$ by Lemma 6.1, and so it remains to bound the second term in $L^{2}$. Let $s \geq 0$, and let $\mu^{N}$ be a measure obtained from $\mu_{s}^{N}$ by a collision, as in (1.3). Then, using the estimate (3.9), we bound

$$
\begin{align*}
\left\|\psi\left(t-s, \mu_{s}^{N}, \mu^{N}\right)\right\|_{\mathrm{TV}+2} & =\left\|\phi_{t-s}\left(\mu^{N}\right)-\phi_{t-s}\left(\mu_{s}^{N}\right)-\mathcal{D} \phi_{t-s}\left(\mu_{s}^{N}\right)\right\|_{\mathrm{TV}+2} \\
& \lesssim e^{-\lambda_{0}(t-s) / 2}\left\|\mu^{N}-\mu_{2}^{N}\right\|_{\mathrm{TV}}^{2-\epsilon} \Lambda_{k}\left(\mu^{N}, \mu_{s}^{N}\right)^{\frac{1}{2}} \tag{6.8}
\end{align*}
$$

By Lemma 2.1, we know that $\Lambda_{k}\left(\mu^{N}\right) \lesssim \Lambda_{k}\left(\mu_{s}^{N}\right)$. Moreover, from the form (1.3) of possible $\mu^{N}$, we know that

$$
\begin{equation*}
\left\|\mu^{N}-\mu_{s}^{N}\right\|_{\mathrm{TV}} \leq \frac{4}{N} \quad \text { for } \mathcal{Q}_{N}\left(\mu_{s}^{N}, \cdot\right) \text {-almost all } \mu^{N} \tag{6.9}
\end{equation*}
$$

Therefore, almost surely, for all $s$ and $\mathcal{Q}_{N}\left(\mu_{s}^{N}, \cdot\right)$-almost all $\mu^{N}$, we have the bound

$$
\begin{equation*}
\left\|\psi\left(t-s, \mu_{s}^{N}, \mu^{N}\right)\right\|_{\mathrm{TV}+2} \lesssim e^{-\lambda_{0}(t-s) / 2} N^{\epsilon-2} \Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}} \tag{6.10}
\end{equation*}
$$

where the implied constants are independent of $s, \mu_{s}^{N}$. Integrating with respect to $\mathcal{Q}_{N}\left(\mu_{s}^{N}, d \mu^{N}\right)$, we obtain an upper bound for $\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle$ :

$$
\begin{align*}
\sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle & \leq \int_{\mathcal{S}_{N}}\left\|\psi\left(t-s, \mu_{s}^{N}, \mu^{N}\right)\right\|_{\mathrm{TV}+2} \mathcal{Q}_{N}\left(\mu_{s}^{N}, d \mu^{N}\right)  \tag{6.11}\\
& \lesssim e^{-\lambda_{0}(t-s) / 2} N^{\epsilon-1} \Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}}
\end{align*}
$$

We now take the $L^{2}$ norm of the second term in (6.5). Using Proposition 2(i) to control the moments $\Lambda_{k}$ appearing in the integral, we obtain

$$
\begin{align*}
& \left\|\int_{0}^{t} \sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s\right\|_{L^{2}(\mathbb{P})} \\
& \quad \leq \int_{0}^{t}\left\|\sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle\right\|_{L^{2}(\mathbb{P})} d s  \tag{6.12}\\
& \quad \lesssim \int_{0}^{t} e^{-\lambda(t-s) / 2} N^{\epsilon-1}\left\|\Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}}\right\|_{L^{2}(\mathbb{P})} d s \\
& \quad \lesssim N^{\epsilon-1} a^{1 / 2} .
\end{align*}
$$

Noting that the exponent $\epsilon-1<\epsilon-\frac{1}{d}$, we combine this with Lemma 6.1, and keep the worse asymptotics.

Proof of Lemma 6.1. We begin by reviewing the following estimates for 1 -Lipschitz functions from [33]. Following [33], we use angle brackets $\langle f\rangle_{C}$ to denote the average of a bounded function $f$ over a Borel set $C$ of finite, nonzero measure.

Let $f$ be $1-$ Lipschitz, and consider $B=\left[0,2^{-j}\right]^{d}$. Then, for some numerical constant $c_{d}$, we have

$$
\begin{equation*}
\forall v \in B, \quad\left|f(v)-\langle f\rangle_{B}\right| \leq c_{d} 2^{-j} ; \quad\left|\langle f\rangle_{B}-\langle f\rangle_{2 B}\right| \leq c_{d} 2^{-j} \tag{6.13}
\end{equation*}
$$

We note that both of these bounds are linear in the length scale $2^{-j}$ of the box. We deal with the case $N \geq 2^{2 d}$.

The proof is based on the following 'hierarchical' partition of $\mathbb{R}^{d}$, given in the proof [33], Proposition 7.1.

- For $j \in \mathbb{Z}$, we take $B_{j}=\left(-2^{j}, 2^{j}\right]$.
- Set $A_{0}=B_{0}$ and, for $j \geq 1, A_{j}=B_{j} \backslash B_{j-1}$.
- For $j \geq 1$ and $l \geq 2$, there is a unique partition $\mathcal{P}_{j, l}$ of $A_{j}$ by $2^{l d}-2^{(l-1) d}$ translates of $B_{j-l}$.
- Similarly, write $\mathcal{P}_{0, l}$ for the unique partition of $A_{0}$ by $2^{d l}$ translates of $B_{-l}$.
- For $l \geq 3$ and $k \in \mathbb{Z}$, let $B \in \mathcal{P}_{j, l}$. We write $\pi(B)$ for the unique element of of $\mathcal{P}_{j, l-1}$ such that $B \subset \pi(B)$.

We deal first with the case $d \geq 3$. Fix discretisation parameters $L, J \geq 1$. Given a test function $f \in \mathcal{A}$, we can decompose

$$
\begin{equation*}
f=\sum_{j=0}^{J} \sum_{l=2}^{L} \sum_{B \in \mathcal{P}_{j, l}} a_{B}(f)\left(1+|v|^{2}\right) 1_{B}+\beta(f) \tag{6.14}
\end{equation*}
$$

where we define

$$
a_{B}(f)= \begin{cases}\langle\hat{f}\rangle_{B} & \text { if } B \in \mathcal{P}_{j, 2}, \text { for some } j \geq 0  \tag{6.15}\\ \langle\hat{f}\rangle_{B}-\langle\hat{f}\rangle_{\pi(B)} & \text { if } B \in \mathcal{P}_{j, l}, \text { for some } j \geq 0, l \geq 3\end{cases}
$$

and the equation serves to define the remainder term $\beta(f)$. Write $h_{B}=2^{2 j}(1+$ $\left.|v|^{2}\right) 1_{B}$, for $B \in \mathcal{P}_{j, l}$, and write $M_{t}^{N ; B}=M_{t}^{N, h_{B}}$. We can now write

$$
\begin{align*}
M_{t}^{N, f} & =\sum_{j=0}^{J} \sum_{l=2}^{L} \sum_{B \in \mathcal{P}_{j, l}} 2^{-2 j} a_{B}(f) M_{t}^{N ; B}+R_{t}^{N, f}  \tag{6.16}\\
R_{t}^{N, f} & =\int_{(0, t] \times \mathcal{S}_{N}}\left\langle\beta(f), \Delta\left(s, t, \mu^{N}\right)\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu^{N}\right) \tag{6.17}
\end{align*}
$$

and where $\Delta, m^{N}$ and $\bar{m}^{N}$ are defined in Section 5. This is the key decomposition in the proof. Roughly speaking:

- The martingales $M^{N ; B}$ are controlled by a bound (A.2) from the general theory of Markov chains, independently of $f$.
- The coefficients $a_{B}$ depend on $f$, but are bounded, uniformly over $f \in \mathcal{A}$.
- On $B_{J}, \beta(f)$ will be small, uniformly in $f$, due to the Lipschitz bound on $f$ and the estimate (6.13). This may be viewed as a relative compactness argument, as discussed in Remark 3.1: given $\epsilon>0$, one could use this construction to produce a finite $\epsilon$-net for $\left.\mathcal{A}\right|_{B_{J}}$ in the norm of $\left.\mathcal{A}_{0}\right|_{B_{J}}$.
- $|\beta(f)| \leq 1$ is bounded on $\mathbb{R}^{d} \backslash B_{J}$, and the contribution from this region will be controlled by the moment bounds.

To control the martingale term uniformly in $f$, observe that for $B \in \mathcal{P}_{j, l}$, the bound (6.13) gives $2^{-2 j}\left|a_{B}(f)\right| \lesssim 2^{-j-l}$, and $\# \mathcal{P}_{j, l} \leq 2^{d l}$. Hence, independently of $f \in$ $\mathcal{A}$,

$$
\begin{equation*}
\left(\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left(a_{B}(f) 2^{-2 j}\right)^{2}\right) \lesssim 2^{(d-2) l} \tag{6.18}
\end{equation*}
$$

Now, by Cauchy-Schwarz,

$$
\begin{align*}
& \sup _{f \in \mathcal{A}}\left|\sum_{j=0}^{J} \sum_{l=2}^{L} \sum_{B \in \mathcal{P}_{j, l}} 2^{-2 j} a_{B}(f) M_{t}^{N ; B}\right| \\
& \quad \lesssim \sum_{l=2}^{L}\left(\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{M_{t}^{N ; B}\right\}^{2}\right)^{1 / 2} 2^{(d / 2-1) l} . \tag{6.19}
\end{align*}
$$

Let $\left(M_{s}^{N ; B ; t}\right)_{s \leq t}$ be the martingale

$$
\begin{equation*}
M_{s}^{N ; B ; t}=\int_{(0, s] \times \mathcal{S}_{N}}\left\langle h_{B}, \Delta\left(u, t, \mu^{N}\right)\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d u, d \mu^{N}\right) . \tag{6.20}
\end{equation*}
$$

We can control the remaining martingale term pointwise in $L^{2}$ by applying the martingale bound (A.2) at the terminal time $t$ :

$$
\left\|M_{t}^{N ; B}\right\|_{L^{2}(\mathbb{P})}^{2}=\mathbb{E} \int_{(0, t] \times \mathcal{S}_{N}}\left\langle\left(1+|v|^{2}\right) 2^{2 j} 1_{B},\left.\Delta\left(s, t, \mu^{N}\right)\right|^{2} \bar{m}^{N}\left(d s, d \mu^{N}\right)\right.
$$

$$
\begin{equation*}
\lesssim \mathbb{E}\left[\int_{(0, t] \times \mathcal{S}_{N}}\left\langle\left(1+|v|^{4}\right) 1_{B},\right| \Delta\left(s, t, \mu^{N}\right)| \rangle^{2} \bar{m}^{N}\left(d s, d \mu^{N}\right)\right] . \tag{6.21}
\end{equation*}
$$

Summing over $B \in \mathcal{P}_{j, l}$ and $j=0, \ldots, J$, we Minkowski's inequality to move the sum inside the integral against $\Delta$, and note that $\sum_{j} \sum_{B \in \mathcal{P}_{j, l}} h_{B} \lesssim\left(1+|v|^{4}\right)$. This produces the bound

$$
\begin{align*}
& \sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\|M_{t}^{N ; B}\right\|_{L^{2}(\mathbb{P})}^{2} \\
& \quad \lesssim \mathbb{E}\left[\left.\int_{(0, t] \times \mathcal{S}_{N}}\left\langle\left(1+|v|^{4}\right),\right| \Delta\left(s, t, \mu^{N}\right)\right|^{2} \bar{m}^{N}\left(d s, d \mu^{N}\right)\right]  \tag{6.22}\\
& \quad=\mathbb{E}\left[\int_{(0, t] \times \mathcal{S}_{N}}\left\|\phi_{t-s}\left(\mu^{N}\right)-\phi_{t-s}\left(\mu_{s-}^{N}\right)\right\|_{\mathrm{TV}+4}^{2} \bar{m}^{N}\left(d s, d \mu^{N}\right)\right],
\end{align*}
$$

where the second line follows by the definition of $\Delta$ in (5.2). Using the stability estimates in Corollary 3.2 with $q=4$, we find

$$
\begin{align*}
& \sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\|M_{t}^{N ; B}\right\|_{L^{2}(\mathbb{P})}^{2}  \tag{6.23}\\
& \quad \lesssim \mathbb{E}\left[\int_{(0, t] \times \mathcal{S}_{N}} e^{-\lambda(t-s)} \Lambda_{k}\left(\mu_{s}^{N}, \mu^{N}\right) N^{2(\epsilon-1)} \bar{m}^{N}\left(d s, d \mu^{N}\right)\right]
\end{align*}
$$

For $\bar{m}^{N}$-almost all $\left(s, \mu^{N}\right)$, we bound $\Lambda_{k}\left(\mu_{s}^{N}, \mu^{N}\right) \lesssim \Lambda_{k}\left(\mu_{s}^{N}\right)$ by Lemma 2.1, and $\bar{m}^{N}\left(d s, \mathcal{S}_{N}\right) \leq 2 N d s$, to bound the right hand side by

$$
\begin{align*}
\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\|M_{t}^{N ; B}\right\|_{L^{2}(\mathbb{P})}^{2} & \lesssim \int_{0}^{t} e^{-\lambda(t-s)} N^{2 \epsilon-1} \mathbb{E}\left[\Lambda_{k}\left(\mu_{s}^{N}\right)\right] d s  \tag{6.24}\\
& \lesssim N^{2 \epsilon-1} a^{\frac{1}{2}}
\end{align*}
$$

where the second line follows using the moment estimates for the Kac process, established in Proposition 2. Therefore, (6.19) gives

$$
\begin{align*}
\left\|\sup _{f \in \mathcal{A}}\left|\sum_{j=0}^{J} \sum_{l=2}^{L} \sum_{B \in \mathcal{P}_{j, l}} a_{B}(f) M_{t}^{N ; l}\right|\right\|_{L^{2}(\mathbb{P})} & \lesssim N^{\epsilon-1 / 2} a^{1 / 2} \sum_{l=2}^{L} 2^{(d / 2-1) l}  \tag{6.25}\\
& \lesssim N^{\epsilon-1 / 2} 2^{(d / 2-1) L} a^{1 / 2}
\end{align*}
$$

The remaining points are a control on $\beta(f)$, uniformly in $f \in \mathcal{A}$, dealing with $B_{J}$ and $\mathbb{R}^{d} \backslash B_{J}$ separately. Fix $f \in \mathcal{A}$ and let $B \in \mathcal{P}_{j, L}$ with $j \leq J$. The definition gives $\hat{\beta}(f)=\hat{f}-\langle\hat{f}\rangle_{B}$ on $B$, and so

$$
\begin{equation*}
\text { On } B, \quad|\beta(f)|=\left(1+|v|^{2}\right)\left|\hat{f}-\langle\hat{f}\rangle_{B}\right| \lesssim\left(1+|v|^{2}\right) 2^{j-L} \tag{6.26}
\end{equation*}
$$

Since $|v| \geq 2^{j-1}$ on $B$, and $B \in \mathcal{P}_{j, L}$ is arbitrary, we see that

$$
\begin{equation*}
\text { On } B_{J}, \quad|\beta(f)| \lesssim 2^{-L}\left(1+|v|^{4}\right) \tag{6.27}
\end{equation*}
$$

On the other hand, the uniform bound $\|\hat{f}\|_{\infty} \leq 1$ implies that

$$
\begin{equation*}
\text { On } B_{J}^{c}, \quad|\beta(f)| \leq\left(1+|v|^{2}\right) \leq 2^{-2 J}\left(1+|v|^{4}\right) \tag{6.28}
\end{equation*}
$$

Combining, we have the global bound for all $f \in \mathcal{A}$ :

$$
\begin{equation*}
\forall v \in \mathbb{R}^{d}, \quad|\beta(f)| \lesssim\left(2^{-2 J}+2^{-L}\right)\left(1+|v|^{4}\right) \tag{6.29}
\end{equation*}
$$

Recalling the definition (5.2) of $\Delta$, we use the stability estimate in Corollary 3.2, with $q=4$, and the moment increase bound Lemma 2.1, as above to see that almost surely, for $m^{N}+\bar{m}^{N}$-almost all $\left(s, \mu^{N}\right)$, we have the bound

$$
\begin{align*}
\left.\sup _{f \in \mathcal{A}}\left|\langle\beta(f),| \Delta\left(s, t, \mu^{N}\right)\right|\right\rangle \mid & \lesssim\left(2^{-2 J}+2^{-L}\right)\left\|\Delta\left(s, t, \mu^{N}\right)\right\|_{\mathrm{TV}+4} \\
& \lesssim\left(2^{-2 J}+2^{-L}\right) e^{-\lambda(t-s) / 2} N^{\epsilon-1} \Lambda_{k}\left(\mu_{s-}^{N}\right)^{\frac{1}{2}}  \tag{6.30}\\
& =: H_{s},
\end{align*}
$$

where we introduced the shorthand $H_{s}$ for the final expression, for simplicity. We now use the trivial observation that

$$
\begin{align*}
\sup _{f \in \mathcal{A}}\left|R_{t}^{N, f}\right| \leq & \int_{(0, t] \times \mathcal{S}_{N}}\left\{\sup _{f \in \mathcal{A}}\langle | \beta(f)\left|,\left|\Delta\left(s, t, \mu^{N}\right)\right|\right\rangle\right\}  \tag{6.31}\\
& \times\left(m^{N}+\bar{m}^{N}\right)\left(d s, d \mu^{N}\right) .
\end{align*}
$$

We split the measure $m^{N}+\bar{m}^{N}=\left(m^{N}-\bar{m}^{N}\right)+2 \bar{m}^{N}$ to obtain a uniform bound for the error terms $R_{t}^{N, f}$ defined in (6.17):

$$
\begin{align*}
\left\|\sup _{f \in \mathcal{A}} R_{t}^{N, f}\right\|_{L^{2}(\mathbb{P})} & \lesssim\left\|\int_{0}^{t} H_{s}\left(m^{N}+\bar{m}^{N}\right)\left(d s, \mathcal{S}_{N}\right)\right\|_{L^{2}(\mathbb{P})}  \tag{6.32}\\
& \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon-1}\left[\mathcal{T}_{1}+\mathcal{T}_{2}\right]
\end{align*}
$$

where we have written

$$
\begin{align*}
& \mathcal{T}_{1}=\left\|\int_{0}^{t} e^{-\lambda(t-s) / 2} \Lambda_{k}\left(\mu_{s-}^{N}\right)^{\frac{1}{2}} \bar{m}^{N}\left(d s, \mathcal{S}_{N}\right)\right\|_{L^{2}(\mathbb{P})},  \tag{6.33}\\
& \mathcal{T}_{2}=\left\|\int_{0}^{t} e^{-\lambda(t-s) / 2} \Lambda_{k}\left(\mu_{s-}^{N}\right)^{\frac{1}{2}}\left(m^{N}-\bar{m}^{N}\right)\left(d s, \mathcal{S}_{N}\right)\right\|_{L^{2}(\mathbb{P})} . \tag{6.34}
\end{align*}
$$

$\mathcal{T}_{1}$ is controlled by dominating $\bar{m}^{N}\left(d s, \mathcal{S}_{N}\right) \leq 2 N d s$ to obtain

$$
\begin{align*}
\mathcal{T}_{1} & \lesssim N\left\|\int_{0}^{t} e^{-\lambda(t-s) / 2} \Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}} d s\right\|_{L^{2}(\mathbb{P})} \\
& \lesssim N \int_{0}^{t} e^{-\lambda(t-s) / 2}\left\|\Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}}\right\|_{L^{2}(\mathbb{P})} d s  \tag{6.35}\\
& \lesssim N a^{1 / 2}
\end{align*}
$$

We control $\mathcal{T}_{2}$ by Itô's isometry for $m^{N}-\bar{m}^{N}$, which is reviewed in (A.3):

$$
\begin{align*}
\mathcal{T}_{2}^{2} & =\mathbb{E}\left\{\int_{0}^{t} e^{-\lambda(t-s)} \Lambda_{k}\left(\mu_{s-}^{N}\right) \bar{m}^{N}\left(d s, \mathcal{S}_{N}\right)\right\} \\
& \lesssim N \int_{0}^{t} e^{-\lambda(t-s)} \mathbb{E}\left\{\Lambda_{k}\left(\mu_{s-}^{N}\right)\right\} d s  \tag{6.36}\\
& \lesssim N a .
\end{align*}
$$

Combining (6.32), (6.35), (6.36), we obtain

$$
\begin{equation*}
\left\|\sup _{f \in \mathcal{A}} R_{t}^{N, f}\right\|_{L^{2}(\mathbb{P})} \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon-1} a^{1 / 2} \tag{6.37}
\end{equation*}
$$

Finally, we combine (6.16), (6.25), (6.37) to obtain

$$
\begin{equation*}
\left\|\sup _{f \in \mathcal{A}}\left|M_{t}^{N, f}\right|\right\|_{L^{2}(\mathbb{P})} \lesssim N^{\epsilon} a^{1 / 2}\left(N^{-1 / 2} 2^{(d / 2-1) L}+2^{-L}+2^{-2 J}\right) . \tag{6.38}
\end{equation*}
$$

Taking $L=\left\lfloor\log _{2}(N) / d\right\rfloor$ and $J \uparrow \infty$ produces the claimed result. For $d=2$, we replace $2^{(d / 2-1) L}$ by $L$ in (6.25), and optimise as before, absorbing the factors of $(\log N)$ to make the exponent of $N$ slightly larger.
7. Proof of Theorem 1.2. We now adapt the ideas of Theorem 6.1 to a local uniform setting, and working in $L^{p}$, to prove the local uniform approximation result Theorem 1.2. As in the proof above, most of the work is in controlling the martingale term $\left(M_{t}^{N, f}\right)_{f \in \mathcal{A}}$ defined in (5.5), uniformly in $f$; for a pathwise local uniform estimate, we wish to control an expression of the form

$$
\begin{equation*}
\left\|\sup _{f \in \mathcal{A}} \sup _{t \leq t_{\mathrm{fin}}}\left|M_{t}^{N, f}\right|\right\|_{L^{p}(\mathbb{P})} \tag{7.1}
\end{equation*}
$$

Since we will frequently encounter suprema of processes on compact time intervals, we introduce notation. For any stochastic process $M$, we write

$$
\begin{equation*}
M_{\star, t}=\sup _{s \leq t}\left|M_{t}\right| \tag{7.2}
\end{equation*}
$$

Proving the sharpest asymptotics in the time horizon $t_{\mathrm{fin}}$ requires working in $L^{p}$ instead of $L^{2}$, for large exponents $p$. This leads to a weaker exponent in $N$ : we obtain only $N^{\epsilon-p^{\prime} / 2 d}$ instead of $N^{\epsilon-1 / d}$, where $p^{\prime} \leq 2$ is the Hölder conjugate to $p$. However, by making $p$ large, we are able to obtain estimates which degrade slowly in the time horizon $t_{\mathrm{fin}}$, with only a factor of $\left(1+t_{\mathrm{fin}}\right)^{1 / p}$. The exponent for $t_{\mathrm{fin}}$ can thus be made arbitrarily small, while the resulting exponent for $N$ is bounded away from 0 as we make $p$ large.

The key result required for the local uniform estimate is the following control of the expression (7.1), in analogy to Lemma 6.1.

Lemma 7.1. Let $\epsilon>0, a \geq 1$ and $p \geq 2$, and let $1<p^{\prime} \leq 2$ be the Hölder conjugate to $p$. Let $k$ be large enough that Corollary 3.2 holds for $q=5$, with Hölder exponent $1-\epsilon$, and with some $0<\lambda<\lambda_{0}$.

Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process on $N \geq 2$ particles, with initial moment $\Lambda_{k p}\left(\mu_{0}^{N}\right) \leq a^{p}$. Let $M_{t}^{N, f}$ be the processes given by (5.5), and $M_{\star, t}^{N, f}$ their local suprema, as in (7.2). Then, for any time horizon $t_{\mathrm{fin}} \in[0, \infty)$, we have the control

$$
\begin{equation*}
\left\|\sup _{f \in \mathcal{A}} M_{\star, \text { fin }}^{N, f}\right\|_{L^{p}(\mathbb{P})} \lesssim a^{1 / 2} N^{-\alpha}(\log N)^{1 / p^{\prime}}\left(1+t_{\mathrm{fin}}\right)^{\frac{3 p+1}{2 p}}, \tag{7.3}
\end{equation*}
$$

where $\alpha=\frac{p^{\prime}}{2 d}-\epsilon$.
The proof of this lemma follows the same ideas as the proof of the equivalent result, Lemma 6.1, for the pointwise bound. However, in this case, we must modify the argument to work in $L^{p}$ rather than $L^{2}$, and also to control all terms uniformly on the compact time interval $\left[0, t_{\mathrm{fin}}\right]$. This will be deferred until the end of this section.

Following the argument of the pointwise bound in Theorem 1.1, we can now produce an initial pathwise, local uniform estimate for the case $\mu_{0}=\mu_{0}^{N}$, with worse long-time behaviour. From this, we will 'bootstrap' to the desired long-time behaviour in Theorem 1.2.

Lemma 7.2. Let $\epsilon>0, a \geq 1$ and $p \geq 2$, with Hölder conjugate $p^{\prime} \leq 2$. Choose $k$ large enough that Proposition 3 holds with exponent $1-\epsilon$, and that Corollary 3.2 holds with exponent $1-\epsilon$ and $q=5$. Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a Kac process on $N \geq 2$ particles, with initial moment $\Lambda_{k p}\left(\mu_{0}^{N}\right) \leq a^{p}$. Then, for any time horizon $t_{\text {fin }} \geq 0$, we have the control

$$
\begin{equation*}
\left\|\sup _{t \leq t_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})} \lesssim a^{1 / 2} N^{\epsilon-\frac{p^{\prime}}{2 d}}(\log N)^{1 / p^{\prime}}\left(1+t_{\mathrm{fin}}\right)^{\frac{3 p+1}{2 p}} \tag{7.4}
\end{equation*}
$$

Proof of Lemma 7.2. As in Theorem 1.1, it remains to control the supremum of the integral term in Formula 5.1

$$
\begin{equation*}
\sup _{t \leq t_{\text {fin }}} \int_{0}^{t} \sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s \tag{7.5}
\end{equation*}
$$

where $\rho^{N}$ is given by (5.6). Following the previous calculation (6.11), we majorise, for $s \leq t \leq t_{\text {fin }}$,

$$
\begin{equation*}
\sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle \lesssim N^{\epsilon-1} \sup _{u \leq t_{\text {fin }}}\left\{\Lambda_{k}\left(\mu_{u}^{N}\right)^{\frac{1}{2}}\right\} \tag{7.6}
\end{equation*}
$$

from which it follows that

$$
\begin{equation*}
\sup _{t \leq t_{\text {fin }}} \int_{0}^{t} \sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s \lesssim N^{\epsilon-1} t_{\text {fin }} \sup _{u \leq t_{\text {fin }}}\left\{\Lambda_{k}\left(\mu_{u}^{N}\right)^{\frac{1}{2}}\right\} . \tag{7.7}
\end{equation*}
$$

From the local uniform moment bound established in Proposition 2(i), and the initial moment bound on $\mu_{0}^{N}$,

$$
\begin{align*}
\left\|\sup _{u \leq t_{\text {fin }}}\left\{\Lambda_{k}\left(\mu_{u}^{N}\right)^{\frac{1}{2}}\right\}\right\|_{L^{p}(\mathbb{P})} & \leq\left\|\sup _{u \leq t_{\text {fin }}}\left\{\Lambda_{k}\left(\mu_{u}^{N}\right)^{\frac{1}{2}}\right\}\right\|_{L^{2 p}(\mathbb{P})} \\
& \leq \mathbb{E}\left[\sup _{u \leq t_{\text {fin }}} \Lambda_{p k}\left(\mu_{u}^{N}\right)^{\frac{1}{2}}\right]^{1 / 2 p}  \tag{7.8}\\
& \lesssim a^{1 / 2}\left(1+t_{\text {fin }}\right)^{1 / 2 p} .
\end{align*}
$$

Combining the estimates (7.7), (7.8), we see that

$$
\begin{equation*}
\left\|\sup _{t \leq t_{\mathrm{fin}}} \int_{0}^{t} \sup _{f \in \mathcal{A}}\left\langle f, \rho^{N}\left(t-s, \mu_{s}^{N}\right)\right\rangle d s\right\|_{L^{p}(\mathbb{P})} \lesssim N^{\epsilon-1} a^{1 / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{2 p+1}{2 p}} . \tag{7.9}
\end{equation*}
$$

We combine this with Lemma 7.1 and keep the worse asymptotics.
We will now show how to 'bootstrap' to better dependence of the time horizon $t_{\mathrm{fin}}$. Heuristically, the proof allows us to replace powers of $t_{\mathrm{fin}}$ in the initial bound with the same power of $\log N$, and introduce an additional factor of $\left(1+t_{\text {fin }}\right)^{1 / p}$. As
was remarked below Proposition 3, we could derive Theorem 1.1 and Lemma 7.2 under the milder assumptions

$$
\begin{align*}
&\left\|\phi_{t}(v)-\phi_{t}(\mu)\right\|_{\mathrm{TV}+5} \leq F(t) \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{\eta}  \tag{7.10}\\
&\left\|\phi_{t}(v)-\phi_{t}(\mu)-\xi_{t}\right\|_{\mathrm{TV}+2} \leq G(t) \Lambda_{k}(\mu, v)^{\frac{1}{2}}\|\mu-v\|_{\mathrm{TV}}^{1+\eta} \tag{7.11}
\end{align*}
$$

for functions $F, G$ such that

$$
\begin{equation*}
\left(\int_{0}^{\infty} F^{2} d t\right)^{1 / 2}<\infty ; \quad \int_{0}^{\infty} G d t<\infty \tag{7.12}
\end{equation*}
$$

If we also assume that $F \rightarrow 0$ as $t \rightarrow \infty$, we can use an identical bootstrap argument, with $\log N$ replaced by a power of

$$
\begin{equation*}
\tau_{N}:=\sup \left\{t: F(t)>N^{-\alpha}\right\} \tag{7.13}
\end{equation*}
$$

which produces a potentially larger loss. Hence, the the full strength of exponential decay in Proposition 3 is used to control the asymptotic loss due to the bootstrap.

Proof of Theorem 1.2. As in the proof of Theorem 1.1, it is sufficient to prove the case $\mu_{0}^{N}=\mu_{0}$. Then, making $k$ larger if necessary, we may use Theorem 1.6 to control $\sup _{t \geq 0} W\left(\phi_{t}\left(\mu_{0}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right)$, which proves the general result.

Let $0<\epsilon^{\prime}<\epsilon$, and choose $k$ such that Lemma 7.2 holds for $\epsilon^{\prime}$. Let $\alpha^{\prime}<\alpha$ be the exponent of $N$ obtained with $\epsilon^{\prime}$ in place of $\epsilon$. From the stability estimate Proposition 3, we have

$$
\begin{equation*}
\forall \mu, v \in \mathcal{S}_{a}^{k}, \quad\left\|\phi_{t}(\mu)-\phi_{t}(v)\right\|_{\mathrm{TV}+2} \lesssim \Lambda_{k}(\mu, v)^{\frac{1}{2}} e^{-\lambda_{0} t / 2} \tag{7.14}
\end{equation*}
$$

Define $\tau=\tau_{N}=-2 \lambda_{0}^{-1} \log \left(N^{-\alpha^{\prime}}\right)$ and consider $t_{\text {fin }}>\tau+1$. Fix a positive integer $n$, and partition the interval $\left[0, t_{\mathrm{fin}}\right]$ as $I_{1} \cup I_{1} \cup \cdots \cup I_{n}$ :

$$
\begin{equation*}
I_{0}=[0, \tau] ; \quad I_{r}=\left[\tau+(r-1) \frac{t_{\mathrm{fin}}-\tau}{n}, \tau+r \frac{t_{\mathrm{fin}}-\tau}{n}\right]=:\left[s_{r}+\tau, t_{r}\right] . \tag{7.15}
\end{equation*}
$$

Write also $H_{r}=\left[s_{r}, t_{r}\right] \supset I_{r}$. Since the norm $\|\cdot\|_{\mathrm{TV}+2}$ dominates the Wasserstein distance $W$, we have the bound

$$
\begin{align*}
& \sup _{t \in I_{r}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right) \\
& \quad \lesssim \sup _{t \in H_{r}} W\left(\mu_{t}^{N}, \phi_{t-s_{r}}\left(\mu_{s_{r}}^{N}\right)\right)+e^{-\lambda \tau} \Lambda_{k}\left(\mu_{s_{r}}^{N}, \phi_{s_{r}}\left(\mu_{0}^{N}\right)\right)^{\frac{1}{2}} . \tag{7.16}
\end{align*}
$$

We bound the two terms in (7.16) separately. Denote $\left(\mathcal{F}_{t}^{N}\right)_{t \geq 0}$ the natural filtration of $\left(\mu_{t}^{N}\right)_{t \geq 0}$. We control the first term by Lemma 7.2, applied to the restarted
process $\left(\mu_{t}^{N}\right)_{t \geq s_{r}}$ :

$$
\begin{align*}
& \left\|\sup _{t \in H_{r}} W\left(\mu_{t}^{N}, \phi_{t-s_{r}}\left(\mu_{s_{r}}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})}^{p} \\
& \quad=\mathbb{E}\left\{\mathbb{E}\left(\left[\sup _{s_{r} \leq t \leq t_{r}} W\left(\mu_{t}^{N}, \phi_{t-s_{r}}\left(\mu_{s_{r}}^{N}\right)\right)\right]^{p} \mid \mathcal{F}_{s_{r}}^{N}\right)\right\}  \tag{7.17}\\
& \quad \lesssim \mathbb{E}\left\{\Lambda_{p k}\left(\mu_{s_{r}}^{N}\right)^{1 / p}\right\}\left(1+\tau+\frac{t-\tau}{n}\right)^{\frac{3 p+1}{2}} N^{-p \alpha^{\prime}}(\log N)^{\frac{p}{p^{\prime}}} .
\end{align*}
$$

We control the moment in the usual way, using Proposition 2(i), to obtain

$$
\begin{align*}
& \left\|\sup _{t \in H_{r}} W\left(\mu_{t}^{N}, \phi_{t-s_{r}}\left(\mu_{s_{r}}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})}^{p} \\
& \quad \lesssim a^{p}\left(1+\tau+\frac{t-\tau}{n}\right)^{\frac{3 p+1}{2}} N^{-p \alpha^{\prime}}(\log N)^{\frac{p}{p^{\prime}}} \tag{7.18}
\end{align*}
$$

We now turn to the second term in (7.16). Using the definition of $\tau$ and the moment estimates (2.1), (2.3) in Proposition 2,

$$
\begin{equation*}
\left\|e^{-\lambda \tau / 2} \Lambda_{k}\left(\mu_{s_{r}}^{N}, \phi_{s_{r}}\left(\mu_{0}^{N}\right)\right)^{\frac{1}{2}}\right\|_{L^{p}(\mathbb{P})} \lesssim N^{-\alpha^{\prime}} a^{1 / 2} \tag{7.19}
\end{equation*}
$$

Combining the estimates (7.18), (7.19), and absorbing powers of $\tau$ into the powers of $(\log N)$, we obtain

$$
\begin{align*}
& \left\|\sup _{t \in I_{r}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})} \\
& \quad \lesssim a^{1 / 2}\left(1+\frac{t_{\mathrm{fin}}-\tau}{n}\right)^{\frac{3 p+1}{2 p}}\left(N^{-\alpha^{\prime}}(\log N)^{\frac{3 p+1}{2 p}+\frac{1}{p^{\prime}}}\right) . \tag{7.20}
\end{align*}
$$

Observe that

$$
\begin{equation*}
\left\{\sup _{\tau \leq t \leq t_{\text {fin }}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\}^{p} \leq \sum_{r=1}^{n}\left\{\sup _{t \in I_{r}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\}^{p} . \tag{7.21}
\end{equation*}
$$

Taking expectations and $p$ th root, we find that

$$
\begin{aligned}
& \left\|\sup _{\tau \leq t \leq t_{\text {fin }}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})} \\
& \quad \lesssim n^{\frac{1}{p}} a^{1 / 2}\left(1+\frac{t_{\mathrm{fin}}-\tau}{n}\right)^{\frac{3 p+1}{2 p}}\left(N^{-\alpha^{\prime}}(\log N)^{\frac{3 p+1}{2 p}+\frac{1}{p^{\prime}}}\right) .
\end{aligned}
$$

This is optimised at $n \sim\left(t_{\mathrm{fin}}-\tau\right)$, where we obtain the estimate

$$
\begin{align*}
\left\|\sup _{\tau \leq t \leq t_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})} & \lesssim a^{1 / 2}\left(t_{\mathrm{fin}}-\tau\right)^{\frac{1}{p}}\left(N^{-\alpha^{\prime}}(\log N)^{\frac{3 p+1}{2 p}+\frac{1}{p^{\prime}}}\right)  \tag{7.23}\\
& \leq a^{1 / 2} t_{\mathrm{fin}}^{\frac{1}{p}}\left(N^{-\alpha^{\prime}}(\log N)^{\frac{3 p+1}{2 p}+\frac{1}{p^{\prime}}}\right) .
\end{align*}
$$

From Lemma 7.2 applied up to time $\tau=\tau_{N}$, we have

$$
\begin{aligned}
& \left\|\sup _{0 \leq t \leq \tau_{N}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\|_{L^{2}(\mathbb{P})} \\
& \quad \lesssim a^{1 / 2} N^{-\alpha^{\prime}}\left(1+\frac{2 \alpha}{\lambda} \log (N)\right)^{\frac{3 p+1}{2 p}}(\log N)^{\frac{1}{p^{\prime}}} \\
& \quad \lesssim a^{1 / 2}\left(N^{-\alpha}(\log N)^{\frac{3 p+1}{2 p}+\frac{1}{p^{\prime}}}\right) .
\end{aligned}
$$

Combining (7.23), (7.24), and absorbing the powers of $(\log N)$ into $N^{\epsilon-\epsilon^{\prime}}$, we have

$$
\begin{equation*}
\left\|\sup _{0 \leq t \leq t_{\mathrm{fin}}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right)\right\|_{L^{p}(\mathbb{P})} \lesssim a^{1 / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{1}{p}} N^{-\alpha} \tag{7.25}
\end{equation*}
$$

The case where $t_{\text {fin }} \leq \tau+1$ is essentially identical to (7.24).
REMARK 7.3. We note that this 'bootstrap' argument would produce the same result with any polynomial time dependence in Lemma 7.2. As a result, the precise time dependence of Lemmas 7.1, 7.2 is uninteresting, and we do not attempt to optimise it. We also remark that this method produces the same long-time behaviour even starting from an exponential estimate, at the cost of a fractional power of $N$.

It remains to prove Lemma 7.1. We draw attention to the fact that $M^{f, N}$ are not themselves martingales, despite the general construction (A.1), since the integrand $\phi_{t-s}\left(\mu^{N}\right)-\phi_{t-s}\left(\mu_{s-}^{N}\right)$ depends on the terminal time $t$. We address this by computing an associated family of martingales.

Lemma 7.4. Let $\left(M_{t}^{N, f}\right)_{t \geq 0}$ be the processes defined in Formula 5.1. Recalling the notation $Q_{t}=Q \circ \phi_{t}$, define

$$
\begin{equation*}
\chi\left(s, t, \mu^{N}\right)=Q_{t-s}\left(\mu^{N}\right)-Q_{t-s}\left(\mu_{s-}^{N}\right) . \tag{7.26}
\end{equation*}
$$

Suppose $f$ satisfies a growth condition $|f(v)| \leq\left(1+|v|^{q}\right)$, for some $q \geq 0$. Consider the martingales $Z_{t}^{N, f}$ given by

$$
\begin{equation*}
\left.Z_{t}^{N, f}=\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f, \mu^{N}-\mu_{s-}^{N}\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu^{N}\right)\right\rangle \tag{7.27}
\end{equation*}
$$

Then we have the equality

$$
\begin{align*}
Z_{t}^{N, f} & =M_{t}^{N, f}-C_{t}^{N, f} \\
& =M_{t}^{N, f}-\int_{0}^{t} d s \int_{(0, s] \times \mathcal{S}_{N}}\left\langle f, \chi\left(u, s, \mu^{N}\right)\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d u, d \mu^{N}\right) \tag{7.28}
\end{align*}
$$

Proof. First, we note that $Z_{t}^{N, f}$ are martingales by standard results from Markov chains, (A.1). Observe that the integrand in the definition of $C_{t}^{N, f}$ is bounded, since whenever $0 \leq u \leq s$, and $\mu^{N}$ is obtain from $\mu_{u-}^{N}$ by collision, we use the estimate (3.17) with $\eta=\frac{1}{2}$, to obtain for some $k$,

$$
\begin{align*}
\left|\left\langle f, \chi\left(u, s, \mu^{N}\right)\right\rangle\right| & \leq\left\|Q_{s-u}\left(\mu^{N}\right)-Q_{s-u}\left(\mu_{u-}^{N}\right)\right\|_{\mathrm{TV}+q} \\
& \lesssim \Lambda_{k}\left(\mu^{N}, \mu_{u-}^{N}\right)^{\frac{1}{2}} N^{-\frac{1}{2}} \lesssim N^{\frac{k-2}{4}}<\infty . \tag{7.29}
\end{align*}
$$

Moreover, for initial data $\mu^{N} \in \mathcal{S}_{N}$, the Boltzmann flow $\left(\phi_{s}\left(\mu^{N}\right)\right)_{s=0}^{t}$ has uniformly bounded $(q+1)$ th moments and so, by approximation, the Boltzmann dynamics (BE) extend to $f$. Now, we apply Fubini to the integral:

$$
\begin{align*}
C_{t}^{N, f}= & \int_{(0, t] \times \mathcal{S}_{N}} \int_{0}^{t} d s\left\langle f, Q_{s-u}\left(\mu^{N}\right)-Q_{s-u}\left(\mu_{u-}^{N}\right)\right\rangle \\
& \times 1[u \leq s \leq t]\left(m^{N}-\bar{m}^{N}\right)\left(d u, d \mu^{N}\right) \\
= & \int_{(0, t] \times \mathcal{S}_{N}}\left\{\int_{u}^{t}\left(\left\langle f, Q_{s-u}\left(\mu^{N}\right)\right\rangle-\left\langle f, Q_{s-u}\left(\mu_{u-}^{N}\right)\right\rangle\right) d s\right\} \\
& \times\left(m^{N}-\bar{m}^{N}\right)\left(d u, d \mu^{N}\right)  \tag{7.30}\\
= & \int_{(0, t] \times \mathcal{S}_{N}}\left\{\left\langle f, \phi_{t-u}\left(\mu^{N}\right)-\phi_{t-u}\left(\mu_{u-}^{N}\right)\right\rangle-\left\langle f, \mu^{N}-\mu_{u-}^{N}\right)\right\} \\
& \times\left(m^{N}-\bar{m}^{N}\right)\left(d u, d \mu^{N}\right) \\
= & M_{t}^{N, f}-Z_{t}^{N, f}
\end{align*}
$$

where the third equality is precisely the (extended) Boltzmann dynamics (BE) in the variable $s \in[u, t]$.

To prove Lemma 7.1, we return to the decomposition (6.16) used in the proof of Lemma 6.1. Our first point is to establish a control on

$$
\begin{equation*}
\mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{M_{\star, \text { fin }}^{N ; B}\right\}^{p}\right], \tag{7.31}
\end{equation*}
$$

where $\star$ denotes the local supremum (7.2). We will do so by breaking the supremum into two parts, each of which can be controlled by elementary martingale estimates. Let $\left(J_{s}^{N ; B ; t}\right)_{0 \leq s \leq t}$ be the process

$$
\begin{equation*}
J_{s}^{N ; B ; t}=\int_{(0, s] \times \mathcal{S}_{N}}\left\langle h_{B}, Q_{t-u}\left(\mu^{N}\right)-Q_{t-u}\left(\mu_{u-}^{N}\right)\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d u, d \mu^{N}\right), \tag{7.32}
\end{equation*}
$$

where, as in the proof of Theorem 1.1,

$$
\begin{equation*}
h_{B}=2^{2 j}\left(1+|v|^{2}\right) 1_{B} ; \quad B \in \mathcal{P}_{j, l} . \tag{7.33}
\end{equation*}
$$

Each process $\left(J_{s}^{N ; B: t}\right)_{0 \leq s \leq t}$ is a martingale, by standard results for Markov chains (A.1). Writing $Z^{N ; B}=Z^{N, h_{B}}$, Lemma 7.4 gives

$$
\begin{equation*}
Z_{t}^{N ; B}=M_{t}^{N ; B}+\int_{0}^{t} J_{s}^{N ; B ; s} d s \tag{7.34}
\end{equation*}
$$

Lemma 7.5. Let $p \geq 2$, and let $p^{\prime}$ be the Hölder conjugate to $p$. In the notation above, we have the comparison

$$
\begin{align*}
& \mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|M_{\star,, t_{\text {fin }}}^{N ; B}\right|\right\}^{p}\right]  \tag{7.35}\\
& \quad \lesssim \mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|M_{t_{\mathrm{fin}}}^{N ; B}\right|^{p}+t_{\mathrm{fin}}^{p / p^{\prime}} \int_{0}^{t_{\mathrm{fin}}}\left|J_{t}^{N ; B ; t}\right|^{p} d t\right\}\right]
\end{align*}
$$

Proof. For each $B$, we observe that

$$
\begin{equation*}
\sup _{t \leq t_{\mathrm{fin}}}\left|M_{t}^{N ; B}-Z_{t}^{N ; B}\right| \leq \int_{0}^{t_{\mathrm{fin}}}\left|J_{s}^{N ; B ; s}\right| d s \tag{7.36}
\end{equation*}
$$

which implies the two bounds

$$
\begin{equation*}
M_{\star, t_{\mathrm{fin}}}^{N ; B} \leq Z_{\star, t_{\mathrm{fin}}}^{N ; B}+\int_{0}^{t_{\mathrm{fin}}}\left|J_{s}^{N ; B ; s}\right| d s ; \quad Z_{t_{\mathrm{fin}}}^{N ; B} \leq M_{t_{\mathrm{fin}}}^{N ; B}+\int_{0}^{t_{\mathrm{fin}}}\left|J_{s}^{N ; B ; s}\right| d s \tag{7.37}
\end{equation*}
$$

By Doob's $L^{p}$ inequality, we have

$$
\begin{equation*}
\left\|Z_{\star, t_{\text {fin }}}^{N ; B}\right\|_{L^{p}(\mathbb{P})} \leq p^{\prime}\left\|Z_{t_{\text {fin }}}^{N ; B}\right\|_{L^{p}(\mathbb{P})} \tag{7.38}
\end{equation*}
$$

Combining (7.37), (7.38), we obtain

$$
\begin{equation*}
\left\|M_{\star, t_{\text {fin }}}^{N ; B}\right\|_{L^{p}(\mathbb{P})} \lesssim\left\|M_{t_{\text {fin }}}^{N ; B}\right\|_{L^{p}(\mathbb{P})}+\left\|\int_{0}^{t_{\text {fin }}}\left|J_{s}^{N ; B ; s}\right| d s\right\|_{L^{p}(\mathbb{P})} \tag{7.39}
\end{equation*}
$$

Using Hölder's inequality on the integral,

$$
\begin{align*}
\mathbb{E}\left[\left\{M_{\star, t t_{\text {fin }}}^{N ; B}\right\}^{p}\right] & \lesssim \mathbb{E}\left[\left|M_{t_{\text {fin }}}^{N ; B}\right|^{p}\right]+\mathbb{E}\left[\left\{\int_{0}^{t_{\mathrm{fin}}}\left|J_{s}^{N ; B ; s}\right| d s\right\}^{p}\right]  \tag{7.40}\\
& \lesssim \mathbb{E}\left[\left|M_{t_{\text {fin }}}^{N ; B}\right|^{p}\right]+t_{\mathrm{fin}}^{p / p^{\prime}} \int_{0}^{t_{\mathrm{fin}}} \mathbb{E}\left[\left|J_{t}^{N ; B ; t}\right|^{p}\right] d s
\end{align*}
$$

Summing over $B \in \mathcal{P}_{j, l}$ and $j=0,1, \ldots, J$, we obtain the desired comparison.

Proof of Lemma 7.1. We begin by controlling the integral term in Lemma 7.5. The quadratic variation is given by

$$
\begin{align*}
{\left[J^{N ; B ; t}\right]_{s} } & =\int_{(0, s] \times \mathcal{S}_{N}}\left\langle h_{B}, \chi\left(u, t, \mu^{N}\right)\right\rangle^{2} m^{N}\left(d u, d \mu^{N}\right)  \tag{7.41}\\
& \leq \int_{(0, s] \times \mathcal{S}_{N}}\left\langle h_{B},\right| \chi\left(u, t, \mu^{N}\right)| \rangle^{2} m^{N}\left(d u, d \mu^{N}\right),
\end{align*}
$$

where $h_{B}$ is as in (7.33) and $\chi$ is as in (7.26). Hence, using Burkholder's inequality (A.1) we see that, for all $t \leq t_{\text {fin }}$,

$$
\begin{align*}
\mathbb{E} & {\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|J_{t}^{N ; B ; t}\right|\right\}^{p}\right] }  \tag{7.42}\\
& \lesssim \mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\int_{(0, t] \times \mathcal{S}_{N}}\left\langle h_{B},\right| \chi\left(u, t, \mu^{N}\right)| \rangle^{2} m^{N}\left(d u, d \mu^{N}\right)\right\}^{p / 2}\right]
\end{align*}
$$

Using Minkowski's inequality to move the double sum inside the parentheses, and recalling that $\sum_{j} \sum_{B \in \mathcal{P}_{j, l}} h_{B} \lesssim\left(1+|v|^{4}\right)$, we obtain the bound

$$
\begin{align*}
& \mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|J_{t}^{N ; B ; t}\right|\right\}^{p}\right] \\
& \left.\quad \lesssim \mathbb{E}\left[\left\{\left.\int_{(0, t] \times \mathcal{S}_{N}}\langle 1+| v\right|^{4},\left|\chi\left(u, t, \mu^{N}\right)\right|\right\rangle^{2} m^{N}\left(d u, d \mu^{N}\right)\right\}^{p / 2}\right]  \tag{7.43}\\
& \quad=\mathbb{E}\left[\left\{\int_{(0, t] \times \mathcal{S}_{N}}\left\|Q_{t-u}\left(\mu^{N}\right)-Q_{t-u}\left(\mu_{u-}^{N}\right)\right\|_{\mathrm{TV}+4}^{2} m^{N}\left(d u, d \mu^{N}\right)\right\}^{p / 2}\right]
\end{align*}
$$

where the equality is the definition of $\chi$ (7.26).
Using the continuity estimate for $Q$ established in (3.17), and arguing as in the proof of Lemma 6.1, we see that almost surely, for $m^{N}$ almost all $\left(u, \mu^{N}\right)$, we have

$$
\begin{equation*}
\left\|Q_{t-u}\left(\mu^{N}\right)-Q_{t-u}\left(\mu_{u-}^{N}\right)\right\|_{\mathrm{TV}+4} \lesssim N^{\epsilon-1} \Lambda_{k}\left(\mu_{u-}^{N}\right) \tag{7.44}
\end{equation*}
$$

Therefore, using Cauchy-Schwarz, (7.43) gives the bound

$$
\begin{align*}
\mathbb{E} & {\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|J_{t}^{N ; B ; t}\right|\right\}^{p}\right] }  \tag{7.45}\\
& \lesssim N^{p(\epsilon-1)} \mathbb{E}\left[\sup _{t \leq t_{\mathrm{fin}}} \Lambda_{k p}\left(\mu_{t}^{N}\right)\right]^{1 / 2}\left\|m^{N}\left(\left(0, t_{\mathrm{fin}}\right] \times \mathcal{S}_{N}\right)\right\|_{L^{p}(\mathbb{P})}^{p / 2} .
\end{align*}
$$

The moment term is controlled by the initial moment bound and Proposition 2:

$$
\begin{equation*}
\mathbb{E}\left[\sup _{t \leq t_{\mathrm{fin}}} \Lambda_{k p}\left(\mu_{t}^{N}\right)\right] \lesssim\left(1+t_{\mathrm{fin}}\right) \Lambda_{k p}\left(\mu_{0}^{N}\right) \leq\left(1+t_{\mathrm{fin}}\right) a^{p} \tag{7.46}
\end{equation*}
$$

Since the rates of the Kac process are bounded by $2 N$, we can stochastically dominate $m^{N}\left(d t \times \mathcal{S}_{N}\right)$ by a Poisson random measure $\mathfrak{m}^{N}(d t)$ of rate $2 N$. By the additive property of Poisson processes, it follows that

$$
\begin{equation*}
\left\|m^{N}\left(\left(0, t_{\mathrm{fin}}\right] \times \mathcal{S}_{N}\right)\right\|_{L^{p}(\mathbb{P})} \leq\left\|\mathfrak{m}^{N}\left(0, t_{\mathrm{fin}}\right]\right\|_{L^{p}(\mathbb{P})} \lesssim N\left(1+t_{\mathrm{fin}}\right) \tag{7.47}
\end{equation*}
$$

Combining (7.45), (7.46), (7.47), we have the control of the integrand:

$$
\begin{equation*}
\sup _{t \leq t_{\text {fin }}} \mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|J_{t}^{N ; B ; t}\right|\right\}^{p}\right] \lesssim N^{p(\epsilon-1 / 2)} a^{p / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{p+1}{2}} . \tag{7.48}
\end{equation*}
$$

This gives the following control of the integral term in Lemma 7.5:

$$
\begin{equation*}
t_{\mathrm{fin}}^{p / p^{\prime}} \mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}} \int_{0}^{t_{\mathrm{fin}}}\left\{\left|J_{t}^{N ; B ; t}\right|\right\}^{p} d t\right] \lesssim N^{p(\epsilon-1 / 2)} a^{p / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{p+3}{2}+\frac{p}{p^{\prime}}} \tag{7.49}
\end{equation*}
$$

Using the definition of $p^{\prime}$ as the Hölder conjugate to $p$, it is straightforward to see that the exponent of $\left(1+t_{\mathrm{fin}}\right)$ is $\frac{3 p+1}{2}$.

We now perform a similar analysis for the terms $M_{t_{\text {fin }}}^{N ; B}$ in Lemma 7.5. Let $\left(M_{s}^{N ; B ; t}\right)_{s \leq t}$ be the martingale defined in (6.20). The quadratic variation is

$$
\begin{align*}
{\left[M^{N ; B ; t}\right]_{s} } & =\int_{(0, s] \times \mathcal{S}_{N}}\left\langle h_{B}, \phi_{t-u}\left(\mu^{N}\right)-\phi_{t-u}\left(\mu_{u-}^{N}\right)\right\rangle^{2} m^{N}\left(d u, d \mu^{N}\right) \\
& \leq \int_{(0, s] \times \mathcal{S}_{N}}\left\langle h_{B},\right| \phi_{t-u}\left(\mu^{N}\right)-\phi_{t-u}\left(\mu_{u-}^{N}\right)| \rangle^{2} m^{N}\left(d u, d \mu^{N}\right) \tag{7.50}
\end{align*}
$$

Arguing using Burkholder and the stability estimate Corollary 3.2, an identical calculation to the above shows that

$$
\begin{equation*}
\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\|M_{\mathrm{ffin}^{N ; B}}^{N}\right\|_{L^{p}(\mathbb{P})}^{p} \lesssim N^{p(\epsilon-1 / 2)} a^{p / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{p+1}{2}} \tag{7.51}
\end{equation*}
$$

Hence, by Lemma 7.5, we obtain

$$
\begin{equation*}
\mathbb{E}\left[\sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{\left|M_{\star, t_{\text {fin }}}^{N ; B}\right|\right\}^{p}\right] \lesssim N^{p(\epsilon-1 / 2)} a^{p / 2}\left(1+t_{\text {fin }}\right)^{\frac{3 p+1}{2}} \tag{7.52}
\end{equation*}
$$

We control the coefficients $2^{-2 j} a_{B}(f)$ as in the argument of Lemma 6.1. Using Hölder's inequality in place of Cauchy-Schwarz, we obtain

$$
\begin{align*}
&\left\|\sup _{f \in \mathcal{A}} \sup _{t \leq t_{\text {fin }}}\left|\sum_{j=0}^{J} \sum_{l=2}^{L} \sum_{B \in \mathcal{P}_{j, l}} 2^{-2 j} a_{B}(f) M_{t}^{N ; B}\right|\right\|_{L^{p}(\mathbb{P})} \\
& \lesssim \sum_{l=2}^{L}\left[\mathbb{E} \sum_{j=0}^{J} \sum_{B \in \mathcal{P}_{j, l}}\left\{M_{\star, l_{\text {fin }}}^{N ; B}\right\}^{p}\right]^{1 / p} 2^{\left(d / p^{\prime}-1\right) l} J^{1 / p^{\prime}}  \tag{7.53}\\
& \lesssim \sum_{l=2}^{L} N^{\epsilon-\frac{1}{2}} a^{1 / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{3 p+1}{2 p}} 2^{\left(d / p^{\prime}-1\right) l} J^{1 / p^{\prime}} \\
& \lesssim N^{\epsilon-\frac{1}{2}} a^{1 / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{3 p+1}{2 p}} 2^{\left(d / p^{\prime}-1\right) L} J^{1 / p^{\prime}}
\end{align*}
$$

Following the argument of Lemma 6.1, we wish to control the error terms $R_{t}^{N, f}$ given by (6.17), locally uniformly in time. As in (6.30), we majorise, for $m^{N}+\bar{m}^{N}$ almost all $\left(s, \mu^{N}\right)$,

$$
\begin{align*}
\sup _{f \in \mathcal{A}}\left|\left\langle\beta(f), \phi_{t-s}\left(\mu^{N}\right)-\phi_{t-s}\left(\mu_{s-}^{N}\right)\right\rangle\right| & \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon-1} \Lambda_{k}\left(\mu_{s-}^{N}\right)^{\frac{1}{2}}  \tag{7.54}\\
& =: H_{s}^{\prime} .
\end{align*}
$$

As in (6.31), we may bound
(7.55) $\left\|\sup _{f \in \mathcal{A}} \sup _{t \leq t_{\mathrm{fin}}}\left|R_{t}^{N, f}\right|\right\|_{L^{p}(\mathbb{P})} \leq\left\|\int_{0}^{t_{\mathrm{fin}}} H_{s}^{\prime}\left(m^{N}+\bar{m}^{N}\right)\left(d s, \mathcal{S}_{N}\right)\right\|_{L^{p}(\mathbb{P})} \leq \mathcal{T}_{1}+\mathcal{T}_{2}$,
where the two error terms are

$$
\begin{equation*}
\mathcal{T}_{1}=\left\|\int_{0}^{t_{\mathrm{fin}}} H_{s}^{\prime} m^{N}\left(d s, \mathcal{S}_{N}\right)\right\|_{L^{p}(\mathbb{P})} \tag{7.56}
\end{equation*}
$$

and

$$
\begin{equation*}
\mathcal{T}_{2}=\left\|\int_{0}^{t_{\mathrm{fin}}} H_{s}^{\prime} \bar{m}^{N}\left(d s, \mathcal{S}_{N}\right)\right\|_{L^{p}(\mathbb{P})} \tag{7.57}
\end{equation*}
$$

We now deal with the two terms separately. For the $\mathcal{T}_{1}$, we dominate $\bar{m}^{N}\left(d s, \mathcal{S}_{N}\right) \leq$ $2 N d s$ to see that

$$
\begin{equation*}
\int_{0}^{t_{\mathrm{fin}}} H_{s}^{\prime} \bar{m}^{N}\left(d s, \mathcal{S}_{N}\right) \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon} t_{\mathrm{fin}}\left(\sup _{s \leq t_{\mathrm{fin}}} \Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}}\right) \tag{7.58}
\end{equation*}
$$

Using the monotonicity of $L^{p}$ norms, and using the moment control in the usual way,

$$
\begin{align*}
\mathcal{T}_{1} & \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon} t_{\mathrm{fin}} \mathbb{E}\left[\sup _{s \leq f_{\mathrm{fin}}} \Lambda_{p k}\left(\mu_{s}^{N}\right)\right]^{\frac{1}{2 p}}  \tag{7.59}\\
& \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon} a^{1 / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{2 p+1}{2 p}}
\end{align*}
$$

For $\mathcal{T}_{2}$, we dominate $m^{N}\left(d s, \mathcal{S}_{N}\right)$ by a Poisson random measure $\mathfrak{m}^{N}(d s)$ of rate $2 N$, as above. Controlling $\mathfrak{m}^{N}$ as in (7.47), we obtain

$$
\begin{align*}
\mathcal{T}_{2} & \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon-1}\left\|\int_{0}^{t_{\mathrm{fin}}} \Lambda_{k}\left(\mu_{s-}^{N}\right)^{\frac{1}{2}} \mathfrak{m}^{N}(d s)\right\|_{L^{p}(\mathbb{P})} \\
& \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon-1}\left\|\left(\sup _{s \leq t_{\mathrm{fin}}} \Lambda_{k}\left(\mu_{s}^{N}\right)^{\frac{1}{2}}\right)\right\|_{L^{2 p}(\mathbb{P})}\left\|\mathfrak{m}^{N}\left(\left(0, t_{\mathrm{fin}}\right)\right)\right\|_{L^{2 p}(\mathbb{P})}  \tag{7.60}\\
& \lesssim\left(2^{-2 J}+2^{-L}\right) N^{\epsilon}\left(1+t_{\mathrm{fin}}\right)^{\frac{2 p+1}{2 p}} .
\end{align*}
$$

Combining the local uniform estimates (7.53), (7.55), (7.59), (7.60) of the terms in the decomposition (6.16), we find that

$$
\left\|\sup _{f \in \mathcal{A}} M_{\star, \text { fin }}^{N, f}\right\|_{L^{p}(\mathbb{P})} \lesssim N^{\epsilon} a^{1 / 2}\left(1+t_{\mathrm{fin}}\right)^{\frac{3 p+1}{2 p}}\left(N^{-1 / 2} 2^{(d / q-1) L} J^{1 / p^{\prime}}+2^{-2 J}+2^{-L}\right)
$$

Taking $J=\left\lfloor\frac{p^{\prime}}{4 d} \log _{2}(N)\right\rfloor$ and $L=\left\lfloor\frac{p^{\prime}}{2 d} \log _{2}(N)\right\rfloor$ proves the result claimed.
8. Proof of Theorem 1.3. We now turn to the proof of Theorem 1.3, which establishes a convergence estimate in the presence of a $k$ th moment bound, for any $k>2$. Our strategy will be to use the ideas of [33], which work well with few moments, to prove convergence on a small initial time interval $\left[0, u_{N}\right]$, for some $u_{N}$ to be chosen later. Then, thanks to the moment production property recalled in Proposition 2, we may use Theorems 1.1, 1.2 to control the behaviour at times $t \geq u_{N}$. The argument is similar to the final argument in the proof of Theorem 1.6 given in Section 4, which may be read as a warm-up to this proof.

Throughout, let $k, a,\left(\mu_{t}^{N}\right), \mu_{0}$ be as in the statement of the theorem.
We begin by recalling the representation formula established in [33], Proposition 4.2 , which is a noisy version of Proposition 5.

Proposition 6. Let $\mu \in \mathcal{S}^{k}$ for some $k>2$, and let $\mu_{t}^{N}$ be a Kac process on $N$ particles. Let $\rho_{t}=\left(\phi_{t}\left(\mu_{0}\right)+\mu_{t}^{N}\right) / 2$, and for $f \in \mathcal{A}, 0 \leq s \leq t$, let $f_{\text {st }}$ be the propagation described in Definition 4.1 in this environment. Then, for all $t \geq 0$, we have the equality

$$
\begin{align*}
\left\langle f, \mu_{t}^{N}-\phi_{t}\left(\mu_{0}\right)\right\rangle= & \left\langle f_{0 t}, \mu_{0}^{N}-\mu_{0}\right\rangle \\
& +\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f_{\mathrm{st}}, \mu^{N}-\mu_{s-}^{N}\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu^{N}\right) \tag{8.1}
\end{align*}
$$

where $m^{N}, \bar{m}^{N}$ are as defined in Section 5.
The major difficulty in using this representation formula is the appearance of an exponentiated random moment in the quantity $z_{t}$ parametrising the continuity of $f_{\mathrm{st}}$. We will use the following proposition, which controls the stochastic integrals on the right-hand side, modulo this difficulty.

PROPOSITION 7. Let $\rho_{t}$ be a potentially random environment such that, for some $\beta>0$,

$$
\begin{equation*}
w=\left\|\sup _{t \leq 1}\left(\frac{\Lambda_{3}\left(\rho_{t}\right)}{\beta t^{\beta-1}+1}\right)\right\|_{L^{\infty}(\mathbb{P})}<\infty \tag{8.2}
\end{equation*}
$$

For $f \in \mathcal{A}$ and $0 \leq s \leq t \leq 1$, let $f_{\mathrm{st}}[\rho]$ denote the propagation in this environment, as described in Definition 4.1.

Let $k>2$ and $a \geq 1$, and let $\mu_{t}^{N}$ be a Kac process with initial moment $\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a$, and let $\bar{m}^{N}, \bar{m}^{N}$ be as in Section 5 . We write

$$
\begin{equation*}
\widetilde{M}_{t}^{N, f}[\rho]=\int_{(0, t] \times \mathcal{S}_{N}}\left\langle f_{\mathrm{st}}[\rho], \mu^{N}-\mu_{s-}^{N}\right\rangle\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu^{N}\right) . \tag{8.3}
\end{equation*}
$$

In this notation, we have the bound

$$
\begin{equation*}
\left\|\sup _{t \leq 1} \sup _{f \in \mathcal{A}} \widetilde{M}_{t}^{N, f}[\rho]\right\|_{1} \leq C a N^{-\eta} \tag{8.4}
\end{equation*}
$$

for some $C=C(d, k, \beta)$ and $\eta=\eta(d, \beta)>0$. Here, we emphasise that $\|\cdot\|_{L^{1}(\mathbb{P})}$ refers to the $L^{1}$ norm with simultaneous expectation over $\mu_{t}^{N}$ and the environment $\rho$.

This largely follows from the proof of [33], Theorem 1.1, and the argument follows a similar pattern to Lemmas 6.1, 7.1, using the continuity estimate recalled in Proposition 4 and a similar estimate for the dependence on the initial time $s$. The key difference is that the hypotheses on the environment $\rho$ guarantee an $L^{\infty}(\mathbb{P})$ control on the quantities

$$
\begin{align*}
& z_{1}=\exp \left(8 \int_{0}^{1} \Lambda_{3}\left(\rho_{u}\right) d u\right) \\
& y_{\beta}=z_{1} \sup _{0 \leq s \leq s^{\prime} \leq 1}\left[\left(s^{\prime}-s\right)^{-\beta} \int_{s}^{s^{\prime}} \Lambda_{3}\left(\rho_{u}\right) d u\right] \tag{8.5}
\end{align*}
$$

which describe the continuity of $f_{\text {st }}(v)$ in $v$ and $s$, respectively. By contrast, these are only controlled in probability in [33], Theorem 1.1; correspondingly, we obtain an $L^{1}(\mathbb{P})$ estimate rather than an estimate in probability. With this estimate, we turn to the proof of Theorem 1.3.

Proof of Theorem 1.3. We first introduce a localisation argument, following the argument in Section 4, which allows us to guarantee that (8.2) holds for the environment $\rho=\left(\mu_{t}^{N}+\phi_{t}\left(\mu_{0}\right)\right) / 2$. Let $\beta=\frac{k-2}{2}$, and let $u_{N} \leq 1$ be chosen later. Now, define $T_{N}$ to be the stopping time

$$
\begin{equation*}
T_{N}=\inf \left\{t \leq u_{N}: \Lambda_{3}\left(\rho_{t}\right)>\frac{\left(\beta t^{\beta-1}+1\right)}{8 \sqrt{2}}\right\} \tag{8.6}
\end{equation*}
$$

We use the convention that $\inf \varnothing=\infty$, so that if $T_{N}>u_{N}$, then $T_{N}=\infty$. Let $\rho^{T}$ be the stopped environment $\rho_{t}^{T}=\rho_{t \wedge T_{N}}$, and write $f_{\mathrm{st}}^{T}$ for the propagation in the stopped environment.

We observe first that on the event $T_{N}=\infty$, we have the equality $f_{\mathrm{st}}^{T}=f_{\mathrm{st}}$ for all $f \in \mathcal{A}, s \leq t \leq u_{N}$. Moreover, since $\Lambda_{3}\left(\rho_{t}\right)$ increases by a factor of at most $4 \sqrt{2}$ at jumps by Lemma 2.1, we have the bound, almost surely for all $t \geq 0$,

$$
\begin{equation*}
\Lambda_{3}\left(\rho_{t}^{T}\right) \leq \frac{\left(\beta t^{\beta-1}+1\right)}{2} \tag{8.7}
\end{equation*}
$$

Therefore, the stopped environment $\rho^{T}$ satisfies the bound 8.2 with $w=\frac{1}{2}$. Now, we write $\widetilde{M}_{t}^{N, f}=\widetilde{M}_{t}^{N, f}\left[\rho^{T}\right]$ as in the proposition above, and by the representation formula in Proposition 6, we have the bound for all $t \leq u_{N}$,

$$
\begin{equation*}
W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right) 1\left[T_{N}=\infty\right] \leq C W\left(\mu_{0}^{N}, \mu_{0}\right)+\sup _{f \in \mathcal{A}} \widetilde{M}_{t}^{N, f} \tag{8.8}
\end{equation*}
$$

for some absolute constant $C$. By Proposition 7, we obtain the estimate

$$
\begin{equation*}
\left\|\sup _{t \leq u_{N}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right) 1\left[T_{N}=\infty\right]\right\|_{1} \lesssim W\left(\mu_{0}^{N}, \mu_{0}\right)+a N^{-\eta} . \tag{8.9}
\end{equation*}
$$

Let $k_{0}=k_{0}(d)$ be large enough that Theorem 1.1 holds with $\epsilon=\frac{1}{2 d}$. By applying Theorem 1.1, restarted at time $u_{N}$, and the moment production property, we obtain

$$
\begin{align*}
\sup _{t \geq u_{N}}\left\|W\left(\mu_{t}^{N}, \phi_{t-u_{N}}\left(\mu_{u_{N}}^{N}\right)\right)\right\|_{2} & \lesssim N^{\epsilon-1 / d} \mathbb{E}\left[\Lambda_{k_{0}}\left(\mu_{u_{N}}^{N}\right)\right]^{1 / 2}  \tag{8.10}\\
& \lesssim N^{\epsilon-1 / d} u_{N}^{1-k_{0} / 2}
\end{align*}
$$

Using our continuity estimate Theorem 1.6, we have the bound for some $\zeta=\zeta(d)$

$$
\begin{align*}
& \sup _{t \geq u_{N}} W\left(\phi_{t-u_{N}}\left(\mu_{u_{N}}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right)  \tag{8.11}\\
& \quad \lesssim W\left(\mu_{u_{N}}^{N}, \phi_{u_{N}}\left(\mu_{0}\right)\right)^{\zeta} \Lambda_{k_{0}}\left(\mu_{u_{N}}^{N}, \phi_{u_{N}}\left(\mu_{0}\right)\right)
\end{align*}
$$

and, considering the cases $\left\{T_{N} \leq u_{N}\right\},\left\{T_{N}=\infty\right\}$ separately, we see that

$$
\begin{align*}
& \sup _{t \geq u_{N}} W\left(\phi_{t-u_{N}}\left(\mu_{u_{N}}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right) \\
& \quad \lesssim W\left(\mu_{u_{N}}^{N}, \phi_{u_{N}}\left(\mu_{0}\right)\right)^{\zeta} \Lambda_{k_{0}}\left(\mu_{u_{N}}^{N}, \phi_{u_{N}}\left(\mu_{0}\right)\right) 1\left[T_{N}=\infty\right]  \tag{8.12}\\
& \quad+1\left[T_{N} \leq u_{N}\right]
\end{align*}
$$

To ease notation, we will write $\mathcal{T}_{1}, \mathcal{T}_{2}$ for the two terms, respectively. We estimate the expectation of $\mathcal{T}_{1}$ using Hölder's inequality: for some $k_{1}>k_{0}$,

$$
\begin{align*}
\left\|\mathcal{T}_{1}\right\|_{L^{1}(\mathbb{P})} & \lesssim \mathbb{E}\left(W\left(\mu_{u_{N}}^{N}, \phi_{u_{N}}\left(\mu_{0}\right)\right) 1\left[T_{N}=\infty\right]\right)^{\zeta} \mathbb{E}\left(\Lambda_{k_{1}}\left(\mu_{u_{N}}^{N}, \phi_{u_{N}}\left(\mu_{0}\right)\right)\right) \\
& \lesssim\left(N^{-\eta}+W\left(\mu_{0}^{N}, \mu_{0}\right)\right)^{\zeta} u_{N}^{1-k_{1} / 2} \tag{8.13}
\end{align*}
$$

where $\eta$ is as in (8.9) with our choice of $\beta$. In order to deal with $\mathcal{T}_{2}$, we now estimate $\mathbb{P}\left(T_{N} \leq u_{N}\right)$. Let $Z_{N}$ be given by

$$
\begin{equation*}
\left.Z_{N}=\left.\sum_{l: 2^{-l} \leq u_{N}} 2^{(\beta-1) l+1} \beta^{-1} \sup _{t \in\left[2^{-l}, 2^{1-l}\right]}\langle 1+| v\right|^{3}, \rho_{t}\right\rangle \tag{8.14}
\end{equation*}
$$

and observe that, for all $t \leq u_{N}$, we have the bound

$$
\begin{equation*}
\left.\left.\langle 1+| v\right|^{3}, \rho_{t}\right\rangle \leq \frac{\left(\beta t^{\beta-1}+1\right) Z_{N}}{2} \tag{8.15}
\end{equation*}
$$

Therefore,

$$
\begin{equation*}
\mathbb{P}\left(T_{N} \leq u_{N}\right) \leq \mathbb{P}\left(Z_{N}>1 / 8\right) \leq 8 \mathbb{E}\left[Z_{N}\right] \tag{8.16}
\end{equation*}
$$

Using the moment production property of the Kac process and Boltzmann equation in Proposition 2, we compute

$$
\begin{equation*}
\mathbb{E}\left(Z_{N}\right) \leq \sum_{l: 2^{-l} \leq u_{N}} 2^{(\beta-1) l+1} 2^{-l(k-3)} \beta^{-1} a \lesssim a u_{N}^{\beta} \tag{8.17}
\end{equation*}
$$

and so

$$
\begin{align*}
& \left\|\sup _{t \geq u_{N}} W\left(\phi_{t-u_{N}}\left(\mu_{u_{N}}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right)\right\|_{L^{1}(\mathbb{P})}  \tag{8.18}\\
& \quad \lesssim\left(N^{-\eta}+W\left(\mu_{0}^{N}, \mu_{0}\right)\right)^{\zeta} u_{N}^{1-k_{1} / 2}+a u_{N}^{\beta} .
\end{align*}
$$

We now return to (8.9) and observe that

$$
\begin{align*}
& \left\|\sup _{t \leq u_{N}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)\right\|_{1} \\
& \quad \lesssim\left\|\sup _{t \leq u_{N}} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right) 1\left[T_{N}=\infty\right]\right\|_{1}+\mathbb{P}\left(T_{N} \leq u_{N}\right)  \tag{8.19}\\
& \quad \lesssim W\left(\mu_{0}^{N}, \mu_{0}\right)+a N^{-\eta}+a u_{N}^{\beta} .
\end{align*}
$$

Combining (8.10), (8.18), (8.19) and keeping the worst terms, we have shown that

$$
\begin{equation*}
\sup _{t \geq 0} \| W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right) \|_{L^{1}(\mathbb{P})} \lesssim\left(N^{-\eta}+W\left(\mu_{0}^{N}, \mu_{0}\right)\right)^{\delta} u_{N}^{-\alpha}+a u_{N}^{\beta}\right. \tag{8.20}
\end{equation*}
$$

for some $\eta, \delta, \alpha, \beta>0$, depending on $d, k$. If we choose

$$
\begin{equation*}
u_{N}=\left(N^{-\eta}+W\left(\mu_{0}^{N}, \mu_{0}\right)\right)^{\delta /(\alpha+\beta)} \tag{8.21}
\end{equation*}
$$

then we finally obtain

$$
\begin{align*}
\sup _{t \geq 0}\left\|W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}\right)\right)\right\|_{L^{1}(\mathbb{P})} & \lesssim a\left(N^{-\eta}+W\left(\mu_{0}^{N}, \mu_{0}\right)\right)^{-\beta \delta /(\alpha+\beta)} \\
& \lesssim a\left(N^{-\eta \beta \delta /(\alpha+\beta)}+W\left(\mu_{0}^{N}, \mu_{0}\right)^{\beta \delta /(\alpha+\beta)}\right)  \tag{8.22}\\
& \lesssim a\left(N^{-\epsilon}+W\left(\mu_{0}^{N}, \mu_{0}\right)^{\epsilon}\right)
\end{align*}
$$

as desired, for sufficiently small $\epsilon=\epsilon(d, k)>0$. The case for the local uniform estimate is similar, using Theorem 1.2 in place of Theorem 1.1.
9. Proof of Theorem 1.4. The proof of Theorem 1.4 is based on the following heuristic argument.

Heuristic. Fix $N$, and consider a Kac process ( $\mu_{t}^{N}$ ) on $N$ particles. As $t \rightarrow \infty$, its law relaxes to the equilibrium distribution $\pi_{N}$, which is known to be the uniform distribution $\sigma^{N}$ on $\mathcal{S}_{N}$. Since this measure assigns nonzero probability to regions $R_{N}$ at macroscopic distance from the fixed-point $\gamma$, given by

$$
\begin{equation*}
\gamma(d v)=\frac{e^{-\frac{d}{2}|v|^{2}}}{\left(2 \pi d^{-1}\right)^{d / 2}} d v \tag{9.1}
\end{equation*}
$$

the process will almost surely hit $R_{N}$ on an unbounded set of times. Meanwhile, the Boltzmann flow $\phi_{t}\left(\mu_{0}\right)$ will converge to $\gamma$. Therefore, at some large time, the particle system $\mu_{t}^{N}$ will have macroscopic distance from the Boltzmann flow $\phi_{t}\left(\mu_{0}^{N}\right)$.

The regions $R_{N}$ which we construct in the proof are those where the energy is concentrated in only a few particles, which might naïvely be considered 'highly ordered, and so low-entropy'. This appears to contradict the principle that entropy should increase; this apparent paradox is explained in the discussion section at the beginning of the paper.

We recall that a labelled Kac process is the Markov process of velocities $\left(v_{1}(t), \ldots, v_{N}(t)\right)$ corresponding to the particle dynamics. The state space is the set $\mathbb{S}^{N}=\left\{\left(v_{1}, \ldots, v_{N}\right) \in\left(\mathbb{R}^{d}\right)^{N}: \sum_{i=1}^{N} v_{i}=0, \sum_{i=1}^{N}\left|v_{i}\right|^{2}=N\right\}$, which we call the labelled Boltzmann Sphere. We denote $\theta_{N}$ the map taking $\left(v_{1}, \ldots, v_{N}\right)$ to its empirical measure in $\mathcal{S}_{N}$ :

$$
\begin{equation*}
\theta_{N}: \mathbb{S}^{N} \rightarrow \mathcal{S}_{N} ; \quad\left(v_{1}, \ldots, v_{N}\right) \mapsto \frac{1}{N} \sum_{i=1}^{N} \delta_{v_{i}} \tag{9.2}
\end{equation*}
$$

Moreover, if $\mathcal{V}_{t}^{N}$ is a labelled Kac process, then the empirical measures $\mu_{t}^{N}:=$ $\theta_{N}\left(\mathcal{V}_{t}^{N}\right)$ are a Kac process in the sense defined in the Introduction.

Considered as a $((N-1) d-1)$-dimensional sphere, $\mathbb{S}^{N}$ has a uniform (Hausdorff) distribution $\gamma^{N}$. We define the 'uniform distribution' $\sigma^{N}$ on $\mathcal{S}_{N}$ to be the push forward of $\gamma^{N}$ by $\theta_{N}$ :

$$
\begin{equation*}
\sigma^{N}(A):=\gamma^{N}\left\{\left(v_{1}, \ldots, v_{N}\right) \in \mathbb{S}^{d}: \theta_{N}\left(v_{1}, \ldots, v_{N}\right) \in A\right\} . \tag{9.3}
\end{equation*}
$$

We will use this definition to transfer the positivity of the measure $\gamma^{N}$ forward to $\sigma^{N}$.

As discussed in the literature review, the problem of relaxation to equilibrium for the Kac process is a subtle problem, and has been extensively studied. For our purposes, the following $L^{2}$ convergence is sufficient.

Proposition 8. Suppose that $\left(\mu_{t}^{N}\right)_{t \geq 0}$ is a hard-spheres Kac process, where the law of the initial data $\mathcal{L} \mu_{t}^{N}$ has a density $h_{0}^{N} \in L^{2}\left(\sigma^{N}\right)$ with respect to $\sigma^{N}$. Then at all positive times $t \geq 0$, the law $\mathcal{L} \mu_{t}^{N}$ has a density $h_{t}^{N} \in L^{2}\left(\sigma^{N}\right)$ with respect to $\sigma^{N}$, and for some universal constant $\lambda_{0}>0$, we have

$$
\begin{equation*}
\left\|h_{t}^{N}-1\right\|_{L^{2}\left(\sigma^{N}\right)} \leq e^{-\lambda_{0} t}\left\|h_{0}^{N}-1\right\|_{L^{2}\left(\sigma^{N}\right)} \tag{9.4}
\end{equation*}
$$

A version of this, for the labelled Kac process, appears as [29], Theorem 6.8 and corollary; the result stated above follows by a push forward argument. This is sufficient to prove the following weak ergodic theorem.

Lemma 9.1. Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a hard-spheres Kac process on $N$ particles, started from $\mu_{0}^{N} \sim \sigma^{N}$. Let $R_{N} \subset \mathcal{S}_{N}$ be such that $p=\sigma^{N}\left(R_{N}\right)>0$. Then

$$
\begin{equation*}
\frac{1}{t} \int_{0}^{t} 1\left(\mu_{s}^{N} \in R_{N}\right) d s \rightarrow p \tag{9.5}
\end{equation*}
$$

in $L^{2}$. In particular, almost surely, $\mu_{t}^{N}$ visits $R_{N}$ on an unbounded set of times.
Proof. Observe that

$$
\begin{equation*}
\mathbb{E}\left[\frac{1}{t} \int_{0}^{t} 1\left(\mu_{s}^{N} \in R_{N}\right) d s\right]=\frac{1}{t} \int_{0}^{t} \mathbb{P}\left(\mu_{s}^{N} \in R_{N}\right) d s=p \tag{9.6}
\end{equation*}
$$

so our claim reduces to bounding the variance.
For times $t \geq 0$, write $A(t)$ as the event $A(t)=\left\{\mu_{t}^{N} \in R_{N}\right\}$; we will compute the covariance of $1_{A\left(s_{1}\right)}$ and $1_{A\left(s_{2}\right)}$, for $0 \leq s_{1} \leq s_{2}$. Observe that

$$
\begin{equation*}
\mathbb{E}\left[1_{A\left(s_{1}\right)}\left(1_{A\left(s_{2}\right)}-p\right)\right]=p\left(\mathbb{P}\left(A\left(s_{2}\right) \mid A\left(s_{1}\right)\right)-p\right) \tag{9.7}
\end{equation*}
$$

Conditional on $A\left(s_{1}\right)$, the law of $\mu_{s_{1}}^{N}$ has a conditional density $h_{s_{1}}^{N} \propto 1_{R_{N}}$ with respect to $\sigma^{N}$. By Proposition 8 , conditional on $A\left(s_{1}\right), \mu_{s_{2}}^{N}$ has a density $h_{s_{2}}^{N}$, and we can bound

$$
\begin{align*}
\left|\mathbb{P}\left(A\left(s_{2}\right) \mid A\left(s_{1}\right)\right)-p\right| & \leq\left\|h_{s_{2}}^{N}-1\right\|_{L^{1}\left(\sigma^{N}\right)}  \tag{9.8}\\
& \leq\left\|h_{s_{2}}^{N}-1\right\|_{L^{2}\left(\sigma^{N}\right)} \leq C\left(R_{N}\right) e^{-\lambda_{0}\left(s_{2}-s_{1}\right)}
\end{align*}
$$

for some constant $C\left(R_{N}\right)$ independent of time. Hence

$$
\begin{align*}
\mathbb{E}\left[\left(1_{A\left(s_{1}\right)}-p\right)\left(1_{A\left(s_{2}\right)}-p\right)\right] & =p\left(\mathbb{P}\left(A\left(s_{2}\right) \mid A\left(s_{1}\right)\right)-p\right) \\
& \leq p C\left(R_{N}\right) e^{-\lambda_{0}\left(s_{2}-s_{1}\right)} \tag{9.9}
\end{align*}
$$

We can now integrate to bound the variance:

$$
\begin{align*}
\operatorname{Var}\left(\frac{1}{t} \int_{0}^{t} 1\left(\mu_{s}^{N} \in R_{N}\right) d s\right) & =\frac{2}{t^{2}} \int_{0}^{t} d s_{1} \int_{s_{1}}^{t} d s_{2} \mathbb{E}\left[\left(1_{A\left(s_{1}\right)}-p\right)\left(1_{A\left(s_{2}\right)}-p\right)\right] \\
& \leq \frac{2 p C}{t^{2}} \int_{0}^{t} d s_{1} \int_{s_{1}}^{\infty} d s_{2} e^{-\lambda_{0}\left(s_{2}-s_{1}\right)}  \tag{9.10}\\
& \leq \frac{2 p C}{\lambda_{0} t} \rightarrow 0 .
\end{align*}
$$

An immediate corollary is that the long-run deviation must be bounded below by the essential supremum of the deviation under the invariant measure.

Corollary 9.2. Let $\left(\mu_{t}^{N}\right)_{t \geq 0}$ be a $N$-particle Kac process in equilibrium. Then, almost surely,

$$
\begin{aligned}
\limsup _{t \rightarrow \infty} W\left(\mu_{t}^{N}, \gamma\right) & \geq\|W(\cdot, \gamma)\|_{L^{\infty}\left(\sigma^{N}\right)} \\
& =\underset{\sigma^{N}(d \mu)}{\operatorname{ess} \sup } W(\mu, \gamma) .
\end{aligned}
$$

Proof. For ease of notation, write $W^{*}$ as the essential supremum appearing on the right-hand side. For any $\epsilon>0$, let $R_{N, \epsilon}=\left\{\mu \in \mathcal{S}_{N}: W(\mu, \gamma)>W^{*}-\epsilon\right\}$; it is immediate that $\sigma^{N}\left(R_{N, \epsilon}\right)>0$. By the remark in Lemma 9.1, almost surely, $\mu_{t}^{N}$ visits $R_{N, \epsilon}$ on an unbounded set of times, and so

$$
\begin{equation*}
\limsup _{t \rightarrow \infty} W\left(\mu_{t}^{N}, \gamma\right) \geq W^{*}-\epsilon \tag{9.12}
\end{equation*}
$$

The conclusion now follows on taking an intersection over some sequence $\epsilon_{n} \downarrow 0$.

To prove Theorem 1.4, it now only remains to show a lower bound on the essential supremum.

Lemma 9.3. Let $f$ be given by

$$
\begin{equation*}
f(v)=\left(1+|v|^{2}\right) \min \left(\frac{|v|}{\sqrt{N / 2}}, 1\right) \tag{9.13}
\end{equation*}
$$

Then $f \in \mathcal{A}$, and

$$
\begin{equation*}
\|\langle f, \mu-\gamma\rangle\|_{L^{\infty}\left(\sigma^{N}\right)} \geq 1-\frac{C}{\sqrt{N}} \tag{9.14}
\end{equation*}
$$

for some constant $C=C(d)$. In particular, this is a lower bound for the essential supremum $W^{*}$, and so for the long-run deviation.

Proof. It is easy to see that $f \in \mathcal{A}$. Moreover, the region

$$
\begin{equation*}
\widetilde{R}_{N}=\left\{\left(v_{1}, \ldots, v_{N}\right) \in \mathbb{S}^{N}:\left\langle f, \theta_{N}\left(v_{1}, \ldots, v_{N}\right)\right\rangle>1\right\} \tag{9.15}
\end{equation*}
$$

is an open subset of $\mathbb{S}^{N}$, containing $\left(\sqrt{\frac{N}{2}} e_{1},-\sqrt{\frac{N}{2}} e_{1}, 0, \ldots, 0\right)$. By positivity of the uniform measure $\gamma^{N}$ on $\mathbb{S}^{N}$, it follows that $\gamma^{N}\left(\widetilde{R}_{N}\right)>0$. The corresponding region in $\mathcal{S}_{N}$ :

$$
\begin{equation*}
R_{N}=\left\{\mu^{N} \in \mathcal{S}_{N}:\left\langle f, \mu^{N}\right\rangle>1\right\} \supset \theta_{N}\left(\widetilde{R}_{N}\right) \tag{9.16}
\end{equation*}
$$

By definition (9.3) of $\sigma^{N}$, we have

$$
\begin{equation*}
\sigma^{N}\left(R_{N}\right) \geq \gamma^{N}\left(\widetilde{R}_{N}\right)>0 . \tag{9.17}
\end{equation*}
$$

For all $\mu^{N} \in R_{N}$, we have

$$
\begin{equation*}
W\left(\mu^{N}, \gamma\right) \geq\left\langle f, \mu^{N}-\gamma\right\rangle \geq 1-N^{-1 / 2}\left\langle\left(1+|v|^{2}\right)\right| v|, \gamma\rangle . \tag{9.18}
\end{equation*}
$$

Since $R_{N}$ has positive measure, taking $C=\left\langle\left(1+|v|^{2}\right)\right| v|, \gamma\rangle$, we can conclude that

$$
\begin{equation*}
W^{*} \geq 1-\frac{C}{\sqrt{N}} \tag{9.19}
\end{equation*}
$$

Proof of Theorem 1.4. From the previous two lemmas, we know that for all $N \geq 2$, and for $\sigma^{N}$ almost all $\mu^{N}$,

$$
\begin{equation*}
\mathbb{P}_{\mu^{N}}\left(\limsup _{t \rightarrow \infty} W\left(\mu_{t}^{N}, \gamma\right) \geq 1-\frac{C}{\sqrt{N}}\right)=1 \tag{9.20}
\end{equation*}
$$

where $\mathbb{P}_{\mu^{N}}$ denotes the law of a Kac process started at $\mu^{N}$.
Let $N \geq 2, k>2$ and $a>1$. The region $R_{\star, N}$ of the labelled sphere such that $\Lambda_{k}\left(\theta_{N}(\mathcal{V})\right)<a$ is an open set; to conclude that it has positive $\sigma^{N_{-}}$measure, it suffices to show that it is nonempty.

Let $r$ be a rotation by $\frac{2 \pi}{N}$ in the plane corresponding to the first two axes $\left(e_{1}, e_{2}\right)$. Then the data

$$
\begin{equation*}
\mathcal{V}_{\star}=\left(e_{1}, r e_{1}, \ldots, r^{N-1} e_{1}\right) \tag{9.21}
\end{equation*}
$$

belongs to $\mathbb{S}^{N}$, and has $\Lambda_{k}\left(\theta_{N}\left(\mathcal{V}_{\star}\right)\right)=\frac{1}{N} \sum_{i=1}^{N} 1^{s}=1$. Hence $\mathcal{V}_{\star} \in R_{\star, N}$ is open and nonempty, so $\gamma^{N}\left(R_{\star, N}\right)>0$. The positivity transfers to the corresponding region of $\mathcal{S}_{N}$ :

$$
\begin{equation*}
\sigma^{N}\left\{\mu^{N} \in \mathcal{S}_{N}: \Lambda_{k}\left(\mu^{N}\right)<a\right\}=\gamma^{N}\left(R_{N, \star}\right)>0 . \tag{9.22}
\end{equation*}
$$

Hence, for any $N \geq 2$, we can choose an initial datum $\mu_{0}^{N}=\mu^{N}$, with $\Lambda_{k}\left(\mu_{0}^{N}\right)<$ $a$, such that (9.20) holds. Observing that

$$
\begin{equation*}
W\left(\phi_{t}\left(\mu_{0}^{N}\right), \gamma\right) \leq\left\|\phi_{t}\left(\mu_{0}^{N}\right)-\phi_{t}(\gamma)\right\|_{\mathrm{TV}+2} \rightarrow 0 \tag{9.23}
\end{equation*}
$$

it follows that, $\mathbb{P}_{\mu^{N}}$ almost surely

$$
\begin{equation*}
\limsup _{t \rightarrow \infty} W\left(\mu_{t}^{N}, \gamma\right)=\limsup _{t \rightarrow \infty} W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right) \geq 1-\frac{C}{\sqrt{N}} \tag{9.24}
\end{equation*}
$$

## REMARK 9.4.

(i) The proof of Lemma 9.1 leaves open the possibility that there is a nonempty 'exceptional set' of initial data $\mu^{N}$ where (9.20) does not hold. A stronger assertion would be positive Harris recurrence, as defined in [20], which allows a similar ergodic theorem for any initial data $\mu^{N}$. This is not necessary for our purposes.
(ii) In principle, one could use this compute the typical time scales necessary for these deviations to occur, and sharper estimates may be obtained by using more detailed forms of relaxation, such as the entropic relaxation considered by [6]. This is not necessary for our arguments.
10. Proof of Theorem 1.7. Finally, we show that Theorems 1.1, 1.3 implies the claimed chaoticity estimates in Theorem 1.7. The following proof largely follows that of [29], Theorem 3.1, using the estimates derived in this paper. As remarked in the Introduction, the novelty is the use of the Hölder estimate (1.36) to control the term $\mathcal{T}_{3}$.

In the following proof, we will use estimates from Theorem 1.3, which allow us to minimise the moment conditions required on the initial data. Better results can be obtained using Theorem 1.1 at the cost of requiring a stronger moment estimate, although these still do not obtain optimal rates.

Proof of Theorem 1.7. Let $k>2$, and $\epsilon=\epsilon(d, k)>0$ be the resulting exponent from Theorem 1.3. Let $\mu_{0}^{N} \in \mathcal{S}_{N}$ satisfy $\Lambda_{k}\left(\mu_{0}^{N}\right) \leq a$.

Recall that we wish to estimate

$$
\begin{equation*}
\frac{\mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{P}_{t}^{N}\left(\mu_{0}^{N}, \cdot\right), \phi_{t}\left(\mu_{0}^{N}\right)^{\otimes l}\right)\right.}{l} \tag{10.1}
\end{equation*}
$$

uniformly in $t \geq 0$ and $l=1, \ldots, N$, and where $\mathcal{W}_{1, l}$ is the Wasserstein ${ }_{1}$ distance on laws, given by (1.21). Let $\mathcal{V}_{t}^{N}$ be a labelled Kac process, and let $\mu_{t}^{N}$ be the associated process of empirical measures. Fixing a test function $f \in B_{X}^{\otimes l}$, we break up the difference as

$$
\begin{align*}
& \int_{\left(\mathbb{R}^{d}\right)^{N}} f(V)\left(\Pi_{l}\left[\mathcal{P}_{t}^{N}\left(\mu_{0}^{N}, \cdot\right)\right]-\left(\phi_{t}\left(\mu_{0}^{N}\right)\right)^{\otimes l}\right)(d V) \\
& \quad=\mathbb{E}_{\mu_{0}^{N}}\left[\prod_{j=1}^{l} f_{j}\left(v_{j}(t)\right)\right]-\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle  \tag{10.2}\\
& \quad=\mathcal{T}_{1}+\mathcal{T}_{2}
\end{align*}
$$

where $\mathbb{E}_{\mu_{0}^{N}}$ denotes expectation under the law $\mathcal{P}_{t}^{N}\left(\mu_{0}^{N}, \cdot\right)$, and where the two error terms are

$$
\begin{align*}
& \mathcal{T}_{1}:=\mathbb{E}_{\mu_{0}^{N}}\left[\prod_{j=1}^{l} f_{j}\left(v_{j}(t)\right)-\prod_{j=1}^{l}\left\langle f_{j}, \mu_{t}^{N}\right\rangle\right]  \tag{10.3}\\
& \mathcal{T}_{2}:=\mathbb{E}_{\mu_{0}^{N}}\left[\prod_{j=1}^{l}\left\langle f_{j}, \mu_{t}^{N}\right\rangle-\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle\right] . \tag{10.4}
\end{align*}
$$

Now, $\mathcal{T}_{1}$ is a purely combinatorial term, based on the use of empirical measures, and $\mathcal{T}_{2}$ may be controlled using the pointwise estimates Theorems 1.1, 1.3. We will indicate how these terms may be controlled for the simple case $l=2$, and use this to show the full, 'infinite dimensional' chaos estimate claimed.

Step 1: Estimate on $\mathcal{T}_{1}$ Since the law $\mathcal{P}_{t}^{N}\left(\mu_{0}^{N}, \cdot\right)$ is symmetric, we may rewrite

$$
\begin{equation*}
\mathbb{E}_{\mu_{0}^{N}}\left[f_{1}\left(v_{1}(t)\right) f_{2}\left(v_{2}(t)\right)\right]=\mathbb{E}_{\mu_{0}^{N}}\left[\frac{1}{N(N-1)} \sum_{i \neq j} f_{1}\left(v_{i}(t)\right) f_{2}\left(v_{j}(t)\right)\right], \tag{10.5}
\end{equation*}
$$

where $N(N-1)$ counts the number of ordered pairs of indexes $(i, j)$. Similarly, the second term may be written

$$
\begin{equation*}
\mathbb{E}_{\mu_{0}^{N}}\left[\left\langle f_{1}, \mu_{t}^{N}\right\rangle\left\langle f_{2}, \mu_{t}^{N}\right\rangle\right]=\mathbb{E}_{\mu_{0}^{N}}\left[\left(\frac{1}{N} \sum_{i=1}^{N} f_{1}\left(v_{i}(t)\right)\right)\left(\frac{1}{N} \sum_{j=1}^{N} f_{1}\left(v_{j}(t)\right)\right)\right] . \tag{10.6}
\end{equation*}
$$

Comparing the two terms, and using the bound $\left\|f_{j}\right\|_{L^{\infty}} \leq\left\|f_{j}\right\|_{X} \leq 1$ for $j=1,2$, we obtain the estimate

$$
\begin{equation*}
\left|\mathcal{T}_{1}\right| \leq \sum_{i \neq j}\left|\frac{1}{N(N-1)}-\frac{1}{N^{2}}\right|+\sum_{i=1}^{N} \frac{1}{N^{2}} \tag{10.7}
\end{equation*}
$$

Therefore, we have the bound $\left|\mathcal{T}_{1}\right| \leq \frac{2}{N}$, uniformly in $f$ and $t$.
Step 2: Estimate on $\mathcal{T}_{2}$ For the case $l=2$, we break up the product as

$$
\begin{align*}
& \prod_{j=1}^{2}\left\langle f_{j}, \mu_{t}^{N}\right\rangle-\prod_{j=1}^{2}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle  \tag{10.8}\\
& \quad=\left\langle f_{1}, \mu_{t}^{N}-\phi_{t}\left(\mu_{0}^{N}\right)\right\rangle\left\langle f_{2}, \mu_{t}^{N}\right\rangle+\left\langle f_{1}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle\left\langle f_{2}, \mu_{t}^{N}-\phi_{t}\left(\mu_{0}^{N}\right)\right\rangle
\end{align*}
$$

In each case, the difference term is dominated by a multiple of the Wasserstein distance $W\left(\mu_{t}^{N}, \phi_{t}(\mu)\right)$, where $W$ is as in (1.15), and the remaining term is absolutely bounded, by the boundedness of $f_{j}, j=1,2$. Therefore, we estimate

$$
\begin{equation*}
\left|\prod_{j=1}^{l}\left\langle f_{j}, \mu_{t}^{N}\right\rangle-\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle\right| \lesssim W\left(\mu_{t}^{N}, \phi_{t}\left(\mu_{0}^{N}\right)\right) \tag{10.9}
\end{equation*}
$$

Now, the right-hand side is precisely the term controlled by Theorems 1.1, 1.3, in the special case $\mu_{0}=\mu_{0}^{N}$. By the choice of $\epsilon$ and $k$ above, we obtain the control

$$
\begin{equation*}
\mathcal{T}_{2} \lesssim \Lambda_{k}\left(\mu_{0}^{N}\right)^{\frac{1}{2}} N^{-\epsilon} \lesssim a N^{-\epsilon} \tag{10.10}
\end{equation*}
$$

for some explicit $\epsilon=\epsilon(d, k)>0$.
We also remark here that this implication, given Theorems 1.1, 1.3 is immediate. However, attempting to reverse this implication, and deduce a theorem similar to 1.1 from a control of $\mathcal{T}_{2}$, requires moving the supremum over test functions $f$ inside the expectation. This corresponds to the most technical step in our proof (Lemmas 6.1, 7.1). Therefore, while it may be possible to deduce a version Theorem 1.1 from the control of $\mathcal{T}_{2}$ given by [29], this would scarcely be less technical than the proof given, and would not lead to a proof of Theorem 1.2.

Step 3: Deduction of infinite-dimensional chaos Combining the two estimates for the case $l=2$ above, we deduce that there exists $\epsilon=\epsilon(d, k)>0$ such that

$$
\begin{equation*}
\sup _{t \geq 0} \mathcal{W}_{1,2}\left(\Pi_{2}\left[\mathcal{P}_{t}^{N}\left(\mu_{0}^{N}, \cdot\right)\right], \phi_{t}\left(\mu_{0}^{N}\right)^{\otimes l}\right) \lesssim a N^{-\epsilon} \tag{10.11}
\end{equation*}
$$

To deduce the full statement, we appeal to the following result from [19], which may also be found in [28], Theorem 2.1. For any probability measure $\mu$ on $\mathbb{R}^{d}$, and any symmetric distribution $\mathcal{L}^{N}$ on $\left(\mathbb{R}^{d}\right)^{N}$, we may estimate

$$
\begin{equation*}
\max _{l \leq N} \frac{\mathcal{W}_{1, l}\left(\Pi_{l}\left[\mathcal{L}^{N}\right], \mu^{\otimes l}\right)}{l} \leq C\left(\mathcal{W}_{1,2}\left(\Pi_{2}\left[\mathcal{L}^{N}\right], \mu^{\otimes 2}\right)^{\alpha_{1}}+N^{-\alpha_{2}}\right) \tag{10.12}
\end{equation*}
$$

for some explicit constants $C, \alpha_{1}, \alpha_{2}>0$ depending on the dimension $d$. The claimed result (1.38) now follows.

We now turn to the two consequences claimed as a result.
(i) Chaotic case. Let $\mu_{0} \in \mathcal{S}$ have an $k$ th moment $\Lambda_{k}\left(\mu_{0}\right) \leq a$, and construct $\mathcal{V}_{0}^{N}=\left(v_{1}(0), \ldots, v_{N}(0)\right)$ be as described in the statement of the theorem with associated empirical measure $\mu_{0}^{N}$. It is straightforward to show that this construction preserves moments up to a constant, that is, $\mathbb{E}\left(\Lambda_{k}\left(\mu_{0}^{N}\right)\right) \lesssim a$.

For a fixed test function $f \in B_{X}^{\otimes l}$, we return to the decomposition (10.2). For this case, where $\mu_{0}^{N} \neq \mu_{0}$, we have a third error term:

$$
\begin{equation*}
\int_{\left(\mathbb{R}^{d}\right)^{N}} f(V)\left(\Pi_{l}\left[\mathcal{L} \mathcal{V}_{t}^{N}\right]-\left(\phi_{t}\left(\mu_{0}\right)\right)^{\otimes l}\right)(d V)=\mathcal{T}_{1}+\mathcal{T}+\mathcal{T}_{3} \tag{10.13}
\end{equation*}
$$

Here, $\mathcal{T}_{1}$ and $\mathcal{T}_{2}$ are as above, replacing $\mathbb{E}_{\mu_{0}^{N}}$ by the full expectation $\mathbb{E}$, and $\mathcal{T}_{3}$ is an additional error term, from approximating $\mu_{0}$ by $\mu_{0}^{N}$ :

$$
\begin{equation*}
\mathcal{T}_{3}:=\mathbb{E}\left[\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle-\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}\right)\right\rangle\right] . \tag{10.14}
\end{equation*}
$$

As in the case above, we consider first the case $l=2$. The first two terms $\mathcal{T}_{1}, \mathcal{T}_{2}$ may be estimated as above, by conditioning on $\left(v_{1}(0), \ldots, v_{N}(0)\right)$ to conclude that

$$
\begin{equation*}
\mathcal{T}_{1}+\mathcal{T}_{2} \lesssim a N^{-\epsilon} \tag{10.15}
\end{equation*}
$$

for some $\epsilon>0$, uniformly in $f \in B_{X}^{\otimes l}$ and $t \geq 0$.
Arguing as in (10.9), we bound

$$
\begin{equation*}
\mathcal{T}_{3} \lesssim \mathbb{E} W\left(\phi_{t}\left(\mu_{0}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right) \tag{10.16}
\end{equation*}
$$

We estimate this term using the continuity estimate Theorem 1.6. Let $k^{\prime} \in(2, k)$, and let $\zeta>0$ be the resulting exponent using Theorem 1.6; by making $\zeta$ smaller if necessary, we assume that

$$
\begin{equation*}
\frac{\zeta k}{k-k^{\prime}} \leq 1 \tag{10.17}
\end{equation*}
$$

From Theorem 1.6, we have the estimate

$$
\begin{equation*}
\sup _{t \geq 0} W\left(\phi_{t}\left(\mu_{0}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right) \lesssim \Lambda_{k^{\prime}}\left(\mu_{0}^{N}, \mu_{0}\right) W\left(\mu_{0}^{N}, \mu_{0}\right)^{\zeta} \tag{10.18}
\end{equation*}
$$

and we use Hölder's inequality to obtain, uniformly in $t \geq 0$,

$$
\begin{align*}
\mathbb{E}\left[W\left(\phi_{t}\left(\mu_{0}^{N}\right) \phi_{t}\left(\mu_{0}\right)\right)\right] & \lesssim \mathbb{E}\left[\Lambda_{k}\left(\mu_{0}^{N}\right)\right]^{k^{\prime} / k} \mathbb{E}\left[W\left(\mu_{0}^{N}, \mu_{0}\right)^{\frac{\zeta k}{k-k^{\prime}}}\right]^{\frac{k-k^{\prime}}{k}}  \tag{10.19}\\
& \lesssim a^{k^{\prime} / k} \mathbb{E}\left[W\left(\mu_{0}^{N}, \mu_{0}\right)\right]^{\zeta}
\end{align*}
$$

From [33], Proposition 9.2, there is a constant $\beta=\beta(d, k)>0$ such that $\mathbb{E} W\left(\mu_{0}^{N}, \mu_{0}\right) \lesssim N^{-\beta}$, so we obtain

$$
\begin{equation*}
\mathbb{E}\left[W\left(\phi_{t}\left(\mu_{0}^{N}\right), \phi_{t}\left(\mu_{0}\right)\right)\right] \lesssim a N^{-\beta \zeta} \tag{10.20}
\end{equation*}
$$

Combining, and since all of our estimates are uniform in $f$ and $t$, we have shown that

$$
\begin{equation*}
\mathcal{W}_{1,2}\left(\Pi_{2}\left[\mathcal{L} \mathcal{V}_{t}^{N}\right], \phi_{t}\left(\mu_{0}\right)^{\otimes 2}\right) \lesssim a N^{-\alpha} \tag{10.21}
\end{equation*}
$$

for some $\alpha=\alpha(d, k)>0$. The improvement to infinite-dimensional chaos is exactly as above.
(ii) General case. The general case follows from the first case, by taking expectations over the initial data $\mu_{0}^{N}$. Indeed, for all $l \leq N$, all $f \in B_{X}^{\otimes l}$ and $t \geq 0$, and for any initial data $\left(v_{1}(0), \ldots, v_{N}(0)\right)$ with associated measure $\mu_{0}^{N}$, we have the bound

$$
\begin{equation*}
\frac{1}{l} \mathbb{E}_{\mu_{0}^{N}}\left[f_{1}\left(v_{1}(t)\right) \cdots f_{l}\left(v_{l}(t)\right)-\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle\right] \lesssim \Lambda_{k}\left(\mu_{0}^{N}\right) N^{-\epsilon} \tag{10.22}
\end{equation*}
$$

Taking expectation over the random initial data $\left(v_{1}(0), \ldots, v_{N}(0)\right)$ produces a full expectation on the left-hand side, and by definition of $\mathcal{L}_{t}^{l}$ in (1.44),

$$
\begin{equation*}
\mathbb{E}\left[\prod_{j=1}^{l}\left\langle f_{j}, \phi_{t}\left(\mu_{0}^{N}\right)\right\rangle\right]=\int_{\left(\mathbb{R}^{d}\right)^{l}} f(V) \mathcal{L}_{t}^{l}(d V) \tag{10.23}
\end{equation*}
$$

Optimising over $f \in B_{X}^{\otimes l}, l \leq N$ and $t \geq 0$ proves the claimed result.

## APPENDIX: CALCULUS OF MARTINGALES

We also review some basic facts and inequalities for martingales associated to the Kac process. All of these facts are true for general Markov chains; see [9].

Let $\mu_{t}^{N}$ be a Kac process, and write $m^{N}, \bar{m}^{N}$ for the jump measure and compensator defined in Section 5. Then, for any bounded and measurable $F^{N}$ : $[0, T] \times \mathcal{S}_{N} \rightarrow \mathbb{R}$, the process

$$
\begin{align*}
\mathcal{M}_{t}^{N} & =\int_{(0, t] \times \mathcal{S}_{N}}\left\{F_{s}^{N}\left(\mu^{N}\right)-F_{s}^{N}\left(\mu_{s-}^{N}\right)\right\}\left(m^{N}-\bar{m}^{N}\right)\left(d s, d \mu^{N}\right),  \tag{A.1}\\
0 & \leq t \leq T
\end{align*}
$$

is a martingale for the natural filtration $\left(\mathcal{F}_{t}^{N}\right)_{t \geq 0}$ of the process. We have the $L^{2}$ control

$$
\begin{equation*}
\left\|\mathcal{M}_{t}^{N}\right\|_{2}^{2}=\mathbb{E}\left\{\int_{(0, t] \times \mathcal{S}_{N}}\left\{F_{s}^{N}\left(\mu^{N}\right)-F_{s}^{N}\left(\mu_{s-}^{N}\right)\right\}^{2} \bar{m}^{N}\left(d s, d \mu^{N}\right)\right\} \tag{A.2}
\end{equation*}
$$

We will also use another special case of Itô's isometry for the measure $m^{N}-\bar{m}^{N}$ for a similar form of martingale. If $F^{N}$ is bounded and measurable on $[0, T] \times \mathcal{S}_{N}$, then for $t \leq T$,

$$
\begin{equation*}
\left\|\int_{0}^{t} F_{s}^{N}\left(\mu_{s-}^{N}\right)\left(m^{N}-\bar{m}^{N}\right)\left(d s, \mathcal{S}_{N}\right)\right\|_{2}^{2}=\mathbb{E}\left\{\int_{0}^{t} F_{s}^{N}\left(\mu_{s}^{N}\right)^{2} \bar{m}^{N}\left(d s, \mathcal{S}_{N}\right)\right\} \tag{A.3}
\end{equation*}
$$

For the local uniform case, Theorem 1.2, it will be necessary to control martingales of the form (A.1) in general $L^{p}$ spaces, rather than simply $L^{2}$. Since $\mathcal{M}^{N}$ of this form are finite variation martingales, the quadratic variation is given by

$$
\begin{equation*}
\left[\mathcal{M}^{N}\right]_{t}=\int_{(0, t] \times \mathcal{S}_{N}}\left\{F_{s}^{N}\left(\mu^{N}\right)-F_{s}^{N}\left(\mu_{s-}^{N}\right)\right\}^{2} m^{N}\left(d s, d \mu^{N}\right), \quad 0 \leq t \leq T \tag{A.4}
\end{equation*}
$$

Our analysis in $L^{p}$ is based on Burkholder's inequality for càdlàg martingales, which we state here for the class of martingales constructed above.

Lemma A.1. Suppose that $\left(\mathcal{M}_{t}^{N}\right)_{t=0}^{T}$ is the process given by (A.1), and let $p \geq 2$. Then there exists a constant $C=C(p)<\infty$ such that for all $t \leq T$, we have the $L^{p}$ control
(A.5) $\left\|\sup _{s \leq t}\left|\mathcal{M}_{s}^{N}\right|\right\|_{p}^{p} \leq C(p) \mathbb{E}\left[\left(\int_{0}^{t}\left\{F_{s}^{N}\left(\mu^{N}\right)-F_{s}^{N}\left(\mu_{s-}^{N}\right)\right\}^{2} m^{N}\left(d s, d \mu^{N}\right)\right)^{p / 2}\right]$.

Acknowledgements. I am grateful to my supervisor, James Norris, for the suggestion of this project and for several useful remarks which allowed me to strengthen the results, and to Clément Mouhot, for a useful conversation concerning the interpretation of our results. I would also like to express my gratitude to the two anonymous reviewers, whose suggestions and comments over the course of two iterations led to several substantial strengthening of the results.

## REFERENCES

[1] Bobylev, A. V. (1997). Moment inequalities for the Boltzmann equation and applications to spatially homogeneous problems. J. Stat. Phys. 88 1183-1214. MR1478067
[2] Bolley, F., Guillin, A. and Malrieu, F. (2010). Trend to equilibrium and particle approximation for a weakly selfconsistent Vlasov-Fokker-Planck equation. ESAIM Math. Model. Numer. Anal. 44 867-884. MR2731396
[3] Bolley, F., Guillin, A. and Villani, C. (2007). Quantitative concentration inequalities for empirical measures on non-compact spaces. Probab. Theory Related Fields 137 541-593. MR2280433
[4] BoltZmann, L. (1872). Weitere Studien über das Wärmegleichgewicht unter Gasmolekülen Vorgelegt in der Sitzung am 10. October 1872. K. und k. Hof-und Staatsdr.
[5] Carlen, E., Carvalho, M. C. and Loss, M. (2001). Many-body aspects of approach to equilibrium. In Séminaire: Équations aux Dérivées Partielles, 2000-2001. Sémin. Équ. Dériv. Partielles Exp. No. XIX, 12. École Polytech., Palaiseau. MR1860691
[6] Carlen, E. A., Carvalho, M. C., Le Roux, J., Loss, M. and Villani, C. (2010). Entropy and chaos in the Kac model. Kinet. Relat. Models 3 85-122. MR2580955
[7] Carlen, E. A., Carvalho, M. C. and Loss, M. (2003). Determination of the spectral gap for Kac's master equation and related stochastic evolution. Acta Math. 191 1-54. MR2020418
[8] Cortez, R. and Fontbona, J. (2018). Quantitative uniform propagation of chaos for Maxwell molecules. Comm. Math. Phys. 357 913-941. MR3769742
[9] Darling, R. W. R. and Norris, J. R. (2008). Differential equation approximations for Markov chains. Probab. Surv. 5 37-79. MR2395153
[10] DeSvillettes, L. (1993). Some applications of the method of moments for the homogeneous Boltzmann and Kac equations. Arch. Ration. Mech. Anal. 123 387-404. MR1233644
[11] Dobrušin, R. L. (1979). Vlasov equations. Funct. Anal. Appl. 13 115-123.
[12] Durmus, A., Eberle, A., Guillin, A. and Zimmer, R. (2018). An elementary approach to uniform In time propagation of chaos. Preprint. Available at arXiv:1805.11387.
[13] Elmroth, T. (1983). Global boundedness of moments of solutions of the Boltzmann equation for forces of infinite range. Arch. Ration. Mech. Anal. 82 1-12. MR0684411
[14] Fournier, N. and Guillin, A. (2015). On the rate of convergence in Wasserstein distance of the empirical measure. Probab. Theory Related Fields 162 707-738. MR3383341
[15] Fournier, N. and Méléard, S. (2001). A Markov process associated with a Boltzmann equation without cutoff and for non-Maxwell molecules. J. Stat. Phys. 104 359-385. MR1851392
[16] Gottlieb, A. D. (1998). Markov transitions and the propagation of chaos. Preprint. Available at arXiv:math/0001076.
[17] Graversen, S. E. and Peskir, G. (2000). Maximal inequalities for the Ornstein-Uhlenbeck process. Proc. Amer. Math. Soc. 128 3035-3041. MR1664394
[18] Grünbaum, F. A. (1971). Propagation of chaos for the Boltzmann equation. Arch. Ration. Mech. Anal. 42 323-345. MR0334788
[19] Hauray, M. and Mischler, S. (2014). On Kac's chaos and related problems. J. Funct. Anal. 266 6055-6157. MR3188710
[20] Hernández-Lerma, O. and Lasserre, J. B. (2001). Further criteria for positive Harris recurrence of Markov chains. Proc. Amer. Math. Soc. 129 1521-1524. MR1712909
[21] Holding, T. (2016). Propagation of chaos for Hölder continuous interaction kernels via Glivenko-Cantelli. Preprint. Available at arXiv:1608.02877.
[22] Janvresse, E. (2001). Spectral gap for Kac's model of Boltzmann equation. Ann. Probab. 29 288-304. MR1825150
[23] Kac, M. (1956). Foundations of kinetic theory. In Proceedings of the Third Berkeley Symposium on Mathematical Statistics and Probability, 1954-1955, Vol. III 171-197. Univ. California Press, Berkeley and Los Angeles. MR0084985
[24] Lu, X. and Mouhot, C. (2012). On measure solutions of the Boltzmann equation, part I: Moment production and stability estimates. J. Differential Equations 252 3305-3363. MR2871802
[25] Malrieu, F. (2003). Convergence to equilibrium for granular media equations and their Euler schemes. Ann. Appl. Probab. 13 540-560. MR1970276
[26] Maslen, D. K. (2003). The eigenvalues of Kac's master equation. Math. Z. 243 291-331. MR1961868
[27] McKean, H. P. Jr. (1967). An exponential formula for solving Boltmann's equation for a Maxwellian gas. J. Combin. Theory 2 358-382. MR0224348
[28] Mischler, S. (2012). Sur le programme de Kac concernant les limites de champ moyen. In Seminaire: Equations aux Dérivées Partielles, 2009-2010. Sémin. Équ. Dériv. Partielles Exp. No. XXXIII, 19. École Polytech., Palaiseau. MR3098637
[29] Mischler, S. and Mouhot, C. (2013). Kac's program in kinetic theory. Invent. Math. 193 1-147. MR3069113
[30] Mischler, S., Mouhot, C. and Wennberg, B. (2015). A new approach to quantitative propagation of chaos for drift, diffusion and jump processes. Probab. Theory Related Fields 161 1-59. MR3304746
[31] Mischler, S. and Wennberg, B. (1999). On the spatially homogeneous Boltzmann equation. Ann. Inst. H. Poincaré Anal. Non Linéaire 16 467-501. MR1697562
[32] Norris, J. (2016). Private communication.
[33] Norris, J. (2016). A consistency estimate for Kac's model of elastic collisions in a dilute gas. Ann. Appl. Probab. 26 1029-1081. MR3476632
[34] Povzner, A. J. (1962). On the Boltzmann equation in the kinetic theory of gases. Mat. Sb .58 65-86. MR0142362
[35] Rousset, M. (2014). A $N$-uniform quantitative Tanaka's theorem for the conservative Kac's $N$-particle system with Maxwell molecules. Preprint. Available at arXiv:1407.1965.
[36] Sznitman, A.-S. (1984). Équations de type de Boltzmann, spatialement homogènes. Z. Wahrsch. Verw. Gebiete 66 559-592. MR0753814
[37] Sznitman, A.-S. (1991). Topics in propagation of chaos. In École D'Été de Probabilités de Saint-Flour XIX—1989. Lecture Notes in Math. 1464 165-251. Springer, Berlin. MR1108185
[38] Talagrand, M. (1992). Matching random samples in many dimensions. Ann. Appl. Probab. 2 846-856. MR1189420
[39] Talagrand, M. (1994). The transportation cost from the uniform measure to the empirical measure in dimension $\geq 3$. Ann. Probab. 22 919-959. MR1288137
[40] TANAKA, H. (1978/79). Probabilistic treatment of the Boltzmann equation of Maxwellian molecules. Z. Wahrsch. Verw. Gebiete 46 67-105. MR0512334
[41] TANAKa, H. (2002). On the uniqueness of Markov process associated with the Boltzmann equation of Maxwellian molecules. In Stochastic Processes: Selected Papers of Hiroshi Tanaka 101-117.
[42] Wennberg, B. (1997). Entropy dissipation and moment production for the Boltzmann equation. J. Stat. Phys. 86 1053-1066. MR1450762
[43] Zermelo, E. (1896). Über einen Satz der Dynamik und die mechanische Wärmetheorie. Ann. Phys. 293 485-494.

Centre for Mathematical Sciences University of Cambridge
Wilberforce Road
Cambridge CB3 0WA
United Kingdom
E-MAIL: dh489@cam.ac.uk


[^0]:    Received January 2018; revised August 2018.
    ${ }^{1}$ Supported by the UK Engineering and Physical Sciences Research Council (EPSRC) Grant EP/L016516/1 for the University of Cambridge Centre for Doctoral Training, the Cambridge Centre for Analysis.

    MSC2010 subject classifications. Primary 60J25, 60K35; secondary 35Q20.
    Key words and phrases. Kac process, law of large numbers, Wasserstein distance, Boltzmann equation.

