# Kinetic Approach to Long time Behavior of Linearized Fast Diffusion Equations 

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#### Abstract

We show that the rate of convergence towards the self-similar solution of certain linearized versions of the fast diffusion equation can be related to the number of moments of the initial datum that are equal to the moments of the self-similar solution at a fixed time. As a consequence, we find an improved rate of convergence to self-similarity in terms of a Fourier based distance between two solutions. The results are based on the asymptotic equivalence of a collisional kinetic model of Boltzmann type with a linear Fokker-Planck equation with nonconstant coefficients, and make use of methods first applied to the reckoning of the rate of convergence towards equilibrium for the spatially homogeneous Boltzmann equation for Maxwell molecules.


KEY WORDS: fast diffusion equation, Maxwell models

## 1. INTRODUCTION

This paper concerns some aspects of the rate of convergence to equilibrium for solutions to the Cauchy problem of the fast diffusion equation posed on the whole space $\mathbb{R}^{N}$

$$
\begin{align*}
\frac{\partial v}{\partial \tau} & =\Delta v^{m}, \quad y \in \mathbb{R}^{N}, \quad \tau>0  \tag{1.1}\\
v(y, 0) & =v_{0}(y) \tag{1.2}
\end{align*}
$$

where $\frac{N}{N+2}<m<1$.

[^0]The long time asymptotics for Eq. (1.1) is described by the family of self similar source type Barenblatt-Pattle solutions

$$
\begin{equation*}
B_{C}(y, \tau)=\tau^{-\frac{1}{m+1}}\left(C+\frac{1-m}{2 m}|y|^{2} \tau^{-\frac{2}{m+1}}\right)^{\frac{1}{m-1}} \tag{1.3}
\end{equation*}
$$

where the constant $C$ must be chosen in order to match the initial mass. Recent results on the subject allow to assert that the rate of convergence is sensitive to the choice of the initial datum, and it is believed that the values of its moments play an important role.

As a matter of fact, in the linear case described by the heat equation the rate of convergence towards the fundamental solution is improved by fixing the center of mass of the initial data. Moreover, one can fix a certain number of initial moments in order to recover a higher order asymptotic approximation for the solutions with a faster rate of convergence. ${ }^{(16,18)}$ In the same linear case it has been shown that mass-centering also speeds up the entropy decay. ${ }^{(13)}$ These results suggest the possibility of detecting a more accurate asymptotic description even in the nonlinear case by fixing suitably a certain number of initial moments. ${ }^{(32)}$

This phenomenon was first established by J.L. Vazquez in Ref. 29 for porous medium equations. He was able to prove that, while the support of a general solution takes the shape of the support of its corresponding shifted Barenblatt profile for large times, a faster convergence rate in $L^{\infty}$ towards such profile could be obtained in case of solutions with radial symmetry. This convergence rate result is valid for the fast diffusion range as well.

In the fast diffusion case, the spatial translation invariance is related to an explicit eigenvalue of the linearization of the scaled equation in self-similar variables. ${ }^{(11,14)}$ Therefore, one expects an improvement of the rate of convergence for the fast diffusion equation by mass-centering to the rate $t^{-1}$ in $L^{1}$ in the translation dominated range, which is given by the next eigenvalue of the linearization. A nearly optimal convergence rate in the fast diffusion range, based on mass-centering, has been recently announced in Refs. 19 and 22.

For the porous medium case, mass-centering has been recently used in Ref. 6 to show that the Euclidean Wasserstein distance between two compactly supported solutions of the one-dimensional porous medium equation having the same center of mass decays to zero for large times with a computable rate. As a consequence, mass-centering allows to detect an improved rate of convergence of solutions of the one-dimensional porous medium equation towards well centered self-similar Barenblatt profiles.

Let us denote the total mass of $v_{0}$ by

$$
M_{0}=\int_{\mathbb{R}^{N}} v_{0}(y) d y>0
$$

and its center of mass by

$$
M_{1}=\frac{1}{M_{0}} \int_{\mathbb{R}^{N}} y v_{0}(y) d y .
$$

Both quantities $M_{0}$ and $M_{1}$ are invariant with respect to time (see Ref. 28). Hence, mass-centering means that we choose the initial datum with the same mass and first moment of the Barenblatt solution. Except in the linear case, however, higher order moments of the solution do not follow the same law of evolution of the higher order moments of a shifted in time Barenblatt profile, even if the same are equal at time $\tau=0$. As far as the second moment is concerned, a recent result establishes however an asymptotic equivalence between them. ${ }^{(26)}$

Taking into account this asymptotic equivalence, it can be reasonably conjectured that information on the rate of decay in terms of moments can be obtained by considering linear or linearized versions of Eq. (1.1). Among others, we shall consider in the sequel a linear version of the fast diffusion Eq. (1.1),

$$
\begin{align*}
\frac{\partial v(y, \tau)}{\partial \tau} & =\operatorname{div}\left(m B_{C}(y, \tau)^{m-1} \nabla v(y, \tau)\right), \quad y \in \mathbb{R}^{N}, \quad \tau>0,  \tag{1.4}\\
v(y, 0) & =v_{0}(y), \tag{1.5}
\end{align*}
$$

where $B_{C}(y, \tau)$ is a Barenblatt solution suitably shifted in time. The rate of convergence to equilibrium for Eq. (1.4) will be studied in terms of the number of moments initially equal to thats of the Barenblatt solution.

Using the spatial-temporal scaling given in Ref. 10, the initial value problem (1.4) can be rewritten as the initial value problem for the Fokker-Planck equation

$$
\begin{align*}
\frac{\partial f}{\partial t} & =\operatorname{div}\left[x f+m B_{C}(x)^{m-1} \nabla f\right] \quad x \in \mathbb{R}^{N}, \quad t>0,  \tag{1.6}\\
f(x, t=0) & =f_{0}(x) \geq 0 \tag{1.7}
\end{align*}
$$

where the initial value for the Fokker-Planck equation coincides with the initial value for the fast diffusion equation $\left(f_{0}=v_{0}\right)$. In Eq. (1.6) $B_{C}(x)$ is the Barenblatt solution (1.3) evaluated at time $\tau=1$,

$$
\begin{equation*}
B_{C}(x)=\left(C+\frac{1-m}{2 m}|x|^{2}\right)^{-\frac{1}{1-m}} \tag{1.8}
\end{equation*}
$$

Direct computations show that, for any fixed mass, the Fokker-Planck Eq. (1.6) has a unique stationary state given by (1.8).

A different linearization was first considered in Ref. 8. There, the large-time asymptotic of linearized very fast diffusion equations with and without potential confinements were studied, by reckoning estimates for the spectral gap and drawing conclusions on the time decay of the solution. The results in Ref. 8 hold for arbitrary algebraically large diffusion speeds, provided the solutions have the
mass-conservation property. Hence, the problem we will be dealing with, was not considered.

Likewise, this linearized fast diffusion equation was deeply investigated in Ref. 14 by Denzler and McCann, who were able to analyze its spectrum extracting sharp rates of asymptotic convergence to the Barenblatt profile. Our result recovers this rate of convergence (see Theorem 4.2 for more details). And we must point out, as in Ref. 14 the authors did, that in our case, although our kinetic analysis is rigorous the linearizations are formal calculations, obtaining in this way only conjectures for the nonlinear equation. For any given Barenblatt profile, characterized by the exponent $m$, higher moments of the solution stay uniformly bounded in time up to a critical exponent linked to the value of $m$. In addition to mass and momentum conservation, due to linearity, the higher moments of the solutions to Eq. (1.6) evolve in time in a closed form in terms of the lower order moments. This imply that the moments corresponding to two different initial data which are equal initially, remain equal to any subsequent time. Moreover, since the Barenblatt profile is a steady state to Eq. (1.6), if the initial datum has moments up some natural number $n>2$ equal to thats of the Barenblatt function (1.8), these moments remain constant in time. This property allows to compute precise rates of convergence to the stationary state in terms of the number of moments of the initial datum which are initially equal to thats of the equilibrium solution. The convergence rate, of the linearized equations, will be derived in terms of a Fourier based metric which has been proven very useful in the finding of rates of convergence towards equilibrium in kinetic theory of rarefied gases, both in the conservative case ${ }^{(17)}$ and in nonconservative one. ${ }^{(2,23)}$

The Fourier-based metrics $d_{s}$, for any $s>0$, are defined as

$$
\begin{equation*}
d_{s}(f, g)=\sup _{\xi \in \mathbb{R}^{N}} \frac{|\hat{f}(\xi)-\hat{g}(\xi)|}{|\xi|^{s}} \tag{1.9}
\end{equation*}
$$

for any pair of probability measures in $\mathcal{P}_{s}\left(\mathbb{R}^{N}\right)$, where $\mathcal{P}_{s}\left(\mathbb{R}^{N}\right)$ is the set of probability measures with bounded $s$-moment and as usual, $\hat{f}$ is the Fourier transform of the density $f(x)$,

$$
\hat{f}(\xi)=\int_{\mathbb{R}^{N}} f(x) \mathrm{e}^{-i x \cdot \xi} d x
$$

One shows that the distance is well-defined by simple Taylor expansion and finite for any pair of probability measures with equal moments up to order $[s]$, where [ $s$ ] denotes the integer part of $s$. Moreover, in case $s \geq 1$ be an integer, it suffices equality of moments up to order $s-1$ for being $d_{s}$ finite. In this way, the initial bounded distance in our framework means initial data with same moments of its corresponding steady state. In fact, $d_{s}$ with $s \geq 2$ topology is equivalent to the
weak-star topology for measures plus convergence of moments up to order [s], and can be related to the Wasserstein distance between probability measures. ${ }^{(27)}$

The study of the convergence in terms of this distance implies the analysis of the evolution of it. Therefore the first tentative could be study the derivative on time of the distance using the linearized equation in Fourier transform. But, at least for us, it looks like not easy since we must control the term with Laplacian, which comes from the diffusion part. Then, to overcome this problem we "invent" a kinetic model which will be our tool to obtain convergence on Fourier metrics for the linearizations equation of the fast diffusion equation. The key idea of our analysis will be the reckoning of rates of decay towards equilibrium for some kinetic equations of Boltzmann type, with a constant rate of collision. The results obtained for the kinetic model are subsequently shown to pass to the FokkerPlanck equation in a suitable asymptotic procedure reminiscent of the so-called grazing collision asymptotic. In kinetic theory of rarefied gases this asymptotic procedure became popular after the studies by C. Villani, ${ }^{(30,31)}$ who established a rigorous connection between the elastic Boltzmann equation ${ }^{(12)}$ and the Landau equation. ${ }^{(21)}$ This procedure, which corresponds to concentrate collisions on the grazing ones, namely collisions which leave velocities unaffected, allows to recover in the limit Fokker-Planck (or more generally Landau-Fokker-Planck) equations. Other applications to the one-dimensional dissipative Boltzmann equation can be found in Ref. 25. In the rest of the paper we will assume the initial datum of unit mass, zero center of mass and unit second moment

$$
\begin{equation*}
\int_{\mathbb{R}^{N}} f_{0} d x=1, \quad \int_{\mathbb{R}^{N}} x f_{0} d x=0, \quad \int_{\mathbb{R}^{N}}|x|^{2} f_{0} d x=1 \tag{1.10}
\end{equation*}
$$

Since the mass and momentum are conserved in time, the solution to Eq. (1.6), which is initially a probability density of zero mean, so remains at any subsequent time $t>0$. Our main results deal with the rate of decay to zero of solutions to different linearizations of the fast diffusion Eq. (1.1) in terms of the $d_{2+\delta}$ distance, where $\delta=\delta(m)$ is given in terms of the exponent $m$ of the fast diffusion. If Eq. (1.6) is considered, the typical result in one dimension reads

Theorem 1.1. Let $f(v, t)$ be the solution to Eq. (1.6), with initial datum $f_{0}(v)$. Suppose in addition that

$$
\begin{equation*}
d_{2+\delta}\left(f_{0}, B_{C}\right) \leq M<+\infty, \quad \delta<(3 m-1) /(1-m) \tag{1.11}
\end{equation*}
$$

Then, the solution $f(v, t)$ converges exponentially to $B_{C}(x)$ in $d_{2+\delta-m e t r i c, ~ a n d ~}^{\text {a }}$

$$
\begin{equation*}
d_{2+\delta}\left(f(t), B_{C}\right) \leq d_{2+\delta}\left(f_{0}, B_{C}\right) \exp \left\{\frac{2+\delta}{2}\left(\frac{1-m}{m} \delta-\frac{3 m-1}{m}\right) t\right\} \tag{1.12}
\end{equation*}
$$

Moreover, provided $3 / 5<m<1$ the rate of convergence in (1.12) increases in the interval $0<\delta \leq(5 m-3) /(2-2 m)$.

A careful reading of Theorem 1.1 clarifies the role of moments in the rate of convergence to equilibrium for solutions of Eq. (1.6). In fact, by definition, the $d_{s}(f, g)$ distance is finite as soon as the moments of $f$ and $g$ are equal up to the entire part $[s]$ of $s$. Hence, while $\delta=0$ in Theorem 1.1 only implies that there is mass-centering of the initial data, $\delta_{0}=(5 m-3) /(2-2 m) \leq 1$ implies than if at least the second moments of the initial data are equal, convergence to equilibrium is faster than before. Note that $\delta_{0}$ increases as soon as $m$ increases, and the convergence result is better if we are close to the linear case of the heat equation.

Instead of working on the Fokker-Planck equation directly, we will introduce a nonlinear kinetic model of Maxwell type ${ }^{(4)}$ for which the recovering of the rate of decay in terms of the Fourier based distance is immediate. Then, the result for the Fokker-Planck equation will follow by a well-established asymptotic analysis recently introduced for analogous nonconservative Boltzmann equation. ${ }^{(23)}$

## 2. PRELIMINARY RESULTS

The main object of our analysis will be the study of the large-time asymptotic of the solution to the Fokker-Planck equation

$$
\begin{align*}
\frac{\partial f}{\partial t} & =\operatorname{div}\left[x f+\nabla\left(\left(\alpha^{2}+\beta^{2}|x|^{2}\right) f(x)\right)\right], \quad\left(x \in \mathbb{R}^{N}, t>0\right),  \tag{2.1}\\
f(x, t=0) & =f_{0}(x) \geq 0 . \quad\left(x \in \mathbb{R}^{N}\right) \tag{2.2}
\end{align*}
$$

In Eq. (2.1) the constants $\alpha$ and $\beta$ are such that

$$
\alpha^{2}+\beta^{2}=\frac{1}{N}
$$

This condition implies that, if the initial datum $f_{0}(x)$ satisfies conditions (1.10), so does the solution at any subsequent time $t \geq 0$. It can be easily verified that Eq. (2.1) has a steady state with overpopulated tails, given by

$$
\begin{equation*}
B_{\infty}(x)=C\left(\alpha^{2}+\beta^{2}|x|^{2}\right)^{-1-\left(2 \beta^{2}\right)^{-1}} \tag{2.3}
\end{equation*}
$$

The constant $C$ in (2.3) has to be chosen to fit unit mass. In the one-dimensional case, Eq. (2.1), with $\alpha^{2}=\beta^{2}=1 / 2$ has been recently considered in Ref. 23 as the grazing limit of the one-dimensional dissipative Boltzmann equation studied in Ref. 1.

### 2.1. Fokker-Planck Equations as Linearization of Fast Diffusion Equations

The linear Fokker-Planck Eq. (2.1) is strictly connected to the linearization of the fast diffusion Eq. (1.1). Consider in fact Eq. (1.1), with initial datum (1.2). It
is well-known that (1.1)-(1.2) can be transformed into the fast diffusion equation with harmonic convection

$$
\begin{align*}
\frac{\partial u}{\partial t} & =\operatorname{div}\left(x u+\nabla u^{m}\right), \quad\left(x \in \mathbb{R}^{N}, t>0\right)  \tag{2.4}\\
u(x, t=0) & =u_{0}(x) \geq 0, \quad\left(x \in \mathbb{R}^{N}\right) \tag{2.5}
\end{align*}
$$

The transformation follows by the spatial-temporal scaling

$$
\begin{align*}
x & =\frac{y}{R(\tau)}, \quad R(\tau)=((2-N+N m) \tau+1)^{1 /(2-N+N m)}  \tag{2.6}\\
t(\tau) & =\frac{1}{2-N+N m} \ln ((2-N+N m) \tau+1)  \tag{2.7}\\
v(y, \tau) & =R(\tau)^{-N} u\left(\frac{y}{R(\tau)}, t(\tau)\right) \tag{2.8}
\end{align*}
$$

Applying (2.6)-(2.8) the similarity solution of the diffusion Eq. (1.1), given by (1.3) yields a equilibrium solution $B_{C}(x)$ of the Fokker-Planck equation, given by (1.8).

It is a simple exercise to show that the scaling (2.6)-(2.7) with

$$
\begin{equation*}
v(y, \tau)=R(\tau)^{-N} f\left(\frac{y}{R(\tau)}, t(\tau)\right) \tag{2.9}
\end{equation*}
$$

transforms the linear diffusion Eq. (1.4) into (1.6) and $v_{0}(x)=f_{0}(y)$.
The linear Fokker-Planck Eq. (1.6) is obtained from (2.4), but other possible linearizations can be done from (2.4) as we will show now.

The linear Fokker-Planck Eq. (1.6) is a linearization of (2.4).- We recall (2.4)

$$
\frac{\partial u}{\partial t}=\operatorname{div}\left(x u+\nabla u^{m}\right), \quad\left(x \in \mathbb{R}^{N}, t>0\right)
$$

and consider a state $u(x, t) \geq 0$ close to the equilibrium $B_{C}(x)$, i.e., $u-B_{C}=$ $O(\epsilon)$ in an appropriate topology. Thus we approximate the nonlinear term as

$$
\nabla u^{m}=m u^{m-1} \nabla f \cong m B_{C}^{m-1} \nabla u
$$

Identifying $u$ with $f$, we just obtain (1.6)

$$
\frac{\partial f}{\partial t}=\operatorname{div}\left[x f+m B_{C}(x)^{m-1} \nabla f\right] \quad x \in \mathbb{R}^{N}, t>0
$$

Since

$$
\begin{equation*}
B_{C}(x)^{m-1}=C+\frac{1-m}{2 m}|x|^{2}, \tag{2.10}
\end{equation*}
$$

one concludes with

$$
m \nabla\left[B_{C}(x)^{m-1} f\right]=m B_{C}(x)^{m-1} \nabla f+(1-m) x f
$$

Hence (1.6) can be rewritten as

$$
\begin{equation*}
\frac{\partial f}{\partial t}=\operatorname{div}\left[m x f+m \nabla\left(B_{C}(x)^{m-1} f\right)\right] \tag{2.11}
\end{equation*}
$$

Other possible linearization of (2.4).- Let us set

$$
\begin{equation*}
f(x, t)=\frac{1}{1+\epsilon}\left(B_{C}(x)+\epsilon \rho(x, t)\right), \tag{2.12}
\end{equation*}
$$

for $\epsilon>0$ small and " $\rho=O(1)$ ". Assuming mass-conservation

$$
\int_{\mathbb{R}^{N}} B_{C}(x) d x=\int_{\mathbb{R}^{N}} f(x, t) d x=M, t \geq 0
$$

we have

$$
\begin{equation*}
\int_{\mathbb{R}^{N}} \rho(x, t) d x=M, t \geq 0 . \tag{2.13}
\end{equation*}
$$

Inserting (2.12) into (2.4), computing the derivative with respect to the parameter $\epsilon$ and setting $\epsilon=0$ gives the linearization:

$$
\begin{align*}
\frac{\partial \rho}{\partial t} & =\operatorname{div}\left[x \rho+m \nabla\left(B_{C}(x)^{m-1} \rho\right)\right] \quad x \in \mathbb{R}^{N}, t>0,  \tag{2.14}\\
\rho(x, t=0) & =\rho_{0}(x) \geq 0 . \tag{2.15}
\end{align*}
$$

with

$$
\begin{equation*}
\int_{\mathbb{R}^{N}} \rho_{0} d x=M \tag{2.16}
\end{equation*}
$$

Note that both linearizations (2.14) and (2.11) have the same diffusion term, and differ only for the presence of the constant in front of the convection term. This difference reflects into the stationary solution, which coincides with the equilibrium Barenblatt solution (1.8) in Eq. (2.11), while for (2.14) is given by

$$
\rho_{\infty}(x)=D B_{C}^{2-m}(x),
$$

where the constant $D$ is fixed in terms of the initial mass. Due to their similar form, both Eqs. (2.14) and (2.11) can be studied by a unified treatment.

Both linearizations of (2.4) can be written as the Fokker-Planck Eq. (2.1).- By a simple scaling argument, one can finally show that both Eqs. (2.14) and (2.11) fall into the form of Eq. (2.1). Hence we will focus in the large-time behavior of solutions to this last equation, and in consequence we obtain the long-time behavior of both linearizations of (2.1).

First of all, we consider Eq. (2.14). Using equality (2.10)

$$
\begin{equation*}
m B_{C}(x)^{m-1}=m C+\frac{1-m}{2}|x|^{2} . \tag{2.17}
\end{equation*}
$$

Note that, provided $m>(N-2) / N$, the coefficient of $|x|^{2}$ is less than $1 / N$. If $\rho(x, t)$ solves (2.14), let us set $\rho(x, t)=\delta^{N} g(\delta x, t)$, where $\delta$ is a constant to be chosen. Then if $y=\delta x, g(y, t)$ solves

$$
\begin{align*}
\frac{\partial g}{\partial t} & =\operatorname{div}\left[y g+\nabla\left(\left(\delta^{2} m C+\frac{1-m}{2}|y|^{2}\right) g(y)\right)\right], \quad\left(y \in \mathbb{R}^{N}, t>0\right)  \tag{2.18}\\
g(y, t & =0)=\delta^{-N} f_{0}\left(\delta^{-1} y\right) \geq 0 \quad\left(y \in \mathbb{R}^{N}\right) \tag{2.19}
\end{align*}
$$

Therefore, $g(y, t)$ satisfies Eq. (2.1) for a convenient $\delta$ in such a way that $\alpha^{2}+\beta^{2}=$ $\frac{1}{N}$, where $\alpha^{2}=\delta^{2} m C$ and $\beta^{2}=\frac{1-m}{2}$. By a simple computation we show that $\delta$ has to be chosen in the following way

$$
\begin{equation*}
0<\delta^{2}=\frac{1}{C m}\left(\frac{1}{N}-\frac{1-m}{2}\right) \tag{2.20}
\end{equation*}
$$

The same strategy applies to Eq. (2.11). The only difference is that here we have to scale only the time. Let us set $f(x, t)=h(x, m t)$. Then, if $\tau=m t, h(x, \tau)$ solves

$$
\begin{gather*}
\frac{\partial h}{\partial \tau}=\operatorname{div}\left[x h+\nabla\left(\left(C+\frac{1-m}{2 m}|x|^{2}\right) h(x)\right)\right], \quad\left(x \in \mathbb{R}^{N}, \tau>0\right)  \tag{2.21}\\
h(x, \tau=0)=f_{0}(x) \geq 0 . \quad\left(x \in \mathbb{R}^{N}\right) \tag{2.22}
\end{gather*}
$$

In this case,

$$
\begin{equation*}
\alpha^{2}=\frac{1}{N}-\frac{1-m}{2 m} ; \quad \beta^{2}=\frac{1-m}{2 m} \tag{2.23}
\end{equation*}
$$

Note that the bound $\beta^{2}<1 / N$ implies $m>N /(N+2)$. By resorting to the definition of the Fourier based distance (1.9), any result on the large-time behavior of the distance between the solution to the Fokker-Planck Eq. (2.1) and its stationary state (2.3) can be easily translated into the corresponding result on the large-time behavior of the distance between the solution to the Fokker-Planck Eq. (2.14) or (2.11) and their stationary Barenblatt solutions. In fact, the relation $\rho(x, t)=\delta^{N} g(\delta x, t)$ implies $\hat{g}(\xi)=\hat{\rho}(\delta \xi)$, and

$$
\begin{equation*}
d_{s}\left(g, B_{\infty}\right)=\sup _{\xi \in \mathbb{R}^{N}} \frac{\left|\hat{\rho}(\delta \xi)-\hat{B_{\infty}}(\delta \xi)\right|}{|\xi|^{s}}=\frac{\left|\hat{\rho}(\xi)-\hat{\rho_{\infty}}(\xi)\right|}{|\xi|^{s}} \delta^{s} \tag{2.24}
\end{equation*}
$$

where $\delta$ is given by (2.20). Analogous result holds for Eq. (2.11).
Remark 2.1. The previous computation emphasize a noticeable difference between the linearized fast diffusion Eq. (2.14) and the linear one given by Eq. (2.11). In the former case, the condition $\beta^{2}<1 / N$ is satisfied if

$$
\begin{equation*}
m>\frac{N-2}{N} \tag{2.25}
\end{equation*}
$$

while, in the latter case, the same condition on $\beta$ holds if

$$
\begin{equation*}
m>\frac{N}{N+2} . \tag{2.26}
\end{equation*}
$$

Thus, finite mass for the Barenblatt solution is enough to study the linearized Eq. (2.14), while finite second moment for the Barenblatt profile is needed to study the linear Eq. (2.11).

### 2.2. On Moments of the Barenblatt Solution

The previous analysis shows that the Fokker-Planck Eq. (2.1) contains various linear or linearized versions of scaled fast diffusion equations. This fact depends clearly on the structure of the steady Barenblatt solution. For the sake of completeness, we shall recall in the following some of these properties. First, we recover the relationship between the zero order moment (mass) of the Barenblatt solution and its second one (the temperature) (see Ref. 26). Let $u(x, t)$ be a solution to Eq. (2.4). For any given $r \in \mathbb{N} \cup\{0\}$, let us set

$$
M_{r}(t)=\int_{\mathbb{R}^{N}}|x|^{r} u(x, t) d x .
$$

Direct computations show that

$$
\begin{equation*}
\frac{d}{d t} M_{2}(t)=-2 M_{2}(t)+2 N \int_{\mathbb{R}^{N}} u^{m} d x \tag{2.27}
\end{equation*}
$$

Since the Barenblatt is a steady solution to Eq. (2.27), its second moment $M_{2}$ satisfies

$$
M_{2}+2 N \int_{\mathbb{R}^{N}} B_{C}^{m} d x
$$

Using the expression of the Barenblatt solution (1.8) we find

$$
\int_{\mathbb{R}^{N}} B_{C}^{m} d x=\int_{\mathbb{R}^{N}}\left(C+\frac{1-m}{2 m}|x|^{2}\right) B_{C} d x=C M_{0}+\frac{1-m}{2 m} M_{2} .
$$

Hence

$$
M_{2}\left(\frac{N(1-m)}{m}-2\right)+2 N C M_{0}=0
$$

and therefore

$$
\begin{equation*}
C M_{0}=\left(\frac{1}{N}-\frac{1-m}{2 m}\right) M_{2} . \tag{2.28}
\end{equation*}
$$

If we fix the values of the moments, $M_{0}=M_{2}=1$, we then obtain for $C$ the value

$$
C=\frac{1}{N}-\frac{1-m}{2 m}
$$

found in (2.23).

Remark 2.2. Since both $M_{0}$ and $M_{2}$ are nonnegative, formula (2.28) implies the restriction

$$
\frac{1}{N}-\frac{1-m}{2 m}>0
$$

Consequently the second moment is finite if and only if $m>\frac{N}{N+2}$.
The precise values of higher order (and bounded) moments of the Barenblatt solution can be evaluated recursively starting from the equality

$$
\begin{equation*}
\frac{d}{d t} M_{n}(t)=-n M_{n}(t)+n(n-2+N) \int_{\mathbb{R}^{N}}|x|^{n-2} u^{m} d x \tag{2.29}
\end{equation*}
$$

where $2<n \in \mathbb{N}$. As before, we use the fact the Barenblatt is a steady solution to Eq. (2.29), which implies

$$
M_{n}=(n-2+N) \int_{\mathbb{R}^{N}} B_{C}^{m} d x
$$

We find

$$
\int_{\mathbb{R}^{N}}|x|^{n-2} B_{C}^{m} d x=\int_{\mathbb{R}^{N}}|x|^{n-2}\left(C+\frac{1-m}{2 m}|x|^{2}\right) B_{C} d x=C M_{n-2}+\frac{1-m}{2 m} M_{n}
$$

Finally, using the expression of $C$ found in (2.23),

$$
\left(\frac{1}{n-2+N}-\frac{1-m}{2 m}\right) M_{n}=\left(\frac{1}{N}-\frac{1-m}{2 m}\right) M_{n-2} .
$$

This relation also gives the maximum moment of the Barenblatt solution which is bounded. In fact, the procedure can be iterated up to the value of $n$ for which the coefficient of $M_{n}$ is greater than zero. A further interesting property of Barenblatt solutions is related to the scaling of variables. We used this scaling in the previous section to normalize the constants $\alpha$ and $\beta$. We have

Lemma 2.3. (Relation between Barenblatt solutions) Let B a Barenblatt solution with constant $C$ and mass and energy $M_{0}$ and $M_{2}$ respectively. Then, if $b^{m-1} a^{2}=1$, the function $\bar{B}(x)=b B(a x)$ is a Barenblatt solution with

$$
\begin{equation*}
\bar{M}_{0}=\frac{b}{a^{N}} M_{0}, \bar{M}_{2}=\frac{b}{a^{N+2}} M_{2} \quad \text { and } \quad \bar{C}=b^{m-1} C . \tag{2.30}
\end{equation*}
$$

Proof: Expressions for mass and energy are direct computations. We will focus only in proving that $\bar{B}$ is really a Barenblatt solution. We have

$$
\bar{B}(x)^{m-1}=b^{m-1}\left(C+\frac{1-m}{2 m} a^{2}|x|^{2}\right)=b^{m-1} C+b^{m-1} \frac{1-m}{2 m} a^{2}|x|^{2} .
$$

Then, since $b^{m-1} a^{2}=1$ the previous equality can be written as

$$
\bar{B}(x)=\left(b^{m-1} C+\frac{1-m}{2 m}|x|^{2}\right)^{\frac{-1}{1-m}} .
$$

### 2.3. Remark on the Heat Equation

In order to understand the role of the Fourier based distance in the reckoning of the large-time behavior of the linear Fokker-Planck Eq. (2.1) we recall briefly how this distance can be used in the case of the linear heat equation. The result that follows is contained in Ref. 18, where it appears as a simple example to justify application of the distance $d_{s}$ to the rate of convergence towards the Gaussian density in the central limit theorem. Another application of the distance (1.9) to the classical Fokker-Planck equation can be found in Ref. 9. These results follow in consequence of the fact that in these cases there exists an explicit solution, which allows for exact computations. Due to its importance in applications to the linear Fokker-Planck equations, we briefly describe how this distance works.

It is well known that a solution of the heat equation

$$
\begin{align*}
\frac{\partial u}{\partial t} & =\Delta u, \quad x \in \mathbb{R}^{N},  \tag{2.31}\\
u(x, 0) & =u_{0}(x), \tag{2.32}
\end{align*}
$$

where $u_{0}(x)$ is a probability density function satisfying conditions (1.10), behaves asymptotically in time as the fundamental solution of (2.31), which is given by the Gaussian density

$$
\begin{equation*}
\omega_{2 t}(x)=\frac{1}{\sqrt{4 \pi t}} \exp \left\{-\frac{|x|^{2}}{4 t}\right\} . \tag{2.33}
\end{equation*}
$$

The rate of convergence, that in $L^{2}$ norm is governed by $t^{-N / 2}$, can be improved as soon as we have more information on the moments of the initial data. A simple proof of this result follows by using the Fourier based distance (1.9). More precisely, let $u_{1}$ and $u_{2}$ be solutions of (2.31), corresponding to initial data $u_{0,1}$ and $u_{0,2}$ respectively, and let $s>0$ be such that $d_{s}\left(u_{0,1}, u_{0,2}\right)$ is finite. By Fourier transform, we get for $i=1,2$

$$
\partial_{t} \widehat{u_{i}}(t, \xi)=-|\xi|^{2} \widehat{u_{i}}(t, \xi) .
$$

Solving the ordinary differential equation leads to

$$
\left|\widehat{u_{1}}(t, \xi)-\widehat{u_{2}}(t, \xi)\right|^{2}=\exp \left(-2 t|\xi|^{2}\right)\left|\widehat{u_{0,1}}(\xi)-\widehat{u_{0,2}}(\xi)\right|^{2} .
$$

It follows that the $L^{2}$ norm is estimated by using the Parseval equality

$$
\begin{aligned}
\int_{\mathbb{R}^{N}}\left|u_{1}(t, x)-u_{2}(t, x)\right|^{2} d x & =(2 \pi)^{-N} \int_{\mathbb{R}^{N}}\left|\widehat{u_{1}}(t, \xi)-\widehat{u_{2}}(t, \xi)\right|^{2} d \xi \\
& =(2 \pi)^{-N} \int_{\mathbb{R}^{N}} \exp \left(-2 t|\xi|^{2}\right)\left|\widehat{u_{0,1}}(\xi)-\widehat{u_{0,2}}(\xi)\right|^{2} d \xi \\
& \leq d_{s}\left(u_{0,1}, u_{0,2}\right)^{2} F(t)
\end{aligned}
$$

The function of time $F(t)$, that is equal to

$$
F(t)=(2 \pi)^{-N} \int_{\mathbb{R}^{N}} \exp \left(-2 t|\xi|^{2}\right)|\xi|^{2 s} d \xi
$$

is bounded and can be computed explicitly to give (2.34). In fact, if

$$
\begin{aligned}
C & =(2 \pi)^{-N} \int_{\mathbb{R}^{N}} \exp \left(-2|\xi|^{2}\right)|\xi|^{2 s} d \xi, \\
F(t) & =C t^{-(s+N / 2)}
\end{aligned}
$$

Therefore, it follows that

$$
\begin{equation*}
\left\|\left(u_{1}-u_{2}\right)(t, \cdot)\right\|_{L^{2}\left(\mathbb{R}^{N}\right)}^{2} \leq d_{s}\left(u_{0,1}, u_{0,2}\right)^{2} C t^{-(s+N / 2)} \tag{2.34}
\end{equation*}
$$

where the constant $C$ is explicitly computable. ${ }^{(18)}$ In particular, as $s \rightarrow 0$, we recover from (2.34) the usual $t^{-N / 2}$ decay rate. The decay rate in (2.34) can be rephrased in terms of moments. Let $s=m+\delta$, where $m$ is an integer and $0 \leq \delta<1$. Since the Gaussian density (2.33) is a solution to the heat equation, looking at the definition of the metric $d_{s}$, if the initial datum $u_{0}(x)$ has the same moments of the fundamental solution $\omega_{2}(x)$ up to $p=m$, the $d_{s}$-distance $d_{s}\left(u_{0}, \omega_{2}\right)$ is bounded, and the rate of convergence becomes $t^{-(s+N / 2)}$.

As the previous example shows, one of the interesting features of the distance (1.9) is that it can be used through interpolation to obtain convergence in stronger spaces. This property has been discovered in Ref. 5, and subsequently used in various applications to the dissipative Boltzmann equation. ${ }^{(2)}$

### 2.4. A Property of the Fourier-Based Distance

Various properties of the distance (1.9) were collected in papers, ${ }^{(17,27)}$ where the interested reader can achieve information on the relationship of this norm to other more familiar equivalent norms used both in probability theory and in mass transportation. For the purposes of this paper, however, further properties are needed. We prove the following

Lemma 2.4. Let $\left\{f_{n}(x)\right\}_{n \geq 0}$ and $\left\{g_{n}(x)\right\}_{n \geq 0}, x \in \mathbb{R}^{N}$, be two sequences satisfying conditions (1.10) and such that $f_{n} \rightharpoonup f, g_{n} \rightharpoonup g$. Suppose in addition that, for some $r>2$

$$
\begin{equation*}
\int_{\mathbb{R}^{n}}|x|^{r} f_{n}(x)<+\infty, \quad \int_{\mathbb{R}^{n}}|x|^{r} g_{n}(x)<+\infty \tag{2.35}
\end{equation*}
$$

and, in addition

$$
d_{r}\left(f_{n}, g_{n}\right)<+\infty .
$$

Then, for all $s<r$,

$$
\begin{equation*}
d_{s}(f, g) \leq \liminf d_{s}\left(f_{n}, g_{n}\right) \tag{2.36}
\end{equation*}
$$

Proof: By definition, $d_{r}\left(f_{n}, g_{n}\right)$ is bounded if all moments of order less or equal than $[r]$ of $f_{n}$ and $g_{n}$ are equal. Since $f_{n} \rightharpoonup f$ and $g_{n} \rightharpoonup g$, condition (2.35) implies that, for any $s<r$

$$
\begin{equation*}
\int_{\mathbb{R}^{n}}|x|^{s} f_{n}(x) \rightarrow \int_{\mathbb{R}^{n}}|x|^{s} f(x), \quad \int_{\mathbb{R}^{n}}|x|^{s} g_{n}(x) \rightarrow \int_{\mathbb{R}^{n}}|x|^{s} g(x) . \tag{2.37}
\end{equation*}
$$

Thus, $d_{r}(f, g)=D$ is bounded. Now, consider that

$$
\begin{aligned}
& \sup _{\xi \in \mathbb{R}^{N}} \frac{|\widehat{f}(\xi)-\widehat{g}(\xi)|}{|\xi|^{s}} \\
& \quad \leq \inf _{n \geq k} \sup _{\xi \in \mathbb{R}^{N}} \frac{\left|\widehat{f}_{n}(\xi)-\widehat{g}_{n}(\xi)-(\widehat{f}(\xi)-\widehat{g}(\xi))\right|}{|\xi|^{s}}+\inf _{n \geq k} \sup _{\xi \in \mathbb{R}^{N}} \frac{\left|\widehat{f}_{n}(\xi)-\widehat{g}_{n}(\xi)\right|}{|\xi|^{s}} .
\end{aligned}
$$

Now, from the inequality

$$
\sup _{|\xi| \leq \delta} \frac{\left|\widehat{f_{n}}(\xi)-\widehat{g}_{n}(\xi)\right|}{|\xi|^{s}} \leq \sup _{|\xi| \leq \delta} \frac{\left|\widehat{f_{n}}(\xi)-\widehat{g}_{n}(\xi)\right|}{|\xi|^{r}}|\delta|^{r-s},
$$

for any given $\epsilon>0$ there exists $\delta>0$ such that, for all $n$

$$
\sup _{|\xi| \leq \delta} \frac{\left|\widehat{f}_{n}(\xi)-\widehat{g}_{n}(\xi)-(\widehat{f}(\xi)-\widehat{g}(\xi))\right|}{|\xi|^{s}} \leq(C+D)|\delta|^{r-s} \leq \epsilon .
$$

Likewise, there exists $R>0$ such that, for all $n$

$$
\sup _{|\xi| \geq R} \frac{\left|\widehat{f_{n}}(\xi)-\widehat{g}_{n}(\xi)-(\widehat{f}(\xi)-\widehat{g}(\xi))\right|}{|\xi|^{s}} \leq 4 / R^{s} \leq \epsilon
$$

Thus, for any $\epsilon>0$ we can find $\delta$ and $R$ such that

$$
\sup _{\xi \in \mathbb{R}^{N}} \frac{\left|\widehat{f}_{n}(\xi)-\widehat{g}_{n}(\xi)-(\widehat{f}(\xi)-\widehat{g}(\xi))\right|}{|\xi|^{s}}
$$

$$
\begin{align*}
& \leq \max \left\{\epsilon, \sup _{\delta \leq|\xi| \leq R} \frac{\left|\widehat{f}_{n}(\xi)-\widehat{g}_{n}(\xi)-(\widehat{f}(\xi)-\widehat{g}(\xi))\right|}{|\xi|^{s}}\right\} \\
& \leq \max \left\{\epsilon, \sup _{\delta \leq|\xi| \leq R} \frac{1}{\delta^{s}}\left(\left|\widehat{f}_{n}(\xi)-\widehat{f}(\xi)\right|+\left|\widehat{g}_{n}(\xi)-\widehat{g}(\xi)\right|\right)\right\} . \tag{2.38}
\end{align*}
$$

Now, since $f_{n} \rightharpoonup f$ and $g_{n} \rightharpoonup g$ we have

$$
\forall \xi, \widehat{f}_{n}(\xi) \longrightarrow \widehat{f}(\xi)
$$

and since $f_{n}$ satisfies conditions (1.10), $\left|D^{2} \widehat{f_{n}}(\xi)\right| \leq 1$ and $D \widehat{f_{n}}(0)=0$. Thus $\left\{\widehat{f}_{n}\right\}_{n \geq 0}$ is uniformly equi-continuous on the compact set $\{\delta \leq|\xi| \leq R\}$. By Ascoli's theorem, this entails that $\sup _{\delta \leq|\xi| \leq D}\left|\widehat{f}_{n}(\xi)-\widehat{f}(\xi)\right|$ goes to 0 . Same conclusion for the sequence $g_{n}$. Finally, there exists $n_{0} \geq 0$ such that for $n \geq n_{0}$,

$$
\begin{equation*}
\max \left\{\epsilon, \sup _{\delta \leq|\xi| \leq R} \frac{1}{\delta^{s}}\left(\left|\widehat{f_{n}}(\xi)-\widehat{f}(\xi)\right|+\left|\widehat{g}_{n}(\xi)-\widehat{g}(\xi)\right|\right)\right\} \leq \epsilon \tag{2.39}
\end{equation*}
$$

This concludes the proof.

## 3. THE ONE-DIMENSIONAL FOKKER-PLANCK EQUATION

### 3.1. A Nonconservative Kinetic Model

This section concerns the introduction of some nonlinear kinetic model of Boltzmann type, which is related to the Fokker-Planck Eq. (2.1) in the so-called grazing collision limit procedure. ${ }^{(30,31)}$ Similar ideas were introduced in Ref. 23 in order to obtain information on the self-similar profile of one-dimensional nonlinear Boltzmann equations of Maxwell type, in the case of lack of conservation of the energy. The main advantage in working with a Boltzmann type equation relies in the possibility to obtain in a relatively easy way the rate of decay of the Fourier based distance.

Let us consider a binary interaction between particles governed by the law:

$$
\begin{equation*}
v^{*}=v+n(\lambda v+\mu w), \quad w^{*}=w+n(\lambda w+\mu v) \tag{3.1}
\end{equation*}
$$

where $(v, w)$ are the pre-collisional velocities which generate the post-collisional ones $\left(v^{*}, w^{*}\right)$. In (3.1) $\lambda$ and $\mu$ are positive constants and $n$ is a parameter which varies in $[-1,1]$.

Let $f(v, t)$ denote the distribution of particles with velocity $v \in \mathbb{R}$ at time $t \geq 0$. A kinetic model governed by binary collisions among particles can be easily derived by standard methods of kinetic theory, considering that the change in time of $f(v, t)$ depends on a balance between the gain and loss of particles
with velocity $v$ due to collisions. This leads to the following integro-differential equation of Boltzmann type, ${ }^{(23)}$

$$
\begin{equation*}
\frac{\partial f}{\partial t}=\int_{-1}^{1} \int_{\mathbb{R}} B_{\epsilon}(n)\left(\frac{1}{J} f\left(v_{*}\right) f\left(w_{*}\right)-f(v) f(w)\right) d w d n \tag{3.2}
\end{equation*}
$$

In (3.2) $\left(v_{*}, w_{*}\right)$ are the pre-collisional velocities that generate the couple $(v, w)$ after the interaction. $J$ is the Jacobian of the transformation of $(v, w)$ into $\left(v^{*}, w^{*}\right)$, while $B_{\epsilon}(n)$ represents the collision frequency. The collision frequency depends on a small parameter $\epsilon>0$. As $\epsilon \rightarrow 0$, binary collisions concentrate on collisions which are grazing, namely the post-collision velocities are close to the pre-collision ones. This phenomenon is well-known for the Boltzmann equation, where the grazing limit allows to recover in the limit the Landau equation. ${ }^{(30,31)}$ In what follows, we will assume that $B_{\epsilon}(n)$ satisfies the following properties:
[P.1] For all $\epsilon, B_{\epsilon}(n)$ is a symmetric function of $n$, so that $\int_{-1}^{1} B_{\epsilon}(n) n d n=0$
[P.2] For all $\epsilon$, the second moment of $B_{\epsilon}(n)$ is bounded, $A_{\epsilon}:=\int_{-1}^{1} B_{\epsilon}(n) n^{2} d n$. Moreover $\lim _{\epsilon \rightarrow 0} A_{\epsilon}=1$
[P.3] For any $r>2$, with $r \in \mathbb{N}$, it holds $\lim _{\epsilon \rightarrow 0} \int_{-1}^{1}|n|^{r} B_{\epsilon}(n) d n=0$.
The kinetic Eq. (3.2) is the analogous of the Boltzmann equation for Maxwell molecules. ${ }^{(4)}$ Also, it presents several similarities with the one-dimensional Kac model without cut-off introduced by Desvillettes. ${ }^{(15)}$ Among the possible choices of the collision frequency, one can consider the following

$$
B_{\epsilon}(n)= \begin{cases}\frac{1}{2 \epsilon\left(n^{2}+\epsilon^{4}\right)} & \text { if } n \in[-\epsilon, \epsilon] \\ 0 & \text { otherwise }\end{cases}
$$

When necessary, we will denote by $f_{\epsilon}(v, t)$ the solution to the initial value problem for Eq. (3.2). In this way, we emphasize the dependence of this solution on the $\epsilon$-parameter. Moreover, without loss of generality, we can fix the nonnegative initial density $f_{0}(v)$ to satisfy conditions (1.10).

Remark 3.1. At a first view, the choice of a nonnegative initial density satisfying conditions (1.10), while quite natural for the study of a kinetic model of Boltzmann type, appears quite restrictive at least for the study of the linearized Eq. (2.14), where the solution is a perturbation of the Barenblatt profile, and there are no reasons to pretend the perturbation to be positive. To clarify why it is enough to get results on the kinetic Eq. (3.2) with nonnegative initial data, consider that thanks to mass conservation, this equation can be rewritten as

$$
\begin{equation*}
\frac{\partial f}{\partial t}=Q_{\epsilon}(f, f)-C_{\epsilon} f(v) \tag{3.3}
\end{equation*}
$$

where

$$
\begin{equation*}
Q_{\epsilon}(f, f)(v)=\int_{-1}^{1} \int_{\mathbb{R}} B_{\epsilon}(n) \frac{1}{J} f\left(v_{*}\right) f\left(w_{*}\right) d w d n \tag{3.4}
\end{equation*}
$$

and

$$
\begin{equation*}
C_{\epsilon}=\int_{-1}^{1} B_{\epsilon}(n) d n . \tag{3.5}
\end{equation*}
$$

The solution to Eq. (3.3) can be expressed as

$$
\begin{equation*}
f(v, t)=f_{0}(v) \exp \left\{-C_{\epsilon} t\right\}+\int_{0}^{t} Q_{\epsilon}(f, f)(v, s) \exp \left\{-C_{\epsilon}(t-s)\right\} d s \tag{3.6}
\end{equation*}
$$

Of course, representation (3.6) holds independently of the sign of the initial datum $f_{0}(v)$. On the other hand, in consequence of the form of the bilinear operator $Q_{\epsilon}$,

$$
\begin{equation*}
|f(v, t)| \leq\left|f_{0}(v)\right| \exp \left\{-C_{\epsilon} t\right\}+\int_{0}^{t} Q_{\epsilon}(|f|,|f|)(v, s) \exp \left\{-C_{\epsilon}(t-s)\right\} d s \tag{3.7}
\end{equation*}
$$

Since the solution to the Boltzmann Eq. (3.3) is nonnegative in correspondence to a nonnegative initial density, choosing $\left|f_{0}(x)\right|$ as initial density, the solution $f_{+}(v, t)$ to Eq. (3.3) is nonnegative, and

$$
\begin{equation*}
f_{+}(v, t)=\left|f_{0}(v)\right| \exp \left\{-C_{\epsilon} t\right\}+\int_{0}^{t} Q_{\epsilon}\left(f_{+}, f_{+}\right)(v, s) \exp \left\{-C_{\epsilon}(t-s)\right\} d s \tag{3.8}
\end{equation*}
$$

This implies $|f(v, t)| \leq f_{+}(v, t)$. Thus, any upper bound related to the absolute value of the solution to the Boltzmann equation with a general initial value (with no sign!) can be derived from the same equation taking as initial value a nonnegative density (the absolute value of the initial value).

On the other hand, since the linearized Eq. (2.14) is a linear equation the solution can be split in positive part and negative one and solve the equation in separated form. Then the solution is the sum of both. Namely, given a initial data $f_{0}$ we split in $f_{0}=f_{0_{+}}+f_{0_{-}}$and we obtain the solution as $f=f_{+}+f_{-}$. Therefore, we can reduce us to nonnegative initial data.

### 3.2. A Related Conservative Kinetic Model

Thanks to the remark of the previous section, without loss of generality, in what follows we will work with nonnegative initial values, satisfying conditions (1.10). To avoid the presence of the Jacobian, and to study approximation to the collision operator it is extremely convenient to write Eq. (3.2) in weak form. It
corresponds to consider, for all smooth functions $\phi(v)$, the equation
$\frac{d}{d t} \int_{\mathbb{R}} \phi(v) f_{\epsilon}(v, t) d v=\int_{-1}^{1} \int_{\mathbb{R}^{2}} B_{\epsilon}(n)\left(\phi\left(v^{*}\right)-\phi(v)\right) f_{\epsilon}(v) f_{\epsilon}(w) d w d v d n$,

One can alternatively use the symmetric form

$$
\begin{align*}
\frac{d}{d t} \int_{\mathbb{R}} f_{\epsilon}(v) \phi(v) d v= & \frac{1}{2} \int_{-1}^{1} \int_{\mathbb{R}^{2}} B_{\epsilon}(n) f_{\epsilon}(v) f_{\epsilon}(w) \\
& \times\left(\phi\left(v^{*}\right)+\phi\left(w^{*}\right)-\phi(v)-\phi(w)\right) d v d w d n . \tag{3.10}
\end{align*}
$$

Existence of solutions to (3.9) or, equivalently, to (3.10), can be proven by using arguments like in Ref. 15. Evolution of moments then follows by a suitable choice of the function $\phi$. If $\phi=1$ we obtain from (3.9)

$$
\begin{equation*}
\frac{d}{d t} \int_{\mathbb{R}} f_{\epsilon}(v, t) d v=0 \tag{3.11}
\end{equation*}
$$

namely mass conservation. Conservation of momentum is obtained by taking $\phi(v)=v$, and making use of assumption P.1. We have

$$
\begin{align*}
\frac{d}{d t} \int_{\mathbb{R}} v f_{\epsilon}(v, t) d v & =\int_{-1}^{1} \int_{\mathbb{R}^{2}} B_{\epsilon}(n)(n(\lambda v+\mu w)) f_{\epsilon}(v) f_{\epsilon}(w) d w d v d n \\
& =(\lambda+\mu) \int_{-1}^{1} B_{\epsilon}(n) n d n \int_{\mathbb{R}} v f_{\epsilon}(v) d v=0 . \tag{3.12}
\end{align*}
$$

Taking finally $\phi(v)=v^{2}$ we obtain

$$
\begin{align*}
& \frac{d}{d t} \int_{\mathbb{R}} v^{2} f_{\epsilon}(v, t) d v \\
& \quad=\int_{-1}^{1} \int_{\mathbb{R}^{2}} B_{\epsilon}(n)\left((v+n(\lambda v+\mu w))^{2}-v^{2}\right) f_{\epsilon}(v) f_{\epsilon}(w) d w d v d n \\
& \quad=\int_{-1}^{1} \int_{\mathbb{R}^{2}} B_{\epsilon}(n)\left(n^{2}(\lambda v+\mu w)^{2}+2 n v(\lambda v+\mu w)\right) f_{\epsilon}(v) f_{\epsilon}(w) d w d v d n \\
& \quad=\left(\lambda^{2}+\mu^{2}\right) \int_{-1}^{1} B_{\epsilon}(n) n^{2} d n \int_{\mathbb{R}} v^{2} f_{\epsilon}(v) d v=\left(\lambda^{2}+\mu^{2}\right) A_{\epsilon} \int_{\mathbb{R}} v^{2} f_{\epsilon}(v) d v . \tag{3.13}
\end{align*}
$$

Since the initial density satisfies (1.10), we obtain that the second moment is exponentially increasing, and

$$
\begin{equation*}
M_{2}(t)=\exp \left\{A_{\epsilon}\left(\lambda^{2}+\mu^{2}\right) t\right\} . \tag{3.14}
\end{equation*}
$$

Due to (3.14), it follows clearly that stationary solutions of finite energy do not exist, and the large-time behavior of the kinetic equation can at best be described by self-similar solutions. The standard way to look for self-similar profiles is to scale the solution according to the role

$$
\begin{equation*}
g_{\epsilon}(v, t)=\sqrt{M_{2}(t)} f_{\epsilon}\left(v \sqrt{M_{2}(t)}, t\right) . \tag{3.15}
\end{equation*}
$$

This scaling implies that $\int v^{2} g_{\epsilon}(v, t)=1$ for all $t \geq 0$. Assuming $\phi$ smooth and of bounded support, by elementary computations one concludes that $g_{\epsilon}=g_{\epsilon}(v, t)$ satisfies the equation

$$
\begin{align*}
& \frac{d}{d t} \\
& \int_{\mathbb{R}} \phi(v) g_{\epsilon}(v, t) d v+\frac{\left(\lambda^{2}+\mu^{2}\right) A_{\epsilon}}{2} \int_{\mathbb{R}} v g_{\epsilon} \phi^{\prime}(v) d v  \tag{3.16}\\
& \quad=\int_{-1}^{1} \int_{\mathbb{R}^{2}} B_{\epsilon}(n)\left(\phi\left(v^{*}\right)-\phi(v)\right) g_{\epsilon}(v) g_{\epsilon}(w) d w d v d n
\end{align*}
$$

Let us consider a third-order Taylor expansion of $\phi\left(v^{*}\right)$ around $v$

$$
\phi\left(v^{*}\right)-\phi(v)=\phi^{\prime}(v)\left(v^{*}-v\right)+\frac{\phi^{\prime \prime}(v)}{2}\left(v^{*}-v\right)^{2}+\frac{\phi^{\prime \prime \prime}(\tilde{v})}{3!}\left(v^{*}-v\right)^{3}
$$

or equivalently

$$
\begin{aligned}
\phi\left(v^{*}\right)-\phi(v)= & \phi^{\prime}(v)(n(\lambda v+\mu w))+\frac{\phi^{\prime \prime}(v)}{2} n^{2}(\lambda v+\mu w)^{2} \\
& +\frac{\phi^{\prime \prime \prime}(\tilde{v})}{3!} n^{3}(\lambda v+\mu w)^{3},
\end{aligned}
$$

where $\tilde{v}$ is a value between $v$ and $v^{*}$. Substituting the Taylor expansion into the integral on the right-hand side of (3.16), taking the limit $\epsilon \rightarrow 0$ and using property P. 3 we obtain that the integral containing the third order term vanishes. Moreover, provided $g_{\epsilon}(v, t)$ converges to $g(v, t)$ as $\epsilon \rightarrow 0$, using mass and momentum conservation into the weak form (3.16) we obtain that $g(v, t)$ satisfies

$$
\begin{align*}
& \frac{d}{d t} \int_{\mathbb{R}} \phi(v) g(v, t) d v+\frac{\left(\lambda^{2}+\mu^{2}\right)}{2} \int_{\mathbb{R}} v g \phi^{\prime}(v) d v \\
& \quad=\int_{\mathbb{R}} \frac{\left(\lambda^{2} v^{2}+\mu^{2}\right)}{2} \phi^{\prime \prime}(v) g(v) d v \tag{3.17}
\end{align*}
$$

which is the weak formulation of the following Fokker-Planck equation

$$
\begin{equation*}
\frac{\partial}{\partial t} g(v, t)-\frac{\left(\lambda^{2}+\mu^{2}\right)}{2} \frac{\partial}{\partial v}(v g)=\frac{\partial^{2}}{\partial v^{2}}\left(\frac{\left(\lambda^{2} v^{2}+\mu^{2}\right)}{2} g(v)\right) . \tag{3.18}
\end{equation*}
$$

Since in our linearized equations the coefficient in the convection term is one we take $\lambda$ and $\mu$ such that $\lambda^{2}+\mu^{2}=2$ and therefore

$$
\begin{equation*}
\frac{\partial}{\partial t} g=\frac{\partial}{\partial v}(v g)+\frac{\partial^{2}}{\partial v^{2}}\left(\left(\frac{\lambda^{2}}{2} v^{2}+\frac{\mu^{2}}{2}\right) g\right) . \tag{3.19}
\end{equation*}
$$

Note that Eq. (3.19) is nothing but the Fokker-Planck Eq. (2.1), where

$$
\begin{equation*}
\beta^{2}=\frac{\lambda^{2}}{2} \quad \text { and } \quad \alpha^{2}=\frac{\mu^{2}}{2} \tag{3.20}
\end{equation*}
$$

The passage to the limit outlined above is largely formal. The derivation, however, can be made rigorous by using the same arguments like in Ref. 23. Consider in fact the class of probability densities $\left\{g_{\epsilon}(v, t)\right\}_{\epsilon \geq 0}$. For all $\epsilon$, the solution to (3.16) satisfies conditions (1.10). By virtue of Prokhorov theorem (cfr. ${ }^{(20)}$ ) the existence of a uniform bound on the second moment implies that this class is tight, so that any sequence $\left\{g_{\epsilon_{n}}(v, t)\right\}_{n \geq 0}$ contains an infinite subsequence which converges weakly to some probability measure $g(v, t)$. Therefore, the weak solution to the Boltzmann Eq. (3.16) converges, up to extraction of a subsequence, to a probability density $g(v, t)$. This density is a weak solution of the Fokker-Planck Eq. (3.19).

### 3.3. The Evolution of Higher Moments

A detailed calculation, along the lines of Ref. 23, allows to recognize how many moments of the solution to (3.2) remain uniformly bounded in time with respect to $\epsilon$.

Suppose that the initial density $g_{0}(v)=f_{0}(v)$ is such that

$$
\begin{equation*}
\int_{\mathbb{R}}|v|^{2+\delta} g_{0}(v) d v=M_{2+\delta}<\infty . \tag{3.21}
\end{equation*}
$$

Then, since the contribution due to the term $\frac{\partial}{\partial v}\left(v g_{\epsilon}(v)\right)$ can be evaluated integrating by parts,

$$
\int_{\mathbb{R}}|v|^{2+\delta} \frac{\partial}{\partial v}\left(v g_{\epsilon}(v)\right) d v=-(2+\delta) \int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v
$$

we obtain

$$
\begin{align*}
& \frac{d}{d t} \int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v+(2+\delta) \frac{\lambda^{2}+\mu^{2}}{2} A_{\epsilon} \int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v \\
& \quad=\int_{-1}^{1} B_{\epsilon}(n) d n \int_{\mathbb{R}^{2}} d v d w\left(|(1+\lambda n) v+\mu n w|^{2+\delta}-|v|^{2+\delta}\right) g_{\epsilon}(v) g_{\epsilon}(w) . \tag{3.22}
\end{align*}
$$

Let us recover a suitable upper bound for the last integral in (3.22). Given any two constants $a, b$, and $0<\delta \leq 1$ the following inequality holds

$$
\begin{equation*}
(|a|+|b|)^{\delta} \leq|a|^{\delta}+|b|^{\delta} . \tag{3.23}
\end{equation*}
$$

Hence, choosing $a=|1+\lambda n||v|$ and $b=|\mu n||w|$,

$$
|(1+\lambda n) v+\mu n w|^{2+\delta} \leq((1+\lambda n) v+\mu n w)^{2}\left(|1+\lambda n|^{\delta}|v|^{\delta}+|\mu n|^{\delta}|w|^{\delta}\right) .
$$

Substituting into the right-hand side of (3.22), recalling that the mean value of $g$ is equal to zero, and the second moment of $g$ equal to one, gives

$$
\begin{aligned}
& \int_{-1}^{1} B_{\epsilon}(n) d n \int_{\mathbb{R}^{2}}|(1+\lambda n) v+\mu n w|^{2+\delta} g_{\epsilon}(v) g_{\epsilon}(w) d v d w \\
& \leq \int_{\mathbb{R}^{2}}((1+\lambda n) v+\mu n w)^{2}\left(|1+\lambda n|^{\delta}|v|^{\delta}+|\mu n|^{\delta}|w|^{\delta}\right) g_{\epsilon}(v) g_{\epsilon}(w) d v d w \\
&= \int_{-1}^{1} B_{\epsilon}(n) d n\left(|1+\lambda n|^{2+\delta}+|\mu n|^{2+\delta}\right) \int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v) d v \\
&+\int_{-1}^{1} B_{\epsilon}(n) d n\left(|1+\lambda n|^{2}|\mu n|^{\delta}+|\mu n|^{2}|1+\lambda n|^{\delta}\right) \int_{\mathbb{R}}|v|^{\delta} g_{\epsilon}(v) d v .
\end{aligned}
$$

Grouping all these inequalities, we obtain

$$
\begin{equation*}
\frac{d}{d t} \int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v \leq \mathcal{S}_{\epsilon}(\delta) \int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v+D_{\delta} \tag{3.24}
\end{equation*}
$$

where

$$
\begin{equation*}
\mathcal{S}_{\epsilon}(\delta)=\int_{-1}^{1} B_{\epsilon}(n) d n\left(|1+\lambda n|^{2+\delta}+|\mu n|^{2+\delta}-1\right)-(2+\delta) \frac{\lambda^{2}+\mu^{2}}{2} A_{\epsilon}, \tag{3.25}
\end{equation*}
$$

and, by Hölder inequality

$$
\begin{equation*}
D_{\delta} \leq \int_{-1}^{1} B_{\epsilon}(n) d n\left(|1+\lambda n|^{2}|\mu n|^{\delta}+|\mu n|^{2}|1+\lambda n|^{\delta}\right) . \tag{3.26}
\end{equation*}
$$

By property $\mathbf{P} 3$ it follows that $D_{\delta}$ is uniformly bounded with respect to $\epsilon$. On the other hand, using a Taylor development up to the order three we obtain

$$
\begin{equation*}
|1+\lambda n|^{2+\delta}=1+(2+\delta) \lambda n+\frac{(2+\delta)(1+\delta)}{2} \lambda^{2} n^{2}+O\left(n^{3}\right) \tag{3.27}
\end{equation*}
$$

Substituting into (3.25) we obtain

$$
\begin{equation*}
\mathcal{S}_{\epsilon}=\frac{2+\delta}{2}\left(\lambda^{2} \delta-\mu^{2}\right) A_{\epsilon}+0(\epsilon) . \tag{3.28}
\end{equation*}
$$

Hence, $\mathcal{S}_{\epsilon} \leq 0$ as soon as $\lambda^{2} \delta-\mu^{2}<0$ and $\epsilon$ is suitably small. In this case, inequality (3.24) gives an upper bound for the moment, that reads

$$
\begin{equation*}
\int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v \leq M_{2+\delta}+\frac{D_{\delta}}{\left|\mathcal{S}_{\epsilon}(\delta)\right|}<\infty \tag{3.29}
\end{equation*}
$$

Let us set $\bar{\delta}=\mu^{2} / \lambda^{2}$. In the case $\bar{\delta}>3$ we can easily iterate our procedure to obtain that any moment of order $2+\delta$, with $\delta<\bar{\delta}$ which is bounded initially, remains uniformly bounded with respect to $\epsilon$ at any subsequent time. The only difference now is that the explicit expression of the bound is more and more involved. We proved

Lemma 3.2. Let $g_{\epsilon}(v, t)$ be the solution to the initial value problem for the kinetic Eq. (3.16), where the initial datum $g_{0}(v)$ satisfies

$$
\begin{equation*}
\int_{\mathbb{R}}|v|^{2+\delta} g_{0}(v) d v=M_{2+\delta}<\infty . \tag{3.30}
\end{equation*}
$$

Then, if $\delta<\bar{\delta}=\mu^{2} / \lambda^{2}$,

$$
\begin{equation*}
\int_{\mathbb{R}}|v|^{2+\delta} g_{\epsilon}(v, t) d v \tag{3.31}
\end{equation*}
$$

is uniformly bounded in time for all $\epsilon$. Moreover, if $\delta<\bar{\delta}$ the function $\mathcal{S}_{\epsilon}$ defined by (3.25) satisfies

$$
\begin{equation*}
\lim _{\epsilon \rightarrow 0} \mathcal{S}_{\epsilon}(\delta)<0 \tag{3.32}
\end{equation*}
$$

### 3.4. Rate of Convergence in Fourier Based Metrics

Let $\mathcal{M}_{0}$ denote the space of all probability measures in $\mathbb{R}$ and

$$
\mathcal{M}_{r}=\left\{\mu \in \mathcal{M}_{0}: \int_{\mathbb{R}}|v|^{r} \mu(d v)<\infty, r \geq 0\right\}
$$

the space of all Borel probability measures with finite momentum of order $r$ equipped with the topology of the weak convergence of the measures. By a weak solution of the initial value problem for Eq. (3.2), corresponding to the initial probability density $f_{0}(w) \in \mathcal{M}_{r} r>2$ we shall mean any probability density $f \in C^{1}\left(\mathbb{R}, \mathcal{M}_{r}\right)$ satisfying the weak form of the Eq. (3.9) for $t>0$ and all smooth functions $\phi$, and such that for all $\phi$

$$
\begin{equation*}
\lim _{t \rightarrow 0} \int_{\mathbb{R}} \phi(v) f(v, t) d v=\int_{\mathbb{R}} \phi(v) f_{0}(v) d v \tag{3.33}
\end{equation*}
$$

In the rest of this section, we shall study the weak form of Eq. (3.2), with the normalization conditions (1.10). It is equivalent to use the Fourier transform of
the Eq. ${ }^{(4)}$ :

$$
\begin{equation*}
\frac{\partial \widehat{f}(\xi, t)}{\partial t}=\widehat{Q}(\widehat{f}, \widehat{f})(\xi, t) \tag{3.34}
\end{equation*}
$$

where $\widehat{f}(\xi, t)$ is the Fourier transform of $f(v, t)$ and

$$
\begin{equation*}
\widehat{Q}(\widehat{f}, \widehat{f})(\xi)=\int_{-1}^{1} B_{\epsilon}(n)[\hat{f}(\xi(1+\lambda n), t) \hat{f}(n \mu \xi, t)-\hat{f}(\xi, t)] d n \tag{3.35}
\end{equation*}
$$

The initial conditions (1.10) turn into

$$
\widehat{f}(0)=1, \quad \widehat{f}^{\prime}(0)=0, \quad \widehat{f}^{\prime \prime}(0)=-1,
$$

where $\widehat{f_{0}} \in C^{2}(\mathbb{R})$. Hence Eq. (3.34) can be rewritten (recall (3.5)) as

$$
\begin{equation*}
\frac{\partial \widehat{f}(\xi, t)}{\partial t}+C_{\epsilon} \widehat{f}(\xi, t)=\int_{-1}^{1} B_{\epsilon}(n) \hat{f}(\xi(1+\lambda n), t) \hat{f}(n \mu \xi, t) d n \tag{3.36}
\end{equation*}
$$

We work with the metric introduced in ${ }^{(17)}$

$$
\begin{equation*}
d_{s}(f, g)=\sup _{\xi \in \mathbb{R}} \frac{|\hat{f}(\xi)-\hat{g}(\xi)|}{|\xi|^{s}} \tag{3.37}
\end{equation*}
$$

It is easy to show that (see Refs. 23 and 24), for any given collision integral of form (3.35), given two initial densities $f_{0,1}$ and $f_{0,2}$ belonging to $\mathcal{M}_{r}, r>2$, and such that $d_{r}\left(f_{0,1}, f_{0,2}\right)$ is bounded, $d_{r}\left(f_{1}(t), f_{2}(t)\right)$ is bounded at any subsequent time $t>0$, and the following inequality holds

$$
\begin{equation*}
d_{r}\left(f_{1}(t), f_{2}(t)\right) \leq d_{r}\left(f_{0,1}, f_{0,2}\right) \cdot \exp \left\{\left(\int_{-1}^{1} B_{\epsilon}(n)\left[|1+\lambda n|^{r}+(\mu n)^{r}-1\right] d n\right) t\right\} . \tag{3.38}
\end{equation*}
$$

Inequality (3.38) implies a similar inequality for the scaled function $g$, which we recall is defined through the scaling $g(v, t)=\sqrt{M_{2}(t)} f\left(v \sqrt{M_{2}(t)}, t\right)$. The Fourier transform of $g$ is $\hat{g}(\xi)=\hat{f}\left(\xi / \sqrt{M_{2}(t)}\right)$. From the definition of $d_{r}$ it follows

$$
d_{r}\left(g_{1}(t), g_{2}(t)\right)=\left(\frac{1}{\sqrt{M_{2}(t)}}\right)^{r} d_{r}\left(f_{1}(t), f_{2}(t)\right)
$$

Therefore

$$
\begin{align*}
& d_{r}\left(g_{1}(t), g_{2}(t)\right) \leq d_{r}\left(g_{0,1}, g_{0,2}\right) \\
& \quad \cdot \exp \left\{\left(\int_{-1}^{1} B_{\epsilon}(n)\left[|1+\lambda n|^{r}+|\mu n|^{r}-1-\frac{r}{2} n^{2}\left(\lambda^{2}+\mu^{2}\right)\right] d n\right) t\right\} . \tag{3.39}
\end{align*}
$$

The exponent in (3.39) coincides with the function $\mathcal{S}_{\epsilon}(\cdot)$ defined in (3.25), and evaluated at $\delta=r-2$.

Applying Lemma 2.4, and taking the limit $\epsilon \rightarrow 0$ on (3.39) and calling again $g$ the solution, for all $\delta<r-2$ we find that the solution to the Fokker-Planck Eq. (3.18) satisfies

$$
\begin{equation*}
d_{2+\delta}\left(g_{1}(t), g_{2}(t)\right) \leq d_{2+\delta}\left(g_{0,1}, g_{0,2}\right) \exp \left\{\frac{2+\delta}{2}\left(\lambda^{2} \delta-\mu^{2}\right) t\right\} . \tag{3.40}
\end{equation*}
$$

The exact rate of convergence in $\delta$-metric follows by analyzing the parabola $h(\delta)=\frac{2+\delta}{2}\left(\lambda^{2} \delta-\mu^{2}\right)$, that is nonpositive in the interval $\delta \in\left[0, \bar{\delta}=\mu^{2} / \lambda^{2}\right]$. The minimum value in this interval is achieved at $\delta_{0}=\bar{\delta} / 2-1$ where $h\left(\delta_{0}\right)=$ $-\frac{1}{8}\left(\mu^{2}+2 \lambda^{2}\right)$. Thus, if $\delta_{0}>0$, in the interval $\delta \in\left[0, \delta_{0}\right]$ the rate of convergence in $d_{2+\delta}$-metric is increasing, and we have a gain in the rate of convergence if the initial data have a big number of equal moments. Now, $\delta_{0}>0$ if $\mu^{2}>2 \lambda^{2}$, or, what is the same (see (3.20)),

$$
\begin{equation*}
\alpha^{2}=\frac{\mu^{2}}{2}>2 \beta^{2}=\frac{2 \lambda^{2}}{2} \tag{3.41}
\end{equation*}
$$

This implies that we have an improved rate of convergence in $d_{2+\delta}$-metric for Eq. (2.18) in the range $1 / 3<m<1$, while the same holds for Eq. (2.21) in the range $3 / 5<m<1$. We include these results into the following.

Theorem 3.3. Let $f(v, t)$ be the solution to Eq. (2.14), corresponding to the initial datum $f_{0}(v)$ satisfying (1.10). Suppose in addition that

$$
\begin{equation*}
d_{2+\delta}\left(f_{0}, D B_{C}^{2-m}\right)<+\infty, \quad \delta<(1+m) /(1-m) \tag{3.42}
\end{equation*}
$$

Then, the solution $f(v, t)$ converges exponentially to $B_{C}(v)^{2-m}$ in $d_{2+\delta}-m e t r i c$, and
$d_{2+\delta}\left(f(t), D B_{C}^{2-m}\right) \leq d_{2+\delta}\left(f_{0}, D B_{C}^{2-m}\right) \exp \left\{\frac{2+\delta}{2}((1-m) \delta-(1+m)) t\right\}$.

Moreover, provided $1 / 3<m<1$ the rate of convergence in (3.43) increases in the interval $0<\delta \leq(3 m-1) /(2-2 m)$. Likewise, let $f(v, t)$ be the solution to Eq. (1.6), with initial datum $f_{0}(v)$. Suppose in addition that

$$
\begin{equation*}
d_{2+\delta}\left(f_{0}, B_{C}\right)<+\infty, \quad \delta<(3 m-1) /(1-m) \tag{3.44}
\end{equation*}
$$

Then, the solution $f(v, t)$ converges exponentially to $B_{C}(x)$ in $d_{2+\delta}-m e t r i c$, and

$$
\begin{equation*}
d_{2+\delta}\left(f(t), B_{C}\right) \leq d_{2+\delta}\left(f_{0}, B_{C}\right) \exp \left\{\frac{2+\delta}{2}\left(\frac{1-m}{m} \delta-\frac{3 m-1}{m}\right) t\right\} \tag{3.45}
\end{equation*}
$$

Moreover, provided $3 / 5<m<1$ the rate of convergence in (3.45) increases in the interval $0<\delta \leq(5 m-3) /(2-2 m)$.

Proof: The proof of this theorem is a direct consequence of the previous computations and we can obtain it by identification of coefficients, as follows. First of all, we recall that both linear equations ((2.14) and (1.6)) can be written using different scaling as Eq. (2.1) (see Sec. 2.1). For Eq. (2.14) we have

$$
\alpha^{2}=\frac{\mu^{2}}{2}=\frac{1+m}{2} \quad \text { and } \quad \beta^{2}=\frac{\lambda^{2}}{2}=\frac{1-m}{2}
$$

thus (3.41) holds if and only if $m>1 / 3$. On the other hand, since

$$
\mu^{2}=1+m \quad \text { and } \quad \lambda^{2}=1-m
$$

we obtain that the rate of convergence $\frac{2+\delta}{2}\left(\lambda^{2} \delta-\mu^{2}\right)$ is written as

$$
\frac{2+\delta}{2}((1-m) \delta-(1+m))
$$

And the minimum value is achieved in $\delta_{0}=\frac{3 m-1}{2-2 m}$ (see previous computations) and this value is

$$
-\frac{1}{2^{3}} \frac{(3-m)^{2}}{1-m} .
$$

For Eq. (1.6) we have

$$
\alpha^{2}=\frac{\mu^{2}}{2}=\frac{3 m-1}{2 m} \quad \text { and } \quad \beta^{2}=\frac{\lambda^{2}}{2}=\frac{1-m}{2 m}
$$

and in this case (3.41) holds if and only if $m>3 / 5$ and the rate of convergence is written as

$$
\frac{2+\delta}{2}\left(\frac{1-m}{2 m} \delta-\frac{3 m-1}{2 m}\right)
$$

And the minimum value is achieved in $\delta_{0}=\frac{5 m-3}{2-2 m}$ (see previous computations) and this value is

$$
-\frac{1}{2^{3}} \frac{(m+1)^{2}}{m} .
$$

The result of Theorem 3.3 introduces new bounds into the game. As a matter of fact, while our method of proof improves the rate of convergence in both cases as soon as we are sufficiently close from the linear heat equation, while we are not able to show convergence to equilibrium as soon as we are too far from the linear one.

## 4. THE $N$-DIMENSIONAL FOKKER-PLANCK EQUATION

The results of Sec. 3 clarify the strategy we used to reckon precise rates of convergence towards the steady state for the one-dimensional Fokker-Planck Eq. (2.4). In one dimension of the velocity variable the connection between nonconservative models of the Boltzmann equation and the Fokker-Planck equation with non constant diffusion coefficients was known for binary collisions of different type. ${ }^{(23)}$ To our knowledge, however, in higher dimensions, at present no studies on this connection have been considered so far. Clearly, higher dimensional models have been studied in the grazing limit in the conservative case by Villani, ${ }^{(30,31)}$ but there, starting from the nonlinear Boltzmann equation with elastic collisions, the limit equation results in the nonlinear Landau equation. To obtain in the limit a linear Fokker-Planck type equation requires the introduction of a particular collision dynamics. We deal with this extension in the forthcoming Section. Since many arguments will be taken from the one-dimensional analysis of Sec. 3, when not necessary to the understanding, we will only briefly outline them.

## 4.1. $N$-Dimensional Binary Collisions and Nonconservative Kinetic Models

The starting point of our analysis will be the introduction of a particular collision dynamics, that in the grazing limit is able to produce the target Eq. (2.4). Clearly, as the one-dimensional analysis suggests, this is not the only possible model of binary collision that produces the correct result. Nevertheless, this law is in our opinion the best possible generalization of (3.1) to higher dimensions. The binary interaction between particles will be governed by the law:

$$
\begin{equation*}
v^{*}=v+n(\lambda e \cdot v+\mu e \cdot w), \quad w^{*}=w+n(\lambda e \cdot w+\mu e \cdot v) . \tag{4.1}
\end{equation*}
$$

In (4.1) $(v, w) \in \mathbb{R}^{2 N}$ are the pre-collisional velocities which generate the postcollisional ones $\left(v^{*}, w^{*}\right) \in \mathbb{R}^{2 N}, \lambda$ and $\mu$ are positive constants, $n \in[-1,1]^{N}$ and $e \in \mathbb{R}^{N}$ is a unit vector. Consequently, the integro-differential equation of Boltzmann type (3.2) is written as

$$
\begin{equation*}
\frac{\partial f}{\partial t}=\int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{N}} \frac{B_{\epsilon}(n)}{B_{N}}\left(\frac{1}{J} f\left(v_{*}\right) f\left(w_{*}\right)-f(v) f(w)\right) d w d e d n \tag{4.2}
\end{equation*}
$$

where $\left(v_{*}, w_{*}\right)$ are the pre-collisional velocities that generate the couple ( $\left.v, w\right)$ after the interaction, $J$ is the Jacobian of the transformation of $(v, w)$ into $\left(v^{*}, w^{*}\right)$, $C_{N}$ denotes the cube in $\mathbb{R}^{N}$ centered on the origin with dimensions $[-1,1]^{N}, S^{N-1}$ represents the unit sphere in $\mathbb{R}^{N}, B_{N}$ is the measure of the surface of $S^{N-1}$ and $B_{\epsilon}(n)$ represents the collision frequency. Proceeding as in dimension one, we will
assume that $B_{\epsilon}(n)$ satisfies the following properties:
$\mathbf{P}_{\mathbf{N}} .1$ For all $\epsilon, B_{\epsilon}(n)$ is a symmetric function of all the components of $n$, so that for all $i \int_{C_{N}} B_{\epsilon}(n) n_{i} d n=0, \quad$ for all $i \neq j, \int_{C_{N}} B_{\epsilon}(n) n_{i} n_{j} d n=0$
$\mathbf{P}_{\mathrm{N}} .2$ For all $\epsilon$, the second moment of $B_{\epsilon}(n)$ is bounded,

$$
A_{\epsilon}:=\int_{C_{N}} B_{\epsilon}(n)|n|^{2} d n<\infty . \text { Moreover } \lim _{\epsilon \rightarrow 0} A_{\epsilon}=1
$$

$\mathbf{P}_{\mathrm{N}} .3$ For any $r>2$, with $r \in \mathbb{N}$, it holds $\lim _{\epsilon \rightarrow 0} \int_{C_{N}}|n|^{r} B_{\epsilon}(n) d n=0$.
The weak formulation of (4.2) now reads

$$
\begin{align*}
\frac{d}{d t} \int_{\mathbb{R}^{N}} f_{\epsilon}(v) \phi(v) d v= & \frac{1}{2} \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w) \\
& \times\left(\phi\left(v^{*}\right)+\phi\left(w^{*}\right)-\phi(v)-\phi(w)\right) d v d w d e d n \tag{4.3}
\end{align*}
$$

Conservation of mass follows directly considering $\phi(v)=1$ and using $\mathbf{P}_{\mathbf{N}} \cdot \mathbf{1}$. Likewise, conservation of momentum follows by taking $\phi(v)=v$ and using $\mathbf{P}_{\mathbf{N}} \cdot \mathbf{1}$ again. The evaluation of the law of variation of the second momentum requires computations that we present below. Precisely we obtain

Lemma 4.1. Let $f_{\epsilon}$ be a solution of (4.2) corresponding to an initial value satisfying conditions (1.10). Then the moment of order two is exponentially increasing and it holds

$$
\begin{equation*}
M_{2}(t)=\exp \left\{\frac{A_{\epsilon}\left(\lambda^{2}+\mu^{2}\right)}{N} t\right\} . \tag{4.4}
\end{equation*}
$$

Proof: Choosing $\phi(v)=|v|^{2}$ into the weak form (4.3) we find

$$
\begin{aligned}
\frac{d}{d t} \int_{\mathbb{R}^{N}} f_{\epsilon}(v)|v|^{2} d v=\frac{1}{2} & \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w) \\
& \times\left(\left|v^{*}\right|^{2}+\left|w^{*}\right|^{2}-|v|^{2}-|w|^{2}\right) d v d w d e d n
\end{aligned}
$$

By the collision rule (4.1) we obtain

$$
\begin{aligned}
& \frac{d}{d t} \int_{\mathbb{R}^{N}} f_{\epsilon}(v)|v|^{2} d v=\int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w) \\
& \quad \times\left(|n|^{2}|\lambda e \cdot v+\mu e \cdot w|^{2}+2(\lambda e \cdot v+\mu e \cdot w) v \cdot n\right) d v d w d e d n .
\end{aligned}
$$

By property $\mathbf{P}_{\mathbf{N}} .1$ follows

$$
\int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w)(\lambda e \cdot v+\mu e \cdot w) v \cdot n d v d w d e d n=0 .
$$

Therefore,

$$
\begin{aligned}
\frac{d}{d t} \int_{\mathbb{R}^{N}} f_{\epsilon}(v)|v|^{2} d v= & \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w) \\
& \cdot|n|^{2}|\lambda e \cdot v+\mu e \cdot w|^{2} d v d w d e d n,
\end{aligned}
$$

which can be written as

$$
\begin{aligned}
\frac{d}{d t} \int_{\mathbb{R}^{N}} f_{\epsilon}(v)|v|^{2} d v= & \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w) \\
& \times|n|^{2}\left(\lambda^{2}(e \cdot v)^{2}+\mu^{2}(e \cdot w)^{2}\right. \\
& +2 \lambda \mu(e \cdot v)(e \cdot w)) d v d w d e d n .
\end{aligned}
$$

Finally, since mass and momentum are conserved, we obtain

$$
\begin{aligned}
\frac{d}{d t} \int_{\mathbb{R}^{N}} f_{\epsilon}(v)|v|^{2} d v= & \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v) f_{\epsilon}(w) \\
& \times|n|^{2}\left(\lambda^{2}(e \cdot v)^{2}+\mu^{2}(e \cdot w)^{2}\right) d v d w d e d n \\
= & \left(\lambda^{2}+\mu^{2}\right) \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{N}} \frac{B_{\epsilon}(n)}{B_{N}} f_{\epsilon}(v)|n|^{2}(e \cdot v)^{2} d v d e d n \\
= & \frac{A_{\epsilon}\left(\lambda^{2}+\mu^{2}\right)}{N} \int_{\mathbb{R}^{N}} f_{\epsilon}(v)|v|^{2} d v .
\end{aligned}
$$

The last equality follows by symmetry, since from the equality

$$
\begin{equation*}
\frac{1}{B_{N}} \int_{S^{N-1}} e_{i}^{2} d e=\frac{1}{B_{N}} \int_{S^{N-1}} e_{j}^{2} d e, \quad i \neq j \tag{4.5}
\end{equation*}
$$

it follows

$$
\frac{1}{B_{N}} \int_{S^{N-1}} e_{i}^{2} d e=\frac{1}{N} \quad i=1, \ldots n,
$$

and

$$
\frac{1}{B_{N}} \int_{S^{N-1}}(e \cdot v)^{2} d e=\frac{1}{N}|v|^{2} .
$$

As discussed in the one-dimensional case, the exponential growth of the second moment implies that do not exist stationary solutions with finite energy. To look for self-similar profiles, we scale the solution according to

$$
\begin{equation*}
g_{\epsilon}(v, t)=\left(M_{2}(t)\right)^{N / 2} f_{\epsilon}\left(v \sqrt{M_{2}(t)}, t\right) . \tag{4.6}
\end{equation*}
$$

By this scaling, we obtain that the second moment of $g_{\epsilon}$ remains equal to one (for all $t \geq 0$ ). Thus, if the initial value to the kinetic Eq. (4.2) satisfies condition (1.10), so does $g_{\epsilon}(v, t)$ at any subsequent time. The weak formulation of Eq. (4.2) for the scaled function $g_{\epsilon}$ is written as

$$
\begin{align*}
& \frac{d}{d t} \int_{\mathbb{R}^{N}} \phi(v) g_{\epsilon}(v, t) d v+\frac{\left(\lambda^{2}+\mu^{2}\right) A_{\epsilon}}{2 N} \int_{\mathbb{R}^{N}} \nabla \phi \cdot v g_{\epsilon} d v \\
& \quad=\int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}}\left(\phi\left(v^{*}\right)-\phi(v)\right) g_{\epsilon}(v) g_{\epsilon}(w) d w d v d e d n \tag{4.7}
\end{align*}
$$

If we consider a third-order Taylor development of $\phi\left(v^{*}\right)$ around $v$ we obtain

$$
\begin{aligned}
\phi\left(v^{*}\right)-\phi(v)= & \nabla_{v} \phi(v) \cdot\left(v^{*}-v\right)+\frac{1}{2}\left(v^{*}-v\right) \operatorname{Hess}(\phi)(v)\left(v^{*}-v\right)^{t} \\
& +R\left(\phi, \tilde{v}, v^{*}, v\right)
\end{aligned}
$$

where we denoted by Hess the Hessian matrix of a function, and by $w^{t}$ the transpose of a vector $w$. Moreover $R$ represents the remainder evaluated in correspondence to a vector $\tilde{v}=\lambda v+(1-\lambda) v^{*}$, where $\lambda$ is a suitable constant, $0 \leq \lambda \leq 1$. Resorting to the expression of $v^{*}$ in terms of $v$ we find

$$
\begin{aligned}
\phi\left(v^{*}\right)-\phi(v)= & \nabla_{v} \phi(v) \cdot n(\lambda e \cdot v+\mu e \cdot w) \\
& +\frac{1}{2} n(\lambda e \cdot v+\mu e \cdot w) \operatorname{Hess}(\phi)(v)(n(\lambda e \cdot v+\mu e \cdot w))^{t} \\
& +R\left(\phi, \tilde{v}, v^{*}, v\right) .
\end{aligned}
$$

Substituting the above expression into (4.7), and owing to property $\mathbf{P}_{\mathbf{N}} . \mathbf{3}$ shows that the term corresponding to the remainder vanishes when $\epsilon$ goes to zero. Therefore, we will focus on the remaining two terms. The integral corresponding to the first term in the Taylor expansion vanishes thanks to property $\mathbf{P}_{\mathbf{N}}$. 1. In fact

$$
\begin{aligned}
& \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} \nabla_{v} \phi(v) \cdot n(\lambda e \cdot v+\mu e \cdot w) g_{\epsilon}(v) g_{\epsilon}(w) d w d v d e d n \\
& =\int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{1}{B_{N}} \nabla_{v} \phi(v) \cdot\left(\int_{C_{N}} B_{\epsilon}(n) n d n\right)(\lambda e \cdot v+\mu e \cdot w) \\
& \quad \times g_{\epsilon}(v) g_{\epsilon}(w) d w d v d e=0 .
\end{aligned}
$$

On the other hand, by property $\mathbf{P}_{\mathbf{N}} \cdot \mathbf{1}$ the second term can we written as follow

$$
\begin{aligned}
& \frac{1}{2} \int_{C_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{B_{\epsilon}(n)}{B_{N}} g_{\epsilon}(v) g_{\epsilon}(w) n(\lambda e \cdot v+\mu e \cdot w) \\
& \quad \times \operatorname{Hess}\left(g_{\epsilon}\right)(v)(n(\lambda e \cdot v+\mu e \cdot w))^{t} d w d v d e d n
\end{aligned}
$$

$$
\begin{aligned}
= & \frac{1}{2} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{1}{B_{N}} g_{\epsilon}(v) g_{\epsilon}(w)\left(\sum_{i} \frac{\partial^{2} \phi(v)}{\partial v_{i}^{2}} \int_{C_{N}} B_{\epsilon}(n) n_{i}^{2} d n\right) \\
& \times(\lambda e \cdot v+\mu e \cdot w)^{2} d w d v d e \\
= & \frac{1}{2} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} \frac{1}{B_{N}} g_{\epsilon}(v) g_{\epsilon}(w)\left(\sum_{i} \frac{\partial^{2} \phi(v)}{\partial v_{i}^{2}} \frac{A_{\epsilon}}{N}\right) \\
& \times\left(\lambda^{2}(e \cdot v)^{2}+\mu^{2}(e \cdot w)^{2}+2 \lambda \mu(e \cdot v)(e \cdot w)\right) d w d v d e \\
= & \frac{A_{\epsilon}}{2 N B_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} g_{\epsilon}(v) g_{\epsilon}(w) \Delta \phi(v) \\
& \times\left(\lambda^{2}(e \cdot v)^{2}+\mu^{2}(e \cdot w)^{2}+2 \lambda \mu(e \cdot v)(e \cdot w)\right) d w d v d e \\
= & \frac{A_{\epsilon}}{2 N B_{N}} \int_{S^{N-1}} \int_{\mathbb{R}^{2 N}} g_{\epsilon}(v) g_{\epsilon}(w) \Delta \phi(v)\left(\lambda^{2}(e \cdot v)^{2}+\mu^{2}(e \cdot w)^{2}\right) d w d v d e \\
= & \frac{A_{\epsilon}}{2 N^{2}} \int_{\mathbb{R}^{2 N}} g_{\epsilon}(v) g_{\epsilon}(w) \Delta \phi(v)\left(\lambda^{2}|v|^{2}+\mu^{2}|w|^{2}\right) d w d v \\
= & \frac{A_{\epsilon}}{2 N^{2}} \int_{\mathbb{R}^{N}} g_{\epsilon}(v) \Delta \phi(v)\left(\lambda^{2}|v|^{2}+\mu^{2}\right) d v .
\end{aligned}
$$

Therefore, collecting terms, and taking the limit $\epsilon \rightarrow 0$ we find that the limit function $g(v, t)$ satisfies the equation (in weak form)

$$
\begin{gather*}
\frac{d}{d t} \int_{\mathbb{R}^{N}} \phi(v) g(v, t) d v+\frac{\left(\lambda^{2}+\mu^{2}\right) A_{\epsilon}}{2 N} \int_{\mathbb{R}^{N}} \nabla \phi(v) \cdot v g d v \\
=\frac{A_{\epsilon}}{2 N^{2}} \int_{\mathbb{R}^{N}} \Delta \phi g(v)\left(\lambda^{2}|v|^{2}+\mu^{2}\right) d v \tag{4.8}
\end{gather*}
$$

which is the weak formulation of the following Fokker-Planck equation

$$
\begin{equation*}
\frac{\partial}{\partial t} g(v, t)-\frac{\left(\lambda^{2}+\mu^{2}\right)}{2 N} \operatorname{div}(v g)=\Delta\left(\frac{\lambda^{2} v^{2}+\mu^{2}}{2 N^{2}} g(v)\right) \tag{4.9}
\end{equation*}
$$

As in the 1-dimensional case the coefficient in the convection term must to be one since we consider $\lambda^{2}+\mu^{2}=2 N$ and therefore

$$
\begin{equation*}
\frac{\partial}{\partial t} g(v, t)=\operatorname{div}(v g)+\Delta\left[\left(\frac{\lambda^{2}}{2 N^{2}} v^{2}+\frac{\mu^{2}}{2 N^{2}}\right) g(v)\right] . \tag{4.10}
\end{equation*}
$$

Again, as in the one dimensional case, it follows that the limit function $g(v, t)$ satisfies the Fokker-Planck Eq. (2.1), by taking

$$
\alpha^{2}=\frac{\mu^{2}}{2 N^{2}} \quad \text { and } \quad \beta^{2}=\frac{\lambda^{2}}{2 N^{2}}
$$

which satisfy $\alpha^{2}+\beta^{2}=1 / N$. We can now repeat the reasoning we did in onedimension of space, to conclude that the sequence $\left\{g_{\epsilon}(v, t)\right\}_{\epsilon \geq 0}$ contains a subsequence that converges weakly to a function $g$ satisfying (4.8).

### 4.2. Convergence in Fourier Based Metrics

We use the same notation as in one dimensional case and we denote by $\mathcal{M}_{0}$ the space of all probability measures in $\mathbb{R}^{N}$ and

$$
\mathcal{M}_{r}=\left\{\mu \in \mathcal{M}_{0}: \int_{\mathbb{R}^{N}}|v|^{r} \mu(d v)<\infty, r \geq 0\right\}
$$

the space of all Borel probability measures with finite momentum of order $r$ equipped with the topology of the weak convergence of the measures. A weak solution of the initial value problem for Eq. (4.2), corresponding to the initial probability density $f_{0}(w) \in \mathcal{M}_{r}, r>2$ we shall mean any probability density $f \in C^{1}\left(\mathbb{R}^{N}, \mathcal{M}_{r}\right)$ satisfying the weak form of the Eq. (4.3) for $t>0$ and all smooth functions $\phi$, and such that for all $\phi$

$$
\begin{equation*}
\lim _{t \rightarrow 0} \int_{\mathbb{R}} \phi(v) f(v, t) d v=\int_{\mathbb{R}} \phi(v) f_{0}(v) d v \tag{4.11}
\end{equation*}
$$

In the rest of this section, we shall study the weak form of Eq. (4.2), with the normalization conditions (1.10). It is equivalent to use the Fourier transform of the equation ${ }^{(4)}$ :

$$
\begin{equation*}
\frac{\partial \widehat{f}(\xi, t)}{\partial t}=\widehat{Q}(\widehat{f}, \widehat{f})(\xi, t) \tag{4.12}
\end{equation*}
$$

where

$$
\begin{align*}
& \widehat{Q}(\widehat{f}, \widehat{f})(\xi) \\
& \quad=\int_{C_{N}} \int_{S^{N-1}} \frac{B_{\epsilon}(n)}{B_{N}}[\hat{f}(\xi+(n \cdot \xi) \lambda e, t) \hat{f}((n \cdot \xi) \mu e, t)-\hat{f}(\xi, t)] d e d n \\
& \quad \times \int_{C_{N}} \int_{S^{N-1}} \frac{B_{\epsilon}(n)}{B_{N}} \hat{f}(\xi+(n \cdot \xi) \lambda e, t) \hat{f}((n \cdot \xi) \mu e, t) d e d n-C_{\epsilon}^{N} \hat{f}(\xi, t) \tag{4.13}
\end{align*}
$$

where

$$
C_{\epsilon}^{N}=\int_{C_{N}} B_{\epsilon}(n) d n
$$

In the $N$-dimensional case, the initial conditions (1.10) take the form

$$
\widehat{f_{0}}(\xi=0)=1, \quad \nabla \widehat{f_{0}}(\xi=0)=0, \quad \Delta \widehat{f_{0}}(\xi=0)=-1
$$

where $\widehat{f_{0}} \in C^{2}\left(\mathbb{R}^{N}\right)$.

Then, proceeding as in the one-dimensional case, ${ }^{(17)}$ given two initial data $f_{0,1}$ and $f_{0,2}$ such that $d_{r}\left(f_{0,1}, f_{0,2}\right)$ is bounded for some $r \geq 2$, one obtains for the corresponding solutions to ((4.3)) the estimate

$$
\begin{equation*}
d_{r}\left(f_{1}(t), f_{2}(t)\right) \leq d_{r}\left(f_{0,1}, f_{0,2}\right) \cdot \exp \left\{R_{N}^{\epsilon}(r) t\right\} \tag{4.14}
\end{equation*}
$$

where

$$
R_{N}^{\epsilon}(r)=\int_{C_{N}} \int_{S^{N-1}} \frac{B_{\epsilon}(n)}{B_{N}}\left[|\gamma+(n \cdot \gamma) \lambda e|^{r}+|(n \cdot \gamma) \mu e|^{r}-1\right] d e d n,
$$

and we denote with $\gamma=\xi /|\xi|$ a unit vector. The estimation (4.14) is translated to the analogous estimation for the scaled function $g_{\epsilon}$. As in the one-dimensional case

$$
d_{r}\left(g_{\epsilon_{1}}(t), g_{\epsilon_{2}}(t)\right)=M_{2}(t)^{-r / 2} d_{r}\left(f_{1}(t), f_{2}(t)\right) \leq d_{r}\left(f_{0,1}, f_{0,2}\right) \cdot \exp \left\{R_{N}^{\epsilon}(r) t\right\}
$$

or equivalently, since $d_{r}\left(f_{0,1}, f_{0,2}\right)=d_{r}\left(g_{0,1}, g_{0,2}\right)$

$$
\begin{equation*}
d_{r}\left(g_{\epsilon_{1}}(t), g_{\epsilon_{2}}(t)\right)=M_{2}(t)^{-r / 2} d_{r}\left(f_{1}(t), f_{2}(t)\right) \leq d_{r}\left(g_{\epsilon_{01}}, g_{\epsilon_{02}}\right) \cdot \exp \left\{S_{N}^{\epsilon}(r) t\right\} \tag{4.15}
\end{equation*}
$$

where now

$$
\begin{align*}
& S_{N}^{\epsilon}(r)=\int_{C_{N}} \int_{S^{N-1}} \frac{B_{\epsilon}(n)}{B_{N}} \\
& \quad \cdot\left[|\gamma+(n \cdot \gamma) \lambda e|^{r}+|(n \cdot \gamma) \mu e|^{r}-1-\frac{r}{2 N} n^{2}\left(\lambda^{2}+\mu^{2}\right)\right] d e d n \tag{4.16}
\end{align*}
$$

In order to pass to the limit $\epsilon \rightarrow 0$ we use a Taylor development for the function $h(v)=|v|^{2+\delta}$ with $v \in \mathbb{R}^{N}$,

$$
\begin{aligned}
h(v) & =h\left(v_{0}\right)+\nabla_{v} h\left(v_{0}\right) \cdot\left(v-v_{0}\right)+\frac{1}{2}\left(v-v_{0}\right) \operatorname{Hess}(h)\left(v_{0}\right)\left(v-v_{0}\right)^{t} \\
& +R\left(h, \tilde{v}, v_{0}, v\right)
\end{aligned}
$$

Since $v=\gamma+(n \cdot \gamma) \lambda e$ and $v_{0}=(n \cdot \gamma) \lambda e$, we use properties $\mathbf{P}_{\mathbf{N}} \cdot \mathbf{1}-\mathbf{P}_{\mathbf{N}} . \mathbf{3}$ and (4.5) into (4.16) to obtain

$$
\begin{aligned}
S_{N}^{\epsilon}(2+\delta)= & \frac{\lambda^{2}}{2} \frac{1}{N} A_{\epsilon}(2+\delta) \delta \frac{1}{N}+\frac{\lambda^{2}}{2 N} A_{\epsilon}(2+\delta)-A_{\epsilon} \frac{2+\delta}{2 N}\left(\lambda^{2}+\mu^{2}\right) \\
& +O\left(n^{2+\delta}\right)
\end{aligned}
$$

Taking the limit $\epsilon \rightarrow 0$ we find

$$
S_{N}(2+\delta)=\frac{2+\delta}{2 N}\left[\frac{\delta}{N} \lambda^{2}-\mu^{2}\right]
$$

Therefore, using Lemma 2.4 we obtain for the limit function $g$ the following inequality

$$
\begin{equation*}
d_{2+\delta}\left(g_{1}(t), g_{2}(t)\right) \leq d_{2+\delta}\left(g_{0_{1}}, g_{0_{2}}\right) \cdot \exp \left\{S_{N}(2+\delta) t\right\} \tag{4.17}
\end{equation*}
$$

which can be rewritten in terms of the general Fokker-Planck Eq. (2.1) for the stationary solution and a general solution as

$$
\begin{equation*}
d_{2+\delta}\left(f(t), B_{\infty}\right) \leq d_{2+\delta}\left(f_{0}, B_{\infty}\right) \cdot \exp \left\{2(2+\delta)\left(\delta \beta^{2}-2 N \alpha^{2}\right) t\right\} \tag{4.18}
\end{equation*}
$$

We are now in a position to prove the main result of this Section, namely the $N$-dimensional version of Theorem 3.3. We can observe that for the linearized Eq. (2.14) also considered in Ref. 14 we obtain the same rate of convergence taking $\delta=0$.

Theorem 4.2. Let $f(v, t)$ be the solution to Eq. (2.14), corresponding to the initial datum $f_{0}(v)$ satisfying (1.10), and let us denote by $\bar{B}=D B_{C}^{2-m}$ the stationary solution to the same equation. Suppose in addition that

$$
\begin{equation*}
d_{2+\delta}\left(f_{0}, \bar{B}\right)<+\infty, \quad \delta<(2-N+N m) /(1-m) \tag{4.19}
\end{equation*}
$$

Then, the solution $f(v, t)$ converges exponentially to $\bar{B}$ in $d_{2+\delta}$-metric, and

$$
\begin{equation*}
d_{2+\delta}(f(t), \bar{B}) \leq d_{2+\delta}\left(f_{0}, \bar{B}\right) \exp \left\{\frac{2+\delta}{2}((1-m) \delta-(2-N+N m)) t\right\} \tag{4.20}
\end{equation*}
$$

Moreover, provided $\frac{N}{N+2}<m<1$ the rate of convergence in (4.20) increases in the interval $0<\delta \leq((N+2) m-N) /(2-2 m)$. Likewise, let $f(v, t)$ be the solution to Eq. (1.6), with initial datum $f_{0}(v)$. Suppose in addition that

$$
\begin{equation*}
d_{2+\delta}\left(f_{0}, B_{C}\right)<+\infty, \quad \delta<(m(N+2)-N) /(1-m) \tag{4.21}
\end{equation*}
$$

Then, the solution $f(v, t)$ converges exponentially to $B_{C}(x)$ in $d_{2+\delta-m e t r i c, ~ a n d ~}^{\text {a }}$

$$
\begin{equation*}
d_{2+\delta}\left(f(t), B_{C}\right) \leq d_{2+\delta}\left(f_{0}, B_{C}\right) \exp \left\{\frac{2+\delta}{2 N}\left(\frac{1-m}{m} \delta-\frac{(2+N) m-N}{m}\right) t\right\} \tag{4.22}
\end{equation*}
$$

Moreover, provided $(N+2) /(N+4)<m<1$ the rate of convergence in (4.22) increases in the interval $0<\delta \leq \frac{(N+4) m-(N+2)}{2-2 m}$.

Proof: The proof of this theorem is an easy generalization of the 1-dimensional case and we can obtain it simply by identifying coefficients, since in Sec. 2.1 we showed that both linear equations ((2.14) and (1.6)) can be written, using different
scaling, as Eq. (2.1). For Eq. (2.14) we have

$$
\alpha^{2}=\frac{\mu^{2}}{2 N^{2}}=\frac{1}{N}-\frac{1-m}{2} \quad \text { and } \quad \beta^{2}=\frac{\lambda^{2}}{2 N^{2}}=\frac{1-m}{2}
$$

$S_{N}(2+\delta)$ is nonpositive if and only if $\delta \leq \frac{N \mu^{2}}{\lambda^{2}}$. And the minimum value is achieved at $\delta_{0}=\frac{N \mu^{2}}{2 \lambda^{2}}-1$ which is positive if and only if $N \mu^{2}>2 \lambda^{2}$ or equivalently if and only if $N \alpha^{2}>2 \beta^{2}$ (analogous N -dimensional version of (3.41)). And finally, this inequality holds if and only if $m>\frac{N}{N+2}$. On the other hand, since

$$
\mu^{2}=2 N \frac{2-N+N m}{2} \quad \text { and } \quad \lambda^{2}=2 N^{2} \frac{1-m}{2}
$$

we obtain that the rate of convergence $\frac{2+\delta}{2 N}\left(\frac{\lambda^{2}}{N} \delta-\mu^{2}\right)$ is written as

$$
\frac{2+\delta}{2}((1-m) \delta-(2-N+N m))
$$

The minimum value is achieved at $\delta_{0}=\frac{(2+N) m-N}{2-2 m}$ (see previous computations) and this value is

$$
-\frac{1}{2^{3}} \frac{(4-N+(N-2) m)^{2}}{1-m} .
$$

For Eq. (1.6) we have

$$
\alpha^{2}=\frac{\mu^{2}}{2 N^{2}}=\frac{1}{N}-\frac{1-m}{2 m} \quad \text { and } \quad \beta^{2}=\frac{\lambda^{2}}{2 N^{2}}=\frac{1-m}{2 m}
$$

and in this case $N \alpha^{2}>2 \beta^{2}$ holds if and only if $m>(N+2) /(N+4)$ and the rate of convergence is written as

$$
\frac{2+\delta}{2 N}\left(\frac{1-m}{m} \delta-\frac{(2+N) m-N}{m}\right) .
$$

The minimum value is now achieved in $\delta_{0}=\frac{(N+4) m-(N+2)}{2-2 m}$ and this value is

$$
-\frac{1}{2^{3}} \frac{((2-N)+N m)^{2}}{N m} .
$$

Remark 4.3. For the Eq. (2.14) $\delta<1$ (which means only equality for the initial mass) if

$$
\frac{N-2}{N}<m<\frac{N-1}{N+1}
$$

and $\delta_{o}<1$ (which means equality of the second moments) if

$$
\frac{N}{N+2}<m<\frac{2+N}{N+4}
$$

While for (1.6) $\delta<1$ if

$$
\frac{N}{N+2}<m<\frac{N+1}{N+3}
$$

and $\delta_{o}<1$ if

$$
\frac{N+2}{N+4}<m<\frac{N+4}{N+6}
$$

Some comments to Theorem 4.2 are needed. Despite the one-dimensional situation, the boundedness of the Fourier based metric for $s>2$ requires that all the cross moments of both the initial value $f_{0}$ and the stationary state are equal up to $[s]$, where $[s]$ is the entire part of the number $s$. When $s>2$, this condition reads very strong. Thus, the conditions of Theorem 4.2 are equivalent to thats of Theorem 3.3 only for radially symmetric initial data, which, as it can be easily seen looking at the kinetic model, imply that the solution is radially symmetric at any subsequent time. In this case, in fact, only the principal moments of both the initial value and the stationary state need to be equal.

In the general situation, however, it has been shown first for the elastic spatially homogeneous Boltzmann equation for Maxwell molecules, ${ }^{(17)}$ and subsequently for the inelastic Boltzmann equation for Maxwell type interactions, ${ }^{(3)}$ that the explicitly computable and fast decay to zero of the cross moments can be used to relax the condition of Theorem 4.2 to all the initial data satisfying (1.10). We will outline this strategy in the next section.

Second, as shown in paper ${ }^{(5)}$ for the elastic spatially homogeneous Boltzmann equation for Maxwell molecules, the decay of the Fourier metric, combined with the uniform boundedness of stronger norms, gives rise to precise rates of convergence towards the stationary state in these strong norms. This idea has been recently applied to the spatially uniform Boltzmann equation for inelastic Maxwell model in presence of a thermal bath in Ref 2. In the purely inelastic setting, however, the lack of an entropy principle, makes it difficult to extend the results of Ref. 5 to this case. ${ }^{(3)}$ In the next Section we will give indications on a possible way to obtain decay rates in stronger norm for the target Fokker-Planck equations.

## 5. RATE OF CONVERGENCE FOR GENERAL INITIAL VALUES

### 5.1. Propagation of Regularity

The goal of this section is to show that the smoothness of the initial data of the Fokker-Planck Eq. (2.1) is propagated, so that we have bounds on the smoothness uniform in time. This result, combined with suitable interpolation inequalities ${ }^{(5)}$ can be used to obtain convergence in stronger space with almost the same rate of convergence of the metric $d_{2}$. This strategy has been recently applied to the dissipative Boltzmann equation in a thermal bath in Ref. 2.

After application of the Fourier transform, the Fokker-Planck Eq. (2.1) takes the form

$$
\begin{equation*}
\frac{\partial \hat{f}}{\partial t}=-\xi \cdot \nabla_{\xi} \hat{f}-\alpha^{2}|\xi|^{2} \hat{f}+\beta^{2}|\xi|^{2} \Delta_{\xi} \hat{f} . \tag{5.23}
\end{equation*}
$$

To simplify computations, thanks to linearity, we can consider separately the real and imaginary parts of $\hat{f}(\xi), \hat{f}(\xi)=a(\xi)+i b(\xi)$. Both the functions $a$ and $b$ then satisfy the same Eq. (5.23). In the following we argument for $a$ but clearly the same holds for $b$. Multiplying the Eq. (5.23) satisfied by the real part $a(\xi)$ by $2|\xi|^{2 k} a$ one has

$$
\frac{\partial a^{2}}{\partial t}|\xi|^{2 k}=-|\xi|^{2 k} \xi \cdot \nabla_{\xi} a^{2}-2 \alpha^{2}|\xi|^{2 k+2} a^{2}+2 \beta^{2}|\xi|^{2 k+2} a \Delta_{\xi} a .
$$

Therefore, integrating over $\mathbb{R}^{N}$ we obtain

$$
\begin{gather*}
\frac{d}{d t} \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2} d \xi=-\int_{\mathbb{R}^{N}}|\xi|^{2 k} \xi \cdot \nabla_{\xi} a^{2} d \xi \\
-2 \alpha^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2} a^{2} d \xi+2 \beta^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2} a \Delta_{\xi} a d \xi \tag{5.24}
\end{gather*}
$$

Direct computations, integrating by parts produce

$$
\int_{\mathbb{R}^{N}}|\xi|^{2 k} \xi \cdot \nabla_{\xi} a^{2} d \xi=-(N+2 k) \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2} d \xi
$$

and
$\int_{\mathbb{R}^{N}}|\xi|^{2 k+2} a \Delta_{\xi} a d \xi=(k+1)(N+2 k) \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2} d \xi-\int_{\mathbb{R}^{N}}|\xi|^{2 k+2}\left(\nabla_{\xi} a\right)^{2} d \xi$.
Hence (5.24) is written as

$$
\begin{align*}
& \frac{d}{d t} \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2} d \xi=(N+2 k)\left[1+2 \beta^{2}(k+1)\right] \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2} d \xi \\
- & 2 \alpha^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2} a^{2} d \xi-2 \beta^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2}\left(\nabla_{\xi} a\right)^{2} d \xi . \tag{5.25}
\end{align*}
$$

Since $|\hat{f}|^{2}=a^{2}+b^{2}$, we obtain

$$
\begin{align*}
& \frac{d}{d t} \int_{\mathbb{R}^{N}}|\xi|^{2 k}|\hat{f}|^{2} d \xi=(N+2 k)\left[1+2 \beta^{2}(k+1)\right] \int_{\mathbb{R}^{N}}|\xi|^{2 k}|\hat{f}|^{2} d \xi \\
& \quad-2 \alpha^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2}|\hat{f}|^{2} d \xi-2 \beta^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2}\left[\left(\nabla_{\xi} a\right)^{2}+\left(\nabla_{\xi} b\right)^{2}\right] d \xi \tag{5.26}
\end{align*}
$$

The next step is to find suitable lower bounds for the two negative terms in (5.26). It is almost clear that the bounds we will derive in what follows are not new, and one can find them in the literature. Due to the simplicity of the proofs, however, we will present them in some detail. We have

Lemma 5.1. Let $f \in L^{1}\left(\mathbb{R}^{N}\right)$ be a probability density function such that, for some $k>0$

$$
\begin{equation*}
\int_{\mathbb{R}^{N}}|\xi|^{2 k+2}\left[\left(\nabla_{\xi} a\right)^{2}+\left(\nabla_{\xi} b\right)^{2}\right] d \xi<\infty . \tag{5.27}
\end{equation*}
$$

Then, $\int|\xi|^{2 k}|\hat{f}(\xi)|^{2}$ is bounded, and the following inequality holds

$$
\begin{equation*}
\int_{\mathbb{R}^{N}}|\xi|^{2 k}|\hat{f}(\xi)|^{2} d \xi \leq \frac{4}{(N+2 k)^{2}} \int_{\mathbb{R}^{N}}|\xi|^{2 k+2}\left[\left(\nabla_{\xi} a\right)^{2}+\left(\nabla_{\xi} b\right)^{2}\right] d \xi \tag{5.28}
\end{equation*}
$$

Proof: Let $\lambda$ be a positive constant. Then

$$
\begin{align*}
0 & \leq \int_{\mathbb{R}^{N}}|\xi|^{2 k}\left|\lambda \frac{\xi}{|\xi|} a(\xi)+|\xi| \nabla a(\xi)\right|^{2} d \xi=\lambda^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2}(\xi) d \xi \\
& +\int_{\mathbb{R}^{N}}|\xi|^{2 k+2}|\nabla a(\xi)|^{2} d \xi+2 \lambda \int_{\mathbb{R}^{N}}|\xi|^{2 k} a(\xi) \xi \cdot \nabla a(\xi) d \xi \tag{5.29}
\end{align*}
$$

The last integral into (5.29) can be integrated by parts to give

$$
\begin{aligned}
2 \lambda \int_{\mathbb{R}^{N}}|\xi|^{2 k} a(\xi) \xi \nabla a(\xi) d \xi & =\lambda \int_{\mathbb{R}^{N}}|\xi|^{2 k} \xi \cdot \nabla a^{2}(\xi) d \xi \\
& =-\lambda(N+2 k) \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2}(\xi) d \xi
\end{aligned}
$$

Substituting into (5.29) gives

$$
\begin{equation*}
\int_{\mathbb{R}^{N}}|\xi|^{2 k+2}|\nabla a(\xi)|^{2} d \xi \geq \lambda(N+2 k) \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2}(\xi) d \xi-\lambda^{2} \int_{\mathbb{R}^{N}}|\xi|^{2 k} a^{2}(\xi) d \xi \tag{5.30}
\end{equation*}
$$

Choosing $\lambda=(N+2 k) / 2$ into (5.30) proves inequality (5.28).
Let us introduce, for $r \geq 0$, the Sobolev space norms $\|\cdot\|_{H^{r}\left(\mathbb{R}^{N}\right)}$ by

$$
\|f\|_{H^{r}\left(\mathbb{R}^{N}\right)}^{2}=\int_{\mathbb{R}^{N}}|\xi|^{2 r}|\hat{f}(\xi)|^{2} d \xi
$$

Using inequality (5.28) into inequality (5.26), we can rewrite it as

$$
\begin{equation*}
\frac{d}{d t}\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2} \leq D\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2}-2 \alpha^{2}\|f\|_{H^{k+1}\left(\mathbb{R}^{N}\right)}^{2} \tag{5.31}
\end{equation*}
$$

where the constant $D$ holds

$$
\begin{equation*}
D=D(N, k, \beta)=(N+2 k)\left(1+2 \beta^{2}(k+1)\right)-\beta^{2} \frac{(N+2 k)^{2}}{2} . \tag{5.32}
\end{equation*}
$$

From the differential inequality (5.31) we deduce the uniform boundedness of $\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}$. The result follows by the following Nash-type inequality.

Lemma 5.2. Let $f \in L^{1}\left(\mathbb{R}^{N}\right)$ be a probability density function such that, for some $q>0$

$$
\int_{\mathbb{R}^{N}}|\xi|^{q}|\hat{f}(\xi)|^{2} d \xi<\infty
$$

Then, if $0<p<q, \int|\xi|^{p}|\hat{f}(\xi)|^{2}$ is bounded, and the following inequality holds

$$
\begin{equation*}
\int_{\mathbb{R}^{N}}|\xi|^{p}|\hat{f}(\xi)|^{2} d \xi \leq C(p, q, N)\left(\int_{\mathbb{R}^{N}}|\xi|^{q}|\hat{f}(\xi)|^{2} d \xi,\right)^{d(p, q, N)} \tag{5.33}
\end{equation*}
$$

where, denoting by $\left|B_{N}\right|$ the measure of the unit ball in $\mathbb{R}^{N}$,

$$
C(p, q, N)=\left|B_{N}\right|^{\frac{q-p}{N+q}}\left(\frac{q-p}{p+N}\right)^{\frac{p+N}{q+N}} \frac{q+N}{q-p}
$$

and

$$
d(p, q, N)=\frac{p+N}{q+N}<1 .
$$

Proof: Since $f$ is a probability density, $|\hat{f}| \leq 1$. Therefore, for any constant $R>0$ the following bound holds,

$$
\begin{aligned}
\int_{\mathbb{R}^{v}}|\xi|^{p}|\hat{f}(\xi)|^{2} d \xi & \leq R^{p} \int_{|\xi| \leq R}|\hat{f}|^{2} d \xi+R^{p-q} \int_{|\xi|>R}|\xi|^{q}|\hat{f}|^{2} d \xi \\
& \leq R^{p} \int_{|\xi| \leq R} d \xi+R^{p-q} \int_{|\xi|>R}|\xi|^{q}|\hat{f}|^{2} d \xi
\end{aligned}
$$

$$
\begin{equation*}
\leq R^{N+p}\left|B_{N}\right|+R^{p-q} \int_{\mathbb{R}^{N}}|\xi|^{q}|\hat{f}|^{2} d \xi \tag{5.34}
\end{equation*}
$$

Optimizing in $R$ we obtain the result.
We can use Lemma 5.2 to obtain a bound in (5.31). To this aim, considering $p=2 k$ and $q=2 k+2$

$$
\begin{aligned}
\frac{d}{d t}\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2} & \leq D\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2}-2 \alpha^{2} C(2 k, 2 k+2, N)^{-1 / d(2 k, 2 k+2, N)}\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2 / d(2 k+2, N)} \\
& \leq D\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2}\left(1-2 \frac{\alpha^{2}}{D C^{1 / d}}\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2(1-d) / d}\right)
\end{aligned}
$$

Since $d<1$, the time derivative can growth only if the right-hand side remains positive, that is when

$$
\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2(1-d) / d} \leq \frac{D C^{1 / d}}{2 \alpha^{2}} .
$$

Clearly this implies a uniform bound on $\|f\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2}$. We proved
Theorem 5.5. Let $f_{0}$ be the initial datum for Eq. (2.1) such that $\left\|f_{0}\right\|_{H^{k}\left(\mathbb{R}^{N}\right)}$ finite. Then, any probability density solution $f(t, v)$ of $(2.1)$ is bounded in $H^{k}\left(\mathbb{R}^{N}\right)$, and there is a universal constant $D$ so that, for all $t>0$,

$$
\|f(t)\|_{H^{k}\left(\mathbb{R}^{N}\right)} \leq \max \left\{\left\|f_{0}\right\|_{\left.H^{k} \mathbb{R}^{N}\right)}, B\right\},
$$

where

$$
B=\left(\frac{D C^{1 / d}}{2 \alpha^{2}}\right)^{\frac{d}{2(1-d)}}
$$

$C(2 k, 2 k+2, N)$ and $d(2 k, 2 k+2, N)$ are defined as in Lemma 5.2, while the constant $D$ is defined in (5.32).

### 5.2. Strong Convergence to Equilibrium

This section is very short. The regularity result proven in the previous section can be used together with Theorem 4.2 to obtain rates of convergence to equilibrium in stronger norms. We only need to collect results and to explain how to compute the constants involved in it. The key point are some interpolation inequalities, recently considered in Ref. 5 to obtain strong convergence results for the elastic Boltzmann equation for Maxwell molecules. These inequalities have been subsequently been used for obtaining sharp rates of convergence to equilibrium for the dissipative Boltzmann equation in presence of a thermal bath. ${ }^{(2)}$ The first of these inequalities reads

Lemma 5.4. [Control of the $\boldsymbol{H}^{k}$-distance] Let $k, p \geq 0$, then

$$
\|f-g\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2} \leq C d_{2}(f, g)^{\beta_{1}}\|f-g\|_{H^{k+p}\left(\mathbb{R}^{N}\right)}^{\beta_{2}},
$$

with

$$
\begin{aligned}
\beta_{1} & =\frac{4 p}{2 k+2 p+N+4}, \\
\beta_{2} & =\frac{2(2 k+N+4)}{2 k+2 p+N+4}, \\
C(k, p, N) & =\left(\frac{B_{N}}{2 k+N+4}\right)^{\frac{\beta_{1}}{2}}(2 p)^{\frac{\beta_{2}}{2}}+\left(\frac{2 p(2 k+N+4)}{B_{N}}\right)^{\frac{-\beta_{1}}{2}}
\end{aligned}
$$

Proof: For $R>0$ we obtain

$$
\begin{aligned}
\|f-g\|_{H^{k}\left(\mathbb{R}^{N}\right)}^{2} & \leq \int_{|\xi| \leq R}|\xi|^{2 k}|\hat{f}-\hat{g}|^{2} d \xi+\frac{1}{R^{2 p}} \int_{|\xi|>R}|\xi|^{2(k+p)}|\hat{f}-\hat{g}|^{2} d \xi \\
& \leq \int_{|\xi| \leq R}|\xi|^{2(k+2)} d_{2}^{2}(f, g) d \xi+\frac{1}{R^{2 p}}\|f-g\|_{H^{k+p}}^{2}
\end{aligned}
$$

and optimizing on $R$ we obtain the result.
This result shows that the weak $d_{2}$ distance coupled with $\|\cdot\|_{H^{r}}$ smoothness, controls the $H^{k}$ distance for $r$ sufficiently larger than $k$. The next inequality shows that control of the sufficiently many moments and control of the $L^{2}$ norm together, control the $L^{1}$ norm.

Lemma 5.5. [Control of the $\boldsymbol{L}^{1}$-distance ${ }^{(5)}$ (Theorem 4.2)] Let $f$ be an integrable function on $\mathbb{R}^{N}$. Then, for all $r>0$
$\int_{\mathbb{R}^{N}}|f(v)| d v \leq C(N, r)\left(\int_{\mathbb{R}^{N}}|f(v)|^{2} d v\right)^{2 r /(N+4 r)}\left(\int_{\mathbb{R}^{N}}|v|^{2 r}|f(v)| d v\right)^{N /(N+4 r)}$ with

$$
C(N, r)=\left[\left(\frac{N}{4 r}\right)^{4 r /(N+4 r)}+\left(\frac{4 r}{N}\right)^{N /(N+4 r)}\right]\left|B^{N}\right|^{2 r /(N+4 r)} .
$$

Using both Theorem 5.5 and 5.4 with $r=2$, we prove

Theorem 5.6. Let $f_{0}$ be any initial datum for the Fokker-Planck Eq. ((2.1)) satisfying conditions ((1.10)). Suppose in addition that the conditions of Theorem
(4.2) are satisfied, so that there exists a positive constant $\lambda$ such that for some $\delta>0$ there is convergence of the solution $f(v, t)$ in $d_{2+\delta}$-norm towards the steady state $f_{\infty}(v)$ at the exponential rate $\lambda$. Let $\epsilon>0$ be given. Then, there is a number $r$ depending only on $\epsilon$ so that whenever

$$
\int_{\mathbb{R}^{N}}|v|^{2} f_{0}(v) d v+\int_{\mathbb{R}^{N}}|k|^{2 r}\left|\hat{f}_{0}(k)\right|^{2} d k<\infty
$$

then it holds that

$$
\begin{equation*}
\left\|f(\cdot, t)-f_{\infty}\right\|_{L^{1}\left(\mathbb{R}^{N}\right)} \leq C_{\epsilon} \exp \{-(1-\epsilon) \lambda t\} \tag{5.35}
\end{equation*}
$$

$C_{\epsilon}$ is explicitly computable in terms of the integrals specified above. Moreover, increasing $r$, we obtain the same result if the $L^{1}$-norm is replaced by any $H^{m}$-norm.

Proof: The result follows from the uniform propagation of regularity shown in Theorem 5.3, and from the interpolation inequalities, Lemmas 5.4 and 5.5 above.

## 6. CONCLUSIONS

We discussed in this paper the approach to equilibrium for various FokkerPlanck equations with variable coefficient of diffusion. These equations separate in a distinctive way from other Fokker-Planck equations with variable coefficients for the reason that they are obtained by linearizing scaled fast diffusion equations. We investigated the approach to equilibrium by using suitable kinetic models of the Boltzmann equation corresponding to nonconservative collision interactions. These interactions are in fact such that the energy is produced by binary collisions. This allows to use for the reckoning of rates of decay Fourier-based metrics which are easy to handle in this kinetic framework. From one side, this seems artificial, since one has to use nonlinear kinetic equations to obtain results on a linear equation. Looking however to previous attempts to the solution of this problem, an approach based on the nonlinear diffusion equation has been used in Ref. 8 to obtain Poincaré inequalities for the linearized fast diffusion equation. Maybe some new approach has to be dealt with, to treat in full generality the problem of convergence to equilibrium for these Fokker-Planck equations.

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