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# When are Poisson and Moyal Brackets equal? 

Didier Robert*


#### Abstract

In this paper we prove that if a given smooth Hamiltonian $H$ on the phase space $\mathbb{R}^{2 d}$, with derivatives of moderate growth, satisfies $\{A, H\}=\{A, H\}_{\circledast}$ for any observable $A$ in the Schwartz space $\mathcal{S}\left(\mathbb{R}^{2 d}\right)$, then, as it is expected, $H$ must be a polynomial of degree at most 2 in $\mathbb{R}^{2 d}$. Here $\{A, B\}$ denotes the Poisson bracket of two smooth classical observables and $\{A, B\}_{\circledast}$ their Moyal bracket, defined as the Weyl symbol of $i[\hat{A}, \hat{B}]$, where $\widehat{A}$ is the Weyl quantization of $A$ and $[\widehat{A}, \widehat{B}]=$ $\widehat{A} \widehat{B}-\widehat{B} \widehat{A}$ (commutator). A related answer to this question is given in the Groenewold-van Hove Theorem [9] concerning quantization of polynomial observables. We revisit the paper [1] where the authors considered an algebraic formal approach of this question for the time dependent evolution of observables.


## 1 Introduction

Let $H, A, B$ be smooth classical observables on $\mathbb{R}^{2 d}$ in the variables $X=$ $(x, \xi)$. The Poisson brackets is defined as $\{A, B\}=\partial_{\xi} A \cdot \partial_{x} B-\partial_{x} A \cdot \partial_{\xi} B$. So the classical time evolution of $A$ determined by the Hamilton equation for $H$ is solution of the equation:

$$
\begin{align*}
\frac{d}{d t} A(t) & =\{A(t), H\}  \tag{1.1}\\
A(0) & =A .
\end{align*}
$$

[^0]The Weyl quantization $\widehat{A}$ of $A$ is defined as the following operator, depending on the Planck constant $\hbar>0$,

$$
\begin{equation*}
\widehat{A} f(x):=\left(\mathrm{Op}_{\hbar}^{w} A\right) f(x)=(2 \pi \hbar)^{-d} \int_{\mathbb{R}^{2 d}} A\left(\frac{x+y}{2}, \xi\right) \mathrm{e}^{i \xi \cdot(x-y) / \hbar} f(y) d y d \xi \tag{1.2}
\end{equation*}
$$

for any $f \in \mathcal{S}\left(\mathbb{R}^{d}\right)$. Recall that $f \in \mathcal{S}\left(\mathbb{R}^{d}\right)$ means that $f \in C^{\infty}\left(\mathbb{R}^{d}\right)$ and for any multiindex $\alpha, \beta, x^{\alpha} \partial_{x}^{\beta} f(x)$ is bounded on $\mathbb{R}^{d}$.
The quantum time evolution of the quantum observable $\widehat{A}$ must satisfy the Heisenberg equation

$$
\begin{align*}
\frac{d}{d t} \widehat{A}(t) & =\frac{i}{\hbar}[\widehat{A}(t), \widehat{H}]  \tag{1.3}\\
\widehat{A}(0) & =\widehat{A} \tag{1.4}
\end{align*}
$$

where $[\hat{A}, \hat{B}]=\hat{A} \hat{B}-\hat{B} \hat{A}$.
The Moyal bracket of the observables $A, H$, is defined such that

$$
\begin{equation*}
\frac{i}{\hbar}[\widehat{A}, \widehat{H}]=\operatorname{Op}_{\hbar}^{w}\left(\{A, H\}_{\circledast}\right) . \tag{1.5}
\end{equation*}
$$

Notice that it results from the Weyl quantization calculus with a small parameter $\hbar$ that we have

$$
\lim _{\hbar \searrow 0}\{A, H\}_{\circledast}=\{A, H\} .
$$

A natural question is to ask when the classical dynamics generated by the Hamiltonian $H$ (1.1) has an exact correspondence with the quantum dynamics generated by $H$ (1.3) (see below the quotation from Van Hove). In the correspondence principle stated by N. Bohr the Planck constant $\hbar$ is supposed to be small. The question discussed here is for fixed $\hbar>0$.

A well known trick to check the correspondence Bohr principle is to compute the time evolution of Gaussian coherent states. Let us denote $\varphi_{Y}=$ $\hat{T}_{Y} \varphi_{0}$ the coherent state center in $Y \in \mathbb{R}^{2 d}$ and $\varphi_{0}(x)=(\pi \hbar)^{-d / 4} \mathrm{e}^{-|x|^{2} / 2 \hbar}$ ( $\hat{T}_{Y}$ is defined in the next section). We have [3]

$$
\lim _{\hbar \searrow 0}\left\langle\varphi_{Y}, \hat{A} \varphi_{Y}\right\rangle=A(Y) .
$$

Hence taking the average of (1.3) on $\varphi_{Y}$ and passing to the limit $\hbar \searrow 0$, we recover (1.1).

To define the Moyal bracket, there is a more explicit definition by introducing the Moyal product $A \circledast B$ (see the next section) such that

$$
\left(\mathrm{Op}_{\hbar}^{w} A\right)\left(\mathrm{Op}_{\hbar}^{w} B\right)=\mathrm{Op}_{\hbar}^{w}(A \circledast B)
$$

Then we have

$$
\{A, B\}_{\circledast}=\frac{i}{\hbar}(A \circledast B-B \circledast A) .
$$

These definitions make sense for $A, B \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$ and can be extended to suitable classes of symbols with moderate growth. To be more explicit we introduce the classes $\mathbb{S}_{\delta}^{\mu}$, for $\delta<1, \mu \in \mathbb{R} . A \in \mathbb{S}_{\delta}^{\mu}$ iff $A \in C^{\infty}\left(\mathbb{R}^{2 d}\right)$ and for any multiindex $\gamma \in \mathbb{N}^{2 d}$ we have:

$$
\left|\partial_{X}^{\gamma} A(X)\right| \leq C_{\gamma}\langle X\rangle^{\mu+\delta|\gamma|}
$$

Using Theorem A. 1 in [2], we can see that $A \circledast H$ is a smooth symbol if $H \in \mathbb{S}_{\delta}^{\mu}$ and $A \in \mathbb{S}_{\delta}^{\nu}$ where $\mu, \nu \in \mathbb{R}$ and $\delta<1 / 2$. Our aim here is to prove the following result.

Theorem 1.1. Assume that $\hbar>0$ is fixed (for example $\hbar=1$ ). Let be $H \in \mathbb{S}_{\delta}^{\mu}$ for some $\mu \in \mathbb{R}$ and $\delta<1 / 2$. Assume that for any $A \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$ we have $\{A, H\}_{\circledast}=\{A, H\}$. Then $H(X)$ must be a polynomial in $X=(x, \xi)$ of degree at most 2.

Remark 1.2. It is well known that if $H$ is a polynomial of degree at most 2 then $\{A, H\}_{\circledast}=\{A, H\}$ for any $A \in \mathbb{S}_{0}^{\nu}$. This is an easy consequence of the Weyl calculus.

Remark 1.3. The usual proofs of the Groenewold-van Hove Theorem on the phase space $\mathbb{R}^{2 d}$ concern more general quantization procedures but are restricted to polynomial symbols $A, H$.
A quotation from [9] p.66-67:
"On établit ensuite qu'une correspondance biunivoque entre grandeurs classiques et quantiques, ayant le caractère d'un isomorphisme entre algèbres de Lie, existe entre les grandeurs représentées par des polynômes de degré 0,1 , 2 en les variables $p_{1}, \cdots p_{N}, q_{1}, \cdots q_{n}$ mais ne peut être étendue sans perdre ses propriétés essentielles à l'ensemble de toutes les grandeurs classiques" The Theorem of Groenewold-van Hove is detailed $p .76$ and the quadratic case p. 87 of [9].

Notice that the quadratic case is related with the metaplectic representation [5].

Let us consider now the time dependent consequence of theorem 1.1. Assume that $H$ is real and generates a classical flow $\Phi_{H}^{t}$ on $\mathbb{R}^{2 d}$ such that for any $\gamma \in \mathbb{N}^{2 d}$ we have

$$
\left|\partial_{X}^{\gamma} \Phi_{H}^{t}(X)\right| \leq C_{\gamma}\langle X\rangle^{m}, \forall X \in \mathbb{R}^{2 d},|t| \leq \epsilon_{0},
$$

for some $m \geq 0$ and $\epsilon_{0}>0$.
The solution of (1.1) is $A_{\mathrm{cl}}(t, X):=A\left(\Phi_{H}^{t}(X)\right)$ and for $|t| \leq \epsilon_{0}$ with $A(t) \in$ $\mathcal{S}\left(\mathbb{R}^{2 d}\right)$ if $A(0) \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$.
On the quantum side assume that $\hat{H}$ is essentially self-adjoint in $L^{2}\left(\mathbb{R}^{d}\right)$. The solution of the Heisenberg equation (1.3) is $\hat{A}_{\text {quant }}(t)=\mathrm{e}^{-\frac{i}{\hbar} \hat{H}} \hat{A} \mathrm{e}^{\frac{i}{\hbar} \hat{H}}$.

Corollary 1.4. Under the conditions of Theorem 1.1, if for some $\hbar>0$ we have

$$
\hat{A}_{\text {qant }}(t)=\mathrm{op}_{\hbar}^{w}\left(A_{\mathrm{cl}}(t)\right), \forall A \in \mathcal{S}\left(\mathbb{R}^{2 d}\right), \quad|t|<\epsilon_{0},
$$

if and only if $H(X)$ is a polynomial in $X$ of degree at most 2.
Remark 1.5. In [1] the authors consider the semi-classical expansion ( $\hbar \searrow$ 0) in the Egorov theorem using a formal algebraic approach, and Bargman transform. In the Appendix $C$ we give a direct proof of the main result of [1] which is stated below (proposition 1.6).

Recall that according the semiclassical Egorov theorem the following asymptotic formula holds true as $\hbar \searrow 0$ :

$$
\begin{equation*}
\hat{A}_{\text {quant }}(t):=\mathrm{e}^{-\frac{i}{\hbar} \hat{H}} \hat{A} \mathrm{e}^{\frac{i}{\hbar} \hat{H}}=\sum_{0 \leq k \leq N} \hbar^{2 j} \hat{A}_{2 j}(t)+O\left(\hbar^{2 N+2}\right), \tag{1.6}
\end{equation*}
$$

where the symbol of the coefficient $A_{2 j}$ can be computed by induction [2] and C.1.

Proposition 1.6. [1] Let us consider the formal serie of symbols in $\hbar$ :
$\sum_{j \geq 0} A_{2 j}(t, X) \hbar^{2 j}$ where the coefficient $A_{2 j}(t, X)$ are defined by (C.2).
If for some $j_{0} \geq 1$ we have $A_{2 j_{0}}(t, X) \equiv 0$ then $A_{2}(t, X) \equiv 0$ and $H$ must be a polynomial of degree at most 2. In particular if in a neigborhood of 0 in $t$ we have $\hat{A}_{\text {quant }}(t)=\hat{A}_{\mathrm{cl}}(t)+O\left(\hbar^{3}\right)$ then $H$ must be a polynomial of degree at most 2.

Remark 1.7. The above Proposition may be a little surprising: if one term of order $\geq 2$ in the asymptotic expansion (1.6) vanishes then $H$ must be a polynomial of degree $\leq 2$ and all the terms vanish excepted the term of order 0 .

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## 2 Weyl calculus

### 2.1 Introduction to the Weyl quantization

In this section, we recall some basic properties of the Weyl calculus (for more details see [7]).
Weyl quantization starts by quantization of exponent of linear forms $L_{Y}(X)=$ $\sigma(Y, X)=\eta \cdot x-y \cdot \xi$ with $X=(x, \xi), Y=(y, \eta)$. Apart the usual properties asked for an admissible quantization, Weyl quantization is uniquely determined by imposing that the Weyl symbol of $\mathrm{e}^{i \widehat{L_{Y}}}$ is $\mathrm{e}^{i L_{Y}}$. Recall that $\widehat{T}(Y):=\mathrm{e}^{-i \widehat{L_{Y}}}$ is the Weyl-Heisenberg translation operator by $Y$ in the phase space $\mathbb{R}^{2 d}$. In other words the Weyl quantization $A \mapsto \hat{A}$ has to satisfy $\mathrm{e}^{i \widehat{L_{Y}}}=\widehat{\left(\mathrm{e}^{i L_{Y}}\right)}$. Then for any observable $A$, using a Fourier transform, the Weyl quantization $A$ is defined for any $\psi \in \mathcal{S}\left(\mathbb{R}^{d}\right)$, as

$$
\begin{equation*}
\widehat{A} \psi=(2 \pi)^{-d} \int_{\mathbb{R}^{2 d}} \tilde{A}_{\sigma}(Y) \hat{T}(Y) \psi d Y \tag{2.1}
\end{equation*}
$$

where $\tilde{A}_{\sigma}(Y)=\int_{\mathbb{R}^{2 d}} A(z) \mathrm{e}^{-i \sigma(Y, z)} d z$ is the symplectic Fourier transform of $A$ (in the sense of distributions). So that the family $\{T(Y)\}_{Y \in \mathbb{R}^{2 d}}$ is an overcomplete basis for operators between the Schwartz spaces $\mathcal{S}\left(\mathbb{R}^{d}\right)$ and $\mathcal{S}^{\prime}\left(\mathbb{R}^{d}\right)$. $\tilde{A}_{\sigma}$ is the covariant symbol of $\hat{A}$ and $A$ the contravariant symbol of $\hat{A}$.

Remark 2.1. Notice that from (2.1) for any symbol $A$ and for any linear form $L_{Z}$ we get

$$
\begin{equation*}
i\left[\hat{A}, \hat{L}_{Z}\right]=\left\{\widehat{A, L_{Z}}\right\} \tag{2.2}
\end{equation*}
$$

It is enough to prove (2.2) for $\hat{A}=\hat{T}(Y)$. This is done using the translation property of the Heisenberg unitary operators $\hat{T}(Y)$ where $Y=(y, \eta), D_{x}=$ $i^{-1} \nabla_{x}$, we have:

$$
\begin{equation*}
\hat{T}(s Y)^{*}\binom{x}{D_{x}} \hat{T}(s Y)=\binom{x-s y}{D_{x}-s \eta}, s \in \mathbb{R}, \tag{2.3}
\end{equation*}
$$

$s \mapsto(x-s y, \xi-s \eta)$ is the classical translation motion for the linear Hamiltonian $L(Y)$.
For quadratic Hamiltonians the classical flow is a time dependent linear symplectic map and the extension of (2.2) and (2.3) to quadratic Hamiltonians can be proved by the same method [3, Theorem 15, p.65].

### 2.2 The Moyal Product

We first recall the formal product rule for quantum observables with Weyl quantization. Let $A, B \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$. The Moyal product $C:=A \circledast B$ is the
observable $C$ such that $\widehat{A} \cdot \widehat{B}=\widehat{C}$. Some computations with the Fourier transform give the following well known formulas [7] (see also [8])

$$
\begin{equation*}
(A \circledast B)(X)=(\pi \hbar)^{-2 d} \iint_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} \mathrm{e}^{\frac{2 i}{\hbar} \sigma(u, v)} A(X+u) B(X+v) d u d v . \tag{2.4}
\end{equation*}
$$

Some more computations with the Fourier transform give the following formula :

$$
\begin{equation*}
(A \circledast B)(x, \xi)=\left.\exp \left(\frac{i \hbar}{2} \sigma\left(D_{q}, D_{p} ; D_{q^{\prime}}, D_{p^{\prime}}\right)\right) A(q, p) B\left(q^{\prime}, p^{\prime}\right)\right|_{(q, p)=\left(q^{\prime}, p^{\prime}\right)=(x, \xi)} \tag{2.5}
\end{equation*}
$$

where $\sigma$ is the symplectic bilinear form $\sigma\left((q, p),\left(q^{\prime}, p^{\prime}\right)\right)=p \cdot q^{\prime}-p^{\prime} \cdot q$ and $D=i^{-1} \hbar \nabla$. By expanding the exponential term in a formal power series in $\hbar$ we get

$$
\begin{equation*}
C(x, \xi)=\left.\sum_{j \geq 0} \frac{\hbar^{j}}{j!}\left(\frac{i}{2} \sigma\left(D_{q}, D_{p} ; D_{q^{\prime}}, D_{p^{\prime}}\right)\right)^{j} A(q, p) B\left(q^{\prime}, p^{\prime}\right)\right|_{(q, p)=\left(q^{\prime}, p^{\prime}\right)=(x, \xi)} \tag{2.6}
\end{equation*}
$$

Hence $C(x, \xi)$ is a formal power series in $\hbar$ with coefficients given by

$$
\begin{equation*}
C_{j}(A, B ; x, \xi)=(2 i)^{-j} \sum_{|\alpha+\beta|=j} \frac{(-1)^{|\alpha|}}{\alpha!\beta!}\left(\nabla_{x}^{\beta} \nabla_{\xi}^{\alpha} A\right) \cdot\left(\nabla_{x}^{\alpha} \nabla_{\xi}^{\beta} B\right)(x, \xi) . \tag{2.7}
\end{equation*}
$$

Furthermore we need remainder estimates for the expansion of the Moyal product.
For every $N \geq 1$ we introduce

$$
\begin{equation*}
R_{N}(A, B ; X):=A \circledast B(X)-\sum_{0 \leq j \leq N} \hbar^{j} C_{j}(X) \tag{2.8}
\end{equation*}
$$

The following estimate is a particular case of Theorem A. 1 in [2] see also Remark A. 3 .

Lemma 2.2. Let be $A \in \mathbb{S}_{\delta}^{\mu_{A}}$ and $B \in \mathbb{S}_{\delta}^{\mu_{B}}, \delta<1 / 2$, then for any $N \geq 1$, $\gamma \in \mathbb{N}^{2 d}, M \geq M_{0}$ there exists $C_{N, \gamma, M}>0$ (independent of $\left.(A, B)\right)$ such that

$$
\begin{gather*}
\left|\partial_{X}^{\gamma} R_{N}(A, B ; X)\right| \leq C_{N, \gamma, M} \hbar^{N+1} \sum_{\substack{|\alpha+\beta|=N+1 \\
|\mu+\nu| \leq M+\gamma \mid}} \sup _{u, v \in \mathbb{R}^{2 d}}\left(1+|u|^{2}+|v|^{2}\right)^{\left(M_{0}-M\right) / 2}\left|\partial_{u}^{(\alpha, \beta)+\mu} A(X+u)\right|\left|\partial_{v}^{(\beta, \alpha)+\nu} B(X+v)\right| \tag{2.9}
\end{gather*}
$$

In particular $R_{N}(A, B, X) \in \mathbb{S}_{\delta}^{\mu_{A B}}$ for some $\mu_{A B} \geq \mu_{A}+\mu_{B}$.

For proving this lemma one assumes first that $A, B \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$. For the general case we put $A_{\varepsilon}(X)=\mathrm{e}^{-\epsilon|X|^{2}} A(X), B_{\varepsilon}(X)=\mathrm{e}^{-\epsilon|X|^{2}} B(X)$ and pass to the limit for $\varepsilon \searrow 0$. In the appendix we give more details. We also need to use the following lemma.

Lemma 2.3. Let $A \in \mathbb{S}_{\delta}^{\mu_{A}}$ and $B \in \mathbb{S}_{\delta}^{\mu_{B}}, \delta<1 / 2$. Then uniformly in every compact of $\mathbb{R}_{X}^{2 d}$, we have

$$
\lim _{\epsilon \searrow 0}\left(A_{\epsilon} \circledast B\right)(X)=\lim _{\epsilon \searrow 0}\left(A \circledast B_{\epsilon}\right)(X)=(A \circledast B)(X) .
$$

In particular we have

$$
\lim _{\epsilon \searrow 0}\left\{A_{\epsilon}, B\right\}_{\circledast}(X)=\{A, B\}_{\circledast}(X) .
$$

For completeness a proof is given in the appendix B.

## 3 Proof of Theorem 1.1

Here $\hbar=1$. Notice first that from Lemma 2.3 we also have for any $A \in \mathbb{S}_{0}^{0}$, $\{A, H\}_{\circledast}=\{A, H\}$. So it is enough to consider the test observables $A=T_{Y}:=\mathrm{e}^{-i L_{Y}}\left(Y \in \mathbb{R}^{2 d}\right)$.
We have

$$
\widehat{T}_{Y} \widehat{H} \widehat{T}_{Y}^{*}=\left[\widehat{T}_{Y}, \widehat{H}\right] \widehat{T}_{Y}^{*}+\widehat{H}
$$

Using the assumption of Theorem 1.1 and Lemma 2.3 we get

$$
\begin{equation*}
\frac{1}{i}\left(\left\{T_{T}^{*}, H\right\} \circledast T_{Y}\right)(X)=H(X+Y)-H(X), \forall X, Y \in \mathbb{R}^{2 d} \tag{3.1}
\end{equation*}
$$

Computing the Poisson bracket in (3.1) gives

$$
\begin{equation*}
\left(\left(\left(y \cdot \partial_{x} H+\eta \cdot \partial_{\xi} H\right) T_{Y}^{*}\right) \circledast T_{Y}\right)(X)=H(X+Y)-H(X), \forall X, Y \in \mathbb{R}^{2 d} \tag{3.2}
\end{equation*}
$$

Our aim is to prove that (3.2) implies that $H(X)$ is a polynomial of degree at most 2. For that purpose we shall compute the asymptotic expansion as $Y \rightarrow 0$ of the left hand side of (3.2) and compare it with the Taylor expansion for $H(X+Y)$ modulo $O\left(|Y|^{4}\right)$. From that we shall conclude that all the third order derivatives of $H$ vanish for $X$ in any bounded subset of $\mathbb{R}^{2 d}$, hence the conclusion will follow.

We have

$$
\partial_{x}^{\alpha} \partial_{\xi}^{\beta} T_{Y}=i^{-|\alpha+\beta|} \eta^{\alpha} y^{\beta} T_{Y}
$$

Let us denote by $C(X, Y)$ the left hand side in (3.2). So using Lemma 2.2 uniformly in every compact in $X \in \mathbb{R}^{2 d}$, we have

$$
C(X, Y)=\sum_{0 \leq j \leq 2} C_{j}(X, Y)+O\left(|Y|^{4}\right),
$$

where

$$
\begin{align*}
& C_{0}(X, Y)=Y \cdot \nabla_{X} H(X)  \tag{3.3}\\
& C_{1}(X, Y)=\frac{1}{2} Y \cdot \nabla_{X}^{2} H(X) Y, \tag{3.4}
\end{align*}
$$

where $\nabla_{X}^{2} H(X)$ is the Hessian matrix of $H$.
Let us compute now $C_{2}(X, Y)$, which is an homogeneous polynomial of degree 3 in $Y$.
For simplicity let us consider the 1-D case. The same computation can clearly be done for $d>1$.
Using (2.6) we get with $Y=(y, \eta)$,

$$
\begin{equation*}
C_{2}(X, Y)=\frac{1}{8}\left(y^{3} \partial_{x}^{3} H+\eta^{3} \partial_{\eta}^{3} H-y^{2} \eta \partial_{\xi} \partial_{x}^{2} H-y \eta^{2} \partial_{\xi}^{2} \partial_{x} H\right) \tag{3.5}
\end{equation*}
$$

According (3.2), $C_{2}(X, Y)$ must coincide with the term of order 3 in $Y$ of the Taylor expansion in $X$ for $H(X+Y)-H(X)$. But this is possible only if $\partial_{x}^{3} H=\partial_{\eta}^{3} H=\partial_{\xi} \partial_{x}^{2} H=\partial_{\xi}^{2} \partial_{x} H=0$ for any $(x, \xi) \in \mathbb{R}^{2}$. So $H$ must be a polynomial of degree $\leq 2$.

## 4 Extension to polynomials of arbitrary degree

The asymptotic expansion in $\hbar$ in the Moyal product suggests to introduce the following semi-classical approximations of the Moyal bracket:

$$
\{A, B\}_{\circledast, m}=\{A, B\}+\hbar^{2}\{A, B\}_{3}+\cdots+\hbar^{2 m}\{A, B\}_{2 m+1},
$$

where $\{A, B\}_{j}=\frac{i}{\hbar}\left(C_{j}(A, B)-C_{j}(B, A)\right)$ (notation of (2.7)). More explicitly we have

$$
\begin{equation*}
\{A, B\}_{j}=(2 i)^{-j} \sum_{|\alpha+\beta|=j} \frac{(-1)^{|\alpha|}-(-1)^{|\beta|}}{\alpha!\beta!}\left(\nabla_{x}^{\beta} \nabla_{\xi}^{\alpha} A\right) \cdot\left(\nabla_{x}^{\alpha} \nabla_{\xi}^{\beta} B\right)(x, \xi) \tag{4.1}
\end{equation*}
$$

Notice that $\{A, B\}_{j}=0$ for $j$ even. In particular if $A=T_{Y}$ we have if $j=2 k+1$, with $Y=(y, \eta)$,

$$
\begin{equation*}
T_{Y}^{*}\left\{T_{Y}, B\right\}_{2 k+1}(x, \xi)=(2)^{-2 k} \sum_{|\alpha+\beta|=2 k+1} \frac{y^{\alpha} \eta^{\beta}}{\alpha!\beta!} \cdot\left(\nabla_{x}^{\alpha} \nabla_{\xi}^{\beta} B\right)(x, \xi) . \tag{4.2}
\end{equation*}
$$

It is clear that if $H$ is a polynomial of degree at most $2 m+2$ then we have $\{A, H\}_{\circledast, m}=\{A, H\}_{\circledast}$ for any $A$. Conversely we have

Theorem 4.1. Assume $\hbar=1$ and $H \in \mathbb{S}_{\delta}^{\mu}, \mu \in \mathbb{R}, \delta<1 / 2$. If for any $A \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$ we have $\{A, H\}_{\circledast, m}=\{A, H\}_{\circledast}$ then $H$ must be a polynomial of degree at most $2 m+2$.

Proof. Here we give a proof different from the case $m=0$, without connection with the Taylor formula, for simpler computations. Using Lemma 2.2 we have, uniformly in every compact in $X \in \mathbb{R}^{2 d}$,

$$
\begin{equation*}
T_{Y}^{*}\left(\left\{T_{Y}, H\right\}_{\circledast}(X)-\left\{T_{Y}, H\right\}_{\circledast, m}(X)\right)=\mathcal{O}\left(|Y|^{2 m+3}\right), Y \rightarrow 0 \tag{4.3}
\end{equation*}
$$

Moreover from (2.7) we get:

$$
\begin{equation*}
T_{Y}^{*}\left\{T_{Y}, H\right\}_{2 j+1}(X)=\frac{1}{2^{j+1}} \sum_{\mid \alpha+\beta=2 j+1} \frac{y^{\alpha} \eta^{\beta}}{\alpha!\beta!} \partial_{x}^{\alpha} \partial_{\xi}^{\beta} H(X) \tag{4.4}
\end{equation*}
$$

Using the assumption of Theorem 4.1. and (4.3) we get that

$$
T_{Y}^{*}\left\{T_{Y}, H\right\}_{2 m+3}(X)=\mathcal{O}_{X}\left(|Y|^{2 m+5}\right)
$$

But $T_{Y}^{*}\left\{T_{Y}, H\right\}_{2 m+3}$ is an homogeneous polynomial of degree $2 m+3$ in $Y$ so we get that this polynomial is 0 and from (4.4) we get that $\partial_{x}^{\alpha} \partial_{\xi}^{\beta} H(X)=0$ for $|\alpha+\beta|=2 m+3$. Then we can conclude that $H(X)$ is a polynomial of degree at most $2 m+2$ in $X \in \mathbb{R}^{2 d}$.

## A Proofs for formulas (2.4) and (2.5)

It is enough to assume that $A, B \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$.
Recall first the relationship between Weyl symbols and integral kernel of $\hat{A}$. We have

$$
K_{\hat{A}}(x, y)=(2 \pi \hbar)^{-d} \int_{\mathbb{R}^{d}} \mathrm{e}^{\frac{i}{\hbar}(x-y) \cdot \eta} A\left(\frac{x+y}{2}, \eta\right) d \eta
$$

and

$$
A(x, \xi)=\int_{\mathbb{R}^{d}} \mathrm{e}^{-\frac{i}{\hbar} \xi \cdot t} K_{\hat{A}}(x+t / 2, x-t / 2) d t
$$

Using these formulas and the relation $K_{\hat{A} \hat{B}}(x, z)=\int_{\mathbb{R}^{d}} K_{\hat{A}}(x, y) K_{\hat{B}}(y, z) d y$ we get

$$
\begin{align*}
(A \circledast B)(X)=(2 \pi \hbar)^{-2 d} & \int_{\mathbb{R}^{4 d}} \exp \left(\frac{i}{\hbar}(-t \cdot \xi+(x-y-t / 2) \cdot \eta+(y-x+t / 2)) \cdot \zeta\right) . \\
& A((x+y) / 2+t / 4, \eta) B((x+y) / 2-t / 4, \zeta) d \zeta d \eta d y d t \tag{A.1}
\end{align*}
$$

Then after the change of variables in the integral $v_{\xi}=\zeta-\xi, u_{\xi}=\eta-\xi$, $u_{x}=(y-x) / 2+t / 4, v_{x}=(y-x) / 2-t / 4$, we get formula (2.4), with $u=\left(u_{x}, u_{\xi}\right), v=\left(v_{x}, v_{\xi}\right)$,

$$
(A \circledast B)(X)=(\pi \hbar)^{-2 d} \int_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} \mathrm{e}^{\frac{2 i}{\hbar} \sigma(u, v)} A(X+u) B(X+v) d u d v
$$

To get formula (2.5) we notice that $(u, v) \mapsto 2 \sigma(u, v)$ is non degenerate and its matrix is $G:=\left(\begin{array}{cc}0 & -J \\ J & 0\end{array}\right)$, so $G^{-1}=G$. Hence using Fourier transform in $(u, v)$ and the Fourier multiplier formula we get (2.5).

## B Proofs for Lemmas 2.2 and 2.3

## B. 1 Proof of lemma 2.3

Using (2.4) for $A_{\epsilon} \circledast B$ we split the integral in two pieces :
$1=\chi_{0}\left(|u|^{2}+|v|^{2}\right)+\chi_{1}\left(|u|^{2}+|v|^{2}\right)$, where $\chi_{0} \in C_{0}^{\infty}(\mathbb{R}), \chi_{0}(t)=1$ for $|t| \leq 1 / 2$. On the support of $\chi_{0}$ we can obviously pass to the limit in $\epsilon$. On the support of $\chi_{1}$ we first perform integrations by parts with the differential operator $L$ to get a uniformly and absolutely convergent integral,

$$
L=\frac{J u \cdot \partial_{v}-J v \cdot \partial_{u}}{|u|^{2}+|v|^{2}}
$$

using that $L \mathrm{e}^{\frac{2 \mathrm{i}}{\hbar} \sigma(\mathrm{u}, \mathrm{v})}=\mathrm{Le} \mathrm{e}^{\frac{2 \mathrm{i}}{\hbar} \sigma(\mathrm{u}, \mathrm{v})}=\frac{2 \mathrm{i}}{\hbar} \mathrm{i}^{\frac{2 \mathrm{i}}{\hbar} \sigma(\mathrm{u}, \mathrm{v})}$. On the support of $\chi_{1}$, performing $4 d+1$ integrations by parts for gaining enough decay to ensure integrability in $(u, v) \in \mathbb{R}^{4 d}$. Then passing to the limit in $\epsilon$ we get $\lim _{\epsilon \backslash 0}\left(A_{\epsilon} \circledast B\right)(X)=$ $(A \circledast B)(X)$ and the same for $\lim _{\epsilon} \bigwedge_{0}\left(A \circledast B_{\varepsilon}\right)(X)=\{A, B\}_{\circledast}(X)$.
The other properties follow.

## B. 2 Proof of Lemma 2.2

From (2.4), by Fourier transform computations and application of the Taylor formula, we get the following formula for the remainder,

$$
\begin{equation*}
R_{N}(A, B, X)=\frac{1}{N!}\left(\frac{i \hbar}{2}\right)^{N+1} \int_{0}^{1}(1-t)^{N} R_{N, t}(X ; \hbar) d t \tag{B.1}
\end{equation*}
$$

where
$R_{N, t}(X ; \varepsilon)=$
$(2 \pi \bar{t})^{-2 d} \iint_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} \exp \left(-\frac{i}{2 t \hbar} \sigma(u, v)\right) \sigma^{N+1}\left(D_{u}, D_{v}\right) A(u+X) B(v+X) d u d v$.

Notice that the integral is an oscillating integral as we shall see below. So we shall use the following lemma :

Lemma B.1. There exists a constant $C_{d}>0$ such that for any $F \in \mathcal{S}\left(\mathbb{R}^{2 d} \times\right.$ $\mathbb{R}^{2 d}$ ) the integral

$$
\begin{equation*}
I(\lambda)=\lambda^{2 d} \iint_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} \exp [-i \lambda \sigma(u, v)] F(u, v) d u d v \tag{B.2}
\end{equation*}
$$

satisfies the following estimate:
for any $M>0$ there exists $C_{M}>0$, independent of $F$ such that

$$
\begin{equation*}
\mid I\left(\lambda\left|\leq C_{M} \sup _{\substack{u, v \in \mathbb{R}^{2} d \\|\alpha+\beta| \leq M+4 d+1}}\left(1+|u|^{2}+|v|^{2}\right)^{(4 d+1-M) / 2}\right| \partial_{u}^{\alpha} \partial_{v}^{\beta} F(u, v) \mid .\right. \tag{B.3}
\end{equation*}
$$

A proof will be given later in the next section.
Using this lemma for $A, B \in \mathcal{S}\left(\mathbb{R}^{2 d}\right)$ with the integrand

$$
F_{N, \gamma}(X ; u, v)=\pi^{-2 d} \partial_{X}^{\gamma}\left(\sigma^{N+1}\left(D_{u}, D_{v}\right) A(u+X) B(v+X)\right)
$$

and the parameter $\lambda=1 /(2 t \hbar)$, we then have that

$$
\left|\partial_{X}^{\gamma} R_{N, t}(X ; \varepsilon)\right| \leq C_{d} \sup _{\substack{u, v \in \mathbb{R}^{2 d} \\|\alpha|+|| | \leq 4 d+1}}\left|\partial_{u}^{\alpha} \partial_{v}^{\beta} F_{N, \gamma}(X ; u, v)\right| .
$$

Moreover, we have the elementary estimate

$$
\begin{equation*}
\left|\sigma^{N+1}\left(D_{u}, D_{v}\right) A(u) B(v)\right| \leq(2 d)^{N+1} \sup _{|\alpha|+|\beta|=N+1}\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} A(x, \xi) \partial_{y}^{\beta} \partial_{\eta}^{\alpha} B(y, \eta)\right| . \tag{B.4}
\end{equation*}
$$

Together with the Leibniz formula, we then get the claimed result with universal constants. For symbols $A \in \mathbb{S}_{0}^{\mu}$ and $B \in \mathbb{S}_{0}^{\nu}$ we argue by localisation. We use $A_{\epsilon}(u)=\mathrm{e}^{-\epsilon u^{2}} A(u)$ and $B_{\epsilon}(v)=\mathrm{e}^{-\epsilon v^{2}} B(v)$ for $\epsilon>0$ and pass to the limit as $\epsilon \rightarrow 0$.

## B. 3 Proof of the Lemma B. 1

We consider the same cut-off $\chi_{0}$ as above. We split $I(\lambda)$ into two pieces and write $I(\lambda)=I_{0}(\lambda)+I_{1}(\lambda)$ with

$$
\begin{aligned}
& I_{0}(\lambda)=\lambda^{2 d} \iint_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} \exp [-i \lambda \sigma(u, v)] \chi_{0}\left(\left(u^{2}+v^{2}\right)\right) F(u, v) d u d v \\
& \left.I_{1}(\lambda)=\lambda^{2 d} \iint_{\mathbb{R}^{2 d} \times \mathbb{R}^{2 d}} \exp [-i \lambda \sigma(u, v)]\left(1-\chi_{0}\right)\left(u^{2}+v^{2}\right)\right) F(u, v) d u d v .
\end{aligned}
$$

We notice that $(u, v) \mapsto \sigma(u, v)$ is a quadratic non-degenerate real form on $\mathbb{R}^{4 d}$.
Let us estimate $I_{1}(\lambda)$. We can integrate by parts with the differential operator

$$
L=\frac{i}{|u|^{2}+|v|^{2}}\left(J u \cdot \frac{\partial}{\partial v}-J v \cdot \frac{\partial}{\partial u}\right)
$$

using that $L \mathrm{e}^{-\mathrm{i} \lambda \sigma(\mathrm{u}, \mathrm{v})}=\mathrm{Le}^{-\mathrm{i} \lambda J \mathrm{u} \cdot \mathrm{v}}=\lambda \mathrm{e}^{-\mathrm{i} \lambda \sigma(\mathrm{u}, \mathrm{v})}$. For $I_{1}(\lambda)$, the integrand is supported outside the ball of radius $1 / \sqrt{2}$ in $\mathbb{R}^{4 d}$. Performing $4 d+1$ integrations by parts for gaining enough decay to ensure integrability in $(u, v) \in \mathbb{R}^{4 d}$, we get a constant $c_{d}$ such that

$$
\begin{equation*}
\left|I_{1}(\lambda)\right| \leq c_{d} \sup _{\substack{u, v \in \mathbb{R}^{2 d} \\|\mu|+\mid \nu \nu \leq 4 d+1}}\left|\partial_{u}^{\mu} \partial_{v}^{\nu} F(u, v)\right| . \tag{B.5}
\end{equation*}
$$

But we need to control the behaviour for $u^{2}+v^{2}$ large, so with $M$ more integrations by parts we get

$$
\begin{equation*}
\left|I_{1}(\lambda)\right| \leq C_{M} \sup _{\substack{u, v \mathbb{R}^{2 d} \\|\alpha+\beta| \leq M}}\left(1+|u|^{2}+|v|^{2}\right)^{(4 d+1-M) / 2}\left|\partial_{u}^{\alpha} \partial_{v}^{\beta} F(u, v)\right| \tag{B.6}
\end{equation*}
$$

To estimate $I_{0}(\lambda)$ we apply the stationary phase theorem. The symmetric matrix of the quadratic form $\sigma(u, v)$ is

$$
A_{\sigma}=\left(\begin{array}{cc}
0 & -J \\
J & 0
\end{array}\right) .
$$

So by the stationary phase theorem ([7], Vol.I, section 7.7), noticing that the leading term in the stationary phase theorem is of order $\lambda^{-2 d}$, we get

$$
\begin{equation*}
\left|I_{0}(\lambda)\right| \leq c_{d}^{\prime} \sup _{\substack{u, v \in \mathbb{R}^{2 d} \\|\mu|+\mid \nu \nu \leq 2 d+3}}\left|\partial_{u}^{\mu} \partial_{v}^{\nu} F(u, v)\right| . \tag{B.7}
\end{equation*}
$$

## C Proof of proposition $\mathbf{1 . 6}$

Let us recall the statement of the semiclassical Egorov theorem (see [2, 8] for a proof)

Theorem C.1. Let us assume that $H$ is a subquadratic Hamiltonian, that means here that for $|\alpha|+|\beta| \geq 2$ there exists $C_{\alpha, \beta}$ such that

$$
\left|\partial_{x}^{\alpha} \partial_{\xi}^{\beta} H(x, \xi)\right| \leq C_{\alpha, \beta}, \forall(x, \xi) \in \mathbb{R}^{2 d}
$$

For any $A \in \mathbb{S}_{0}^{0}$ we have, for the operator norm in $L^{2}\left(\mathbb{R}^{d}\right)$ and for every $N \geq 0$,

$$
\begin{equation*}
\hat{A}_{\text {quant }}(t)=\mathrm{e}^{-\frac{i}{\hbar} \hat{H}} \hat{\mathrm{~A}} \mathrm{e}^{\frac{i}{\hbar} \hat{H}}=\sum_{0 \leq k \leq N} \hbar^{2 j} \hat{A}_{2 j}(t)+O\left(\hbar^{2 N+2}\right) . \tag{C.1}
\end{equation*}
$$

In particular we have $A_{0}(t, X)=A\left(\Phi_{H}^{t}(X)\right):=A_{\mathrm{cl}}(t, X)$ and for $j \geq 1$, $A_{2 j}(t, X)$ is defined by induction on $j$ such that $A_{2 j}(0)=0$ and

$$
\begin{equation*}
\partial_{t} A_{2 j+2}(t)-\left\{A_{2 j+2}, H\right\}_{1}=\sum_{k+\ell=j+1, k \leq j}\left\{A_{2 k}(t), H\right\}_{2 \ell+1} . \tag{C.2}
\end{equation*}
$$

Let us remark that here the expansion is even in $\hbar$ because $H$ is reduced to its principal part (no perturbative term in $\hbar$ ).
$A_{2} \equiv 0$ hence $H$ is a polynomial of degree at most 2 . We can recover easily here these results.
Proof of proposition 1.6.
From (C.2) we get

$$
\begin{equation*}
\partial_{t} A_{2 j}(t)-\left\{A_{2 j}(t), H\right\}_{1}=\sum_{k+\ell=j, k \leq j-1}\left\{A_{2 k}(t), H\right\}_{2 \ell+1} \tag{C.3}
\end{equation*}
$$

If for $j=j_{0}$ we have $\partial_{t} A_{2 j}(t)-\left\{A_{2 j}(t), H\right\}_{1} \equiv 0$, so taking $t=0$ in the r.h.s of (C.3) we get $\{A, H\}_{2 j_{0}+1}=0$. Hence using the observables $A=T_{Y}$, like in (4.4), we get that $H$ is a polynomial of degree at most $2 j_{0}$.
Moreover using that $\partial_{t} A_{2 j}(0)=0$ and $A_{2 k}(0)=0$ for $k \geq 1, A_{0}(0)=A$, we get the formula (4.5) of the paper [1]:

$$
\begin{equation*}
\sum_{k+\ell=j, k \leq j-1}\left\{\{A, H\}_{2 \ell+1}, H\right\}_{2 k+1}=0 \tag{C.4}
\end{equation*}
$$

Then as in [1] from (C.4) one can conclude that $H$ is a polynomial of degree at most 2 .

For the last part of the proposition we notice that from the assumption we get that $A_{2}(t, X) \equiv 0$.

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