# Probabilistic version of Trotter's exponential product formula in Banach algebras 

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## 1. Introduction and results

It is an elementary fact that the exponential function may be defined by the equivalent formulae

$$
\exp (x)=\lim _{n \rightarrow \infty}\left(1+\frac{1}{n} x\right)^{n}=\sum_{n=0}^{\infty} \frac{1}{n!} x^{n}
$$

not only when $x$ is a real or complex number but also when it is a matrix with real or complex entries or a bounded operator acting on a Hilbert space or a Banach space or, even when it is an element of an abstract Banach algebra $\mathfrak{B}$ with identity 1 . (for a definition of a Banach algebra see for instance [1]). If $\mathfrak{B}$ is not commutative then in general $\exp (x) \exp (y) \neq \exp (x+y)$. There is, however, a formula which replaces the addition law of the exponential function, namely

$$
\begin{equation*}
\lim _{n \rightarrow \infty}\left(\exp \left(\frac{x}{n}\right) \exp \left(\frac{y}{n}\right)\right)^{n}=\exp (x+y) \tag{1}
\end{equation*}
$$

and this holds regardless whether $x$ and $y$ commute or not. Formula (1) is capable of further generalization; see Trotter [2]. Specifically, $x$ and $y$ may be unbounded operators of a certain type, namely generators of continuous one-parameter operator semi-groups. In the present paper we are not concerned with Trotter's generalization, but we shall still refer to (1) as the Trotter product formula. The symbols $x, y, \ldots, a, b, \ldots$ shall generally denote elements of the Banach algebra $\mathfrak{B}$. The norm of $x \in \mathfrak{B}$ is written $\|x\|$.

Let $\mathbf{X}=\left(x_{1}, x_{2}, \ldots, x_{m}\right.$ be any finite sequence of elements of $\mathfrak{B}$. With any such sequence we associate the product

$$
T(\mathbf{X})=\exp \left(\frac{x_{1}}{m}\right) \exp \left(\frac{x_{2}}{m}\right) \ldots \exp \left(\frac{x_{m}}{m}\right)
$$

which will be called its Trotter product. Note that it depends essentially on the
order of the factors, i:e. on $\mathbf{X}$ as a sequence, not merely as a set. We also write, for the mean of the elements of $\mathbf{X}, M(\mathbf{X})=\frac{1}{m}\left(x_{1}+x_{2}+\cdots+x_{m}\right)$. Using this notation we can express the Trotter product formula as follows: If $\mathbf{X}_{k}$ (for $k=1,2, \ldots$ ) is the sequence of length $2 k$ whose elements are alternatingly $x$ and $y$, then

$$
\begin{equation*}
\lim _{k \rightarrow \infty} T\left(\mathbf{X}_{k}\right)=\exp \left(\lim _{k \rightarrow \infty} M\left(\mathbf{X}_{k}\right)\right) . \tag{2}
\end{equation*}
$$

We now raise the following question. Under what natural conditions on a sequence $\mathbf{X}_{1}, \mathbf{X}_{2}, \ldots$ (subject to the requirement that their lengths tend to infinity) will (2) hold?

We shall prove two theorems relevant to this question. The first theorem gives a rather general sufficient condition for an infinite sequence $\mathbf{X}_{1}, \mathbf{X}_{2}, \ldots$ to satisfy (2). The original Trotter formula is an obvious consequence of this condition. But our theorem also shows that to a certain extent the order of the factors in the Trotter product may be made subject to considerable rearrangement without destroying the validity of (2).

For any $\mathbf{X}=\left(x_{1}, x_{2}, \ldots, x_{m}\right)$ let us write $\varrho=\varrho(\mathbf{X})=\operatorname{Max}_{1 \leqq j \leqq m}\left\|x_{j}\right\|$. Let $\pi$ denote a partition of the sequence $(1,2, \ldots, m)$ into successive subsequences $\left(1,2, \ldots, m_{1}\right),\left(m_{1}+1, m_{1}+2, \ldots, m_{1}+m_{2}\right), \ldots,\left(m-m_{s}+1, m-m_{s}+2, \ldots, m\right)$, and let $\mathbf{Y}_{1}, \mathbf{Y}_{2}, \ldots, \mathbf{Y}_{s}$ be the corresponding subsequences of $\mathbf{X}$. For any element $g \in \mathfrak{B}$ we introduce a quantity $\delta=\delta(\mathbf{X}, \pi, g)$ which measures the closeness to which $g$ uniformly approximates the "partial means" of $\mathbf{X}$ induced by the partition $\pi$

$$
\begin{equation*}
\delta(\mathbf{X}, \pi, g)=\operatorname{Max}_{1 \leqq j \xi s}\left\|M\left(\mathbf{Y}_{j}\right)-g\right\| . \tag{3}
\end{equation*}
$$

We also define a quantity

$$
\begin{equation*}
\eta=\dot{\eta}(\pi)=\sum_{j=1}^{s}\left(\frac{m_{j}}{m}\right)^{2} . \tag{4}
\end{equation*}
$$

If $\eta_{0}=\max _{1 \equiv j \equiv s}\left(\frac{m_{j}}{m}\right)$, we have clearly

$$
\begin{equation*}
\eta_{0}^{2} \leqq \eta \leqq \eta_{0}, \tag{5}
\end{equation*}
$$

so that $\eta$ is a kind of a measure for the relative fineness of $\pi$.
Theorem 1. Let $\mathbf{X}_{1}, \mathbf{X}_{2}, \ldots$ be an infinite sequence of finite sequences of elements of $\mathfrak{B}$. Suppose that $\varrho\left(\mathbf{X}_{k}\right)$ is bounded. Suppose that $g \in \mathfrak{B}$ exists and a sequence of partitions $\pi_{k}$ of $\mathbf{X}_{k}$ into successive subsequences exists such that $\eta\left(\pi_{k}\right) \rightarrow 0$ and $\delta\left(\mathbf{X}_{k}, \pi_{k}, g\right) \rightarrow 0$ as $k \rightarrow \infty$. Then $M\left(\mathbf{X}_{k}\right) \rightarrow g$ and $T\left(X_{k}\right) \rightarrow \exp (g)$ as $k \rightarrow \infty$.

Note that the quantity $\delta(\mathbf{X}, \pi, g)$ is independent of the order of the elements $x_{i}$ within each of the subsequences $\mathbf{Y}_{j}(j=1,2, \ldots, s)$. Thus our theorem formulates
mathematically the intuitive notion that the order of the factors in the Trotter product $T(\mathbf{X})$ is irrelevant "locally" and only essential "in the large".

Our second Theorem deals with the following problem. Suppose we consider infinite sequences $\mathbf{Z}=\left(x_{1}, x_{2}, \ldots\right)$ whose elements are all taken from a fixed finite subset $\left\{a_{1}, a_{2}, \ldots, a_{\sigma}\right\}$ of $\mathfrak{B}$. With any such infinite sequence we may consider the sequence of its finite sections $\mathbf{X}_{1}=\left(x_{1}\right), \mathbf{X}_{2}=\left(x_{1}, x_{2}\right), \ldots$. Is it possible to characterize the extent of the set of those infinite sequences $\mathbf{Z}$ for which the generalized Trotter product formula (2) holds?

This question leads to measure-theoretic considerations, and can be naturally formulated in probabilistic terminology.

Theorem 2. Let $a_{1}, a_{2}, \ldots, a_{\sigma} \in \mathfrak{B}$, and let $p_{1}, p_{2}, \ldots, p_{\sigma} \geqq 0, \Sigma_{v} p_{v}=1$. Suppose an infinite (random) sequence $\mathbf{Z}=\left(x_{1}, x_{2}, \ldots\right)$ is formed by choosing, independently for each $j \geqq 1, x_{j}=a_{\sigma}$ with probability $p_{\sigma}$. Let $g=p_{1} a_{1}+p_{2} a_{2}+\cdots+p_{\sigma} a_{\sigma}$. Then the probability is unity that $M\left(\mathbf{X}_{k}\right) \rightarrow g$ and $T\left(\mathbf{X}_{k}\right) \rightarrow \exp (g)$ as $k \rightarrow \infty$.

Thus in the sense of the given probability measure defined on the set of sequences $\mathbf{Z}$, for almost every sequence the generalized Trotter product formula holds.

## 2. An auxiliary inequality

Both theorems derive from an elementary estimate formulated as follows:
Lemma. Let $\mathbf{X}$ be a finite sequence of elements of $\mathfrak{B}, \pi$ any partition of it into successive subsequences, and $g \in \mathfrak{B}$. Let $\varrho=\varrho(\mathbf{X}), \eta=\eta(\pi)$ and $\delta=\delta(\mathbf{X}, \pi, g)$ be defined as above. Then

$$
\|T(\mathbf{X})-\exp (g)\| \leqq e^{\|g\|}\left(e^{\delta+\frac{1}{2} \eta e^{2} e^{a}}-2 e^{-\frac{1}{2} \eta\|g\|^{2}}+1\right)
$$

In the proof of the Lemma we shall make use of the following facts:
(A) If $P\left(x_{1}, x_{2}, \ldots, x_{q}\right)$ is a polynomial or a power series with non-negative real coefficients then

$$
\left\|P\left(x_{1}, x_{2}, \ldots, x_{q}\right)\right\| \leqq P\left(\left\|x_{1}\right\|,\left\|x_{2}\right\|, \ldots,\left\|x_{q}\right\|\right)
$$

for all $x_{j} \in \mathfrak{B}$ such that the right hand side is finite.
(B) For any $t \geqq 0, e^{t}-1-t \leqq \frac{1}{2} t^{2} e .^{t}$
(C) For any $t \geqq 0, e^{t-\frac{1}{2} t^{2}} \leqq 1+t \leqq e^{t}$.

Let $\mathbf{Y}_{j}(j=1,2, \ldots, s)$ be the subsequences of $\mathbf{X}$ produced by the partition $\pi$, and let $m_{j}$ be their respective length, $m=\Sigma_{j} m_{j}$. Let $T(\mathbf{X})=y_{1} y_{2} \ldots y_{s}$ where $y_{j}$ is the product of those factors in the product, taken in their proper order, which
involve the elements $x_{k} \in \mathbf{Y}_{j}$. Write $y_{j}=1+\frac{m_{j}}{m} g+r_{j}$, thus defining $r_{j}$. For notational convenience we now consider $j=1$. We have

$$
r_{1}=\left[y_{1}-1-\frac{m_{1}}{m} M\left(\mathbf{Y}_{1}\right)\right]+\frac{m_{1}}{m}\left[M\left(\mathbf{Y}_{1}\right)-g\right]
$$

The norm of the second term is bounded by $\frac{m_{1}}{\dot{m}} \delta$. To obtain a bound on the norm of the first term we note that if

$$
y_{1}-1-\frac{m_{1}}{m} M\left(\mathbf{Y}_{1}\right)=\exp \left(\frac{x_{1}}{m}\right) \exp \left(\frac{x_{2}}{m}\right) \ldots \exp \left(\frac{x_{m_{1}}}{m}\right)-1-\frac{1}{m}\left(x_{1}+x_{2}+\cdots+x_{m_{1}}\right)
$$

is regarded as a power series in the $x_{j}\left(j=1,2, \ldots, m_{1}\right)$ it has non-negative coefficients (the negative terms cancel!). Thus by principle (A) above we may replace $x_{j}$ by $\left\|x_{j}\right\|$, and then taking (B) and the definition of $\varrho$ into account we get

$$
\left\|y_{1}-1-\frac{m_{1}}{m} M\left(\mathbf{Y}_{1}\right)\right\| \leqq \frac{1}{2}\left(\frac{m_{1}}{m}\right)^{2} \varrho^{2} e^{\varrho}
$$

Thus we have for $j=1,2, \ldots, s$

$$
\begin{equation*}
\left\|r_{j}\right\| \leqq \frac{1}{2}\left(\frac{m_{j}}{m}\right)^{2} \varrho^{2} e^{\varrho}+\frac{m_{j}}{m} \delta . \tag{6}
\end{equation*}
$$

Next, let $z_{j}=1+\frac{m_{j}}{m} g$, and consider the difference $T(\mathbf{X})-z_{1} z_{2} \ldots \dot{z}_{s}=$ $=y_{1} y_{2} \ldots y_{s}-z_{1} z_{2} \ldots z_{s}$. As a polynomial in $g$ and the $r_{j}$, it has again non-negative coefficients, so we apply principle (A). The norms of $r_{j}$ are majorized by (6), and so by the inequality (C)

$$
1+\frac{m_{j}}{m}\|g\|+\left\|r_{j}\right\| \leqq e^{\frac{m_{j}}{m}\|g\|+\frac{1}{2}\left(\frac{m_{j}}{m}\right)^{2} e^{2} e^{e}+\frac{m_{j}}{m} \delta}
$$

and

$$
-\left(1+\frac{m_{j}}{m}\|g\|\right) \leqq-e^{\frac{m_{j}}{m}\|g\|-\frac{1}{2}\left(\frac{m_{j}}{m}\right)^{2}\|g\|^{2}}
$$

where in the last step we used $\frac{m_{j}}{m}\|g\| \leqq \sqrt{\eta}\|g\| \leqq 1$. Thus we see that

$$
\left\|y_{1} y_{2} \ldots y_{s}-z_{1} z_{2} \ldots z_{s}\right\| \leqq e^{\|g\|}\left(e^{\delta+\frac{1}{2} \eta \varrho^{2} c c^{2}}-e^{-\frac{1}{2}\| \| g i^{2}}\right)
$$

Arguing analogously, we have also

$$
\left\|\exp (g)-z_{1} z_{2} \ldots z_{s}\right\| \leqq e^{\|g\|}\left(1-e^{-\frac{1}{2}\| \| \|^{2}}\right)
$$

The last two inequalities together imply the conclusion of the lemma.

We note that Theorem 1 is an immediate consequence of the lemma, since the estimate for $\left\|T\left(\mathbf{X}_{k}\right)-\exp (g)\right\|$ supplied by the lemma tends to zero as $k \rightarrow \infty$ if the hypotheses of the theorem are fulfilled.

## 3. Proof of Theorem 2

The idea of the proof of Theorem 2 is to find an appropriate sequence of partitions $\pi_{k}(k=1,2,3, \ldots)$ such that if we let $\delta_{k}=\delta\left(\mathbf{X}_{k}, \pi_{k}, g\right)$ then

$$
\begin{equation*}
\mathbf{P}\left\{\lim _{k \rightarrow \infty} \delta_{k}=0\right\}=1 \tag{7}
\end{equation*}
$$

and at the same time such that

$$
\begin{equation*}
\lim _{k \rightarrow \infty} \eta\left(\pi_{k}\right)=0 \tag{8}
\end{equation*}
$$

Indeed, if (7) and (8) are fulfilled then the conclusion of Theorem 2 follows from Theorem 1.

Let $C_{1}<C_{2}$ and $\beta<1$ be three positive constants. We define the partition $\pi_{k}$ of $(1,2,3, \ldots, k)$ into successive subsequences of lengths $m_{j}=m_{j}(k)(j=1,2, \ldots, s=s(k))$ in such a manner that for all $j$ and $k$

$$
\begin{equation*}
C_{1} k^{\beta}<m_{j}<C_{2} k^{\beta} \tag{9}
\end{equation*}
$$

Since $\Sigma_{j} m_{j}=k$, it follows that

$$
\begin{equation*}
s=s(k)=O\left(k^{1-\beta}\right) \tag{10}
\end{equation*}
$$

and therefore

$$
\begin{equation*}
\eta_{k}=\eta\left(\pi_{k}\right)=\sum_{j=1}^{s(k)}\left(\frac{m_{j}}{k}\right)^{i}=O\left(k^{\beta-1}\right) \tag{11}
\end{equation*}
$$

so that (8) holds. Note that in our probabilistic set-up the partitions $\pi_{k}$ are not random variables (i.e. $\pi_{k}$ is constant over the whole probability space).

Next, we remark that, by virtue of the Borel-Cantelli lemma, in order to prove (7) it is sufficient to prove that

$$
\begin{equation*}
\sum_{k=1}^{\infty} \mathbf{P}\left\{\delta_{k} \geqq \varepsilon\right\}<\infty \tag{12}
\end{equation*}
$$

for any positive $\varepsilon$.
Let $\mathbf{Y}_{1}, \mathbf{Y}_{2}, \ldots, \mathbf{Y}_{s}$ be the successive subsequences of $\mathbf{X}_{k}$ produced by the partition $\pi_{k}$. Let $N_{j v}(j=1,2, \ldots, s ; v=1,2, \ldots, \sigma)$ be the number of occurrences of $a_{v}$ among the elements of the subsequence $\mathbf{Y}_{j}$. The $N_{j v}$ are random variables subject to the multinomial distribution determined by the probabilities $p_{v}$, and for different. $j$ they are independent. We have

$$
\delta_{k}=\operatorname{Max}_{1 \leqq j \leqq s}\left\|\sum_{v=1}^{\sigma}\left(\frac{N_{j v}}{m_{j}}-p_{v}\right) a_{v}\right\| \leqq A M_{k}
$$

where

$$
A=\operatorname{Max}_{v}\left\|a_{v}\right\| \quad \text { and } \quad M_{k}=\operatorname{Max}_{1 \leqq j \leqq s} \sum_{v=1}^{\sigma}\left|\frac{N_{j v}}{m_{j}}-p_{v}\right|
$$

Thus we need to show

$$
\begin{equation*}
\sum_{k=1}^{\infty} \mathbf{P}\left\{M_{k} \geqq \varepsilon\right\}<\infty \tag{13}
\end{equation*}
$$

Since the following inclusion (implication) of events holds

$$
\left\{M_{k} \geqq \varepsilon\right\}=\bigcup_{j=1}^{s(k)}\left\{\sum_{v=1}^{\sigma}\left|\frac{N_{j v}}{m_{j}}-p_{v}\right| \geqq \varepsilon\right\} \cong \bigcup_{j=1}^{s(k)} \bigcup_{v=1}^{\sigma}\left\{\left|\frac{N_{j v}}{m_{j}}-p_{v}\right| \geqq \frac{\varepsilon}{\sigma}\right\},
$$

we obtain for the probabilities of the complementary events

$$
\begin{align*}
& \mathbf{P}\left\{M_{k}<\varepsilon\right\} \geqq \mathbf{P} \bigcap_{j=1}^{s(k)} \bigcap_{v=1}^{\sigma}\left\{\left|\frac{N_{j v}}{m_{j}}-p_{v}\right|<\frac{\varepsilon}{\sigma}\right\}=  \tag{14}\\
& \quad=\prod_{j=1}^{s(k)} \mathbf{P} \bigcap_{v=1}^{\sigma}\left\{\left|\frac{N_{j v}}{m_{j}}-p_{v}\right|<\frac{\varepsilon}{\sigma} \geqq \prod_{j=1}^{s(k)}\left[1-\sum_{v=1}^{\sigma} \mathbf{P}\left\{\left|\frac{N_{j v}}{m_{j}}-p_{j}\right| \geqq \frac{\varepsilon}{\sigma}\right\}\right] .\right.
\end{align*}
$$

The equality in (14) is due to the fact that we are dealing with the intersection (conjunction) of independent events.

Suppose $N$ is the number of "successes" in a sequence of $m$ Bernoulli trials with probability $p$ for success. We have then the following fact [3]: given any $\alpha>1$, for all sufficiently large $m$ (depending only on $\alpha$ and $p$ )

$$
\mathbf{P}\left\{\frac{|N-m p|}{[m p(1-p)]^{1 / 2}} \geqq(2 \alpha \log m)^{1 / 2}\right\}<\frac{1}{m^{\alpha}} .
$$

It follows that for any $\varepsilon>0$

$$
\begin{equation*}
\mathbf{P}\left\{\left|\frac{N}{m}-p\right| \geqq \varepsilon\right\}<\frac{1}{m^{\alpha}} \tag{15}
\end{equation*}
$$

for all sufficiently large $m$ (depending only on $\varepsilon, \alpha$ and $p$ ). If the inequality (15) is used for the probabilities on the right hand side of (14) we obtain

$$
\mathbf{P}\left\{M_{k}<\varepsilon\right\} \geqq \prod_{j=1}^{s(k)}\left[1-\frac{\sigma}{m_{j}^{\alpha}}\right]
$$

which is valid for all sufficiently large $k$, since (9) implies that then all the $m_{j}$ are large enough. But (9) and (10) show that for suitable constants $C$ and $C^{\prime}$

$$
\begin{equation*}
\mathbf{P}\left\{M_{k}<\varepsilon\right\} \geqq\left(1-C k^{-\alpha \beta}\right)^{C^{\prime} k^{1-\beta}}=1-O\left(k^{-\gamma}\right) \tag{16}
\end{equation*}
$$

where $\gamma=\alpha \beta+\beta-1$. Since $\alpha>1$ was arbitrary we may suppose $\gamma>1$, so that $\mathbf{P}\left\{M_{k} \geqq \varepsilon\right\}=O\left(k^{-\gamma}\right)$ and (13) follows. This concludes the proof of Theorem 2.

## References

[1] K. Yosida, Functional Analysis (Heidelberg-New York, 1968), Chapter XI.
[2] H. Trotter, On the product of semi-groups of operators, Proc. Amer. Math. Soc., $10^{\circ}$ (1959), 545-551.
[3] W. Feller, Introduction to Probability Theory and its Applications. I (New York, 1957). See Section VIII. 4, especially equ. (4. 5).

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