# Functional Integral Solution of the Complex Diffusion Equation with Complex Quadratic Potential in a Classical Path Space 

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#### Abstract

Based on our result of constructing the real integral solution of the complex diffusion equation, we successfully generate, via the pull back operation in the space of classical paths $\Gamma^{k l a}$, the solution of the complex diffusion equation with complex quadratic potential as a Feynmannian path integral with built-in Feynmannian measure in $\Gamma^{k l a}$.


## 1. Motivation

This article extends the framework set up by Shaharir [17, 19, 20] to formulate the Feynman integral. Intrinsically, Shaharir's framework aspired to generalise the work done by Albeverio and Hoegh-Krohn [2,3] and to generate a path integral over the Riemannian manifold which possesses the qualitative characteristics of the original Feynman integral, specifically the concept of 'sum over all the classical paths'. The central part of our work encompasses the effort to derive a 'Feynmannian' integral (lending Shaharir's [17] connotation) or simply a Feynman integral type from the path integral solution $\psi$ of the complex diffusion equation in the space of $n$-tuples set of complex number $C^{\boldsymbol{n}}$ :

$$
\begin{equation*}
\psi_{t}=\alpha \Delta \psi+\beta V(\boldsymbol{q}) \psi \tag{1.1}
\end{equation*}
$$

$\alpha, \beta \in \boldsymbol{C}$, complex numbers and $V(\boldsymbol{q})$ a complex quadratic potential. Subsript $t$ represents partial differentiation with respect to $t$ and $\Delta$ is the Laplacean. Essentially this is done via a refinement of the complex Hilbert space
$H=\left\{\gamma \in L^{2,2}\left([\mathbf{0}, t] ; \boldsymbol{C}^{\boldsymbol{n}}\right), \gamma(t)=\mathbf{0}\right.$ almost everywhere, and for any $\gamma$ and $\eta$ in the $L^{2,2}$

$$
\begin{equation*}
\left.\langle\gamma \mid \eta\rangle=\int_{0}^{t}(\dot{\gamma}(s) \cdot \dot{\bar{\eta}}(s)+\gamma(s) \cdot \bar{\eta}(s)) d s\right\} \tag{1.2}
\end{equation*}
$$

into

$$
\begin{align*}
\Gamma^{k l a} & =\left\{\text { classical paths over } \boldsymbol{C}^{\boldsymbol{n}} \text { with the inner product as in } H\right\} \\
& =\left\{\gamma: \mathbf{m} \ddot{\gamma}=-\nabla V(\gamma), \gamma(0)=\boldsymbol{q}_{\mathbf{0}}, \gamma(t)=\boldsymbol{q} \text { and } V \text { a potential }\right\} \tag{1.3}
\end{align*}
$$

$L^{2,2}$ is a space of twice differentiable function and each derivative is a Lebesgue square integrable.

This 'Feynmannian' integral is found to be more in agreement with the qualitative aspects of the original Feynman path integral. This article clearly shows the relation between this integral with the original Feynman path integral and the real integral via path integral solution for the complex diffusion equation with complex quadratic potential in $\boldsymbol{C}^{\boldsymbol{n}}$.

## 2. Principal results

In the following we give the necessary results to produce our main conclusion in Theorem 5.

Lemma 1. (Separable complex Hilbert space, H)
$H$ is a separable complex Hilbert space.
Proof. We'll show that $H$ is a direct sum of two separable complex Hilbert spaces $H_{1}$ and $H_{2}$ (i.e. $H_{1} \oplus H_{2}$ ), where
$H_{1}=\left\{\gamma \in L^{2,2}\left([0, t] ; \boldsymbol{C}^{\boldsymbol{n}}\right), \quad \gamma(t)=\mathbf{0}\right.$ almost everywhere, and for each pair $\gamma_{1}, \gamma_{2}$,

$$
\left.\left\langle\gamma_{1} \mid \gamma_{2}\right\rangle_{1}=\int_{0}^{t} \gamma_{1}(s) \cdot \bar{\gamma}_{2}(s) d s\right\}
$$

and
$H_{2}=\left\{\eta \in L^{2,2}\left([0, t] ; \boldsymbol{C}^{\boldsymbol{n}}\right), \eta(t)=\mathbf{0} \quad\right.$ almost everywhere, and for each pair $\eta_{1}, \eta_{2}$,

$$
\left.\left\langle\eta_{1} \mid \eta_{2}\right\rangle_{2}=\int_{0}^{t} \dot{\eta}_{1}(s) \cdot \dot{\bar{\eta}}_{2}(s) d s\right\} .
$$

The facts that $H_{i}, i=1,2$ are complex vector spaces with an inner product $\langle\cdot \mid \cdot\rangle_{i}$ are trivial. The norms $\|\cdot\|_{i}$ in $H_{i}$, which are determined by the above inner products, are given respectively by

$$
\|\gamma\|_{1}=\left(\int_{0}^{t}|\gamma(s)|^{2} d s\right)^{1 / 2},\|\eta\|_{2}=\left(\int_{0}^{t}|\dot{\eta}(s)|^{2} d s\right)^{1 / 2}
$$

Refering to Weierstrass's theorem (refer Kufner et al. [10]) every path $\gamma, \eta \in L^{2,2}\left([0, t] ; \boldsymbol{C}^{\boldsymbol{n}}\right)$ can be approximated in the normed space by a polynomial sequence, i.e.

$$
\begin{aligned}
& \gamma_{n}(s)=\sum_{|j| \leq n} a_{j} s^{j}, \gamma_{n}(t)=0 \text { and } a_{n+1}=a_{n+2}=\cdots=0, \\
& \eta_{n}(s)=\sum_{|j| \leq n} b_{j} s^{j}, \eta_{n}(t)=0 \text { and } b_{n+1}=b_{n+2}=\cdots=0, \\
& n=0,1,2, \cdots, a_{j}, b_{j} \in \boldsymbol{C}, \text { such that }\left\|\gamma-\gamma_{n}\right\|_{1}<\varepsilon,\left\|\eta-\eta_{n}\right\|_{2}<\varepsilon .
\end{aligned}
$$

Consequently,

$$
\left\|\gamma-\gamma_{n}\right\|_{1}^{2}=\int_{0}^{t}\left|\gamma-\gamma_{n}\right|^{2} d s,\left\|\eta-\eta_{n}\right\|_{2}^{2}=\int_{0}^{t}\left|\dot{\eta}-\left(\frac{j}{s}\right) \eta_{n}\right|^{2} d s,
$$

show that $\gamma$ and $\dot{\eta}$ are respectively limits of the sequences $\left\{\gamma_{n}\right\},\left\{\eta_{n}\right\}$ in $L^{2,2}$. Via the completion of $L^{2,2}$, the sequences $\left\{\gamma_{n}\right\},\left\{\eta_{n}\right\}$ converge separately to $\gamma$ and $\dot{\eta}$, and as a result

$$
\lim _{n}\left\|\gamma-\gamma_{n}\right\|_{1}^{2}=\int_{0}^{t}\left|\gamma-\gamma_{n}\right|^{2} d s, \quad \lim _{n}\left\|\eta-\eta_{n}\right\|_{2}^{2}=\int_{0}^{t}\left|\dot{\eta}-\left(\frac{1}{2}\right) \eta_{n}\right|^{2} d s
$$

These outcomes imply Cauchy sequences in $L^{2,2}$ (i.e. converge in the mean square) and subequently $\quad \gamma \in H_{1}, \dot{\eta} \in H_{2}$, and $H_{1}, H_{2}$ are complete. These conclusions guarantee that both $H_{1}, H_{2}$ are complex Hilbert spaces.

In order to show that $H_{1}, H_{2}$ are separable, we consider the sets $D_{1}, D_{2}$ of all paths $\gamma, \eta$ which originate from $H_{1}, H_{2}$ respectively. With reference to Naimark's theorem [12] (i.e. for each coefficient $a_{j}$ of the polynomial $\gamma_{n}$, there exists a rational number coefficient $u_{j}$ of the polynomial $\Gamma_{n}$ such that
$\left|a_{j}-u_{j}\right|<\left(\varepsilon / x d^{|j|}\right), \varepsilon>0,\left|s^{j}\right|<d^{|j|}, d$ a positive number, and $x$ the number of coefficients $a_{j}$, there exist in the sets $D_{1}, D_{2}$ other polynomials $\Gamma_{n}$ and $\Pi_{n}$, where $u_{j}$ and $v_{j}$ are the respective rational coefficients, such that

$$
\left\|\gamma_{n}-\Gamma_{n}\right\|_{1}<\varepsilon, \quad\left\|\eta_{n}-\Pi_{n}\right\|_{2}<\varepsilon .
$$

Moreover $\gamma_{n}(t)=\eta_{n}(t)=0$, and accordingly $u_{n+1}=u_{n+2}=\cdots=0, v_{n+1}=v_{n+2}=\cdots=0$ for some integer $n, n=0,1,2, \cdots$. As a consequence, all the polynomials $\Gamma_{n}, \Pi_{n} ; n=0,1, \cdots$ with the rational coefficients $u_{j}$ and $v_{j}$, of the sets $D_{1}, D_{2}$ are countable. These are due to the fact that the sets can be numerated by the natural numbers $n$.

The sets $D_{1}, D_{2}$ can be shown as follows to be dense everywhere in $H_{1}, H_{2}$ respectively. Let us consider any paths $\gamma \in H_{1}$ and $\eta \in H_{2}$, where $\Gamma_{n}$ and $\Pi_{n}$ respectively represent the polynomials with rational coefficients $u_{j}$ and $v_{j}$. For any $\varepsilon>0$ and since $\gamma, \eta \in L^{2}$, then there exist $\gamma_{n} \in D_{1}, \quad \eta_{n} \in D_{2}$, such that $\left\|\gamma-\gamma_{n}\right\|_{1}<\varepsilon,\left\|\eta-\eta_{n}\right\|_{2}<\varepsilon$. (via Weierstrass's theorem) and $\left\|\gamma_{n}-\Gamma_{n}\right\|_{1}<\varepsilon,\left\|\eta_{n}-\Pi_{n}\right\|_{2}<\varepsilon$. (via Naimark's theorem).
As a result, we obtain

$$
\begin{gathered}
\left\|y-\Gamma_{n}\right\|_{1} \leq\left\|\gamma-\gamma_{n}\right\|_{1}+\left\|\gamma_{n}-\Gamma_{n}\right\|_{1}=2 \varepsilon \\
\left\|\eta-\Pi_{n}\right\|_{2} \leq\left\|\eta-\eta_{n}\right\|_{2}+\left\|\eta_{n}-\Pi_{n}\right\|_{2}=2 \varepsilon
\end{gathered}
$$

which indicate that $D_{1}$ and $D_{2}$ are respectively dense everywhere in $H_{1}, H_{2}$, and thus demonstrate that $H_{1}, H_{2}$ are separable.

The above results show that $H_{1}, H_{2}$ are separable complex Hilbert spaces. Finally we proceed by referring to a theorem by Prugovecki [14] which states that: if $H_{1}, H_{2}, \cdots$ represent a finite or infinite sequence of countable separable Hilbert spaces, then the direct sum $\oplus_{n} H_{n}=H$ corresponds to a separable Hilbert spaces. As a consequence of this theorem, the abovementioned direct sum of $H_{1} \oplus H_{2}$ will form a separable complex Hilbert space with an inner product $\langle\cdot \mid \cdot\rangle$, as stated in the proposition of Lemma 1.

The form of the inner product (1.2) on the $L^{2,2}(G)$ space, $G$ a domain in $R^{n}$, has been briefly mentioned in Rektorys [15] and Griffel [8].

Corollary 2. (construction of $\Gamma^{\text {kla }}$ )
If

$$
\Gamma^{\text {kla }}=\left\{\begin{array}{l}
\sigma \in L^{2,2}\left([0, t] ; \boldsymbol{C}^{\boldsymbol{n}}\right), \text { m } \ddot{\sigma}=-\nabla(\gamma-\boldsymbol{q}), \text { where } m \text { mass of the } \\
\text { particle which is being influenced by a potential } \\
\boldsymbol{V}(\mathbf{z})=(1 / 2) \mathbf{z}^{\boldsymbol{T}} \Omega \overline{\mathbf{z}}+\boldsymbol{L} \cdot \mathbf{z}+P, \Omega \text { Hermitean matrice, } \boldsymbol{L} \text { vector } \\
\text { P complex constant } ; \sigma(t)=\mathbf{0} \text { almost everywhere, } \\
\text { and for any } \sigma \text { and } \beta \text { in the set } \\
\langle\sigma \mid \rho\rangle=\frac{p}{t} \int_{0}^{t}\left(\dot{\sigma}(s) \cdot \dot{\bar{\rho}}(s)+\sigma(s)^{T}\left(\frac{\Omega}{m}\right) \bar{\rho}(s)\right) d s, p>0,
\end{array}\right\}
$$

then $\Gamma^{k l a}$ is a closed subspace of $H$.
The proof of Corollary 2 is given after the following remarks.
Remarks 3. (form of the classical path $\sigma$ of interest)
The substitution $\sigma=\gamma-\boldsymbol{q}$ is used to obtain $\Gamma^{k l a}$ instead of the space

$$
\left\{\begin{array}{l}
\gamma \in L^{2,2}\left([0, t] ; \boldsymbol{C}^{\boldsymbol{n}}\right), \ddot{\gamma}=-\nabla V(\gamma), \gamma(0)=\boldsymbol{q}_{0}, \gamma(t)=\boldsymbol{q} \text { almost everywhere, } \\
\text { and for any pair of } \gamma \text { and } \eta\langle\gamma \mid \eta\rangle=\frac{p}{t} \int_{0}^{t}\left(\dot{\gamma}(s) \cdot \dot{\bar{\eta}}(s)+\gamma(s)^{T}\left(\frac{\Omega}{m}\right) \bar{\eta}(s)\right) d s .
\end{array}\right\}
$$

This implies that $\Gamma^{k l a}$ corresponds to the space of classical paths $\sigma=\gamma-\boldsymbol{q}$. The path $\sigma$, mentioned in the above Corollary 2 , is formed from the consideration of the complex quadratic potential

$$
V(\boldsymbol{q})=\left(\frac{1}{2}\right) \boldsymbol{q}^{\boldsymbol{T}} \Omega \boldsymbol{q}+\boldsymbol{L} \cdot \boldsymbol{q}+P
$$

together with the Euler-Lagrange equation

$$
\begin{aligned}
\ddot{\gamma} & =-\left(\frac{1}{m}\right)(\Omega \gamma+\boldsymbol{L}), \gamma(0)=\boldsymbol{q}_{0}, \gamma(t)=\boldsymbol{q}, \text { i.e. the path } \sigma=\gamma-\boldsymbol{q} \\
\Rightarrow \ddot{\sigma} & =-\left(\frac{\mathbf{1}}{\mathbf{m}}\right)(\Omega(\sigma+\boldsymbol{q})+\boldsymbol{L}), \sigma(0)=\boldsymbol{q}_{0}-\boldsymbol{q}, \sigma(t)=\mathbf{0}
\end{aligned}
$$

## Proof of Corollary 2

It can be seen that $\Gamma^{k l a} \subset L^{2,2}$ is a vector subspace of $L^{2,2}$, since $\sigma, \rho \in \Gamma^{k l a}$ implies that $a \sigma+b \rho$ is an element of $\Gamma^{k l a}$ for all scalars $a, b \in \boldsymbol{C}$. Moreover $\quad \Gamma^{k l a} \subset L^{2,2}$ forms a closed subspace of $L^{2,2}$ (accordingly of $H$ ): if $\Gamma^{k l a}$ is a vector subspace and $\left\{\sigma_{n}\right\} \subset \Gamma^{k l a}$, then $\sigma_{n} \rightarrow \sigma$ implies $\sigma \in \Gamma^{k l a}$.
With reference to the proof of Lemma 1, evidently we have in $\Gamma^{k l a} \subset L^{2,2}$

$$
\left\langle\sigma-\sigma_{n} \mid \sigma-\sigma_{n}\right\rangle_{\Gamma^{k l a}}=\left\|\sigma-\sigma_{n}\right\|_{\Gamma^{k l a}}^{2},
$$

and

$$
\left\|\sigma-\sigma_{n}\right\|_{\Gamma^{k l a}}^{2}=\frac{p}{t} \int_{0}^{t}\left(\left(\frac{\Omega}{m}\right)\left|\sigma-\sigma_{n}\right|^{2}+\left|\dot{\sigma}-\left(\frac{j}{s}\right) \sigma_{n}\right|^{2}\right) d s .
$$

These show that the convergence of $\sigma_{n} \rightarrow \sigma$ signifies that the simultaneous convergence of the sequence $\left\{\sigma_{n}\right\}$ and the corresponding sequence of derivatives $\left\{\dot{\sigma}_{n}\right\}$ in $\Gamma^{k l a}$ to their respective limits $\sigma_{n}$ and $\dot{\sigma}_{n}$. Consequently, via the completeness of $L^{2,2}$,

$$
\lim _{n} \frac{p}{t} \int_{0}^{t}\left(\left(\frac{\Omega}{m}\right)\left|\sigma-\sigma_{n}\right|^{2}+\left|\dot{\sigma}-\left(\frac{j}{s}\right) \sigma_{n}\right|^{2}\right) d s=0
$$

this result implies Cauchy sequence in $L^{2,2}$ and thus for $\sigma, \dot{\sigma} \in \Gamma^{k l a}$, we have a complete $\Gamma^{k l a}$. This implies that $\Gamma^{k l a}$ represents a closed subspace of $H$.

Lemma 4. (use of the pull-back operation in $\Gamma^{k l a}$ )

1. Let $C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right]$ be a set of twice differentiable smooth mapping from $(0, t)$ to unitary space $\boldsymbol{C}^{\boldsymbol{n}}$ such that $m \ddot{\sigma}=-\nabla W(\sigma), W(\sigma)=V(\gamma-\boldsymbol{q})$. $m$ stands for the particle's mass which is being influenced by the quadratic potential $V(\mathbf{z})=\left(\frac{1}{2}\right) \mathbf{z}^{\boldsymbol{T}} \Omega \overline{\mathbf{z}}+\boldsymbol{L} \cdot \mathbf{z}+P, \Omega$ a nonsingular Hermitean matrix (if $\Omega \neq 0$ ),
$\mathbf{L}$ vector, $P$ a complex constant; $\sigma(t)=\mathbf{0}$, and the inner product in $\boldsymbol{C}^{\boldsymbol{n}}$ is given by $g_{1}: g_{1}(a, b)=a \cdot \bar{b}$.
Let

$$
e v_{s}: C_{0}\left[(0, t) ; C^{\boldsymbol{n}}\right] \rightarrow \boldsymbol{C}^{\boldsymbol{n}}, \quad \sigma \mapsto e v_{s}(\sigma)=\sigma(s),
$$

represents the evaluation map at $s$. As a result $<,>_{1}$ is an inner product on $C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right]$ and is defined by

$$
\langle\sigma, \rho\rangle_{1}=g_{1}(\sigma(s), \rho(s))=\sigma(s) \cdot \bar{\rho}(s)
$$

$g_{1}^{\prime}: \boldsymbol{C}^{\boldsymbol{n}} \times \boldsymbol{C}^{\boldsymbol{n}} \rightarrow \boldsymbol{C},(\sigma, \rho) \mapsto \sigma \cdot \rho \quad$ is the 'pull-back' by the mapping $e v_{s}$,

$$
<,>_{1}=e v_{s}^{*} g_{1}
$$

2. The inner product $<,>_{2}$ on $C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right]$, which is defined by $\langle\sigma, \rho\rangle_{2}=\dot{\sigma}(s) \cdot \dot{\bar{\rho}}(s)$, is the pull-back $g_{2}$ by the mapping

$$
d e v_{s}: C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right] \rightarrow \boldsymbol{C}^{\boldsymbol{n}}, \sigma \mapsto \operatorname{dev}_{s}(\sigma)=\dot{\sigma}(s)
$$

3. The inner product on $\Gamma^{k l a}$, which is defined in the Corollary 2, is a pull-back of the metric $g\left(\right.$ i.e. $\left.g=g_{1} \oplus g_{2}\right)$

$$
g \equiv\left(e v_{s_{0}, a}+d e v_{s_{0}, a}\right): \Gamma^{k l a} \rightarrow C^{n}
$$

for the constant a and fixed $s_{0}$,

$$
g \equiv\left(e v_{s_{0}, a}+d e v_{s_{0}, a}\right)(a)=a\left(\sigma\left(s_{0}\right)+\dot{\sigma}(s)\right)
$$

and also the inner product $g$ is the pull-back <,>, a metric on $\Gamma^{\text {kla }}$, by the mapping

$$
F: C^{n} \rightarrow \Gamma^{k l a}, z \mapsto F(z)=\sigma
$$

$\sigma$ as in Remarks 3.

## Proof.

1. Since $e v_{s}$ is linear and via the definition of the Frechet derivative (see Abraham et al. [1], Lang [11]),

$$
D\left(e v_{s}\right)\left(\rho_{0}\right)(\eta)=\eta(s)
$$

for any $\rho_{0} \in C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right]$ and $\eta \in T_{\rho_{0}} C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right]$, the tangent space of $C_{0}\left[(0, t) ; \boldsymbol{C}^{\boldsymbol{n}}\right]$ at $\rho_{0}$. From the definition of the pull-back (see Abraham et al. [1])

$$
\begin{aligned}
\langle\sigma, \lambda\rangle_{1} & =e v_{s}^{*}\left(g_{1}\right)(\sigma, \lambda) ; \sigma, \lambda \in C_{0}\left[(0, t) ; C^{\boldsymbol{n}}\right] \\
\langle\sigma, \lambda\rangle_{1} & =g_{1}\left(e v_{s}\left(\rho_{0}\right)\right)\left(T_{\rho_{0}} e v_{s}(\sigma), T_{\rho_{0}} e v_{s}(\lambda)\right) \\
& =g_{1}\left(\rho_{0}(s)\right)(\sigma(s), \lambda(s))=\sigma(s) \cdot \bar{\lambda}(s)
\end{aligned}
$$

2. Since $d e v_{s}$ is linear,

$$
T_{\rho_{0}} \operatorname{dev}_{s}(\eta)=\left(\operatorname{dev}_{s}\left(\rho_{0}\right), D\left(\operatorname{dev}_{s}\left(\rho_{0}\right)(\eta)\right)=\left(\operatorname{dev}_{s}\left(\rho_{0}\right), \dot{\eta}(s)\right)\right.
$$

and similarly as part 1 above mentioned,

$$
\begin{aligned}
\langle\sigma, \lambda\rangle_{2}=\operatorname{dev}_{s}{ }^{*} g_{2}(\sigma, \lambda) & =g_{2}\left(\operatorname{dev}_{s}\left(\rho_{0}\right)\right)\left(T_{\rho_{0}} \operatorname{dev}_{s}(\sigma), T_{\rho_{0}} \operatorname{dev}_{s}(\lambda)\right) \\
& =g_{2}\left(\operatorname{dev}_{s}\left(\rho_{0}\right)\right)(\dot{\sigma}(s), \dot{\lambda}(s))=\dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s) .
\end{aligned}
$$

3. (a) The first part: We have to show that

$$
\begin{align*}
\langle\sigma, \lambda\rangle_{2} & =g(\sigma, \lambda) \\
& =\operatorname{dev}_{s_{0}, a}{ }^{*} g_{2}(\sigma, \lambda)+e v_{s_{0}, a}{ }^{*} g_{1}(\sigma, \lambda) \\
& =a^{2}\left(\dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)+\sigma\left(s_{0}\right) \cdot \bar{\lambda}\left(s_{0}\right)\right) ; \quad s_{0} \quad \text { is fixed } \tag{1.4}
\end{align*}
$$

This result is obtained with reference to the fact that for any constant of motion $T$ (energy for a conservative system as similar to this case), we have

$$
T\left(\sigma\left(s_{0}\right)\right)=\frac{1}{t} \int_{0}^{t} T(\sigma(s)) d s
$$

where $T(\sigma(s))=\left(\frac{m}{2}\right)|\dot{\sigma}(s)|^{2}+V(\sigma(s)), m$ is the particle's mass with potential $V$. Accordingly for $\sigma$ and $\lambda$ in $\Gamma^{k l a}$,

$$
T\left(\sigma\left(s_{0}\right)\right)+T\left(\lambda\left(s_{0}\right)\right)=\frac{1}{t} \int_{0}^{t}(T(\sigma(s)+T(\lambda(s)) d s
$$

i.e.,

$$
\begin{aligned}
\left(\frac{m}{2}\right)\left(\left|\dot{\sigma}\left(s_{0}\right)\right|^{2}\right. & \left.+\left|\dot{\lambda}\left(s_{0}\right)\right|^{2}\right)+\left(V\left(\lambda\left(s_{0}\right)\right)+V\left(\sigma\left(s_{0}\right)\right)\right) \\
& =\frac{1}{t} \int_{0}^{t}\left(\frac{m}{2}\right)\left(|\dot{\sigma}(s)|^{2}+|\dot{\lambda}(s)|^{2}+(V(\lambda(s))+V(\sigma(s))) d s .\right.
\end{aligned}
$$

As a consequence of this, for the case of a generalised complex harmonic oscillator, $V(\boldsymbol{q})=\left(\frac{1}{2}\right) \boldsymbol{q}^{\boldsymbol{T}} \Omega \boldsymbol{q}$, we have

$$
\begin{aligned}
& \left(\frac{m}{2}\right)\left|\dot{\sigma}\left(s_{0}\right)+\dot{\lambda}\left(s_{0}\right)\right|^{2}-m \dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)+\left(\frac{1}{2}\right)\left(\sigma\left(s_{0}\right)+\lambda\left(s_{0}\right)\right)^{T} \Omega\left(\bar{\sigma}\left(s_{0}\right)\right. \\
& \left.\quad+\bar{\lambda}\left(s_{0}\right)\right)-\sigma\left(s_{0}\right)^{T} \Omega \bar{\lambda}\left(s_{0}\right) \\
& =\frac{1}{t} \int_{0}^{t}\left(\left(\frac{m}{2}\right)|\dot{\sigma}(s)+\dot{\lambda}(s)|^{2}-m \dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\left(\frac{1}{2}\right)(\sigma(s)+\lambda(s))^{T} \Omega\left(\bar{\sigma}\left(s_{0}\right)+\bar{\lambda}\left(s_{0}\right)\right)\right. \\
& \left.\quad-\sigma(s)^{T} \Omega \bar{\lambda}(s)\right) d s
\end{aligned}
$$

or

$$
\begin{aligned}
& \left(\frac{m}{2}\right)\left|\dot{\gamma}\left(s_{0}\right)\right|^{2}+\left(\frac{1}{2}\right)\left(\gamma\left(s_{0}\right)\right)^{T} \Omega \bar{\gamma}\left(s_{0}\right)-\left\{m \dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)+\sigma\left(s_{0}\right)^{T} \Omega \bar{\lambda}\left(s_{0}\right)\right\} \\
& =\frac{1}{t} \int_{0}^{t}\left\{\left(\frac{m}{2}\right)|\dot{\gamma}(s)|^{2}+\left(\frac{1}{2}\right)(\gamma(s))^{T} \Omega \bar{\gamma}(s)-\left(m \dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\sigma(s)^{T} \Omega \bar{\lambda}(s)\right)\right\} d s,
\end{aligned}
$$

where $\sigma, \lambda, \gamma \in \Gamma^{k l a}$.
Since $\gamma=\alpha+\lambda \in \Gamma^{k l a}$, therefore $T\left(\sigma\left(s_{0}\right)\right)=\frac{1}{t} \int_{0}^{t} T(\gamma(s)) d s$, and accordingly we have the result

$$
\left\{\dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)+\sigma\left(s_{0}\right)^{T}\left(\frac{\Omega}{m}\right) \bar{\lambda}\left(s_{0}\right)\right\}=\frac{1}{t} \int_{0}^{t}\left\{\dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\sigma(s)^{T}\left(\frac{\Omega}{m}\right) \bar{\lambda}(s)\right\} d s
$$

or referring to (1.4), we have instead

$$
\begin{aligned}
a^{2}\left\{\dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)\right. & \left.+\sigma\left(s_{0}\right)^{T}\left(\frac{\Omega}{m}\right) \bar{\lambda}\left(s_{0}\right)\right\} \\
& =\frac{a^{2}}{t} \int_{0}^{t}\left\{\dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\sigma(s)^{T}\left(\frac{\Omega}{m}\right) \bar{\lambda}(s)\right\} d s=\langle\sigma, \lambda\rangle_{1}
\end{aligned}
$$

for the case of a generalised complex harmonic oscillator.
For the complex affine potential $V(\boldsymbol{q})=\boldsymbol{L} \cdot \boldsymbol{q}+P$,

$$
\begin{aligned}
& \left(\frac{m}{2}\right)\left|\dot{\sigma}\left(s_{0}\right)+\dot{\lambda}\left(s_{0}\right)\right|^{2}-m \dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)+\boldsymbol{L} \cdot\left(\sigma\left(\mathrm{s}_{0}\right)+\lambda\left(\mathrm{s}_{0}\right)\right)+C \\
& \quad=\frac{1}{t} \int_{0}^{t}\left\{\left(\frac{m}{2}\right)|\dot{\sigma}(s)+\dot{\lambda}(s)|^{2}-m \dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\boldsymbol{L} \cdot(\sigma(\mathrm{s})+\lambda(\mathrm{s}))+C\right\} d s
\end{aligned}
$$

where $C$ is a constant, or this can be expressed as

$$
\begin{aligned}
&\left(\frac{m}{2}\right)\left|\dot{\gamma}\left(s_{0}\right)\right|^{2}-m \dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)+\boldsymbol{L} \cdot \gamma\left(\mathrm{s}_{0}\right)+C \\
&=\frac{1}{t} \int_{0}^{t}\left\{\left(\frac{m}{2}\right)|\dot{\gamma}(s)|^{2}-m \dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\boldsymbol{L} \cdot \gamma(\mathrm{s})+C\right\} d s
\end{aligned}
$$

$\sigma, \lambda, \gamma \in \Gamma^{k l a}$ and therefore $\sigma+\lambda=\gamma \in \Gamma^{k l a}$. As a result, for complex affine potential, we obtain (again referring to (1.4))

$$
a^{2}\left\{\dot{\sigma}\left(s_{0}\right) \cdot \dot{\bar{\lambda}}\left(s_{0}\right)\right\}=\frac{a^{2}}{t} \int_{0}^{t}\{\dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)\} d s=\langle\sigma, \lambda\rangle_{2} .
$$

Consequently for the generalised complex quadratic potential, $V(\boldsymbol{q})=\left(\frac{1}{2}\right) \boldsymbol{q}^{\boldsymbol{T}} \Omega \overline{\boldsymbol{q}}+\boldsymbol{L} \cdot \boldsymbol{q}+P$, we have the result

$$
\langle\sigma, \lambda\rangle=\frac{a^{2}}{t} \int_{0}^{t}\left\{\dot{\sigma}(s) \cdot \dot{\bar{\lambda}}(s)+\sigma^{T}\left(\frac{\Omega}{m}\right) \bar{\lambda}(s)\right\} d s
$$

as stated in Lemma 1 and Corollary 2.
(b) The second part: this part shows that the inner product $g$ can be constructed from the pull-back <,>a metric on $\Gamma^{k l a}$ by the mapping $F$. According to Shaharir and Zainal [20], for the generalised complex quadratic potential $V(\boldsymbol{q})=\left(\frac{1}{2}\right) \boldsymbol{q}^{\boldsymbol{T}} \Omega \overline{\boldsymbol{q}}+\boldsymbol{L} \cdot \boldsymbol{q}+P, \Omega$ a nonsingular Hermitian matrix (if $\Omega \neq 0$ ) of size $n$, it can be written explicitly that

$$
F(\mathbf{z})(s)=\sigma(s)=\left\{\begin{array}{l}
\frac{-1}{2 m}\left(s^{2}-t^{2}\right) \boldsymbol{L}+(s-t) \underline{\ell} ; \quad \underline{\ell}=\frac{1}{t}\left(\boldsymbol{q}-\boldsymbol{q}_{0}+\frac{t^{2}}{2 m} \boldsymbol{L}\right), t>0 ; \text { if } \Omega=0 \\
\left.\sin ^{-1}[\sqrt{\Omega} t](\sin [\sqrt{\Omega} s]) \boldsymbol{X}-\sin [\sqrt{\Omega}(s-t)] \boldsymbol{Y}-\boldsymbol{X}\right), \\
\text { if } \Omega \neq 0 \text { and nonsingular, } \\
(\sqrt{\Omega})^{2}=\Omega, \boldsymbol{X}=\boldsymbol{q}+(\sqrt{\Omega})^{-1} \boldsymbol{L}, \boldsymbol{Y}=\boldsymbol{q}_{0}+(\sqrt{\Omega})^{-1} \boldsymbol{L} ; \text { whereas } \\
t \in\left(0, t_{\min }\right), t_{\min }=\min _{j}\left[\frac{\pi}{\sqrt{\lambda_{j}}}\right] \text { for } \alpha \beta<0 \quad \text { where } \\
|\sin [\sqrt{\Omega} t]|=\prod_{j=1}^{n} \sin \left[\sqrt{\lambda_{j}} t\right], \quad \text { or } t>0 \text { for } \alpha \beta>0, \text { or } \\
t_{\min }=\min _{j}\left[\frac{\pi}{\left.\operatorname{Im} \sqrt{\lambda_{j}}\right]}\right] \text { for } \alpha \beta \in \boldsymbol{C}
\end{array}\right.
$$

$\boldsymbol{L} \in \boldsymbol{C}^{\boldsymbol{n}}, \Omega=0$ for the complex affine potential and $\Omega \neq 0$ for the complex quadratic potential, $\lambda_{i}$ is the eigenvalue of $\Omega$.
Accordingly $D F\left(\mathbf{z}_{\mathbf{0}}\right)(\mathbf{z})=\sigma$, and as a result

$$
\begin{aligned}
g\left(\mathbf{z}_{\mathbf{0}}\right)(\mathbf{z}, \boldsymbol{w}) & =F^{*}(<,>)(\mathbf{z}, \boldsymbol{w}) \\
& =\left\langle T_{\mathbf{z}_{0}} F(\mathbf{z}), T_{\mathbf{z}_{0}} F(\boldsymbol{w})\right\rangle=\langle\sigma, \lambda\rangle \\
& =\frac{a^{2}}{t} \int_{0}^{t}\left\{\dot{\sigma}(s) \cdot \dot{\bar{\gamma}}(s)+\sigma(s)^{T}\left(\frac{\Omega}{m}\right) \bar{\lambda}(s)\right\} d s,
\end{aligned}
$$

and explicitly,

$$
g\left(\mathbf{z}_{\mathbf{0}}\right)(\mathbf{z}, \mathbf{z})=\left\{\begin{array}{l}
\frac{1}{2}\left|\left[\boldsymbol{q}_{\mathbf{0}}-\left(\boldsymbol{q}-\frac{t^{2} \boldsymbol{L}}{2 m}\right)\right]\right|^{2}+\frac{2}{m}\left|\left[\left(\boldsymbol{q}_{\mathbf{0}}-\boldsymbol{q}\right) \cdot \boldsymbol{L}\right]\right|, \quad \Omega=0 \\
{\left[\boldsymbol{Y}-\cos ^{-1}[\sqrt{\Omega} t] \boldsymbol{X}\right]^{T} \sum^{-1}\left[\boldsymbol{Y}-\cos ^{-1}[\sqrt{\Omega} t] \boldsymbol{X}\right]} \\
\quad+\left(\left(\boldsymbol{q}_{\mathbf{0}}-\boldsymbol{q}\right)^{T} \sqrt{\Omega}\left(\boldsymbol{q}_{\mathbf{0}}-\boldsymbol{q}\right)\right)+2\left|\left[\left(\boldsymbol{q}_{\mathbf{0}}-\boldsymbol{q}\right) \cdot \boldsymbol{L}\right]\right|, \Omega \neq 0, \\
\sum(t)=\frac{2 \alpha}{t}(\Omega)^{-1}(\tan [\Omega t]), t \in\left(0, t_{\min }\right), \text { where } t_{\min }=\min _{j} \frac{\pi}{\sqrt{\lambda_{j}}}
\end{array}\right.
$$

for $\alpha \beta<0$, or $t>0$ for $\alpha \beta>0$, or $t_{\min }=\min _{j} \frac{\pi}{\operatorname{Im}\left(\sqrt{\lambda_{j}}\right)}$ for $\alpha \beta \in \boldsymbol{C}$

Theorem 5. ('Feynmannian' integral)

Let $\varphi: \boldsymbol{C}^{\boldsymbol{n}} \rightarrow \boldsymbol{C}$ be Lebesgue integrable and a Fourier transform of a bounded complex Borel measure on $\boldsymbol{C}^{\boldsymbol{n}}$; then the solution of the complex diffusion equation (1.1) with a quadratic complex potential $V(\boldsymbol{q})=\left(\frac{1}{2}\right) \boldsymbol{q}^{\boldsymbol{T}} \Omega \overline{\boldsymbol{q}}+\boldsymbol{L} \cdot \boldsymbol{q}+P$, where $\Omega a$ nonsingular Hermitian matrix (if $\Omega \neq 0$ ), $L$ vector, $P$ a complex constant; can be reduced to a 'Feynmannian' path integral

$$
\begin{equation*}
R(\boldsymbol{q}, t)=\int_{\Gamma^{k l a}} \exp \left(\beta \int_{0}^{t} L(\sigma(s), \dot{\sigma}(s)) d s\right) d F(\sigma), \tag{1.5a}
\end{equation*}
$$

where $\quad L(\sigma(s), \dot{\sigma}(s))=\left(\frac{1}{2}\right) m|\dot{\sigma}(s)|^{2}-W(\sigma(s)), \quad W(\sigma)=V(\gamma-\boldsymbol{q})$.
The entity L represents the classical Lagrangian of a particle of mass

$$
\begin{aligned}
m= & \left(\frac{-1}{2 \alpha \beta}\right), \\
\Gamma^{k l a}= & \left\{\sigma \in L^{2,2}\left([0, t] ; \boldsymbol{C}^{n}\right), m \ddot{\sigma}=-\nabla W(\sigma), \sigma(t)=\mathbf{0}\right. \text { almost everywhere, } \\
& \text { and for any such functions } \sigma, \beta \text { the following inner product is well- } \\
& \text { defined as } \left.\langle\sigma \mid \rho\rangle=\frac{a^{2}}{t} \int_{0}^{t}\left(\dot{\sigma}(s) \cdot \dot{\bar{\rho}}(s)+\sigma(s)^{T}\left(\frac{\Omega}{m}\right) \bar{\rho}(s)\right) d s\right\} .
\end{aligned}
$$

$\Gamma^{k l a}$ is the particle's space of classical paths, and the Feynmannian measure $F$ is given by

$$
\begin{equation*}
d F=\left[T_{\rho_{0}}(G)\right]^{*}\left(d \mu_{\varphi, q, s}\right) ; \quad \rho_{0} \in \Gamma^{k l a} \tag{1.5b}
\end{equation*}
$$

$G$ is defined as in Lemma 4 part (3), where

$$
d \mu_{\varphi, q, s}=\left(S^{-1}\right)^{*}\left(d \mu_{\varphi, q}\right)
$$

and $S: \boldsymbol{C}^{\boldsymbol{n}} \rightarrow \boldsymbol{C}^{\boldsymbol{n}}, \mathbf{z} \mapsto \boldsymbol{q}_{\mathbf{0}}=\left\{\begin{array}{l}{\left[\frac{1}{2 \alpha t}\right]^{1 / 2} \cdot \mathbf{z}, \text { if } \Omega=0} \\ {[2 \alpha|\sqrt{\Omega} \tan (\sqrt{\Omega} t)|]^{-1 / 2} \cdot \mathbf{z}, \text { if } \Omega \neq 0}\end{array}\right.$
is the bounded complex Borel measure and the scalar mapping on $\mathbf{C l}^{\boldsymbol{n}}$ :

$$
d \mu_{\varphi, q}\left(\boldsymbol{q}_{\mathbf{0}}\right)=\exp \left(k \boldsymbol{q} \cdot \boldsymbol{q}_{\mathbf{0}}\right) d \mu_{\varphi}\left(\boldsymbol{q}_{\mathbf{0}}\right)
$$

or equivalently

$$
\mathfrak{J}(A)=\int_{(\boldsymbol{F} \circ T)^{-1}(A)} \exp \left(k \boldsymbol{q} \cdot \boldsymbol{q}_{\mathbf{0}}\right) d \mu_{\varphi}\left(\boldsymbol{q}_{\mathbf{0}}\right) .
$$

$F: \boldsymbol{C}^{\boldsymbol{n}} \rightarrow \Gamma^{k l a}, \mathbf{z} \mapsto \sigma$, where $\sigma$ is the classical path as in Lemma 4 part (3b) and

$$
T \equiv S^{-1}: \boldsymbol{C}^{n} \rightarrow \boldsymbol{C}^{n}, T(z)= \begin{cases}(2 t|\alpha|)^{1 / 2} \cdot \mathbf{z} & \text { if } \Omega=0 \\ (2 \alpha|\sqrt{\Omega} \tan [\sqrt{\Omega} t]|)^{1 / 2} \cdot \mathbf{z} & , \text { if } \Omega \neq 0\end{cases}
$$

Proof. The proof is given separately for the case $\Omega=0$ (complex affine potential) and $\Omega \neq 0$ (complex quadratic potential).

Case $\Omega=0$. From Shaharir and Zainal [20], the solution of the complex diffusion equation (1.1) is given by

$$
\begin{aligned}
& R(\boldsymbol{q}, t)= \int_{R^{n}}\left(\frac{1}{4 \pi \alpha t}\right)^{n / 2} \exp (-g(\mathbf{z}, \mathbf{z})) \exp \left(2 \beta \int_{0}^{t} V(\sigma(s) d s) \varphi(\boldsymbol{q}+\mathbf{z}) d \mathbf{z}\right. \\
&= \int_{C^{n}}\left(\frac{1}{4 \pi \alpha t}\right)^{n / 2} \exp \left(-\frac{\mathbf{z} \cdot \overline{\mathbf{z}}}{4 \alpha t}+\left(\frac{1}{2}\right) \rho t(\boldsymbol{q}+\mathbf{z}) \cdot \boldsymbol{L}+\frac{L^{2} \beta^{2} \alpha t^{3}}{12}+\rho P t\right) \\
& {\left[\int_{C^{n}} \exp \left(k(\boldsymbol{q}+\mathbf{z}) \cdot \boldsymbol{q}_{\mathbf{0}} d \mu_{\varphi}\left(\boldsymbol{q}_{\mathbf{0}}\right)\right] d \mathbf{z},\right.} \\
&= \int_{C^{n}}\left(\int_{C^{n}}\left(\frac{1}{4 \pi \alpha t}\right)^{n / 2} \exp \left(-\frac{\omega \cdot \bar{\omega}}{4 \alpha t}\right) d \omega\right) \exp \left[\left(-\left(\boldsymbol{q}_{\mathbf{0}}-\boldsymbol{q}_{\mathbf{0}}\right) \alpha t+k \boldsymbol{q} \cdot \overline{\boldsymbol{q}}_{\mathbf{0}}-k \alpha \beta t^{2} \boldsymbol{q}_{\mathbf{0}} \cdot \boldsymbol{L}\right.\right. \\
&\left.\quad+\left(\frac{L^{2} \beta^{2} \alpha t^{3}}{3}\right)+\rho P t\right] d \mu_{\varphi}\left(\boldsymbol{q}_{\mathbf{0}}\right) ; \quad \omega=z-2 \alpha t\left(k \boldsymbol{q}_{\mathbf{0}}+\left(\frac{\beta t}{2}\right) \boldsymbol{L}\right) \\
&= \int_{C^{n}} \exp \left(-g\left(\boldsymbol{q}_{\mathbf{0}}, \boldsymbol{q}_{\mathbf{0}}\right)\right) \exp \left(\beta \int_{0}^{t} V(\sigma(s)) d s\right) \exp \left(k \boldsymbol{q} \cdot \overline{\boldsymbol{q}}_{\mathbf{0}}\right) d \mu_{\varphi}\left(\boldsymbol{q}_{\mathbf{0}}\right),
\end{aligned}
$$

upon consideration of Fubini's theorem and characteristics of a 'multinormal complex distribution' (see Shaharir [17] and Andersen et al. [5]). As a result we obtain

$$
R(\boldsymbol{q}, t)=\int_{C^{n}} \exp \left(-g\left(\boldsymbol{q}_{\mathbf{0}}-\boldsymbol{q}_{\mathbf{0}}\right)\right) \exp \left(\beta \int_{0}^{t} V(\sigma(s)) d s\right) d \mu_{\varphi, q}\left(\boldsymbol{q}_{\mathbf{0}}\right)
$$

and

$$
\begin{aligned}
d \mu_{\varphi, q}\left(\boldsymbol{q}_{\mathbf{0}}\right) & =\exp \left(k \boldsymbol{q} \cdot \boldsymbol{q}_{\mathbf{0}}\right) d \mu_{\varphi}\left(\boldsymbol{q}_{\mathbf{0}}\right), \\
& =\int_{C^{n}} \exp \left(-\frac{e^{k p}}{2} g(\mathbf{z}, \mathbf{z})\right) \exp \left(\beta \int_{0}^{t} V(\sigma(s)) d s\right) d \mu_{\varphi, q, s}(\mathbf{z}),
\end{aligned}
$$

(where $p$ represents the phase of $1 / \alpha$ ) or in terms of the differential form, $d \mu_{\varphi, q, s}=\left(S^{-1}\right)^{*}\left(d \mu_{\varphi, q}\right)$, the pull-back of $d \mu_{\varphi, q}$ by $S^{-1}$.
Meanwhile we can write

$$
\begin{aligned}
R(\boldsymbol{q}, t) & =\int_{T_{b} C^{n}} \exp \left(-a g\left(\mathfrak{J}_{\rho_{0}}(\sigma) \cdot \mathfrak{I}_{\rho_{0}}(\sigma)\right)\right) \exp \left(\beta \int_{0}^{t} V\left(\mathfrak{I}_{\rho_{0}}(\sigma(s)) d s\right) d \mu_{\varphi, q, s}\left(\mathfrak{I}_{\rho_{0}}(\sigma)\right),\right. \\
a & =\left(\frac{1}{2}\right) e^{-k p}, b=\alpha \rho_{0}\left(s_{0}\right), \rho_{0} \in \Gamma^{k l a} .
\end{aligned}
$$

Together with the choice

$$
\mathfrak{I}_{\rho_{0}}(\sigma)=T_{\rho_{0}} G(\sigma),
$$

and $G$ as in Lemma 4 part (3), and via Lemma 4 and the general theory of the transformation of a variable of an integral with respect to a differential form, we obtain

$$
\begin{equation*}
R(\boldsymbol{q}, t)=\int_{T_{b} \Gamma^{k l a}} \exp (-a<\sigma, \sigma>) \exp \left(\beta \int_{0}^{t} V(\sigma(s)) d s\right) d F(\sigma) \tag{1.6a}
\end{equation*}
$$

or

$$
\begin{equation*}
R(\boldsymbol{q}, t)=\int_{T_{b} \Gamma k l a} \exp (-a<\gamma-\boldsymbol{q}, \gamma-\boldsymbol{q}>) \exp \left(\beta \int_{0}^{t} V(\gamma(s)-\boldsymbol{q}) d s\right) d F(\gamma-\boldsymbol{q}) \tag{1.6b}
\end{equation*}
$$

The first part of the theorem is obtained when we identify $T_{\rho_{0}} \Gamma^{k l a}$ as $\Gamma^{k l a}$ and upon considering the formulation of Albeverio and Hoegh-Krohn (as in [2]) which states that the normalised integral (or Fresnel integral) on $H$ is invariant under the group transformation of Euclidean type, and the result by Parthasarathy [13] which states that a measure in the function space (eg. Wiener measure) is invariant under such transformation.

We prove the last part of the theorem via the same arguments as proposed in the results of the second part of Lemma 4 and the theory relating to the Borel measure, together with the fact that the trasformation $T$ and $F$ are proper mappings. The measure $F$ is obtained via theorem of Radon-Nikodym (see Halmos [9]).

Case $\Omega \neq 0$. Referring to results obtained in Shaharir and Zainal [20], the solution of the complex diffusion equation (1.1) is given by

$$
\begin{gathered}
R(\boldsymbol{q}, t)=\int_{R^{n}}\left(\frac{1}{2 \pi\left|\sum(t)\right|}\right)^{n / 2} \exp (-g(\mathbf{z}, \mathbf{z})) \exp \left(2 \beta \int_{0}^{t} V(\sigma(s)) d s\right) \varphi\left(\mathbf{z}+\cos ^{-1}[\sqrt{\Omega} t] \boldsymbol{X}\right) d \mathbf{z}, \\
\boldsymbol{X}=\boldsymbol{q}+(\sqrt{\Omega})^{-1} \boldsymbol{L}, \quad \boldsymbol{Y}=\boldsymbol{q}_{\mathbf{0}}+(\sqrt{\Omega})^{-1} \boldsymbol{L} .
\end{gathered}
$$

or

$$
=\int_{R^{n}}\left(\frac{|\sqrt{\Omega}|}{4 \pi \alpha \mid \tan [\sqrt{\Omega} t]}\right)^{n / 2}(\cos [\sqrt{\Omega} t])^{-n / 2} \exp \left[\frac{-1}{4 \pi \alpha} \mathbf{z}^{T}\left(\sqrt{\Omega} \tan ^{-1}[\sqrt{\Omega} t]\right) \overline{\mathbf{z}}\right.
$$

$$
\begin{aligned}
& \left.+\frac{1}{4 \alpha} \boldsymbol{X}^{\boldsymbol{T}}(\sqrt{\Omega} \tan [\sqrt{\Omega} t]) \overline{\boldsymbol{X}}\right] \exp \left(n \beta\left(P-\frac{1}{2} \boldsymbol{L}^{\boldsymbol{T}}(\sqrt{\Omega})^{-1} \overline{\boldsymbol{L}} t\right)\right) \varphi\left(\mathbf{z}+\cos ^{-1}[\sqrt{\Omega} t] \boldsymbol{X}\right) d \mathbf{z} \\
& =\int_{C^{n}}\left(\int_{C^{n}}\left(\frac{|\sqrt{\Omega}|}{4 \pi \alpha|\tan [\sqrt{\Omega} t]|}\right)^{n / 2} \exp \left[\frac{-1}{4 \pi \alpha} \omega^{\boldsymbol{T}}\left(\sqrt{\Omega} \tan ^{-1}[\sqrt{\Omega} t]\right) \bar{\omega}\right] d \omega(\cos [\sqrt{\Omega} t])^{-n / 2}\right. \\
& \\
& \exp \left(-\alpha \boldsymbol{Y}^{\boldsymbol{T}}\left((\sqrt{\Omega})^{-1} \tan [\sqrt{\Omega} t]\right) \overline{\boldsymbol{Y}}+k \boldsymbol{X}^{\boldsymbol{T}} \cos ^{-1}[\sqrt{\Omega} t] \overline{\boldsymbol{X}}\right. \\
& \left.+\frac{1}{4 \alpha} \boldsymbol{X}^{\boldsymbol{T}}(\sqrt{\Omega} \tan [\sqrt{\Omega} t]) \overline{\boldsymbol{X}}\right) \exp \left(n \beta\left(P-\frac{1}{2} \boldsymbol{L}^{\boldsymbol{T}}(\sqrt{\Omega})^{-1} \overline{\boldsymbol{L}}\right) t\right) d u_{\varphi}(\boldsymbol{Y}), \\
& t \in\left(0, t_{\min }\right), t_{\min }=\min \frac{\pi}{\sqrt{\lambda_{j}}} \operatorname{for} \quad \alpha \beta<0 \quad ; \text { or } t>0 \quad \text { for } \alpha \beta>0 ; \text { or } \\
& t_{\min }=\min _{j} \frac{\pi}{\operatorname{Im}\left(\sqrt{\lambda_{j}}\right)} \text { for } \alpha \beta \in \boldsymbol{C}, \text { where } \omega=\mathbf{z}-2 k \alpha\left((\sqrt{\Omega})^{-1} \tan [\sqrt{\Omega} t]\right) \boldsymbol{Y},
\end{aligned}
$$

together with certain manipulation with respect to the second factor of the integrand and the condition that $\left(\alpha(\sqrt{\Omega})^{-1} \tan (\sqrt{\Omega} t)\right)$ is a positive semidefinite complex (hermitian) matrice with non negative real eigenvalues (see Gantmacher [7], Andersen et al. [5]);

$$
R(\boldsymbol{q}, t)=\int_{C^{n}} \exp \left(-g(\boldsymbol{Y}, \boldsymbol{Y}) \exp \left(\beta \int_{0}^{t} V(\sigma(s)) d s\right) d \mu_{\varphi, x}(\boldsymbol{Y})\right.
$$

where

$$
d \mu_{\phi, x}(\boldsymbol{Y})=\exp \left(k \boldsymbol{Y}^{\boldsymbol{T}} \cos ^{-1}[\sqrt{\Omega} t] \boldsymbol{Y}\right) d \mu_{\phi}(\boldsymbol{Y})
$$

and taking into account Fubini's theorem and characteristics of the 'complex multinormal distribution'. We can use the same arguments as the case $\Omega=0$ to obtain the result (1.5).

## Remarks 6

(a) $\quad T_{\rho_{0}} \Gamma^{k l a}$ represents the tangent vector set at $\beta_{0} \in \Gamma^{k l a}$. The space $T_{\rho_{0}} \Gamma^{k l a}$ is identified naturally with $\Gamma^{k l a}$ by associating a tangent vector with an element of $\Gamma^{k l a}$. Furthermore, $T_{\rho_{0}} \Gamma^{k l a}$ being a vector space, it can be naturally identified as isomorphic with respect to $\Gamma^{k l a}$. Nevertheless this natural isomorphism is no longer assumed when we generalise to a manifold. In this case it is necessary to treat the tangent spaces at various points as different (see for example, Crampin and Pirani [6]).
(b) Referring to the formulation of Albeverio and Hoegh-Krohn [2], the potential function $V$ as in Theorem 5 (affine complex and quadratic complex potentials) is derivable from the space $F\left(\boldsymbol{C}^{\boldsymbol{n}}\right)$, Fresnel integrable function space on $\boldsymbol{C}^{\boldsymbol{n}}$. Briefly $V$ is a Fourier transformation of $\mu$, which is a bounded complex Borel measure on $\boldsymbol{C}^{\boldsymbol{n}}$, i.e.

$$
\begin{equation*}
V(\boldsymbol{q})=\int_{C^{n}} \exp (k \boldsymbol{q} \cdot \overline{\boldsymbol{r}}) d \mu(\boldsymbol{r}) \tag{1.7}
\end{equation*}
$$

With reference to Parthasarathy [13],
(i) $\quad V(\boldsymbol{q})$ is uniformly continuous in a normed topology,
(ii) if $V_{1}(\boldsymbol{q})=V_{2}(\boldsymbol{q})$ for all $\boldsymbol{q} \in \boldsymbol{C}^{\boldsymbol{n}}$, then $\mu_{1}=\mu_{2}$,
(iii) convolution $\left(V_{1} * V_{2}\right)(\boldsymbol{q})=V_{1}(\boldsymbol{q}) V_{2}(\boldsymbol{q})$ for all $\boldsymbol{q} \in \boldsymbol{C}^{\boldsymbol{n}}$ and $\mu_{1}, \mu_{2}$ are bounded Borel measures on $C^{\boldsymbol{n}}$, and
(iv) $\bar{V}(\boldsymbol{q})=\overline{V(\boldsymbol{q})}$.

In addition, $V: \boldsymbol{C}^{\boldsymbol{n}} \rightarrow \boldsymbol{C}$ is a function $S^{\infty}$ whereby the first differential is bounded and all the higher order derivatives are most likely experiencing linear growth, i.e. (see Albeverio and Brzezniak [4])

$$
\left|D^{v} V(\boldsymbol{q})\right| \leq\left\{\begin{array}{cll}
M & , & |v|=1  \tag{1.8}\\
m(1+|\boldsymbol{q}|) & , & |v|>1
\end{array}\right.
$$

## 3. Conclusions

Theorem 5 generates a solution for the complex diffusion equation (1.1) with a complex quadratic potential $V(\boldsymbol{q})=\left(\frac{1}{2}\right) \boldsymbol{q}^{\boldsymbol{T}} \Omega \boldsymbol{q}+\boldsymbol{L} \cdot \boldsymbol{q}+P(\Omega$ is nonsingular Hermitean matrix, if $\Omega \neq 0$ ). This solution is in the form of a Feynmannian path integral (1.5a) with the measure ( 1.5 b ) on the classical path space $\Gamma^{k l a}$. This result enables the extension of our framework to formulate the Feynman integral in order to include complex quadratic potential in unitary space $C^{\boldsymbol{n}}$. As a matter of fact, this result shows qualitatively the connection between the Feynmannian integral with the original Feynman path integral and its real integral form. Clearly this is shown via the pull-back operation in $\Gamma^{k l a}$ with respect to the path integral solution of the complex diffusion equation with quadratic potential in $\boldsymbol{C}^{\boldsymbol{n}}$. The extension to the quadratic potential in Riemannian manifold can be done similarly.

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