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# Subnormal operators regarded as generalized observables and compound-system-type normal extension related to $\mathfrak{s u}(1,1)$ 

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#### Abstract

In this paper, subnormal operators, not necessarily bounded, are discussed as generalized observables. In order to describe not only the information about the probability distribution of the output data of their measurement but also the framework of their implementations, we introduce a new concept of a compound-system-type normal extension, and we derive the compound-system-type normal extension of a subnormal operator, which is defined from an irreducible unitary representation of the algebra $\mathfrak{s u}(1,1)$. The squeezed states are characterized as the eigenvectors of an operator from this viewpoint, and the squeezed states in multi-particle systems are shown to be the eigenvectors of the adjoints of these subnormal operators under a representation. The affine coherent states are discussed in the same context, as well.


## 1. Introduction

In quantum mechanics, observables are described by self-adjoint operators and the probability distributions of the output data of their measurement are determined by the spectral measures of those self-adjoint operators and the density operators of states.

When a linear operator has its spectral measure, it is a normal operator where its selfadjoint part and its skew-adjoint part commute with each other (lemma 3). In a broader sense, therefore, it can be regarded as a complexified observable. (Note that from this viewpoint, in the following, we will use the expression 'measurement of a normal operator' in this wider sense, even if the normal operator is not always self-adjoint.) However, the measurements in quantum systems, which are not necessarily the measurements of any observables, are described by the positive operator-valued measures (POVM), which are a generalization of spectral measures (definition 5 and lemma 4). In this paper, from these viewpoints, we try to treat the observables generalized even for the class of subnormal operators $\|$, which is known as a wider class including the class of normal operators. A subnormal operator is defined as the restriction of the normal operator into a narrower domain. As far as the authors know, such an idea generalizing observables was introduced by Yuen and Lax [1]. The pair of the normal operator and the wider domain is called its normal extension (definition 2). We can
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$\|$ The concept of subnormality was introduced by Halmos [2, 3].
define the POVM of a subnormal operator uniquely in a similar sense that we can define the spectral measure of a normal operator uniquely under some condition (lemma 19). By this correspondence, we will formulate the measurements of the subnormal operators which are not necessarily bounded. In this paper, we will not only investigate the POVMs of the subnormal operators but also give some examples of the framework of their implementations in a physical sense.

There are many cases where the adjoint operator of a subnormal operator has eigenvectors with continuous potency and an over-complete eigenvector system. In these cases, the POVM constructed from the over-complete eigenvector system is just the POVM of the subnormal operator (lemma 22). Thus the subnormal operator is closely related to eigenvectors with continuous potency and to over-complete function systems, and these relations are important for the discussion on the properties of the subnormal operator. This fact may give us an illusion that the adjoint of any operator with a point spectrum with continuous potency would be a subnormal operator. However, the subnormality is not necessarily guaranteed only by the condition that its adjoint has a point spectrum with continuous potency $\dagger$.

For example, an implementation of the measurement of a subnormal operator is already known for an actual system in quantum optics. Let $Q$ and $P$ be the multiplication operator and the $(-i)$-times differential operator on the Hilbert space $L^{2}(\mathbb{R})$. A POVM is constructed from the over-complete eigenvector system of the boson annihilation operator $a_{b}:=\sqrt{\frac{1}{2}}(Q+\mathrm{i} P)$ (known as the coherent states system). Then this POVM is just the POVM of the boson creation operator $a_{b}^{*}$ which is a subnormal operator. The measurement of this POVM has been implemented as shown in the following (see section 3, in detail), and is called the heterodyne measurement; this implementation is performed by the measurement of a normal operator on the compound system between the basic system (i.e. the system of interest where the measurement is originally discussed) and an additional ancillary system prepared appropriately. Note that this operation, of measuring a normal operator on the compound system by preparing an additional ancillary system, gives a kind of normal extension of the creation operator $a_{b}^{*}$. However, only giving the definition of the normal extension is not sufficient for discussing such a physical operation. For clarifying such a physical operation, in section 3, we will introduce a new concept of a compound-system-type normal extension, which describes not only the normal extension but also a framework of a physical operation (given in definition 26).

In section 4, under the circumstance where an irreducible unitary representation of the algebra $\mathfrak{s u}(1,1)$ is given, we will construct two types of operators which have point spectra with continuous potency, and will investigate what condition guarantees the subnormality of these operators. The coherent states of the algebra $\mathfrak{s u}(1,1)$ introduced by Perelomov [4], will be reinterpreted as the eigenvectors of these operators. Moreover, in section 6, we will derive the compound-system-type normal extensions of these operators when they are subnormal operators.

In section 5.1, from the relationship between the irreducible unitary representations of the algebra $\mathfrak{s u}(1,1)$ and those of the affine group ( $a x+b$ group), we will discuss what subnormal operators are related to the irreducible unitary representations of the affine group. Moreover, we will discuss the correspondence between the eigenvectors of this subnormal operator (or the coherent states of the algebra $\mathfrak{s u}(1,1)$ ) and the coherent states of the affine group. Hence we will show a relationship between our problem and the irreducible unitary representation of the affine group which is closely related to the continuous wavelet transform.

Next, in section 5.2, from the relationship between the representation of the algebra
$\mathfrak{s u}(1,1)$ and the squeezed states, it will be confirmed that the squeezed states can be described as the coherent states of the algebra $\mathfrak{s u}(1,1)$ in our context. In other words, the squeezed states are characterized as the eigenvectors of the operators (with point spectra with continuous potency) which are canonically constructed from an irreducible unitary representations of the algebra $\mathfrak{s u}(1,1)$. However, the adjoints of these operators are not necessarily subnormal operators and are not directly regarded as generalized observables.

We can easily confirm that the squeezed states are the eigenvectors of an operator with a point spectrum with continuous potency as follows; according to Yuen [5], let $b_{\mu, v}:=\mu a_{b}+v a_{b}^{*}$ with $|\mu|^{2}-|\nu|^{2}=1$, and characterize the squeezed state by the eigenvector $|\alpha ; \mu, \nu\rangle$ of the operator $b_{\mu, \nu}$ associated with the eigenvalue $\alpha \in \mathbb{C}$. In the special cases where $\alpha=0$, the vector $|0 ; \mu, \nu\rangle$ can be obtained by operating the action of the group with the generators $\frac{1}{2} Q^{2},-\frac{1}{2} P^{2}$ and $\frac{1}{2}(P Q+Q P)$ upon the boson vacuum vector $|0 ; 1,0\rangle$. The algebra with these generators satisfies the commutation relations of the algebra $\mathfrak{s u}(1,1)$. By operating $Q^{-1}\left(\right.$ or $\left.\left(a_{b}^{*}\right)^{-1}\right)$ upon the characteristic equation $b_{\mu, \nu}|0 ; \mu, \nu\rangle=0$ from the left, we have the characteristic equations

$$
\begin{align*}
& Q^{-1} P|0 ; \mu, v\rangle=\mathrm{i} \frac{\mu+v}{\mu-v}|0 ; \mu, v\rangle  \tag{1}\\
& -\left(a_{b}^{*}\right)^{-1} a_{b}|0 ; \mu, v\rangle=\frac{v}{\mu}|0 ; \mu, v\rangle \tag{2}
\end{align*}
$$

In section 5.2, we will derive these two equations again and reinterpret them from the viewpoint of the representation theory. In this framework, the operators $Q^{-1} P$ and $\left(a_{b}^{*}\right)^{-1} a_{b}$ have point spectra with continuous potency and they are constructed from an irreducible unitary representation of the algebra $\mathfrak{s u}(1,1)$ naturally. While the adjoints of these operators are not subnormal operators in the case of a one-particle system, the adjoints of these operators are subnormal operators in the cases of two- and multi-particle systems. Hence we can characterize a type of physically interpretable states by a tensor product, as the eigenvectors of the adjoints of subnormal operators in the cases of two- and multi-particle systems.

From a more general viewpoint, our investigation in this paper is regarded as a problem of the joint measurement between the self-adjoint part and the skew-adjoint part of a subnormal operator which do not always commute with each other. However, we should be careful about the difference between self-adjoint operators and symmetric operators in these discussions, because there are many delicate problems when unbounded operators are treated (section 6.1).

In this paper, the complex conjugate and the adjoint operator are denoted by *. And the closure is denoted by the overline.

## 2. Subnormal operator and POVM

In this section, we will summarize several well known lemmas and will modify them for the discussion in the following sections. Some of the well known lemmas will be extended for unbounded operators, and the proofs of the extended version will be given, as well. In this paper, only a densely defined linear operator will be discussed. In the following, $\mathcal{D}_{o}(X)$ denotes the domain of a linear operator $X$. A densely defined operator $X$ is called closed if the domain $\mathcal{D}_{o}(X)$ is complete with respect to the graph norm

$$
\|\phi\|_{\mathcal{D}_{o}(X)}:=\sqrt{\|\phi\|^{2}+\|X \phi\|^{2}}
$$

In operator theory, for two densely defined operators $X$ and $Y$, the product $X Y$ is defined as $\phi \mapsto X(Y(\phi))$ for any vector $\phi$ belonging to the domain $\mathcal{D}_{o}(X Y):=\left\{\phi \in \mathcal{D}_{o}(X) \mid X \phi \in\right.$
$\left.\mathcal{D}_{o}(Y)\right\}$. The notation $X \subset Y$ means that $\mathcal{D}_{o}(X) \subset \mathcal{D}_{o}(Y)$ and $X \phi=Y \phi, \phi \in \mathcal{D}_{o}(X)$. The notation $X=Y$, also means that $X \subset Y$ and $Y \subset X$. We will begin with reviewing the definition of the normal operator and that of the subnormal operator in the unbounded case.

Definition 1. A closed operator $T$ on $\mathcal{H}$ is called normal if it satisfies the condition $T^{*} T=$ $T T^{*}$.

Note that the operator $X^{*} X$ is defined on its domain $\mathcal{D}_{o}\left(X^{*} X\right):=\left\{\phi \in \mathcal{D}_{o}(X) \mid X \phi \in \mathcal{D}_{o}\left(X^{*}\right)\right\}$ and it is self-adjoint and non-negative.

Definition 2. A closed operator $S$ is called subnormal if there exists a Hilbert space $\mathcal{K}$ including $\mathcal{H}$ and a normal operator $T$ on $\mathcal{K}$ such that $S=T P_{\mathcal{H}}$, where $P_{\mathcal{H}}$ denotes the projection from $\mathcal{K}$ to $\mathcal{H}$ and we write the operator $S P_{\mathcal{H}}$ on the bigger space $\mathcal{K}$ by $S$. In the following, we call the pair $(\mathcal{K}, T)$ a normal extension of the subnormal operator $S$.

Remark 1. Many papers, for example, Stochel and Szafraniec [6, 7], Szafraniec [8], Ôta [9] and Lahti et al [10], adopt another definition of the subnormality, which substitutes $S \subset T P_{\mathcal{H}}$ for $S=T P_{\mathcal{H}}$. According to Ôta [9], there exists an example which is not subnormal in our definition, but subnormal in their definition.

For a spectral measure (i.e. a resolution of identity by projections) $E$ over $\mathbb{C}, \int_{\mathbb{C}} z E(\mathrm{~d} z)$ denotes the operator

$$
\phi \mapsto \lim _{n \rightarrow \infty}\left(\int_{|z|<n} z E(\mathrm{~d} z) \phi\right)
$$

with the domain

$$
\left\{\left.\phi \in \mathcal{H}\left|\int_{\mathbb{C}}\right| z\right|^{2}\langle\phi, E(\mathrm{~d} z) \phi\rangle<\infty\right\} .
$$

Concerning normal operators, the following lemma is well known. See theorem 13.33 in Rudin [11].

Lemma 3. For a normal operator $T$, there exists uniquely a spectral measure $E_{T}$ over $\mathbb{C}$ such that $T=\int_{\mathbb{C}} z E_{T}(\mathrm{~d} z)$.
Lemma 3 tells us that a normal operator corresponds to a spectral measure by one to one. Next, we will discuss measurements in a quantum system in order to investigate what corresponds to lemma 3 in the case of subnormal operators.

Let $\mathcal{H}$ be a Hilbert space representing a physical system of interest. Then, the state is denoted by a non-negative operator $\rho$ on $\mathcal{H}$ whose trace is 1 . It is called a density operator on $\mathcal{H}$, and the set of density operators on $\mathcal{H}$ is denoted by $\mathcal{S}(\mathcal{H})$. Let $P_{\rho}$ be the probability distribution given by a density $\rho$ and a measurement. Then, the probabilistic property of the measurement is described by the map $P: \rho \mapsto P_{\rho}$. We can naturally assume that the map $P$ satisfies the following condition from the formulation of quantum mechanics:

$$
\begin{equation*}
\lambda P_{\rho_{1}}+(1-\lambda) P_{\rho_{2}}=P_{\lambda \rho_{1}+(1-\lambda) \rho_{2}} \quad 0<\forall \lambda<1 \quad \forall \rho_{1}, \rho_{2} \in \mathcal{S}(\mathcal{H}) . \tag{3}
\end{equation*}
$$

Lemma 4. For a map P satisfying (3), there uniquely exists a positive operator-valued measure $M$ defined in the following which satisfies the condition

$$
P_{\rho}(B)=\operatorname{tr} M(B) \rho \quad \forall B \in \mathcal{F}(\Omega) \quad \forall \rho \in \mathcal{S}(\mathcal{H})
$$

This lemma was proved by Ozawa [12] in a more general framework. For an easy proof of a finite-dimensional case, see section 6 in chapter I of Holevo [13]. This lemma guarantees that we only have to discuss POVMs in order to describe probabilistic properties.

Definition 5. Let $M$ be a map from a $\sigma$-field $\mathcal{F}(\Omega)$ over $\Omega$ to the set $\mathcal{B}_{s a}^{+}(\mathcal{H})$ of bounded, selfadjoint and non-negative operators on $\mathcal{H}$. The map $M$ is called a positive operator-valued measure on $\mathcal{H}$ over $\Omega$ if it satisfies the following:

- $M(\emptyset)=0 \quad M(\Omega)=I \quad$ ( $I$ : indentity op.)
- $\sum_{i} M\left(B_{i}\right)=M\left(\cup_{i} B_{i}\right)$ for $B_{i} \cap B_{j}=\emptyset \quad(i \neq j)$.

A POVM $M$ is a spectral measure if and only if $M(B)$ is a projection for any $B$. The following lemma 6 is called Naǐmark's extension theorem. For a proof, see Nǎ̌mark [14], section 5 in chapter II in Holevo [13] or theorem 6.2.18 in Hiai and Yanagi [15]. It implies that the set of spectral measures is an important class in POVMs.
Lemma 6. Let $M$ be a POVM over a $\sigma$-field $\mathcal{F}(\Omega)$ on a Hilbert space $\mathcal{H}$. There exists a Hilbert space $\mathcal{K}$ including $\mathcal{H}$ and a spectral measure $E$ on the Hilbert space $\mathcal{K}$ such that

$$
M(B)=P_{\mathcal{H}} E(B) P_{\mathcal{H}} \quad \forall B \in \mathcal{F}(\Omega)
$$

where $P_{\mathcal{H}}$ denotes the projection from $\mathcal{K}$ to $\mathcal{H}$. We call such a pair $(\mathcal{K}, E)$ a Naǐmark extension of the POVM M.
In the following, we will treat only POVMs over the complex numbers $\mathbb{C}$ whose $\sigma$-field is a family of Borel sets.
Definition 7. A closed subspace $\mathcal{H}^{\prime}$ of $\mathcal{H}$ is said to reduce a spectral measure $E$ on $\mathcal{H}$, if the projection $P_{\mathcal{H}^{\prime}}$ to $\mathcal{H}^{\prime}$ commutes with $E(B)$ for any Borel set B. A Naimark extension ( $\mathcal{K}, E$ ) of a POVM M on $\mathcal{H}$ is called minimal if $\mathcal{K}$ has no non-trivial subspace which includes $\mathcal{H}$ and reduces the spectral measure $E$.
The following lemma guarantees the uniqueness of the minimal Nairmark extension. It is proved as a corollary of the principal theorem in section 6 of the appendix in Riesz and Sz-Nagy [16].
Lemma 8. Let $\left(\mathcal{K}_{1}, E_{1}\right)$ and $\left(\mathcal{K}_{2}, E_{2}\right)$ be Nǎ̌mark extensions of a POVM $M$ on $\mathcal{H}$. There exists a unitary map $V$ from $\mathcal{K}_{1}$ to $\mathcal{K}_{2}$ such that $U \phi=\phi$ for any $\phi \in \mathcal{H}$ and $V E_{1}(B) V^{*}=E_{2}(B)$ for any Borel B.
We will give the following definition with respect to the inequalities among linear operators which are not necessarily bounded.
Definition 9. For non-negative and self-adjoint operators $X, Y$ on $\mathcal{H}$, we denote $X \geqslant Y$ if they satisfy

$$
\langle\phi, X \phi\rangle \geqslant\langle\phi, Y \phi\rangle \quad \forall \phi \in \mathcal{D}_{f}(q(X)) \subset \mathcal{D}_{f}(q(Y))
$$

where $q(X)$ denotes the closed non-negative quadratic form defined by a non-negative selfadjoint operator $X$ and $\mathcal{D}_{f}(q)$ denotes the domain of a closed non-negative quadratic form $q$.
We introduce the operators $\mathrm{E}(M)$ and $\mathrm{V}(M)$ on $\mathcal{H}$ which formally represent $\int_{\mathbb{C}} z M(\mathrm{~d} z)$ and $\int_{\mathbb{C}}|z|^{2} M(\mathrm{~d} z)$, respectively. Later, by using lemma 10 , we will give more rigorous definition of $\mathrm{E}(M)$ and $\mathrm{V}(M)$. Then, for $\phi \in \mathcal{D}_{f}(\mathrm{q}(M)),\|\phi\|=1$ and a POVM $M$, the expectation of the measurement of the state by the POVM $M$ is $\langle\phi| \mathrm{E}(M)|\phi\rangle$ and the variance of it is $\langle\phi| \mathrm{V}(M)|\phi\rangle-|\langle\phi| \mathrm{E}(M)| \phi\rangle\left.\right|^{2}$. It is sufficient to evaluate the operator $\mathrm{V}(M)$, in order to evaluate the variance. However, when they are unbounded, we should be more careful with
respect to their domains. We define the closed non-negative quadratic form $\mathrm{q}(M)$ with the domain $\mathcal{D}_{f}(\mathrm{q}(M))$ by

$$
\begin{aligned}
& \mathrm{q}(M)(\phi, \phi):=\int_{\mathbb{C}}|z|^{2}\langle\phi, M(\mathrm{~d} z) \phi\rangle \quad \phi \in \mathcal{D}_{f}(\mathrm{q}(M)) \\
& \mathcal{D}_{f}(\mathrm{q}(M)):=\left\{\left.\phi \in \mathcal{H}\left|\int_{\mathbb{C}}\right| z\right|^{2}\langle\phi, M(\mathrm{~d} z) \phi\rangle<\infty\right\}
\end{aligned}
$$

We assume the condition that the set $\mathcal{D}_{f}(\mathrm{q}(M))$ is a dense subset of $\mathcal{H}$. Let $\mathrm{V}(M)$ be the self-adjoint operator defined by the closed non-negative quadratic form $\mathrm{q}(M)$. Next, we will define the operator $\tilde{\mathrm{E}}(M)$. Define $\mathrm{E}_{R}(M):=\int_{|z|<R} z M(\mathrm{~d} z)$. Then, the sequence $\left\{\mathrm{E}_{n}(M) \phi\right\}$ is a Cauchy sequence for any $\phi \in \mathcal{D}_{f}(\mathrm{q}(M))$, because we have

$$
\left\|\mathrm{E}_{n}(M) \phi-\mathrm{E}_{m}(M) \phi\right\|^{2}=\int_{n \leqslant|z|<m}|z|^{2}\langle\phi, M(\mathrm{~d} z) \phi\rangle
$$

for $n<m$. Therefore, we can define the vector $\tilde{\mathrm{E}}(M) \phi:=\lim _{n \rightarrow \infty} \mathrm{E}_{n}(M) \phi$. Thus, we can define the operator $\tilde{\mathrm{E}}(M)$ on the domain $\mathcal{D}_{f}(\mathrm{q}(M))$.
Lemma 10. The operator $\tilde{\mathrm{E}}(M)$ has a closed extension.
From this lemma, we can define the closed operator $\mathrm{E}(M)$ by the closure of the operator $\tilde{\mathrm{E}}(M)$.
Proof. Let $(E, \mathcal{K})$ and $P_{\mathcal{H}}$ be a Naǐmark extension of $M$ and the projection from $\mathcal{K}$ to $\mathcal{H}$. The operator $T:=\int z E(\mathrm{~d} z)$ is normal. From the definition of $T$, we have $\mathcal{D}_{o}(T)=\{\phi \in$ $\left.\left.\mathcal{K}\left|\int\right| z\right|^{2}\langle\phi, E(\mathrm{~d} z) \phi\rangle<\infty\right\}$. Then the domain $\mathcal{D}_{f}(\mathrm{q}(M))$ equals $\mathcal{D}_{o}(T) \cap \mathcal{H}$. Let $T=U|T|$ be a polar decomposition of $T$. Since the operator $T$ is normal, we have $U|T|=|T| U$. This equation implies that the domain of $|T|$ is invariant under the action of $U$.

In general, for a closed operator $X$ on $\mathcal{K}$ and closed subset $\mathcal{H}$ of $\mathcal{K}$, the operator $X P_{\mathcal{H}}$ with the domain $\mathcal{D}_{o}(X) \cap \mathcal{H}$ is closed if $\mathcal{D}_{o}(X) \cap \mathcal{H}$ is dense in $\mathcal{H}$. We can define the closed operator $T^{*} P_{\mathcal{H}}$ on its domain $\mathcal{D}_{o}\left(T^{*} P_{\mathcal{H}}\right):=\mathcal{D}_{o}\left(T^{*}\right) \cap \mathcal{H}=\mathcal{D}_{o}(T) \cap \mathcal{H}=\mathcal{D}_{f}(\mathrm{q}(M))$. Then, we have the relation $\mathcal{D}_{o}\left(\left(T^{*} P_{\mathcal{H}}\right)^{*}\right) \supset \mathcal{D}_{o}(T)$. Define the closed operator $\left(T^{*} P_{\mathcal{H}}\right)^{*} P_{\mathcal{H}}$ on its domain $\mathcal{D}_{o}\left(\left(T^{*} P_{\mathcal{H}}\right)^{*} P_{\mathcal{H}}\right):=\mathcal{D}_{o}\left(\left(T^{*} P_{\mathcal{H}}\right)^{*}\right) \cap \mathcal{H} \supset \mathcal{D}_{o}(T) \cap \mathcal{H}=\mathcal{D}_{f}(\mathrm{q}(M))$. Then, we obtain $\left(T^{*} P_{\mathcal{H}}\right)^{*} P_{\mathcal{H}} \supset \tilde{\mathrm{E}}(M)$. It follows that the operator $\tilde{\mathrm{E}}(M)$ has a closed extension.

Lemma 11. Let $X$ and $M$ be an operator on a Hilbert space $\mathcal{H}$ and a POVM on the Hilbert space $\mathcal{H}$, respectively. If $X \supset \mathrm{E}(M)$, then we have $\mathrm{V}(M) \geqslant X^{*} X$.

Proof. For a vector $\phi \in \mathcal{D}_{f}(\mathrm{q}(M))$, we have

$$
\mathrm{q}(M)(\phi, \phi)-\langle\phi| X^{*} X|\phi\rangle=\int_{\mathbb{C}}\langle\phi|\left(z^{*}-X^{*}\right) M(\mathrm{~d} z)(z-X)|\phi\rangle \geqslant 0
$$

Since the relation $\mathcal{D}_{f}(\mathrm{q}(M)) \subset \mathcal{D}_{o}(\mathrm{E}(M)) \subset \mathcal{D}_{o}(X)$ holds, we obtain lemma 11 .
The bounded version of this lemma is proved by Helstrom [17] from the viewpoint of quantum estimation theory. Its bounded version also follows from Kadison's inequality [18].

Lemma 12. Let $S$ be an operator defined on the dense subset $\mathcal{D}_{o}(S)$ of $\mathcal{H}$. The operator $S$ is subnormal if and only if there exists a POVM M satisfying the conditions

$$
\begin{align*}
& S=\mathrm{E}(M)  \tag{4}\\
& S^{*} S=\mathrm{V}(M) \tag{5}
\end{align*}
$$

Proof. Let $(\mathcal{K}, T)$ and $P_{\mathcal{H}}$ be a normal extension of the operator $S$ and the projection from $\mathcal{K}$ to $\mathcal{H}$, respectively. By defining a POVM $M$ by $M(B):=P_{\mathcal{H}} E_{T}(B) P_{\mathcal{H}}$, equation (4) is trivial. Since the equation $\mathrm{V}(M)=\left(T P_{\mathcal{H}}\right)^{*}\left(T P_{\mathcal{H}}\right)=S^{*} S$ holds, we have equation (5). Assume equations (4) and (5). From Naǐmark's extension theorem (lemma 6) there exists a Naǐmark extension $(\mathcal{K}, E)$ of the POVM $M$. Define a normal operator $T:=\int_{\mathbb{C}} z E(\mathrm{~d} z)$. Then we have $\mathrm{V}(M)=\left(T P_{\mathcal{H}}\right)^{*}\left(T P_{\mathcal{H}}\right), \mathrm{E}(M)=P_{\mathcal{H}}\left(T P_{\mathcal{H}}\right)$. From equations (4), (5) and lemma 14, we can prove that $S$ is subnormal.

The bounded version of this lemma was proved by Bram [19].
Definition 13. A POVM $M$ is called a POVM of a subnormal operator $S$ if $M$ satisfies the preceding conditions (4) and (5).
We will prove lemma 14 applied in the proof of lemma 12.
Lemma 14. Let $S, \mathcal{K}$ and $P_{\mathcal{H}}$ be an operator on a Hilbert space $\mathcal{H}$, a Hilbert space including the Hilbert space $\mathcal{H}$ and the projection from $\mathcal{K}$ to $\mathcal{H}$, respectively. For an operator $T$ on $\mathcal{K}$, the following are equivalent:
(A) $S=T P_{\mathcal{H}}$.
(B) $S^{*} S=\left(T P_{\mathcal{H}}\right)^{*}\left(T P_{\mathcal{H}}\right) \quad S=P_{\mathcal{H}}\left(T P_{\mathcal{H}}\right)$.

Proof. It is easy to derive condition (B) from condition (A). Assume condition (B). We have $\left(T P_{\mathcal{H}}\right)^{*}\left(T P_{\mathcal{H}}\right)=\left(P_{\mathcal{H}}\left(T P_{\mathcal{H}}\right)\right)^{*}\left(P_{\mathcal{H}}\left(T P_{\mathcal{H}}\right)\right)+\left(\left(I-P_{\mathcal{H}}\right) T P_{\mathcal{H}}\right)^{*}\left(\left(I-P_{\mathcal{H}}\right)\left(T P_{\mathcal{H}}\right)\right)$ and $\left(P_{\mathcal{H}}\left(T P_{\mathcal{H}}\right)\right)^{*}\left(P_{\mathcal{H}}\left(T P_{\mathcal{H}}\right)\right)=S^{*} S=\left(T P_{\mathcal{H}}\right)^{*}\left(T P_{\mathcal{H}}\right)$. Therefore, we obtain $\left(I-P_{\mathcal{H}}\right)\left(T P_{\mathcal{H}}\right)=0$. Thus, we obtain condition (A).

Definition 15. A closed subspace $\mathcal{H}^{\prime}$ of $\mathcal{H}$ is said to reduce a normal operator $T$ on $\mathcal{H}$, if the closed subspace $\mathcal{H}^{\prime}$ of $\mathcal{H}$ reduces its spectral measure $E_{T}$. This condition is equivalent to the condition that the projection $P_{\mathcal{H}^{\prime}}$ commutes with the operators $U, U^{*}$ and $\mathrm{e}^{\mathrm{i} t|T|^{2}}$ for any real number $t$, where $T=U|T|$ is the polar decomposition of $T$ with unitary $U$. A normal extension $(T, \mathcal{K})$ of a subnormal operator $S$ on $\mathcal{H}$ is called minimal if $\mathcal{K}$ has no non-trivial subspace which includes $\mathcal{H}$ and reduces the normal operator $T$.
The POVM $M(B):=P_{\mathcal{H}} E_{T}(B) P_{\mathcal{H}}$ can be defined for a normal extension $(T, \mathcal{K})$ of a subnormal operator $S$, and it is a POVM of $S$. Conversely, from lemma 8, if the normal extension $(T, \mathcal{K})$ is minimal, the spectral measure $E_{T}$ is unitarily equivalent to the minimal Naǐmark extension of $M$. Therefore, there exists a one-to-one correspondence between minimal normal extensions of a subnormal operator $S$ and its POVMs.
Lemma 16. A normal extension $(T, \mathcal{K})$ of a subnormal operator $S$ on $\mathcal{H}$ is minimal if and only if $\mathcal{K}=\overline{\mathcal{L}}$, where the subspaces $\mathcal{L}$ and $\mathcal{C}$ of $\mathcal{K}$ is defined as

$$
\begin{aligned}
\mathcal{L} & :=\left\{\sum_{k=1}^{n}\left(U^{*}\right)^{k} \psi_{k} \mid \psi_{k} \in \overline{\mathcal{C}}, n \in \mathbb{N}\right\} \\
\mathcal{C} & :=\left\{\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} t_{k}|T|^{2}} \psi_{k} \mid \psi_{k} \in \mathcal{H}, t_{k} \in \mathbb{R}, n \in \mathbb{N}\right\}
\end{aligned}
$$

where $T=U|T|$ is the polar decomposition of $T$ with unitary $U$.
Proof. Assume that a closed subspace $\mathcal{K}^{\prime}$ of $\mathcal{K}$ including $\mathcal{H}$ reduces the normal operator $T$. Then, for any $h \in \mathcal{H}$, any integer $m$ and any real number $t$, we have $\mathrm{e}^{\mathrm{i} t|T|^{2}} h \in \mathcal{K}^{\prime}$. Since the
closed subspace $\mathcal{K}^{\prime}$ includes $\mathcal{C}$, the closed subspace $\mathcal{K}^{\prime}$ includes $\overline{\mathcal{C}}$. Similarly, we can show that the closed subspace $\mathcal{K}^{\prime}$ includes $\overline{\mathcal{L}}$ from this fact.

Next, we will prove that the closed subspace $\overline{\mathcal{C}}$ is invariant for $U$. It is sufficient to show that $U \phi \in \overline{\mathcal{C}}$ for any $\phi \in \mathcal{H}$. From the definition of $\mathcal{C}$, the closure $\overline{\mathcal{C}}$ reduces the operator $|T|^{2}$. Also, it reduces the operators $|T|$ and $|T|^{-1}$. Since $\mathcal{D}_{o}\left(|T|^{-1}\right) \subset \operatorname{Im} T, U \phi=|T|^{-1} S \phi \in \overline{\mathcal{C}}$ holds for any $\phi \in \mathcal{D}_{o}(S)$. We have $U \mathcal{H} \subset \overline{\mathcal{C}}$ because $U$ is bounded and $\mathcal{D}_{o}(S)$ is dense in $\mathcal{H}$. Thus, $U \mathrm{e}^{\mathrm{i} t|T|^{2}} \phi=\mathrm{e}^{\mathrm{i} t|T|^{2}} U \phi \in \underline{\mathcal{C}}$ for any $\phi \in \mathcal{H}$. It follows that $\mathcal{C}$ is invariant for $U$.

Therefore, we have the relations $U \overline{\mathcal{L}} \subset \overline{\mathcal{L}}, U^{*} \overline{\mathcal{L}} \subset \overline{\mathcal{L}}$ and $\mathrm{e}^{\mathrm{i} t|T|^{2}} \overline{\mathcal{L}} \subset \overline{\mathcal{L}}$ for any real number $t$. These imply that $\left[P_{\overline{\mathcal{L}}}, U\right]=0,\left[P_{\overline{\mathcal{L}}}, U^{*}\right]=0$ and $\left[P_{\overline{\mathcal{L}}}, \mathrm{e}^{\mathrm{i} t|T|^{2}}\right]=0$. It follows that the closed subspace $\overline{\mathcal{L}}$ reduces the normal operator $T$. The lemma is now immediately obvious.

Lemma 17. Let $(T, \mathcal{K})$ be a minimal normal extension of a subnormal operator $S$ on $\mathcal{H}$. A Hilbert space $\mathcal{K}^{\prime}$ including $\mathcal{H}$ and a normal operator $T^{\prime}$ satisfy the condition $S \subset T^{\prime} P_{\mathcal{H}}$. The following three conditions are equivalent.
(A) $\left\langle\phi_{1}, \mathrm{e}^{\mathrm{i} t|T|^{2}} \phi_{2}\right\rangle=\left\langle\phi_{1}, \mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}} \phi_{2}\right\rangle$ holds for any $\phi_{1}, \phi_{2} \in \mathcal{H}$.
(B) $\left\langle\phi_{1}, \mathrm{e}^{\mathrm{i} t|T|} \phi_{2}\right\rangle=\left\langle\phi_{1}, \mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|} \phi_{2}\right\rangle$ holds for any $\phi_{1}, \phi_{2} \in \mathcal{H}$.
(C) There exists an isometric map $V$ from $\mathcal{K}$ to $\mathcal{K}^{\prime}$ such that $V \phi=\phi$ for any $\phi \in \mathcal{H}$ and $V T V^{*}=T^{\prime} P_{\operatorname{Im} V}$.
The condition (C) implies that $T^{\prime} P_{\mathcal{H}}=S$, i.e. the pair $\left(T^{\prime}, \mathcal{K}^{\prime}\right)$ is a normal extension of $S$.

Proof. It is easy to show that the conditions (A) and (B) follow from condition (C). First, we prove that condition (A) implies condition (B). Define the subspace $\mathcal{C}^{\prime}$ of $\mathcal{K}^{\prime}$ by

$$
\mathcal{C}^{\prime}:=\left\{\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} t_{k}\left|T^{\prime}\right|^{2}} \psi_{k} \mid \psi_{k} \in \mathcal{H}, t_{k} \in \mathbb{R}, n \in \mathbb{N}\right\}
$$

Similarly to the proof of lemma 16 , we can prove that the closure $\overline{\mathcal{C}}$ reduces $|T|^{2}$ and the closure $\overline{\mathcal{C}^{\prime}}$ reduces $\left|T^{\prime}\right|^{2}$. Then, the closures $\overline{\mathcal{C}}$ and $\overline{\mathcal{C}^{\prime}}$ reduce the operators $|T|$ and $\left|T^{\prime}\right|$, respectively. From condition (A), $\left\langle\mathrm{e}^{\mathrm{i} t_{1}|T|^{2}} \phi_{1}, \mathrm{e}^{\mathrm{i} t_{2}|T|^{2}} \phi_{2}\right\rangle=\left\langle\mathrm{e}^{\mathrm{i} t_{1}|T|^{2}} \phi_{1}, \mathrm{e}^{\mathrm{i} t_{2}\left|T^{\prime}\right|^{2}} \phi_{2}\right\rangle$ holds for any $\phi_{1}, \phi_{2} \in \mathcal{H}$ and any real numbers $t_{1}, t_{2}$. Therefore, we can define the unitary map $V_{\mathcal{C}}$ from $\overline{\mathcal{C}}$ to $\overline{\mathcal{C}}^{\prime}$ by

$$
V_{\mathcal{C}}\left(\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} \mathrm{i}_{k}|T|^{2}} \phi_{k}\right)=\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} k_{k}\left|T^{\prime}\right|^{2}} \phi_{k}
$$

Thus, we have $V_{\mathcal{C}}|T|^{2} V_{\mathcal{C}}^{*}=\left|T^{\prime}\right|^{2}$ on $\overline{\mathcal{C}^{\prime}}$. It implies that $V_{\mathcal{C}}|T| V_{\mathcal{C}}^{*}=\left|T^{\prime}\right|$ on $\overline{\mathcal{C}^{\prime}}$ because the closures $\overline{\mathcal{C}}$ and $\overline{\mathcal{C}^{\prime}}$ reduce the operators $|T|$ and $\left|T^{\prime}\right|$, respectively. Since $V_{\mathcal{C}} \phi=\phi$ for any $\phi \in \mathcal{H}$ and $V_{\mathcal{C}} \mathrm{e}^{\mathrm{i} t|T|} V_{\mathcal{C}}^{*}=\mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|}$ for any $t \in \mathbb{R}$, condition (B) holds. Similarly, we can prove that condition (B) implies condition (A).

Next, we prove that condition (A) implies condition (C). From the above discussion, we can define the inverses $|T|^{-1}$ and $\left|T^{\prime}\right|^{-1}$ on $\operatorname{Im}|T| \cap \overline{\mathcal{C}}$ and $\operatorname{Im}\left|T^{\prime}\right| \cap \overline{\mathcal{C}^{\prime}}$, respectively. Then we have $V_{\mathcal{C}}|T|^{-1} V_{\mathcal{C}}^{*}=\left|T^{\prime}\right|^{-1}$ on $\operatorname{Im}\left|T^{\prime}\right| \cap \overline{\mathcal{C}^{\prime}}$.

Let $T=U|T|$ and $T^{\prime}=U^{\prime}\left|T^{\prime}\right|$ be the polar decompositions of $T$ and $T^{\prime}$ satisfying that $U$ and $U^{\prime}$ are unitary, respectively. The image $\operatorname{Im}|T|$ is invariant under the unitary transformation $U$, and the image $\operatorname{Im}\left|T^{\prime}\right|$ is invariant under $U^{\prime}$. Then, we have $\operatorname{Im} S \subset \operatorname{Im}|T| \cap \mathcal{H}$. Similarly, we have $\operatorname{Im} S \subset \operatorname{Im}\left|T^{\prime}\right| \cap \mathcal{H}$. Thus, for any $\phi \in \operatorname{Im} S$, we have $V_{\mathcal{C}}|T|^{-1} \phi=\left|T^{\prime}\right|^{-1} \phi$. From
the proof of lemma 16, the closed subspaces $\overline{\mathcal{C}}$ and $\overline{\mathcal{C}}^{\prime}$ are invariant for $U$ and $U^{\prime}$, respectively. For any $\phi_{1}, \phi_{2} \in \mathcal{D}_{o}(S)$, we have

$$
\begin{aligned}
\left\langle\mathrm{e}^{\mathrm{i} t|T|^{2}} \phi_{1}, U \phi_{2}\right\rangle & \left.=\left\langle V_{\mathcal{C}} \mathrm{e}^{\mathrm{i} t|T|^{2}} \phi_{1}, V_{\mathcal{C}} U \phi_{2}\right\rangle=\left.\left\langle V_{\mathcal{C}} \mathrm{e}^{\mathrm{i} t|T|^{2}} \phi_{1}, V_{\mathcal{C}}\right| T\right|^{-1} S \phi_{2}\right\rangle \\
& \left.\left.=\left.\left\langle V_{\mathcal{C}} \mathrm{e}^{\mathrm{i} t|T|^{2}} V_{\mathcal{C}}^{*} V_{\mathcal{C}} \phi_{1}, V_{\mathcal{C}}\right| T\right|^{-1} V_{\mathcal{C}}^{*} V_{\mathcal{C}} S \phi_{2}\right\rangle=\left.\left\langle\mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}} V_{\mathcal{C}} \phi_{1},\right| T^{\prime}\right|^{-1} V_{\mathcal{C}} S \phi_{2}\right\rangle \\
& \left.=\left.\left\langle\mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}} \phi_{1},\right| T^{\prime}\right|^{-1} S \phi_{2}\right\rangle=\left\langle\mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}} \phi_{1}, U^{\prime} \phi_{2}\right\rangle .
\end{aligned}
$$

Since $\mathrm{e}^{\mathrm{i} t|T|^{2}}, \mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}}, U$ and $U^{\prime}$ are bounded,

$$
\left\langle\mathrm{e}^{\mathrm{i} t|T|^{2}} \phi_{1}, U \phi_{2}\right\rangle=\left\langle\mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}} \phi_{1}, U^{\prime} \phi_{2}\right\rangle
$$

holds, for any $\phi_{1}, \phi_{2} \in \mathcal{H}$. Also, we can prove

$$
\left\langle\left(\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} t_{k}|T|^{2}} \psi_{k}\right), U\left(\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} t_{k}^{\prime}|T|^{2}} \psi_{k}^{\prime}\right)\right\rangle=\left\langle\left(\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} t_{k}\left|T^{\prime}\right|^{2}} \psi_{k}\right), U^{\prime}\left(\sum_{k=1}^{n} \mathrm{e}^{\mathrm{i} t_{k}^{\prime}\left|T^{\prime}\right|^{2}} \psi_{k}^{\prime}\right)\right\rangle
$$

for arbitrary $\psi_{k}, \psi_{k}^{\prime} \in \mathcal{H}, t_{k}, t_{k}^{\prime} \in \mathbb{R}$. Therefore,

$$
\left\langle\phi_{1}, U \phi_{2}\right\rangle=\left\langle V_{\mathcal{C}} \phi_{1}, U^{\prime} V_{\mathcal{C}} \phi_{2}\right\rangle
$$

holds for any $\phi_{1}, \phi_{2} \in \overline{\mathcal{C}}$. Since the closed subspace $\overline{\mathcal{C}}$ is invariant for $U$ and the operator $U$ is bounded,

$$
\begin{equation*}
\left\langle\phi_{1}, U^{n} \phi_{2}\right\rangle=\left\langle V_{\mathcal{C}} \phi_{1}, U^{\prime n} V_{\mathcal{C}} \phi_{2}\right\rangle \tag{6}
\end{equation*}
$$

holds for any $\phi_{1}, \phi_{2} \in \overline{\mathcal{C}}$ and any $n \in \mathbb{N}$.
We can define the isometric map $V$ from $\mathcal{K}=\overline{\mathcal{L}}$ to $\mathcal{K}^{\prime}$ by

$$
V\left(\sum_{k=1}^{n}\left(U^{*}\right)^{k} \psi_{k}\right)=\left(\sum_{k=1}^{n}\left(U^{* *}\right)^{k} \psi_{k}\right)
$$

where $\psi_{k}$ is an arbitrary element of $\overline{\mathcal{C}}$. It can be confirmed that this definition is well defined from (6). Now, we can easily check condition (C).

Definition 18. A vector $\phi \in \mathcal{D}^{\infty}(X):=\cap_{n=0}^{\infty} \mathcal{D}_{o}\left(X^{n}\right)$ is called an analytic vector of $X$ if

$$
\begin{equation*}
\sum_{i=0}^{\infty} \frac{t^{i}}{i!}\left\|X^{i} \phi\right\|<\infty \tag{7}
\end{equation*}
$$

for any $t \in \mathbb{R}$. The set of all analytic vectors of $S$ is written as $\mathcal{A}(S)$.
Lemma 19. Assume that the set $\mathcal{A}(S)$ is dense in $\mathcal{H}$ for a subnormal operator $S$ on $\mathcal{H}$. Let $(T, K)$ be a normal extension of $S$. If a Hilbert space $\mathcal{K}^{\prime}$ including $\mathcal{H}$ and a normal operator $T^{\prime}$ satisfy the condition $T^{\prime} P_{\mathcal{H}} \supset S$, there exists an isometric map $V$ from $\mathcal{K}$ to $\mathcal{K}^{\prime}$ such that $V \phi=\phi$ for any $\phi \in \mathcal{H}$ and $V T V^{*}=T^{\prime} P_{\operatorname{Im} V}$. This implies that $T^{\prime} P_{\mathcal{H}}=S$, i.e. the pair $\left(T^{\prime}, \mathcal{K}^{\prime}\right)$ is a normal extension of $S$. Therefore, this assumption guarantees the uniqueness of the minimal normal extension.

This lemma shows that under the assumption, the pair $(T, \mathcal{H})$ is a normal extension of $S$ if a normal operator $T$ on $\mathcal{K}$ including $\mathcal{H}$ satisfies $T P_{\mathcal{H}} \supset S$. For a simple proof in the bounded case, see section 2 in chapter II of Conway [20]. Szafraniec [8] shows the uniqueness of the minimal normal extension under another assumption that any vector $\phi \in \mathcal{D}(S)$ satisfies (7) for some a real number $t>0$. Stochel and Szafraniec [7] discuss different sufficient conditions for the uniqueness of the minimal normal extension.

Proof. It is sufficient to show that condition (A) in lemma 17 holds. Schwarz's inequality guarantees

$$
\sum_{k=0}^{\infty} \frac{t^{k}}{k!}\left\langle S^{k} \phi_{1}, S^{k} \phi_{2}\right\rangle<\sqrt{\sum_{k=0}^{\infty} \frac{t^{k}}{k!}\left\|S^{k} \phi_{1}\right\|^{2}} \sqrt{\sum_{k=0}^{\infty} \frac{t^{k}}{k!}\left\|S^{k} \phi_{2}\right\|^{2}}<\infty .
$$

for any $\phi_{1}, \phi_{2} \in \mathcal{A}(S)$ and any real number $t$. From Fubini's theorem,

$$
\begin{align*}
\left\langle\phi_{1}, \mathrm{e}^{\mathrm{i} t|T|^{2}} \phi_{2}\right\rangle & =\left\langle\phi_{1}, \sum_{k=0}^{\infty} \frac{\left(\mathrm{i} t|T|^{2}\right)^{k}}{k!} \phi_{2}\right\rangle=\left\langle\phi_{1}, \sum_{k=0}^{\infty} \frac{(\mathrm{i} t)^{k}\left(T^{*}\right)^{k} T^{k}}{k!} \phi_{2}\right\rangle \\
& =\sum_{k=0}^{\infty} \frac{(\mathrm{i} t)^{k}}{k!}\left\langle S^{k} \phi_{1}, S^{k} \phi_{2}\right\rangle=\left\langle\phi_{1}, \mathrm{e}^{\left.\mathrm{i}| | T^{\prime}\right|^{2}} \phi_{2}\right\rangle . \tag{8}
\end{align*}
$$

From (8) and the fact that the operators $\mathrm{e}^{\mathrm{i} t|T|^{2}}$ and $\mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}}$ are bounded and $\mathcal{A}(S)$ is dense in $\mathcal{H}$, we have the equation $P_{\mathcal{H}} \mathrm{e}^{\mathrm{i} t|T|^{2}} P_{\mathcal{H}}=P_{\mathcal{H}} \mathrm{e}^{\mathrm{i} t\left|T^{\prime}\right|^{2}} P_{\mathcal{H}}$. Therefore, condition (A) in lemma 19 holds.

From the one-to-one correspondence between POVMs of a subnormal operator $S$ and its minimal normal extensions, we have the following corollary.

Corollary 20. For any subnormal operator $S$ satisfying the assumption of lemma 19, there uniquely exists the POVM M satisfying equations (4) and (5).
Subnormal operators have the following properties:
Lemma 21. Let $S$ be a subnormal operator on $\mathcal{H}$. Then

$$
\begin{equation*}
S^{*} S \geqslant S S^{*} \tag{9}
\end{equation*}
$$

Proof. We have $\left(P_{\mathcal{H}} T\right)\left(P_{\mathcal{H}} T\right)^{*} \geqslant\left(P_{\mathcal{H}}\left(P_{\mathcal{H}} T\right)\right)\left(P_{\mathcal{H}}\left(P_{\mathcal{H}} T\right)\right)^{*}=S S^{*}$. From the normality of $T$, we have $S^{*} S=\left(T P_{\mathcal{H}}\right)^{*}\left(T P_{\mathcal{H}}\right)=\left(T^{*} P_{\mathcal{H}}\right)^{*}\left(T^{*} P_{\mathcal{H}}\right)$. Since $T^{*} P_{\mathcal{H}} \subset\left(P_{\mathcal{H}} T\right)^{*}$ (see theorem 13.2 in Rudin [11]), the inequality $\left(T^{*} P_{\mathcal{H}}\right)^{*}\left(T^{*} P_{\mathcal{H}}\right) \geqslant\left(P_{\mathcal{H}} T\right)\left(P_{\mathcal{H}} T\right)^{*}$ holds. Therefore, $S^{*} S \geqslant\left(P_{\mathcal{H}} T\right)\left(P_{\mathcal{H}} T\right)^{*} \geqslant S S^{*}$.

Operators satisfying (9) are called hyponormal operators $\dagger$ and the class of these operators is important in the operator theory. The following lemma shows a relation between the POVM of a subnormal operator and an over-complete eigenvector system.
Lemma 22. Let $J$ and $K$ be an operator on $\mathcal{H}$ and a subset of complex numbers $\mathbb{C}$, respectively. Assume that there exists a vector $|z\rangle \in \mathcal{D}_{o}(J)$ satisfying $J|z\rangle=z|z\rangle$ for any complex number $z \in K$, and there exists a measure $\mu$ on $K$ satisfying $\int_{K}\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z)=I$. Then, $J^{*}$ is subnormal and the POVM $\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z)$ is the POVM of the subnormal operator $J^{*}$.

Proof. From the assumptions, we have

$$
J^{*}=\int_{K}\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z) J^{*}=\int_{K} z\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z) .
$$

Note that $\left\langle z^{*}\right| J^{*}=z\left\langle z^{*}\right|$. Thus, The POVM $M(\mathrm{~d} z):=\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z)$ satisfies condition (4). Therefore, we obtain

$$
J J^{*}=\int_{K} J\left|z^{*}\right\rangle\left\langle z^{*}\right| J^{*} \mu(\mathrm{~d} z)=\int_{K}|z|^{2}\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z)=\mathrm{V}(M) .
$$

[^0]Then the POVM $M$ satisfies condition (5) and the operator $J^{*}$ is subnormal, as shown from lemma 12. We can confirm that the POVM $\left|z^{*}\right\rangle\left\langle z^{*}\right| \mu(\mathrm{d} z)$ is the POVM of the subnormal operator $J^{*}$.

In the following of this section, we treat a relation between a subnormal operator and its spectrum.

Lemma 23. Let $S$ and $\phi$ be a subnormal operator and an eigenvector of $S$, respectively. Then, a vector $\phi$ is an eigenvector of the adjoint $S^{*}$ operator of $S$.

Proof. Let $(\mathcal{K}, T)$ and $P_{\mathcal{H}}$ be a normal extension of $S$ and the projection from $\mathcal{K}$ to $\mathcal{H}$, respectively. Assume that $\phi \in \mathcal{D}_{o}(S)$ is an eigenvector of $S$ associated with an eigenvalue $c$ such that $\|\phi\|=1$. Since the equation $T \phi=c \phi$ holds, we have $T^{*} \phi=c^{*} \phi$. Thus $S^{*}=P_{\mathcal{H}} T^{*}$. Therefore, we obtain $S^{*} \phi=c^{*} \phi$. Now, we obtain the lemma.

Definition 24. A subnormal operator $S$ is called pure subnormal if it satisfies the following condition; if a subspace $\mathcal{I}$ of $\mathcal{H}$ satisfies that $S P_{\mathcal{I}}$ is subnormal, then the subspace $\mathcal{I}$ is $\{0\}$ or $\mathcal{H}$.

Lemma 25. Any pure subnormal operator $S$ has no point spectrum.
Proof. Let $(\mathcal{K}, T), P_{\mathcal{H}}$ and $\phi$ be defined in the proof of lemma 23. Since we have $S^{*} \phi=c^{*} \phi$, the operator $|\phi\rangle\langle\phi|$ commutes with the pure subnormal operator $S$. The fact contradicts the definition of pure subnormal operators.

According to Conway [20], it is sufficient to assume the purity and hyponormality in lemma 25.

## 3. Compound-system-type normal extension

Now, as an example of a subnormal operator and its normal extension, we will treat the boson creation operator $a_{b}^{*}$ and the heterodyne measurement in quantum optics. The pair $\left(L^{2}(\mathbb{R}) \otimes L^{2}(\mathbb{R}), a_{b}^{*} \otimes I+I \otimes a_{b}\right)$ is a normal extension of the subnormal operator $a_{b}^{*}$ under the isometric embedding $L^{2}(\mathbb{R}) \rightarrow L^{2}(\mathbb{R}) \otimes L^{2}(\mathbb{R})$ defined by $\psi \mapsto \psi \otimes|0 ; 1,0\rangle$, where $|0 ; 1,0\rangle$ denotes the boson vacuum vector. Here, $a_{b}^{*} \otimes I+I \otimes a_{b}$ is a normal operator, and we have $\left(a_{b}^{*} \otimes I+I \otimes a_{b}\right) \phi \otimes|0 ; 1,0\rangle=\left(a_{b}^{*} \phi\right) \otimes|0 ; 1,0\rangle$ for any $\phi \in L^{2}(\mathbb{R})$. By substituting $a_{b}$ for $J$ in lemma 22 and by letting $|\alpha ; 1,0\rangle$ be the boson coherent state, we can confirm that $|\alpha ; 1,0\rangle\langle\alpha ; 1,0| \mathrm{d}^{2} \alpha$ is the POVM of the subnormal operator $a_{b}^{*}$. The set of rapidly decreasing $C^{\infty}$ functions is dense in $L^{2}(\mathbb{R})$ and any rapidly decreasing $C^{\infty}$ function is analytic of $a_{b}^{*}$. Therefore, $a_{b}^{*}$ 's POVM is uniquely determined.

The heterodyne measurement is implemented by the measurement of $a_{b}^{*} \otimes I+I \otimes a_{b}$ (i.e. the simultaneous measurement between $Q \otimes I+I \otimes Q$ and $P \otimes I-I \otimes P$ which commute with each other) under the circumstance where the state of the basic system is $|\phi\rangle\langle\phi|$ and the state of the ancillary system is controlled to be the vacuum states $|0 ; 1,0\rangle\langle 0 ; 1,0|$. In detail, see section 6 in chapter III in Holevo [13] or section 6 in chapter V in Helstrom [17]. We will generalize normal extensions of a similar type to this, by the name of compound-system-type normal extensions, as follows:

Definition 26. Let $S$ be a subnormal operator defined on a dense linear subspace $\mathcal{D}_{o}(S)$ of $\mathcal{H}$ and let $\mathcal{H}^{\prime}, T$ and $\psi$ be a Hilbert space, a normal operator defined on a dense subspace $\mathcal{D}_{o}(T)$ of the Hilbert space $\mathcal{H} \otimes \mathcal{H}^{\prime}$ and an element of $\mathcal{H}^{\prime}$ whose norm is unity, respectively. We call
the triple $\left(\mathcal{H}^{\prime}, T, \psi\right)$ a compound-system-type normal extension of the subnormal operator $S$ if it satisfies the condition
$\mathcal{D}_{o}(S) \otimes \psi \subset \mathcal{D}_{o}(T) \quad(S \phi) \otimes \psi=T(\phi \otimes \psi) \quad$ for any $\quad \phi \in \mathcal{D}_{o}(S)$.
Thus the definition of the compound-system-type normal extension describes not only the probability distribution but also a framework of the concrete implementation process, while the definition of the normal extension given in section 2 describes only the probability distribution. Therefore, a compound-system-type normal extension contains more information than the corresponding POVM.

In the following section, we discuss compound-system-type normal extensions of isometric operators and symmetric operators, where we let $\{|\uparrow\rangle,|\downarrow\rangle\}$ be a CONS of $\mathbb{C}^{2}$.
Lemma 27. An isometric operator $U$ defined on $\mathcal{H}$ is subnormal. Define the operator $T:=U \otimes|\uparrow\rangle\langle\uparrow|+U^{*} \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\uparrow\rangle\langle\downarrow|$, where $\operatorname{Im} U^{\perp}$ denotes the orthogonal complementary space of $\operatorname{Im} U$. Then, the operator $T$ is unitary on $\mathcal{H} \otimes \mathbb{C}^{2}$ and the triple $\left(\mathbb{C}^{2}, T,|\uparrow\rangle\right)$ is a compound-system-type normal extension of $U$.

Proof. From the definition, we have

$$
\begin{aligned}
& T^{*} T=\left(U^{*} \otimes|\uparrow\rangle\langle\uparrow|+U \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\downarrow\rangle\langle\uparrow|\right) \\
&\left(U \otimes|\uparrow\rangle\langle\uparrow|+U^{*} \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\uparrow\rangle\langle\downarrow|\right) \\
&= I_{\mathcal{H}} \otimes|\uparrow\rangle\langle\uparrow|+P_{\operatorname{Im} U} \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\downarrow\rangle\langle\downarrow|=I_{\mathcal{H}} \otimes I_{\mathbb{C}^{2}} \\
& T T^{*}=\left(U \otimes|\uparrow\rangle\langle\uparrow|+U^{*} \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\uparrow\rangle\langle\downarrow|\right) \\
&\left(U^{*} \otimes|\uparrow\rangle\langle\uparrow|+U \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\downarrow\rangle\langle\uparrow|\right) \\
&= P_{\operatorname{Im} U} \otimes|\uparrow\rangle\langle\uparrow|+I_{\mathcal{H}} \otimes|\downarrow\rangle\langle\downarrow|+P_{\operatorname{Im} U^{\perp}} \otimes|\uparrow\rangle\langle\uparrow|=I_{\mathcal{H}} \otimes I_{\mathbb{C}^{2}} .
\end{aligned}
$$

Then, the operator $T$ is unitary. Moreover, we have $T(\phi \otimes|\uparrow\rangle)=(U \phi) \otimes|\uparrow\rangle$. Therefore, the triple $\left(\mathbb{C}^{2}, T,|\uparrow\rangle\right)$ is a compound-system-type normal extension of $U$.

A closed symmetric operator $X$ is called maximally symmetric, if there exists no symmetric operator $Y$ such that $X \varsubsetneqq Y$.
Lemma 28. A closed symmetric operator $X$ is subnormal on $\mathcal{H}$. Define the operator $T:=X^{*} \otimes|-\rangle\langle+|+X \otimes|+\rangle\langle-|$ on the domain $\mathcal{D}_{o}(T):=\mathcal{D}_{o}\left(X^{*}\right) \otimes|+\rangle \oplus \mathcal{D}_{o}(X) \otimes|-\rangle$ with
 operator and the triple $\left(\mathbb{C}^{2}, T,|\uparrow\rangle\right)$ is a compound-system-type normal extension of $X$.
The classification of (second) self-adjoint extensions of symmetric operators is given in section 5 in Naǐmark [21].

Proof. We can confirm that $T$ is self-adjoint. $\mathcal{D}_{o}(T) \cap \mathcal{H} \otimes|\uparrow\rangle=\mathcal{D}_{o}(X) \otimes|\uparrow\rangle$ and $T(\phi \otimes|\uparrow\rangle)=(X \phi) \otimes(|-\rangle\langle+|+|+\rangle\langle-|)|\uparrow\rangle=(X \phi) \otimes|\uparrow\rangle=X \phi \otimes|\uparrow\rangle$ holds for any $\phi \in \mathcal{D}_{o}(X)$. The lemma is immediately obvious.

For example, we apply the inequalities (9) in lemma 21 to a symmetric operator. If $X$ is self-adjoint, we have $X^{*} X=X X^{*}$. However, if the operator $X$ has no self-adjoint extension, we have $X^{*} X \varsubsetneqq X X^{*}$. This fact does not contradict the inequalities (9).

We have the following lemma from the classification by Naǐmark and the following fact; any maximal symmetric operator is unitarily equivalent to $\left(I \otimes P^{+}\right) \oplus Y$ or $\left(I \otimes P^{-}\right) \oplus Y$, where $Y$ is a self-adjoint operator and $P^{+}$and $P^{-}$are the momentum operators on $L^{2}\left(\mathbb{R}^{+}\right)$and $L^{2}\left(\mathbb{R}^{-}\right)$, respectively. This fact follows from section 104 in Ahkiezer and Glazman [22].

Lemma 29. Any minimal normal (self-adjoint) extension of a closed symmetric operator $X$ is unitarily equivalent to each other if and only if $X$ is maximal symmetric.

Remark 2. Lemma 29 gives an example of a subnormal operator such that its minimal normal extension is not unique in the sense of unitary equivalence.

## 4. Irreducible unitary representations of the algebra $\mathfrak{s u}(1,1)$ and their coherent states

In this section, from the minimal-weight-type unitary representations of the algebra $\mathfrak{s u}(1,1)$ (defined in this section), we will construct the corresponding subnormal operators canonically, and will investigate the relationship between the coherent states defined by Perelomov [4] and these subnormal operators.

Definition 30. A triplet $\left(E_{0}, E_{+}, E_{-}\right)$of skew-adjoint operators is called a unitary representation of the algebra $\mathfrak{s u}(1,1)$ if the relations

$$
\begin{equation*}
\left[E_{0}, E_{ \pm}\right]= \pm 2 E_{ \pm} \quad\left[E_{+}, E_{-}\right]=E_{0} \tag{11}
\end{equation*}
$$

hold.
For the reasoning behind this definition see remark 3. However, it is difficult to discuss the unitary representation in this notation because three operators $E_{0}, E_{+}, E_{-}$have no eigenvector. Thus, we define another triplet ( $L_{0}, L_{+}, L_{-}$) by

$$
\begin{equation*}
L_{0}:=\mathrm{i}\left(E_{-}-E_{+}\right) \quad L_{ \pm}:=\frac{1}{2}\left(E_{0} \pm \mathrm{i}\left(E_{+}+E_{-}\right)\right) \tag{12}
\end{equation*}
$$

Then, this triplet satisfies the commutation relations of the same type

$$
\begin{equation*}
\left[L_{0}, L_{ \pm}\right]= \pm 2 L_{ \pm} \quad\left[L_{+}, L_{-}\right]=L_{0} \tag{13}
\end{equation*}
$$

For this triplet,

$$
\begin{equation*}
L_{0}^{*}=L_{0} \quad L_{+}^{*}=-L_{-} \tag{14}
\end{equation*}
$$

hold, where $L_{+}$and $L_{-}$are neither self-adjoint nor skew-adjoint. Conversely, from the triplet ( $L_{0}, L_{+}, L_{-}$) satisfying the conditions (13) and (14), a unitary representation ( $E_{0}, E_{+}, E_{-}$) of the algebra $\mathfrak{s u}(1,1)$ can be constructed by

$$
\begin{equation*}
E_{0}=L_{+}+L_{-} \quad E_{ \pm}= \pm \frac{\mathrm{i}}{2}\left(L_{0} \mp L_{+} \pm L_{-}\right) \tag{15}
\end{equation*}
$$

The Casimir operator is useful for the analysis of the representation. In the case of the algebra $\mathfrak{s u}(1,1)$, it is given by

$$
\begin{equation*}
C:=E_{0}^{2}+2\left(E_{+} E_{-}+E_{-} E_{+}\right)=L_{0}^{2}+2\left(L_{+} L_{-}+L_{-} L_{+}\right) \tag{16}
\end{equation*}
$$

For the general definition, see pp 130-1 of Perelomov [4] or p 45 of Howe and Tan [23]. The relation (16) can be written in another form

$$
\begin{equation*}
C=L_{0}^{2}-2 L_{0}+4 L_{+} L_{-} \tag{17}
\end{equation*}
$$

by using (13). From (11) and (13), the Casimir operator $C$ is commutative with $E_{0}, E_{+}, E_{-}, L_{0}, L_{+}$and $L_{-}$. From the Schur's lemma, in any irreducible representation, the Casimir operator $C$ is constant.

Lemma 31. Non-trivial irreducible unitary representations of $\mathfrak{s u}(1,1)$ are classified into the following three cases:
case 1: $\quad \operatorname{dim} L_{+}=0 \quad$ and $\quad \operatorname{dim} L_{-}=1$
case 2: $\quad \operatorname{dim} L_{+}=1 \quad$ and $\quad \operatorname{dim} L_{-}=0$
case 3: $\quad \operatorname{dim} L_{+}=0 \quad$ and $\quad \operatorname{dim} L_{-}=0$.
Case 2 is reduced to case 1 , by exchanging $L_{-}$for $L_{+}$and by changing the sign of $L_{0}$. We will not treat case 3 in this paper. Thus, only case 1 will be discussed.

Proof. The irreducibility requires that the dimensions of the kernels of $L_{-}$and $L_{+}$are not more than one. Moreover, if the dimensions of both kernels are one, then the representation should be finite dimensional. However, this circumstance is forbidden by the unitarity of the representation. Now, the lemma follows immediately.

Lemma 32. The unit vector $|0\rangle_{N}$ belonging to the Kernel of $L_{-}$is an eigenvector of $L_{0}$. This eigenvalue $\lambda$ is called the lowest weight and specifies the irreducible unitary representation of $\mathfrak{s u}(1,1)$ uniquely and satisfies $\lambda>0$. The equations

$$
\begin{align*}
& L_{0}|n\rangle_{N}=(\lambda+2 n)|n\rangle_{N} \\
& L_{+}|n\rangle_{N}=\sqrt{(n+1)(\lambda+n)}|n+1\rangle_{N}  \tag{18}\\
& L_{-}|n\rangle_{N}=-\sqrt{n(\lambda+n-1)}|n-1\rangle_{N}
\end{align*}
$$

hold, where we define

$$
|n\rangle_{N}:=\frac{1}{\|\left(L_{+}\right)^{n}|0\rangle_{N} \|}\left(L_{+}\right)^{n}|0\rangle_{N} .
$$

Proof. Because the Casimir operator should be scalar-valued, we can show that $|0\rangle_{N}$ is the eigenvector of $L_{0}$, from (17).

Let $v_{n}:=\left(L_{+}\right)^{n}|0\rangle_{N}$. The commutation relations (13) yields the following relations:

$$
\begin{aligned}
& L_{0} v_{n}=(\lambda+2 n) v_{n} \\
& L_{+} v_{n}=v_{n+1} \\
& L_{-} v_{n}=-n(\lambda+n-1) v_{n-1}
\end{aligned}
$$

whence we can confirm that the lowest weight $\lambda$, with which $|0\rangle_{N}$ is associated, specifies the representation uniquely. From the above assumptions, we can confirm that the basis $\left\{v_{n}\right\}_{n=1}^{\infty}$ is complete and orthogonal. From the above relations, the Casimir operator $C$ is calculated to be the scalar $\lambda(\lambda-2)$. From the commutation relations (13), we have

$$
\left\langle v_{n}, v_{n}\right\rangle=n(\lambda+n-1)\left\langle v_{n-1}, v_{n-1}\right\rangle .
$$

Therefore, the equation

$$
|n\rangle_{N}=\sqrt{\frac{\Gamma(\lambda)}{n!\Gamma(\lambda+n)}} v_{n}
$$

holds. Thus, equation (18) follows immediately. The unitarity of the representation guarantees $\lambda>0$. (See theorem 1.1.5 on p 96 of Howe and Tan [23].)

In the following discussions, $\mathcal{H}_{\lambda}$ denotes the representation space of the irreducible unitary representation of $\mathfrak{s u}(1,1)$ characterized by the lowest weight $\lambda$. We call such a representation
(i.e. case 1) lowest-weight-type. The representation of the Lie group $S U(1,1)$ cannot be constructed unless $\lambda$ is an integer as the representation of the Lie group $S O$ (3) cannot unless the total momentum is an integer. (For more details, see remark 3.) In particular, when the lowest weight $\lambda$ is an integer, the representation of the Lie group $S U(1,1)$ is well known as the discrete series [23,24].

Definition 33. The operator $N:=\frac{1}{2}\left(L_{0}-\lambda\right)$ is called the $\mathfrak{s u}(1,1)$-number operator because of (18). The bounded operator $a:=\frac{1}{2} L_{+}^{-1}\left(L_{0}-\lambda\right)=L_{+}^{-1} N$ is called the $\mathfrak{s u}(1,1)$-annihilation operator. Its definition is well defined because the vector $N|n\rangle_{N}$ belongs to the range of $L_{+}$ for any $n$ and the kernel of $L_{+}$is $\{0\}$. The $\mathfrak{s u}(1,1)$-creation operator is defined by the adjoint $a^{*}$ of $a$.

The equations

$$
\begin{align*}
& a|n\rangle_{N}=\sqrt{\frac{n}{n+\lambda-1}}|n-1\rangle_{N}  \tag{19}\\
& a^{*}|n\rangle_{N}=\sqrt{\frac{n+1}{n+\lambda}}|n+1\rangle_{N} \tag{20}
\end{align*}
$$

hold, where we mean that $a|0\rangle_{N}=0$ by (19) in the exceptional case where $\lambda=1, n=0$, as a convention. From (19), the commutation relation $[a, N]=a$ is derived. From (19) and (20), we have

$$
\begin{align*}
& a^{*} a=(N+\lambda-1)^{-1} N \quad a a^{*}=(N+\lambda)^{-1}(N+1)  \tag{21}\\
& {\left[a, a^{*}\right]=(\lambda-1)(N+\lambda)^{-1}(N+\lambda-1)^{-1}}
\end{align*}
$$

for $\lambda \neq 1$, and

$$
\begin{equation*}
a a^{*}=I \quad a^{*} a=I-|0\rangle_{N}{ }_{N}\langle 0| \quad\left[a, a^{*}\right]=|0\rangle_{N}{ }_{N}\langle 0| \tag{22}
\end{equation*}
$$

instead of (21) for $\lambda=1$. Next, we will construct the $\mathfrak{s u}(1,1)$-coherent state as follows:
Definition 34. Introduce the unitary operator $U(\xi):=\exp \left(\xi L_{+}-\xi^{*} L_{+}^{*}\right)$ for a complex number $\xi$, according to Perelomov [4]. For the complex number $\zeta$ such that $|\zeta|<1$, we define the $\mathfrak{s u}(1,1)$-coherent state $|\zeta\rangle_{a}$ of the algebra $\mathfrak{s u}(1,1)$ by

$$
\begin{equation*}
|\zeta\rangle_{a}:=U\left(\frac{1}{2} \mathrm{e}^{\mathrm{i} \arg \zeta} \ln \frac{1+|\zeta|}{1-|\zeta|}\right)|0\rangle_{N} . \tag{23}
\end{equation*}
$$

Squeezed states are characterized as $\mathfrak{s u}(1,1)$-coherent states, as we discuss in section 5.2.
Lemma 35. The $\mathfrak{s u}(1,1)$-coherent state $|\zeta\rangle_{a}$ is an eigenvector of $a$, i.e. the equation

$$
\begin{equation*}
a|\zeta\rangle_{a}=\zeta|\zeta\rangle_{a} \tag{24}
\end{equation*}
$$

holds.
Proof. From the definition, we have

$$
\begin{aligned}
|\zeta\rangle_{a} & =\exp \left(\zeta L_{+}\right) \exp \left(\frac{1}{2} \ln \left(1-|\zeta|^{2}\right) L_{0}\right) \exp \left(\zeta^{*} L_{-}\right)|0\rangle_{N} \\
& =\left(1-|\zeta|^{2}\right)^{\lambda / 2} \exp \left(\zeta L_{+}\right)|0\rangle_{N}
\end{aligned}
$$

(see pp 73-4 of Perelomov [4] for the derivation of the first equation). Because we can show that $\left[a, L_{+}\right]=I$, we obtain the commutation relation $\left[a, \exp \left(\zeta L_{+}\right)\right]=\zeta \exp \left(\zeta L_{+}\right)$. Moreover, from the relation $\exp \left(\frac{1}{2} \ln \left(1-|\zeta|^{2}\right) L_{0}\right) \exp \left(\zeta^{*} L_{-}\right)|0\rangle_{N}=\left(1-|\zeta|^{2}\right)^{\lambda / 2}|0\rangle_{N}$, we have

$$
\begin{equation*}
a|\zeta\rangle_{a}=\exp \left(\zeta L_{+}\right) a|0\rangle_{N}+\zeta \exp \left(\zeta L_{+}\right)|0\rangle_{N}=\zeta|\zeta\rangle_{a} . \tag{25}
\end{equation*}
$$

Therefore, the coherent states of the algebra $\mathfrak{s u}(1,1)$ are characterized as eigenvectors of the $\mathfrak{s u}(1,1)$-annihilation operator $a$.

Lemma 36. When $1>\lambda>0, a^{*}$ and a are not subnormal. When $\lambda \geqslant 1, a^{*}$ is subnormal and $a$ is not subnormal. For $\lambda>1, a^{*}$ 's POVM is given by $(\lambda-1)\left|\zeta^{*}\right\rangle_{a}{ }_{a}\left\langle\zeta^{*}\right| \mu(\mathrm{d} \zeta)$, where we define

$$
\mu(\mathrm{d} \zeta):=\frac{\mathrm{d}^{2} \zeta}{\pi\left(1-|\zeta|^{2}\right)^{2}}
$$

Proof. When $\lambda>0$, it is shown that $a$ is not subnormal, from lemma 23 and the fact that it has eigenvectors. When $\lambda<1$, it is shown that $a^{*}$ is not subnormal, from lemma 21 and the fact that $\left[a, a^{*}\right] \geqslant 0$ does not hold (see (21)).

Moreover, when $\lambda>1$, we can construct the resolution of the identity by the system of the coherent states:

$$
\begin{equation*}
(\lambda-1) \int_{D}|\zeta\rangle_{a a}\langle\zeta| \mu(\mathrm{d} \zeta)=I \tag{26}
\end{equation*}
$$

where $D$ denotes the unit disc $\{z \in \mathbb{C}||z|<1\}$. From this resolution of the identity and lemma 22, when $\lambda>1$, we can show that $a^{*}$ is a subnormal operator. On the other hand, when $\lambda \leqslant 1$, the integral in (26) diverges. However, equations (22) imply that $a^{*}$ is isometric when $\lambda=1$. Then, $a^{*}$ is subnormal even when $\lambda=1$.

Definition 37. We formally define the operator

$$
A:=-\mathrm{i}(a+1)(a-1)^{-1}
$$

Since this operator is unbounded, we need to pay more attention to this definition. First, define the unbounded operator $\tilde{A}$ by a linear fractional transform (Möbius transform) of $a$, as $\tilde{A}:=-\mathrm{i}(a+1)(a-1)^{-1}$, where the domain $\mathcal{D}_{o}(\tilde{A})$ of $\tilde{A}$ is defined by $\left\langle\left\{|n\rangle_{N}\right\}_{n=0}^{\infty}\right\rangle$. The domain of $\tilde{A}^{*}$ is dense in $\mathcal{H}_{\lambda}$, as will be shown in the last part of remark 6. Therefore, $\tilde{A}$ is closable and we can define the operator $A$ by $A:=\overline{\tilde{A}}=\tilde{A}^{* *}$. (See Reed and Simon [25].)

It is shown that $|\zeta\rangle_{a} \in \mathcal{D}_{o}(A)$ in the last part of remark 6 . Hence we have

$$
A|\zeta\rangle_{a}=-\mathrm{i} \frac{\zeta+1}{\zeta-1}|\zeta\rangle_{a} .
$$

By defining

$$
|\eta\rangle_{A}:=\left|\frac{\eta-\mathrm{i}}{\eta+\mathrm{i}}\right\rangle_{a}
$$

we can show that

$$
\begin{equation*}
A|\eta\rangle_{A}=\eta|\eta\rangle_{A} \tag{27}
\end{equation*}
$$

holds. (Formally, the operator $a$ is the Cayley transform of $A$, with an appropriate discussion on its domain.)

Lemma 38. We have another expression for $A$ :

$$
\begin{equation*}
A=\frac{1}{2} E_{+}^{-1}\left(E_{0}-\lambda\right) . \tag{28}
\end{equation*}
$$

Proof. From the relations $\left[a, L_{+}\right]=I,(13),(15)$ and the definition of $a$, we can show that the relations

$$
\begin{aligned}
2\left(E_{0}-\lambda\right)(a-1) L_{+} & =\left(L_{+}+L_{-}-\lambda\right)\left(L_{0}-(\lambda-2)-2 L_{+}\right) \\
& =\left(L_{0}-L_{+}+L_{-}\right)\left(-\lambda+L_{0}+2+2 L_{+}\right) \\
& =-4 \mathrm{i} E_{+}(a+1) L_{+} \\
\left(E_{0}-\lambda\right)(a-1)|0\rangle_{N} & =-\left(E_{0}-\lambda\right)|0\rangle_{N}=\left(L_{0}-L_{+}\right)|0\rangle_{N} \\
& =-2 \mathrm{i} E_{+}|0\rangle_{N}=-2 \mathrm{i} E_{+}(a+1)|0\rangle_{N}
\end{aligned}
$$

hold on $\mathcal{D}_{o}(\tilde{A})$. Hence, on $\mathcal{D}_{o}(\tilde{A})$, we have

$$
\begin{equation*}
\left(E_{0}-\lambda\right)(a-1)=-2 \mathrm{i} E_{+}(a+1) \tag{29}
\end{equation*}
$$

By using (29), we obtain (28).
From (28), we have

$$
\left[A, A^{*}\right]=-(\lambda-1) E_{+}^{-2}
$$

formally, and
$A^{*} A-A A^{*}=(\lambda-1)\left(E_{+}^{-1}\right)^{*} E_{+}^{-1} \quad \begin{cases}\text { on } \mathcal{D}_{o}\left(A A^{*}\right) & \text { for } \lambda \geqslant 1 \\ \text { on } \mathcal{D}_{o}\left(A^{*} A\right) & \text { for } 0<\lambda<1\end{cases}$
in more precise form. (The proof of this relation will be given in remark 6.)
Lemma 39. When $1>\lambda>0, A$ and $A^{*}$ are not subnormal. When $\lambda \geqslant 1, A^{*}$ is subnormal and $A$ is not subnormal. For $\lambda>1, A^{*}$ 's POVM is given by $(\lambda-1)\left|\eta^{*}\right\rangle_{A}{ }_{A}\left\langle\eta^{*}\right| \nu(\mathrm{d} \eta)$, where we define

$$
v(\mathrm{~d} \eta):=\frac{\mathrm{d}^{2} \eta}{4 \pi(\operatorname{Im} \eta)^{2}}
$$

Proof. For $\lambda>0$, from lemma 23 and the fact that the operator $A$ has eigenvectors, it is shown that $A$ is not subnormal. When $\lambda<1, A^{*}$ is not subnormal because relation (21) shows that the condition $A A^{*} \geqslant A^{*} A$ is not satisfied. Moreover, in a similar manner to the above discussion, the resolution of the identity by the eigenvectors of $A$

$$
(\lambda-1) \int_{\mathrm{H}}|\eta\rangle_{A A}\langle\eta| \nu(\mathrm{d} \eta)=I
$$

holds when $\lambda>1$. Hence, when $\lambda>1$, we can show that $A^{*}$ is subnormal from lemma 22. When $\lambda \leqslant 1$, the integral in (26) diverges. However, as will be proved in the last part of remark 6 , the operator $A^{*}$ is maximal symmetric when $\lambda=1$. From lemma $28, A^{*}$ is subnormal even when $\lambda=1$.

Remark 3 (Relation to unitary representations of $\boldsymbol{S U}(\mathbf{1}, \mathbf{1})$ ). In the following, we discuss definition 30 from the viewpoint of a unitary representation of the group $S U(1,1)$. Any element $g$ in the group $S U(1,1)$ is specified by two complex numbers $\nu(g)=\nu_{1}(g)+v_{2}(g) \mathrm{i}, \mu(g)=$ $\mu_{1}(g)+\mu_{2}(g)$ i satisfying $|\nu(g)|^{2}-|\mu(g)|^{2}=1$ as

$$
g=\left(\begin{array}{cc}
\mu^{*}(g) & v(g) \\
v^{*}(g) & \mu(g)
\end{array}\right)=\left(\begin{array}{cc}
\mu_{1}(g)-\mu_{2}(g) \mathrm{i} & v_{1}(g)+v_{2}(g) \mathrm{i} \\
v_{1}(g)-v_{2}(g) \mathrm{i} & \mu_{1}(g)+\mu_{2}(g) \mathrm{i}
\end{array}\right) .
$$

The group $S U(1,1)$ is isomorphic to the group $S L(2, \mathbb{R})$ by the map
$j:\left(\begin{array}{cc}\mu_{1}(g)-\mu_{2}(g) \mathrm{i} & \nu_{1}(g)+\nu_{2}(g) \mathrm{i} \\ \nu_{1}(g)-v_{2}(g) \mathrm{i} & \mu_{1}(g)+\mu_{2}(g) \mathrm{i}\end{array}\right) \mapsto\left(\begin{array}{cc}\mu_{1}(g)+v_{1}(g) & -\mu_{2}(g)-v_{2}(g) \\ \mu_{2}(g)-v_{2}(g) & \mu_{1}(g)-v_{1}(g)\end{array}\right)$.
The (Lie) algebra $\mathfrak{s u}(1,1)$ associated with $S U(1,1)$ is written as
$\mathfrak{s u}(1,1)=\left\{\left(\begin{array}{ll}a_{1,1} & a_{1,2} \\ a_{2,1} & a_{2,2}\end{array}\right) \left\lvert\,\left(\begin{array}{cc}a_{1,1}^{*} & -a_{2,1}^{*} \\ -a_{1,2}^{*} & a_{2,2}^{*}\end{array}\right)=-\left(\begin{array}{ll}a_{1,1} & a_{1,2} \\ a_{2,1} & a_{2,2}\end{array}\right)\right., a_{1,1}+a_{2,2}=0\right\}$.
The vector space $\mathfrak{s u}(1,1)$ has the following basis $e_{0}, e_{+}, e_{-}$as

$$
e_{0}=\left(\begin{array}{ll}
0 & 1 \\
1 & 0
\end{array}\right) \quad e_{+}=\frac{\mathrm{i}}{2}\left(\begin{array}{cc}
1 & -1 \\
1 & -1
\end{array}\right) \quad e_{-}=\frac{\mathrm{i}}{2}\left(\begin{array}{cc}
-1 & -1 \\
1 & 1
\end{array}\right)
$$

Thus, from the isomorphism (31), we can naturally define the isomorphism $j_{*}$ from the algebra $\mathfrak{s u}(1,1)$ to the algebra $\mathfrak{s l}(2, \mathbb{R})$. Then, the image $j_{*}\left(e_{0}\right), j_{*}\left(e_{-}\right), j_{*}\left(e_{+}\right)$of the basis is written as
$j_{*}\left(e_{0}\right)=\left(\begin{array}{cc}1 & 0 \\ 0 & -1\end{array}\right) \quad j_{*}\left(e_{+}\right)=\left(\begin{array}{cc}0 & 1 \\ 0 & 0\end{array}\right) \quad j_{*}\left(e_{-}\right)=\left(\begin{array}{cc}0 & 0 \\ 1 & 0\end{array}\right)$.
The basis $e_{0}, e_{-}, e_{+}$satisfies the following commutation relation:

$$
\begin{equation*}
\left[e_{0}, e_{ \pm}\right]= \pm 2 e_{ \pm} \quad\left[e_{+}, e_{-}\right]=e_{0} \tag{33}
\end{equation*}
$$

A map $V$ from a group $G$ to the set of unitary operators on a Hilbert space $\mathcal{H}$ is called a unitary representation of the group $G$ on $\mathcal{H}$ if

$$
V\left(g_{1} g_{2}\right)=V\left(g_{1}\right) V\left(g_{2}\right) \quad \forall g_{1}, g_{2} \in G
$$

Let $\mathfrak{g}$ be the Lie algebra associated with a Lie group $G$. From a unitary representation of the group $G$ on a Hilbert space $\mathcal{H}$, we can naturally define the map $V_{*}$ from the Lie algebra $\mathfrak{g}$ to the set of skew-adjoint operators on $\mathcal{H}$, by

$$
V_{*}(X):=\left.\frac{\mathrm{d} V(\exp (t X))}{\mathrm{d} t}\right|_{t=0} .
$$

It satisfies that $V_{*}([X, Y])=\left[V_{*}(X), V_{*}(Y)\right]$. Then, a linear map $f$ from a Lie algebra $\mathfrak{g}$ to the set of skew-adjoint operators on a Hilbert space $\mathcal{H}$ is called a unitary representation of the Lie algebra $\mathfrak{g}$ on $\mathcal{H}$ if

$$
[f(X), f(Y)]=f([X, Y]) \quad \forall X, Y \in \mathfrak{g} .
$$

We can construct the unitary representation $V$ of the universal covering group $\dagger \hat{G}$ associated with a Lie algebra $\mathfrak{g}$ from a unitary representation $f$ of $\mathfrak{g}$, by

$$
V(\exp X):=\exp f(X) \quad \forall X \in \mathfrak{g} .
$$

Since any element of the (Lie) algebra $\mathfrak{s u}(1,1)$ is described by a linear sum of bases $e_{0}, e_{+}, e_{-}$, we can uniquely construct the unitary representation of the algebra $\mathfrak{s u}(1,1)$ from a triplet ( $E_{0}, E_{+}, E_{-}$) of skew-adjoint operators satisfying (11). Thus, we can regard the triplet $\left(E_{0}, E_{+}, E_{-}\right)$satisfying (11) as the unitary representation of the algebra $\mathfrak{s u}(1,1)$.

Remark 4 (Spectra of $\boldsymbol{a}, \boldsymbol{a}^{*}, \boldsymbol{A}, \boldsymbol{A}^{*}$ ). The point spectrum $\sigma_{p}(\boldsymbol{a})$ is the open unit disc $D$, the continuous spectrum $\sigma_{c}(a)$ is the unit circle $S:=\{z \in \mathbb{C}| | z \mid=1\}$ and the residual spectra is the empty set. Moreover, from lemma 45 , the point spectrum $\sigma_{p}\left(a^{*}\right)$, the continuous spectrum $\sigma_{c}\left(a^{*}\right)$ and the residual spectrum $\sigma_{r}\left(a^{*}\right)$ of $a^{*}$ are the empty set, $S$ and $D$, respectively.

[^1]It is shown that the point spectrum $\sigma_{p}(A)$ of $A$ is the upper-half-plane $H$, the continuous spectrum $\sigma_{c}(A)$ is the real axis $\mathbb{R}$ and the residual spectrum $\sigma_{r}(A)$ is the empty set. On the other hand, from lemma 45 , the point spectrum $\sigma_{p}\left(A^{*}\right)$, the continuous spectrum $\sigma_{c}\left(A^{*}\right)$ and the residual spectrum $\sigma_{r}\left(A^{*}\right)$ of $A^{*}$ are the empty set, $\mathbb{R}$ and $H$, respectively.

Remark 5 (Action of the group to operators $\boldsymbol{a}, \boldsymbol{A}$ ). First, we discuss the action to an operator $a$. We let $\pi_{\widehat{S U(1,1)}}$ be the projection from the universal covering group $\widehat{\operatorname{SU(1,1)}}$ to the group $S U(1,1)$, and let $U$ be the maximal Cartan subgroup of $\widehat{S U(1,1)}$, i.e. the oneparameter subgroup generated by $\mathrm{i} L_{0}$.

The homogeneous space $S \widehat{U(1,1)} / U$ is isomorphic to the open unit disc $D$ in the sense that an element $g$ of the group $\widehat{S(1,1)}$ acts on the open unit disc $D$ as

$$
\zeta \mapsto \frac{\mu^{*} \zeta+v}{v^{*} \zeta+\mu} \quad \zeta \in D
$$

where we simply use the notation $\mu$ and $\nu$ instead of the complex numbers $\mu \circ \pi_{\widehat{S U(1,1)}}(g)$ and $v \circ \pi_{\widehat{S U(1,1)}}(g)$ with the functions $\mu$ and $v$ defined at the beginning of remark 3 , respectively. We let $V$ be the representation of the group $\widehat{S(1,1)}$, defined by this representation of $\mathfrak{s u}(1,1)$. Then, we have
$V(g)|\zeta\rangle_{a a}\langle\zeta| V(g)^{*}=\left|\frac{\mu^{*} \zeta+v}{\nu^{*} \zeta+\mu}\right\rangle_{a} \quad\left\langle\frac{\mu^{*} \zeta+v}{\nu^{*} \zeta+\mu}\right| \quad g \in S \widehat{U(1,1)} \quad \zeta \in \mathbb{C}$.
Thus, for any element $g \in \widehat{S(1,1)}$ and any complex number $\zeta$, there exists a real number $\theta(g, \zeta)$ such that

$$
\begin{equation*}
V(g)|\zeta\rangle_{a}=\mathrm{e}^{\mathrm{i} \theta(g, \zeta)}\left|\frac{\mu^{*} \zeta+v}{\nu^{*} \zeta+\mu}\right|_{a} \tag{35}
\end{equation*}
$$

Equations (24) and (35) imply that

$$
\begin{equation*}
V(g)^{*} a V(g)|\zeta\rangle_{a}=\frac{\mu^{*} \zeta+v}{v^{*} \zeta+\mu}|\zeta\rangle_{a} . \tag{36}
\end{equation*}
$$

Since the subspace $\left\langle\left\{|\zeta\rangle_{a}\right\}\right\rangle \dagger$ is dense, we obtain

$$
V(g)^{*} a V(g)=\left(\mu^{*} a+v\right)\left(v^{*} a+\mu\right)^{-1}
$$

where we can define the bounded operator $\left(\nu^{*} a+\mu\right)^{-1}$ by

$$
\left(v^{*} a+\mu\right)^{-1}:=\frac{1}{\mu} \sum_{n=1}^{\infty}\left(-\frac{v^{*}}{\mu} a\right)^{n}
$$

because the norm of the operator $-\frac{v^{*}}{\mu} a$ is less than 1 .
Next, we consider the action to the operator $A$. Similarly to (36), we have

$$
V(g) A V(g)^{*}|\eta\rangle_{A}=\frac{\left(\mu_{1}+v_{1}\right) \eta-\mu_{2}-v_{2}}{\left(\mu_{2}-v_{2}\right) \eta+\mu_{1}-v_{1}}|\eta\rangle_{A}
$$

where we simplify $\mu_{i} \circ \pi_{\widehat{S U(1,1)}}(g)$ and $\nu_{i} \circ \pi_{\widehat{S U(1,1)}}(g)$ as $\mu_{i}$ and $\nu_{i}$, respectively.
$\dagger\langle X\rangle$ denotes the vector space whose elements are finite linear sums of a set $X$.

## 5. Concrete representations of $\mathfrak{s u}(\mathbf{1}, 1)$

### 5.1. Representation associated with irreducible unitary representation of the affine group

Next, we will construct lowest-weight-type irreducible unitary representations of the algebra $\mathfrak{s u}(1,1)$ from an irreducible unitary representation of the affine group ( $a x+b$ group) generated by $E_{+}$and $E_{0}$. The representation which will be constructed in this section is closely related to the continuous wavelet transformation [26,27]. In this representation, the pair $A$ and $|\eta\rangle_{A}$ play a more important role than the pair $a$ and $|\zeta\rangle_{a}$. According to Aslaksen and Klauder [28], there is no irreducible representation of the affine group but the representations equivalent unitarily to the following representation on $L^{2}\left(\mathbb{R}^{+}\right)$or $L^{2}\left(\mathbb{R}^{-}\right)$

$$
\begin{equation*}
E_{0}=\mathrm{i}(P Q+Q P) \quad E_{+}=\mathrm{i} Q \tag{37}
\end{equation*}
$$

where $E_{0}$ and $E_{+}$are shown to be skew adjoint. In this representation, the vector

$$
\sqrt{\frac{(2 \operatorname{Im} \eta)^{2 k+1}}{\Gamma(2 k+1)}} x^{k} \mathrm{e}^{\mathrm{i} \eta x}
$$

is called the affine coherent state $\dagger$, and it is obtained by operating the affine group on the affine vacuum state

$$
\sqrt{\frac{2^{2 k+1}}{\Gamma(2 k+1)}} x^{k} \mathrm{e}^{\mp x}
$$

In the following, we will construct an irreducible unitary representation of the algebra $\mathfrak{s u}(1,1)$ from the above type of unitary representation of the affine group, and will discuss how to interpret the affine coherent states in terms of the unitary representation of the algebra $\mathfrak{s u}(1,1)$. Therefore, in addition to the two generators in (37), we should introduce the representation of another additional generator $E_{-}$. By choosing

$$
\begin{equation*}
\tilde{E}_{-, k}:=-\mathrm{i}\left(P Q P+k^{2} Q^{-1}\right) \quad\left(k>-\frac{1}{2}\right) \tag{38}
\end{equation*}
$$

for this additional generator, we can construct an irreducible unitary representation where the triplet $E_{0}, E_{+}$and $E_{-}$satisfies the commutation relations (11). However, we should be careful about the domain of $\tilde{E}_{-, k}$, as follows; first, define the dense subspace $\mathcal{D}_{o}\left(\tilde{E}_{-, k}\right)$ of $L^{2}\left(\mathbb{R}^{+}\right)$by
$\mathcal{D}_{o}\left(\tilde{E}_{-, k}\right):=\left\{f(x)=x^{k} f_{0}(x) \in L^{2}\left(\mathbb{R}^{+}\right)\right.$

$$
\left.\cap C^{1}\left(\mathbb{R}^{+}\right) \left\lvert\, \begin{array}{l}
(2 k+1) x^{k} f_{0}^{\prime}(x)+x^{k+1} f_{0}^{\prime \prime}(x) \in L^{2}\left(\mathbb{R}^{+}\right) \\
\lim \sup _{s \rightarrow 0} f_{0}(s)<\infty \quad x^{k} f_{0}(x) \rightarrow 0 \quad \text { as } x \rightarrow \infty
\end{array}\right.\right\}
$$

Then $\tilde{E}_{-, k}$ is an operator defined on $\mathcal{D}_{o}\left(\tilde{E}_{-, k}\right)$. We need attention to the domain when $-\frac{1}{2}<k<\frac{1}{2}$.
Lemma 40. The operator $\tilde{E}_{-, k}$ has the skew-adjoint extension, uniquely.
In the following, its skew-adjoint extension is written as $E_{-, k}$.
$\dagger$ The Fourier transform of this affine coherent state is equivalent to the Cauchy wavelet in signal processing, whose basic wavelet function is (constant) $/(t \pm i)^{k+1}$.

Proof. It is confirmed that $\mathrm{i} \tilde{E}_{-, k}=P Q P+k^{2} Q^{-1}$ is a symmetric operator on $\mathcal{D}_{o}\left(\tilde{E}_{-, k}\right)$, from the fact that the difference

$$
\begin{aligned}
\int_{s}^{t}((P Q P+ & \left.\left.k^{2} Q^{-1}\right) f\right)^{*}(x) g(x) \mathrm{d} x-\int_{s}^{t}(f(x))^{*}\left(\left(P Q P+k^{2} Q^{-1}\right) g\right)(x) \mathrm{d} x \\
= & {\left[x\left(f^{\prime}(x)\right)^{*} g(x)-x g^{\prime}(x)(f(x))^{*}\right]_{s}^{t} } \\
= & t\left(f^{\prime}(t)\right)^{*} g(t)-t g^{\prime}(t)(f(t))^{*} \\
& -\left(\operatorname{sg}(s)\left(\left(f^{\prime}(s)\right)^{*}-\frac{k}{s}(f(s))^{*}\right)-s(f(s))^{*}\left(g^{\prime}(s)-\frac{k}{s} g(s)\right)\right) \\
= & t\left(f^{\prime}(t)\right)^{*} g(t)-\operatorname{tg}^{\prime}(t)(f(t))^{*}-\left(\operatorname{sg}(s) s^{k}\left(f_{0}^{\prime}(s)\right)^{*}-s(f(s))^{*} s^{k} g_{0}^{\prime}(s)\right)
\end{aligned}
$$

tends to zero as $s \rightarrow 0, t \rightarrow \infty$. Since $\mathrm{i} \tilde{E}_{-, k}$ is semi-bounded, the Friedrich extension theorem guarantees that there uniquely exists the self-adjoint extension of $\mathrm{i} \tilde{E}_{-, k}$. (See p 177 of Reed and Simon [29].) Now, the proof is complete.

By letting $L_{+, k}, L_{-, k}, L_{0, k}, \tilde{A}_{k}, A_{k}, N_{k},|n\rangle_{N}^{k}$ and $|\eta\rangle_{A}^{k}$ be $L_{+}, L_{-}, L_{0}, \tilde{A}, A, N,|n\rangle_{N}$ and $|\eta\rangle_{A}$ in this representation, respectively, we have

$$
\begin{aligned}
& L_{+, k}=\frac{1}{2}\left(\mathrm{i}(P Q+Q P)-Q+P Q P+k^{2} Q^{-1}\right) \\
& L_{-, k}=\frac{1}{2}\left(\mathrm{i}(P Q+Q P)+Q-P Q P-k^{2} Q^{-1}\right) \\
& L_{0, k}=\left(P Q P+k^{2} Q^{-1}+Q\right) \quad \tilde{A}_{k}=P+\mathrm{i} k Q^{-1} \\
& N_{k}=\frac{1}{2}\left(P Q P+k^{2} Q^{-1}+Q-1-2 k\right) \\
& |n\rangle_{N}^{k}(x)=\sqrt{\frac{2^{2 k+1} n!}{\Gamma(n+2 k+1)}} \mathrm{e}^{-x} x^{k} S_{n}^{2 k}(2 x) \\
& |\eta\rangle_{A}^{k}(x)=\sqrt{\frac{(2 \operatorname{Im} \eta)^{2 k+1}}{\Gamma(2 k+1)}} x^{k} \mathrm{e}^{\mathrm{i} \eta x}
\end{aligned}
$$

when $S_{n}^{l}(x)$ is the Sonine polynomial (or the associated Laguerre polynomial) defined by $\dagger$

$$
S_{n}^{l}(x):=\sum_{m=0}^{n} \frac{(-1)^{m}}{(n-m)!} \frac{\Gamma(n+l+1) x^{m}}{\Gamma(m+l+1) m!}
$$

$|\eta\rangle_{A}^{k}(x)$ is the affine coherent state. Moreover, in this representation, the minimum eigenvalue of $L_{0, k}$ is $\lambda=2 k+1$, and the Casimir operator is $4 k^{2}-1$. Thus, we have the following theorem.
Theorem 41. The representations of lowest-weight-type (defined in section 4), in general, can be concretely constructed by (37) and (38) on $L^{2}\left(\mathbb{R}^{+}\right)$in the correspondence $\lambda=2 k+1$.

Remark 6 (Domains of $A_{k}, A_{k}^{*}$ ). In the following, we will show the properties of $\tilde{A}_{k}$ in order to show the properties of $\tilde{A}$ in the representations of lowest-weight-type. Since the domain of $\tilde{A}_{k}$ is $\left\langle\{|n\rangle\}_{n=0}^{\infty}\right\rangle$ and $\tilde{A}_{k}=P+\mathrm{i} k Q^{-1}$, the relation
$\mathcal{D}_{o}\left(\tilde{A}_{k}^{*}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)=\left\{\begin{array}{l|l}x^{-k} f(x) \in L^{2}\left(\mathbb{R}^{+}\right) \cap C^{1}\left(\mathbb{R}^{+}\right) & \begin{array}{l}x^{-k} f^{\prime}(x) \in L^{2}\left(\mathbb{R}^{+}\right) \\ f(s) \rightarrow 0 \text { as } s \rightarrow 0\end{array}\end{array}\right\}$

[^2]is derived, and hence we can show that $\mathcal{D}_{o}\left(\tilde{A}_{k}^{*}\right)$ is dense in $L^{2}\left(\mathbb{R}^{+}\right)$. Thus $\tilde{A}_{k}$ is shown to be a closable operator. Since $A_{k}=\bar{A}_{k}$, the relation

$\mathcal{D}_{o}\left(A_{k}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)=\left\{\begin{array}{l|l}x^{k} f(x) \in L^{2}\left(\mathbb{R}^{+}\right) \cap C^{1}\left(\mathbb{R}^{+}\right) & \begin{array}{l}x^{k} f^{\prime}(x) \in L^{2}\left(\mathbb{R}^{+}\right) \\ \lim \sup _{s \rightarrow 0} f(s)<\infty\end{array}\end{array}\right\}$
is confirmed. Note that $\bar{X}=X^{* *}$ and $X^{*}=\bar{X}^{*}$ hold for a densely defined linear operator $X$, and that $\lim \sup _{x \rightarrow \infty} f_{0}(x)=0$ for $k>-\frac{1}{2}$ when $x^{k} f_{0}(x) \in L^{2}\left(\mathbb{R}^{+}\right)$. From (40), we can show that $|\zeta\rangle_{a} \in \mathcal{D}_{o}\left(A_{k}\right)$. Thus, the subspaces $\mathcal{D}_{o}\left(A_{k}^{*}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)$and $\mathcal{D}_{o}\left(A_{k}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)$are the cores $\dagger$ of $A_{k}^{*}$ and $A_{k}$. When $-\frac{1}{2}<k<\frac{1}{2}$, the domain of $A_{k}$ is larger than the domain of $A_{-k}^{*}$, though $A_{k}$ and $A_{-k}^{*}$ are the same formally, i.e. $A_{-k}^{*} \varsubsetneqq A_{k}$. In the special case where $\lambda=1$ (i.e. where $k=0$ ), $A_{0}^{*}$ is symmetric. Since $A_{0}=A_{0}^{* *}$ has no spectrum in the lower half-plane, $A_{0}^{*}$ 's deficiency indices are ( 1,0 ). (For the definition of deficiency indices, see p 138 of Reed and Simon [29] or p 360 of Rudin [11].) Therefore, the operator $A_{0}^{*}$ is a maximally symmetric operator.

These relations $\mathcal{D}_{o}\left(A_{k}^{*}\right) \subset \mathcal{D}_{o}\left(A_{k}\right)$ and $\mathcal{D}_{o}\left(A_{k}^{*}\right) \subset \mathcal{D}_{o}\left(E_{+}^{-1}\right)$ are shown in the cases where $\lambda<1(k>0)$, only the relation $\mathcal{D}_{o}\left(A_{k}^{*}\right) \subset \mathcal{D}_{o}\left(A_{k}\right)$ is shown when $\lambda=1(k=0)$, and these relations $\mathcal{D}_{o}\left(A_{k}\right) \subset \mathcal{D}_{o}\left(A_{k}^{*}\right)$ and $\mathcal{D}_{o}\left(A_{k}\right) \subset \mathcal{D}_{o}\left(E_{+}^{-1}\right)$ are shown when $0<\lambda<1\left(-\frac{1}{2}<k<\right.$ 0 ). From theorem 6, these discussions and (28), we obtain (30).

### 5.2. Representation associated with squeezed states

Next, we will discuss the following representation of the algebra $\mathfrak{s u}(1,1)$ on the Hilbert space $L^{2}(\mathbb{R})$; let

$$
\begin{equation*}
E_{0}=\frac{\mathrm{i}}{2}(P Q+Q P) \quad E_{+}=\frac{\mathrm{i}}{2} Q^{2} \quad E_{-}=-\frac{\mathrm{i}}{2} P^{2} \tag{41}
\end{equation*}
$$

then we have

$$
\begin{equation*}
L_{0}=n_{b}+\frac{1}{2} \quad L_{+}=-\frac{1}{2}\left(a_{b}^{*}\right)^{2} \quad L_{-}=\frac{1}{2}\left(a_{b}\right)^{2} \tag{42}
\end{equation*}
$$

where the boson annihilation operator $a_{b}$ and the boson number operator $n_{b}$ are given by

$$
a_{b}=\sqrt{\frac{1}{2}}(Q+\mathrm{i} P)
$$

and

$$
n_{b}=\frac{1}{2}\left(Q^{2}+P^{2}-1\right)=a_{b}^{*} a_{b}
$$

In this representation, the Casimir operator is the scalar $-\frac{3}{4}$. From the fact that the Casimir operator is the scalar $\lambda(\lambda-2)$, the solutions are $\lambda=\frac{1}{2}, \frac{3}{2}$. Under the representation given in $(41), L^{2}(\mathbb{R})$ is not irreducible and it is decomposed into two irreducible subspaces as

$$
L^{2}(\mathbb{R})=L_{\text {even }}^{2}(\mathbb{R}) \oplus L_{\text {odd }}^{2}(\mathbb{R})
$$

where $L_{\text {even }}^{2}(\mathbb{R})$ is the set of square-integrable even functions and $L_{\text {odd }}^{2}(\mathbb{R})$ is the set of squareintegrable odd functions. The solution $\lambda=\frac{1}{2}$ corresponds to the subspace $L_{\text {even }}^{2}(\mathbb{R})$ and the
$\dagger$ A subspace of the domain $\mathcal{D}_{o}(X)$ of a closed operator $X$ is called a core of $X$ if it is dense in $\mathcal{D}_{o}(X)$ with respect to the graph norm of the operator $X$.
solution $\lambda=\frac{3}{2}$ does to the subspace $L_{\text {odd }}^{2}(\mathbb{R})$. In the subspace $L_{\text {even }}^{2}(\mathbb{R})$, the operators $a, A$ and $N$ are written in the forms

$$
\begin{aligned}
& a=-\left(a_{b}^{*}\right)^{-1} a_{b} \quad A=Q^{-1} P \quad N=\frac{1}{2} n_{b} \quad|n\rangle_{N}=(-1)^{n}|2 n\rangle_{n_{b}} \\
& \mathcal{D}_{o}(A) \cap C^{1}(\mathbb{R})=\left\{f \in L_{\text {even }}^{2}(\mathbb{R}) \cap C^{1}(\mathbb{R}) \left\lvert\, \frac{1}{x} f^{\prime}(x) \in L^{2}(\mathbb{R})\right.\right\}
\end{aligned}
$$

where $|n\rangle_{n_{b}}$ denotes the eigenvector in $L^{2}(\mathbb{R})$ of the boson number operator $n_{b}$ associated with the eigenvalue $n$.

Lemma 42. In the action of $\mathfrak{s u}(1,1)$ on $L_{\text {even }}^{2}(\mathbb{R})$, we have
$|0 ; \mu, v\rangle\langle 0 ; \mu, v|=V(g)|0\rangle_{a}{ }_{a}\langle 0| V(g)^{*}=\left|\frac{v}{\mu}\right\rangle_{a}{ }_{a}\left\langle\frac{v}{\mu}\right|=\left|\frac{\mu+v}{\mu-v}\right\rangle_{A}{ }_{A}\left\langle\frac{\mu+v}{\mu-v}\right|$.
Note that the squeezed state $|0 ; \mu, \nu\rangle$ is defined as the unit eigenvector of $b_{\mu, \nu}=\mu a_{b}+v a_{b}^{*}$ associated with the eigenvalue 0 .
Proof. We need the discussion of remarks 3 and 5 for the proof. It is necessary to discuss the action of the group. From this representation of the algebra $\mathfrak{s u}(1,1)$, we can construct the representation of the double-covering group $\widetilde{S(1,1)}$ of the group $S U(1,1)$. In general, we can construct the representation of $\widehat{S U(1,1)}$ in the case where $\lambda$ is a half-integer. Now, we let $\pi_{\widetilde{S U(1,1)}}$ be the projection from $\widehat{S U(1,1)}$ to $S U(1,1)$. From (41), we have

$$
\begin{array}{lll}
\mathrm{e}^{t E_{0}} Q \mathrm{e}^{-t E_{0}}=\mathrm{e}^{t} Q & \mathrm{e}^{t E_{+}} Q \mathrm{e}^{-t E_{+}}=Q & \mathrm{e}^{t E_{-}} Q \mathrm{e}^{-t E_{-}}=Q-P t \\
\mathrm{e}^{t E_{+}} P \mathrm{e}^{-t E_{+}}=P-Q t & \mathrm{e}^{t E_{0}} P \mathrm{e}^{-t E_{0}}=\mathrm{e}^{-t} P & \mathrm{e}^{t E_{0}} P \mathrm{e}^{-t E_{0}}=P .
\end{array}
$$

From (31), (32) and some calculations, we have

$$
\begin{array}{ll}
V(g) Q V(g)^{*}=\left(\mu_{1}+v_{1}\right) Q+\left(v_{2}-\mu_{2}\right) P \\
V(g) P V(g)^{*}=\left(v_{2}+\mu_{2}\right) Q+\left(\mu_{1}-v_{1}\right) P & \forall g \in S \widetilde{U(1,1)}
\end{array}
$$

where the complex numbers $\mu_{i} \circ \pi_{\widetilde{S U(1,1)}}(g)$ and $\nu_{i} \circ \pi_{\widetilde{S U(1,1)}}(g)$ with the functions $\mu_{i}$ and $v_{i}$ defined at the beginning of remark 3 are denoted simply by $\mu_{i}$ and $v_{i}$, respectively, in a similar manner to the previous section. Thus, we have

$$
\begin{aligned}
& V(g) a_{b} V(g)^{*}=\mu a_{b}+v a_{b}^{*} \\
& V(g) a_{b}^{*} V(g)^{*}=v^{*} a_{b}+\mu^{*} a_{b}^{*} \quad \forall g \in \widetilde{S(1,1)}
\end{aligned}
$$

where we simplify $\mu \circ \pi_{\widetilde{S U(1,1)}}(g)$ and $\nu \circ \pi_{\widetilde{S U(1,1)}}(g)$ by $\mu$ and $\nu$, respectively, similarly.
Since $L_{-}=\frac{1}{2}\left(a_{b}\right)^{2}$, the lowest-weight vector $|0\rangle_{a}$ is the boson vacuum vector $|0 ; 1,0\rangle$. The squeezed state $|0 ; \mu, \nu\rangle$ satisfies $\left(\mu a_{b}+\nu a_{b}^{*}\right)|0 ; \mu, \nu\rangle=0$. Assume that $\mu \circ \pi_{\widetilde{S U(1,1)}}(g)=$ $\mu, \nu \circ \pi_{\widetilde{S U(1,1)}}(g)=v$. Then we have $V(g) a_{b} V(g)^{*}|0 ; \mu, \nu\rangle=0$. Hence, we see that the vector $V(g)^{*}|0 ; \mu, \nu\rangle$ equals a scalar times a vacuum vector $|0 ; 1,0\rangle=|0\rangle_{a}$. From these facts and (34), we obtain (43).

From (43), we find the correspondence to the characteristic equations (1) and (2) of squeezed states explained in section 1 . Substituting (43) into (24) and (27), we obtain (2) and (1). In the following, the vector $|\zeta\rangle_{a}$ in $L_{\text {even }}^{2}(\mathbb{R})$ is denoted by $|\zeta\rangle_{a, \text { even }}$. Equations (23) and (43) imply that the squeezed state $|0 ; \mu, \nu\rangle$ equals a scalar times

$$
\exp \left(-\frac{\xi}{2}\left(a_{b}^{*}\right)^{2}+\frac{\xi^{*}}{2}\left(a_{b}\right)^{2}\right)|0 ; 1,0\rangle
$$

corresponding to Caves's notation [30] for squeezed states, where

$$
\xi:=\frac{1}{2} \exp \left[i \arg \frac{\nu}{\mu}\right] \ln \frac{|\mu|+|\nu|}{|\mu|-|\nu|} .
$$

On the other hand, in $L_{\text {odd }}^{2}(\mathbb{R})$, the operators $a, A$ and $N$ are written in the forms
$a=-a_{b}\left(a_{b}^{*}\right)^{-1} \quad A=P Q^{-1} \quad N=\frac{1}{2}\left(n_{b}-1\right) \quad|n\rangle_{N}=(-1)^{n}|2 n+1\rangle_{n_{b}}$.
Next, we will discuss the representation of the algebra $\mathfrak{s u}(1,1)$ in the Hilbert space $L^{2}\left(\mathbb{R}^{n}\right)=\underbrace{L^{2}(\mathbb{R}) \otimes \cdots \otimes L^{2}(\mathbb{R})}_{n}$, for multi-particle systems. In this representation,
$E_{0}=\frac{\mathrm{i}}{2} \sum_{j=1}^{n}\left(P_{j} Q_{j}+Q_{j} P_{j}\right) \quad E_{+}=\frac{\mathrm{i}}{2} \sum_{j=1}^{n} Q_{j}^{2} \quad E_{-}=-\frac{\mathrm{i}}{2} \sum_{j=1}^{n} P_{j}^{2}$
hold, where $Q_{j}$ and $P_{j}$ denote the multiplication operator and the $(-\mathrm{i})$-times differential operator, respectively, with respect to the $j$ th variable. Let $L_{\mathrm{e}}^{2}\left(\mathbb{R}^{n}\right)$ be the closure of the linear space generated by

$$
\{|\zeta\rangle_{a, \text { even }}^{\otimes n}:=\underbrace{|\zeta\rangle_{a, \text { even }} \otimes \cdots \otimes|\zeta\rangle_{a, \text { even }}}_{n}\} .
$$

Then, the Hilbert space $L_{\mathrm{e}}^{2}\left(\mathbb{R}^{n}\right)$ is irreducible under the representation (44) of the algebra $\mathfrak{s u}(1,1)$, and then we have $L_{\mathrm{e}}^{2}\left(\mathbb{R}^{n}\right)=\left\{f \in L^{2}\left(\mathbb{R}^{n}\right) \mid f\right.$ is a function of $\left.\sum_{j=1}^{n} x_{j}^{2}\right\}$, and then the vector $|\zeta\rangle_{a}$ in this representation on $L_{\mathrm{e}}^{2}\left(\mathbb{R}^{n}\right)$ is equivalent to $|\zeta\rangle_{a \text {,even }}^{\otimes n}$.

Letting $A_{n, \mathrm{e}}$ denote the operator $A$ in this representation, we obtain the relation

$$
A_{n, \mathrm{e}}=\left(\sum_{j=1}^{n} Q_{j}^{2}\right)^{-1} \sum_{j=1}^{n} Q_{j} P_{j}=-\mathrm{i}\left(\sum_{j=1}^{n} \frac{2 x_{j}}{r} \frac{\partial}{\partial x_{j}}\right)
$$

with

$$
r:=2 \sum_{j=1}^{n} x_{j}^{2}
$$

Now define the unitary map

$$
U_{n}: L^{2}\left(\mathbb{R}^{n}\right) \rightarrow L^{2}\left(\mathbb{R}^{+}\right) \otimes L^{2}\left(S^{n-1}\right) \cong L^{2}\left(\mathbb{R}^{+} \times S^{n-1}\right)
$$

by

$$
\left(U_{n}(f)\right)\left(r,\left(e_{1}, e_{2}, \ldots, e_{n}\right)\right)=r^{\frac{n-2}{4}} f\left(\sqrt{\frac{r}{2}} e_{1}, \sqrt{\frac{r}{2}} e_{2}, \ldots, \sqrt{\frac{r}{2}} e_{n}\right)
$$

where $S^{n-1}$ denotes the ( $n-1$ )-dimensional spherical surface and $\left(e_{1}, e_{2}, \ldots, e_{n}\right)$ is an element of $S^{n-1}$. Then, the following relations hold:

$$
\begin{aligned}
& U_{n} E_{0} U_{n}^{*}=E_{0, \frac{n-2}{4}} \otimes I \quad U_{n} E_{+} U_{n}^{*}=E_{+, \frac{n-2}{4}} \otimes I \quad U_{n} E_{-} U_{n}^{*}=E_{-, \frac{n-2}{4}} \otimes I \\
& U_{n} A_{n, \mathrm{e}} U_{n}^{*}=\left(P+\mathrm{i}\left(\frac{n}{4}-\frac{1}{2}\right) Q^{-1}\right) \otimes I=-\mathrm{i} \frac{\partial}{\partial r}+\mathrm{i}\left(\frac{n}{4}-\frac{1}{2}\right) \frac{1}{r} \\
& U_{n} L_{\mathrm{e}}^{2}\left(\mathbb{R}^{n}\right)=L^{2}\left(\mathbb{R}^{+}\right) \otimes \psi_{n} \\
& U_{n} \mathcal{D}_{o}\left(A_{n, \mathrm{e}}\right) \cap\left(C^{1}\left(\mathbb{R}^{+}\right) \otimes \psi_{n}\right) \\
& \quad=\left\{\begin{array}{ll}
x^{\frac{n}{4}-\frac{1}{2}} f(x) \in L^{2}\left(\mathbb{R}^{+}\right) \cap C^{1}\left(\mathbb{R}^{+}\right) \left\lvert\, \begin{array}{l}
x^{\frac{n}{4}-\frac{1}{2}} f^{\prime}(x) \in L^{2}\left(\mathbb{R}^{+}\right) \\
f(s)<\infty \quad \text { as } s \rightarrow 0 .
\end{array}\right.
\end{array}\right\} \otimes \psi_{n}
\end{aligned}
$$

where $\psi_{n}$ denotes the constant function on $S^{n-1}$ such that $\left\|\psi_{n}\right\|=1$. The compound-systemtype normal extension of $A_{n, \mathrm{e}}$ in the above relations is reduced to the discussion of $A_{\frac{n}{4}-\frac{1}{2}}$ which will be treated in sections 6.1 and 6.2.

## 6. Construction of a compound-system-type normal extension of $A^{*}$

### 6.1. The case where $\lambda=1$

In this subsection, we will construct a compound-system-type normal extension of $A^{*}$ when $\lambda=1$. Let $\{|\uparrow\rangle,|\downarrow\rangle\}$ be a CONS of $\mathbb{C}^{2}$. From lemma 28 and the fact that $A^{*}$ is maximally symmetric, we obtain the following theorem.
Theorem 43. Define the operator $T:=A \otimes|-\rangle\langle+|+A^{*} \otimes|+\rangle\langle-|$ on the domain $\mathcal{D}_{o}(T):=$ $\mathcal{D}_{o}(A) \otimes|+\rangle \oplus \mathcal{D}_{o}\left(A^{*}\right) \otimes|-\rangle$ with $| \pm\rangle:=\frac{1}{\sqrt{2}}(|\uparrow\rangle \pm|\downarrow\rangle)$. The operator $T$ is a self-adjoint operator. Moreover, the triple $\left(\mathbb{C}^{2}, T,|\uparrow\rangle\right)$ is a compound-system-type normal extension of $A^{*}$.

Similarly, we can construct a compound-system-type normal extension of $a^{*}$ according to lemma 27. The spectrum of the compound-system-type normal extension of $A^{*}$ for $\lambda=1$ appears only on the real axis. That of the compound-system-type normal extension of $a^{*}$ appears only on the unit circle.

### 6.2. The cases where $\lambda>1$

In the following, we will discuss the cases when $\lambda>1$. Let $\{|\uparrow\rangle,|\downarrow\rangle\}$ be CONS of $\mathbb{C}^{2}$. We obtain the following theorem, with $A_{0}\left(A_{k}\right.$ with $\left.k=0\right)$ discussed at the end of section 5.1.
Theorem 44. The pair of $E_{+} \otimes I$ and $E_{0} \otimes I+I \otimes E_{0}$ on $\mathcal{H}_{\lambda} \otimes \mathcal{H}_{\lambda-1}$ satisfies the commutation relation of the generators of the Affine group. This representation of the Affine group is written as follows: there exists a Hilbert space $\mathcal{H}^{\prime}$ and a unitary map $U$ from $\mathcal{H}_{\lambda} \otimes \mathcal{H}_{\lambda-1}$ to $\mathcal{H}^{\prime} \otimes L^{2}\left(\mathbb{R}^{+}\right)$such that $U\left(E_{+} \otimes I\right) U^{*}=I \otimes E_{+}, U\left(E_{0} \otimes I+I \otimes E_{0}\right) U^{*}=I \otimes E_{0}$. Then, the operator $U^{*}\left(I \otimes A_{0}\right) U \otimes|-\rangle\langle+|+U^{*}\left(I \otimes A_{0}^{*}\right) U \otimes|+\rangle\langle-|$ with the domain $\mathcal{D}_{o}\left(U^{*}\left(I \otimes A_{0}\right) U\right) \otimes|+\rangle \oplus \mathcal{D}_{o}\left(U^{*}\left(I \otimes A_{0}\right) U\right) \otimes|-\rangle$ is self-adjoint.

Moreover, the operator $T:=U^{*}\left(I \otimes A_{0}\right) U \otimes|-\rangle\langle+|+U^{*}\left(I \otimes A_{0}^{*}\right) U \otimes|+\rangle\langle-|-\mathrm{i} E_{+}^{-1} \otimes$ $E_{+} \otimes I$ with the domain $\mathcal{D}_{o}(T):=\left(\mathcal{D}_{o}\left(U^{*}\left(I \otimes A_{0}\right) U\right) \otimes|+\rangle \oplus \mathcal{D}_{o}\left(U^{*}\left(I \otimes A_{0}^{*}\right) U\right) \otimes|-\rangle\right) \cap$ $\mathcal{D}_{o}\left(E_{+}^{-1} \otimes E_{+}\right) \otimes \mathbb{C}^{2}$ is normal. The triple $\left(\mathcal{H}_{\lambda}^{\prime}:=\mathcal{H}_{\lambda-1} \otimes \mathbb{C}^{2}, T, \psi:=|0\rangle_{N} \otimes|\uparrow\rangle\right)$ is a compound-system-type normal extension of $A^{*}$.

Proof. We need the discussion of remark 6 for the proof. It is sufficient to prove them under the representations given in section 5.1 because of theorem 41 . Now define the unitary operator $U$ on $L^{2}\left(\mathbb{R}^{+}\right) \otimes L^{2}\left(\mathbb{R}^{+}\right)$by $(U(f))(u, v)=\sqrt{v} f(v, u v)$. Then we have $U\left(E_{+} \otimes I\right) U^{*}=I \otimes E_{+}$, $U\left(E_{0} \otimes I+I \otimes E_{0}\right) U^{*}=I \otimes E_{0}$ and $U\left(-\mathrm{i} E_{+}^{-1} \otimes E_{+}\right) U^{*}=-E^{+} \otimes I$. Because the discussion at the end of section 5.1 shows that $A_{0}$ is closed and symmetric, it follows from the proof of lemma 28 that the operator $A_{0} \otimes|-\rangle\langle+|+A_{0}^{*} \otimes|+\rangle\langle-|$ is self-adjoint and its domain is $\mathcal{D}_{o}\left(A^{*}\right) \otimes|-\rangle \oplus \mathcal{D}_{o}(A) \otimes|+\rangle$. In general, for a self-adjoint operator $X$ on $\mathcal{K}_{1}$ and a skew-adjoint operator $Y$ on $\mathcal{K}_{2}$, we can show that the operator $X \otimes I+I \otimes Y$ with the domain $\mathcal{D}_{o}(X) \otimes \mathcal{D}_{o}(Y)=\mathcal{D}_{o}(X) \otimes \mathcal{K}_{2} \cap \mathcal{K}_{1} \otimes \mathcal{D}_{o}(Y) \subset \mathcal{K}_{1} \otimes \mathcal{K}_{2}$ is normal. Then, the operator $T^{\prime}:=I \otimes\left(A_{0} \otimes|-\rangle\langle+| \oplus A_{0}^{*} \otimes|+\rangle\langle-|\right)-E^{+} \otimes I \otimes I$ with the domain $\mathcal{D}_{o}\left(T^{\prime}\right):=\left(\mathcal{D}_{o}\left(I \otimes A_{0}\right) \otimes|+\rangle \oplus \mathcal{D}_{o}\left(I \otimes A_{0}^{*}\right) \otimes|-\rangle\right) \cap \mathcal{D}_{o}\left(E_{+}\right) \otimes L^{2}\left(\mathbb{R}^{+}\right) \otimes \mathbb{C}^{2}$ is normal. Thus, we have proved that the operator $T\left(=U^{*} T^{\prime} U\right)$ is normal. Now, we will prove that the triple $\left(\mathcal{H}_{\lambda}^{\prime}, T, \psi=|0\rangle_{N}^{k-\frac{1}{2}} \otimes|\uparrow\rangle\right)$ is a compound-system-type normal extension of $A_{k}^{*}$.

Since the set $\mathcal{D}_{o}\left(A_{k}^{*}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)$is a core of the operator $A_{k}^{*}$, it is sufficient to show that $\left(A_{k}^{*} \phi\right) \otimes|0\rangle_{N}^{k-\frac{1}{2}} \otimes|\uparrow\rangle=T\left(\phi \otimes|0\rangle_{N}^{k-\frac{1}{2}} \otimes|\uparrow\rangle\right)$ for any $\phi \in \mathcal{D}_{o}\left(A_{k}^{*}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)$.

From the definitions and (39), some calculations result in

$$
\begin{aligned}
U\left(\left(\mathcal{D}_{o}\left(A_{k}^{*}\right)\right.\right. & \left.\left.\cap C^{1}\left(\mathbb{R}^{+}\right)\right) \otimes|0\rangle_{N}^{k-\frac{1}{2}}\right) \\
& =\left\{f(v) u^{k-1 / 2} \mathrm{e}^{-u v} \in L^{2}\left(\mathbb{R}^{+} \times \mathbb{R}^{+}\right) \mid x^{-k} f^{\prime}(x) \in L^{2}\left(\mathbb{R}^{+}\right), f(s) \rightarrow 0 \text { as } s \rightarrow 0\right\}
\end{aligned}
$$

We can show that a function $u \mapsto u^{k-1 / 2} \mathrm{e}^{-u v}$ is contained by $\mathcal{D}_{o}\left(E_{+}\right) \subset L^{2}\left(\mathbb{R}^{+}\right)$for any $v \in \mathbb{R}^{+}$. If a function $f$ satisfies the condition $x^{-k} f^{\prime}(x) \in L^{2}\left(\mathbb{R}^{+}\right), f(s) \rightarrow 0$ as $s \rightarrow 0$, then a function $v \mapsto f(v) u^{k-1 / 2} \mathrm{e}^{-u v}$ is contained by $\mathcal{D}_{o}\left(A_{0}^{*}\right) \subset L^{2}\left(\mathbb{R}^{+}\right)$for any $u \in \mathbb{R}^{+}$.

Then, the set

$$
U\left(\left(\mathcal{D}_{o}\left(A_{k}^{*}\right) \cap C^{1}\left(\mathbb{R}^{+}\right)\right) \otimes|0\rangle_{N}^{k-\frac{1}{2}}\right)
$$

is included in the set

$$
\mathcal{D}_{o}\left(I \otimes A_{0}^{*}\right) \cap \mathcal{D}_{o}\left(E_{+} \otimes I\right) \cap\left(C^{1}\left(\mathbb{R}^{+}\right) \otimes C^{1}\left(\mathbb{R}^{+}\right)\right)
$$

Hence,

$$
\begin{aligned}
U\left(\left(\mathcal{D}_{o}\left(A_{k}^{*}\right) \cap\right.\right. & \left.\left.C^{1}\left(\mathbb{R}^{+}\right)\right) \otimes|0\rangle_{N}^{k-\frac{1}{2}}\right) \otimes|\uparrow\rangle \\
\subset & \mathcal{D}_{o}\left(I \otimes A_{0}^{*} \otimes I\right) \cap \mathcal{D}_{o}\left(E_{+} \otimes I \otimes I\right) \cap\left(C^{1}\left(\mathbb{R}^{+}\right) \otimes C^{1}\left(\mathbb{R}^{+}\right) \otimes|\uparrow\rangle\right) \\
\subset & \mathcal{D}_{o}\left(I \otimes\left(A_{0}^{*} \otimes|+\rangle\langle-|+A_{0} \otimes|-\rangle\langle+|\right)\right) \cap \mathcal{D}_{o}\left(E_{+} \otimes I \otimes I\right) \\
& \cap\left(C^{1}\left(\mathbb{R}^{+}\right) \otimes C^{1}\left(\mathbb{R}^{+}\right) \otimes|\uparrow\rangle\right) \\
= & \mathcal{D}_{o}\left(T^{\prime}\right) \cap\left(C^{1}\left(\mathbb{R}^{+}\right) \otimes C^{1}\left(\mathbb{R}^{+}\right) \otimes|\uparrow\rangle\right)
\end{aligned}
$$

Thus, for the function $f(x)$ satisfying $\phi(x)=f(x) x^{-k}$, we obtain
$T\left(\phi \otimes|0\rangle_{N}^{k-\frac{1}{2}} \otimes|\uparrow\rangle\right)$

$$
\begin{aligned}
& =-\mathrm{i} \frac{\mathrm{~d}}{\mathrm{~d} v}\left(f(v) u^{k-\frac{1}{2}} \mathrm{e}^{-u v}\right) \otimes(|+\rangle\langle-|+|-\rangle\langle+|)|\uparrow\rangle-\mathrm{i} f(v) u^{k+\frac{1}{2}} \mathrm{e}^{-u v} \otimes|\uparrow\rangle \\
& =-\mathrm{i} \frac{\mathrm{~d} f}{\mathrm{~d} v}(v) u^{k-\frac{1}{2}} \mathrm{e}^{-u v} \otimes|\uparrow\rangle \\
& =\left(A_{k} \phi\right) \otimes|0\rangle_{N}^{k-\frac{1}{2}} \otimes|\uparrow\rangle .
\end{aligned}
$$

The theorem is now immediate.
In the above discussions, it is sufficient only to choose $\mathcal{H}_{\lambda-1}$ instead of $\mathcal{H}_{\lambda}^{\prime}$ in order only to show that the operator $T$ formally satisfies $\left[T, T^{*}\right]=0$ and formally satisfies (10). However, the above definition of $\mathcal{H}_{\lambda}^{\prime}$ is required in order that $T$ may be a normal operator defined in definition 1.

Since the spectrum of the compound-system-type normal extension of $A^{*}$ for $\lambda=1$ appears only in the upper half-plane including the real axis, the spectrum of the compound-system-type normal extension of $a^{*}$ appears only on the unit disc (including the unit circle) if the latter is related to the former by the adjoint of the Cayley transform.

## 7. Conclusions

We have discussed subnormal operators as a class of generalized observables. A POVM of a subnormal operator defined in definition 13 has little information about its implementation. However, in order to describe not only the probability distributions characterized by the POVMs but also a framework of their implementations, we have defined compound-system-type normal extensions in section 3. (The heterodyne measurement known in quantum optics is interpreted as a special case of compound-system-type normal extensions.) In these contexts, we have constructed the compound-system-type normal extensions of two subnormal operators $a^{*}$ and $A^{*}$ canonically introduced from an irreducible unitary representation of $\mathfrak{s u}(1,1)$, when the minimum eigenvalue $\lambda$ of the generator $L_{0}$ is not less than one. The squeezed states are regarded as the coherent states of the algebra $\mathfrak{s u}(1,1)$, and have been characterized as the eigenvectors of an operator defined in this mathematical framework. The squeezed states in two-particle or multi-particle systems have been interpreted as the eigenvectors of the adjoints $a$ and $A$ of the subnormal operators $a^{*}$ and $A^{*}$. The coherent states of the affine group have been interpreted within the same framework, as well. The squeezed states in a one-particle system have been interpreted as the eigenvectors of the operator $a$ and $A$, though the operators $a^{*}$ and $A^{*}$ are not subnormal and their compound-system-type normal extensions do not exist in this case because $\lambda$ is less than one in this case.

The information described by a compound-system-type normal extension is not enough to completely specify the experimental implementation, where the measurement of the normal operator on the compound system is performed by the measurement on each system after some interactions were made between the basic system and the ancillary system. Therefore, the formulation including this specification is a future problem. As another possibility, since the affine group is closely related to the Poincaré group, our results concerning the affine group may be applicable to relativistic quantum mechanics.

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## Appendix

The following lemma about spectra is well known. (See Hiai and Yanagi [15].) In Hiai and Yanagi [15], it is proved in the case of bounded operators. However, it can be easily extended to the case of unbounded operators.

Lemma 45. For a densely defined operator $A$ on $\mathcal{H}$, Let $\sigma_{p}(A), \sigma_{c}(A)$ and $\sigma_{r}(A)$, be the point spectrum, the continuous spectrum and the residual spectrum, respectively. Then we have the following relations:

- $\lambda \in \sigma_{r}(A) \Rightarrow \lambda^{*} \in \sigma_{p}\left(A^{*}\right)$
- $\lambda \in \sigma_{p}(A) \Rightarrow \lambda^{*} \in \sigma_{r}\left(A^{*}\right) \cup \sigma_{p}\left(A^{*}\right)$
- $\lambda \in \sigma_{c}(A) \Rightarrow \lambda^{*} \in \sigma_{c}\left(A^{*}\right)$.


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[^0]:    $\dagger$ This class was introduced by Halmos [2].

[^1]:    $\dagger$ A group is called a universal covering group if it is connected and if its homotopy group is trivial.

[^2]:    $\dagger$ Sometimes another definition with $n+l$ instead of $l$ is used.

