# SUBMANIFOLDS OF A KÄHLERIAN MANIFOLD 

By Mitsue Ako<br>Dedicated to Professor Hitoshi Hombu on his sixtieth birthday

1. Introduction. The following theorem is well known:

Theorem A (Cf. Yano [3]). A holomorphic submanifold of a Kählerian manifold is minimal.

Thus it would be natural to ask whether a minimal submanifold of a Kählerian manifold is holomorphic. As to a very special case of a totally geodesic submanifold we have shown in [1] the following

Theorem B. A totally geodesic submanifold $S$ in a $2 n$-dimensional Fubinian manifold $M$ is holomorphic if we assume that $n \neq 2$ and the codimension of $S$ is 2 . In the exceptional case $n=2, S$ is a holomorphic or an anti-holomorphic submanifold of $M$.

A submanifold $S$ is said to be anti-holomorphic at a point $p \in S$, if $T_{p}(S)$ and $N_{p}(S)$ are transformed under $F$ into each other, where $F$ is an almost complex structure of $M, T_{p}(S)$ and $N_{p}(S)$ denoting respectively the tangent space to $S$ at $p$ and the normal space to $S$ at $p . S$ is called an anti-holomorphic submanifold if it is anti-holomorphic at each point of $S$.

Theorem B shows that the converse of Theorem A is not true in general. Now we shall study, in this paper, submanifolds, especially minimal ones, in a Kählerian manifold. The notations and terminologies are found in [1], but we state some of them at the beginning of the next section for the later use.
2. Let $M$ be a Kählerian manifold of real dimension $2 n$ and $S^{1)}$ a connected orientable submanifold of $M$ whose real dimension is $2 n-2$. It is well known that a Riemannian metric $g$ on $S$ can be induced from the Riemannian metric $G$ of $M$. We denote by $\langle,\rangle_{M}$ the inner product with respect to $G$ and by $\langle,\rangle_{S}$ the inner product with respect to $g$. We now put, for a tangent vector $X$ on $S$,

$$
\begin{equation*}
F\left(\xi_{*} X\right)=T(X)+N(X), \tag{2.1}
\end{equation*}
$$

[^0]where $F$ is the almost complex structure of $M, T(X)$ denotes the tangential part and $N(X)$ the normal part, both of $F\left(\xi_{*} X\right)$. Since $T(X)$ is tangent to $S$, we may put
$$
\left\langle T(X), \xi_{*} Y\right\rangle_{M}=\langle A X, Y\rangle_{S}
$$
where $A$ is a tensor on $S$ of type $(1,1)$ and $Y$ is an arbitrary vector on $S$. If we define a 2 -form $\tilde{A}$ by
$$
\tilde{A}(X, Y)=\langle A X, Y\rangle_{S}
$$
for any pair of vector fields $X$ and $Y$ on $S$, then we have, denoting by $\tilde{F}$ the fundamental 2 -form of $M$,
\[

$$
\begin{equation*}
\tilde{A}(X, Y)=\tilde{F}\left(\xi_{*} X, \xi_{*} Y\right) \tag{2.2}
\end{equation*}
$$

\]

(2.2) shows that $\tilde{A}$ is a skew-symmetric bilinear form.

Now, we restrict ourselves to a sufficiently small neighborhood $\mathcal{U}$ in which there exist two fields of unit normal vectors to $S$. First, we fix in $U$ two unit normal vector fields $C$ and $D$ to $S$ which are mutually orthogonal. Then $N(X)$ defined by (2.1) can be expressed in $U$ as

$$
\begin{equation*}
N(X)=\tilde{\alpha}(X) C+\tilde{\beta}(X) D \tag{2.3}
\end{equation*}
$$

where $\tilde{\alpha}$ and $\tilde{\beta}$ are 1 -forms on $S$. We have

$$
\begin{equation*}
\tilde{\alpha}(X)=\tilde{F}\left(\xi_{*} X, C\right) \tag{2.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\tilde{\beta}(X)=\tilde{F}\left(\xi_{*} X, D\right) \tag{2.5}
\end{equation*}
$$

for any vector fields $X$ on $S$. We define $\|\alpha\|$ and $\|\beta\|$ respectively by

$$
\|\alpha\|=\sqrt{\langle\alpha, \alpha\rangle_{S}} \quad \text { and } \quad\|\beta\|=\sqrt{\langle\beta, \beta\rangle_{S}}
$$

where $\alpha$ and $\beta$ are contravariant tensors of degree 1 defined by $\langle\alpha, X\rangle_{S}=\tilde{\alpha}(X)$ and $\langle\beta, X\rangle_{s}=\tilde{\beta}(X)$ respectively. Then we have, by a direct computation,

$$
\begin{equation*}
\|\alpha\|^{2}=\|\beta\|^{2}=1-\varphi^{2}, \tag{2.6}
\end{equation*}
$$

where

$$
\begin{equation*}
\varphi=\tilde{F}(C, D) . \tag{2.7}
\end{equation*}
$$

$\tilde{F}(C, D)$ seems to depend upon the choice of the pair of unit normal vector fields $C$ and $D$ in $U$, but it is not hard to show that $\tilde{F}(C, D)$ is independent of the choice of $C$ and $D$. That is to say, $\tilde{F}(C, D)$ is left invariant under any orthogonal transformation applied to $C$ and $D$, since $S$ is assumed to be orientable. Thus $\varphi$ is a globally defined function on $S$. On the other hand, (2.6) implies that if $\tilde{\alpha}(X)=0$ at a point $p$ for any vector $X$ on $S$, then $\tilde{\beta}(X)=0$ at $p$ and vice versa. Straightforward computation shows that

$$
\begin{equation*}
\langle\alpha, \beta\rangle_{S}=0 \tag{2.8}
\end{equation*}
$$

and

$$
\begin{equation*}
A^{2}=-I+\tilde{\alpha} \otimes \alpha+\tilde{\beta} \otimes \beta, \tag{2.9}
\end{equation*}
$$

where $I$ is the unit tensor.
The maximal holomorphic subspace $H_{p}$ of the tangent space $T_{p}(S)$ to $S$ at $p$ is defined by

$$
H_{p}=\left\{V \in T_{p}(S) \mid F\left(\xi_{*} V\right) \in T_{p}(S)\right\}
$$

and the anti-holomorphic subspace $H_{p}^{\prime}$ of $T_{p}(S)$ is defined by

$$
H_{p}^{\prime}=\left\{W \in T_{p}(S) \mid F\left(\xi_{*} W\right) \in N_{p}(S)\right\},
$$

where $N_{p}(S)$ denotes the normal space to $S$ at $p$. These definitions show that

$$
H_{p} \oplus H_{p}^{\prime}=T_{p}(S) \quad \text { (direct sum) }
$$

and $H_{p}$ and $H_{p}^{\prime}$ are mutually orthogonal. In fact

$$
\langle V, W\rangle_{S}=\left\langle\xi_{*} V, \xi_{*} W\right\rangle_{M}=\left\langle F\left(\xi_{*} V\right), F\left(\xi_{*} W\right)\right\rangle_{M}=0,
$$

if $V \in H_{p}$ and $W \epsilon H_{p}^{\prime}$. If we restrict ourselves to $U$ and if we take account of (2.3), then we see that a necessary and sufficient condition that $V$ belongs to $H_{p}$ is expressed as

$$
\begin{equation*}
\langle V, \alpha\rangle_{S}=0 \quad \text { and } \quad\langle V, \beta\rangle_{S}=0 . \tag{2.10}
\end{equation*}
$$

The identity (2.9) and the equations (2.10) show that $A$ restricted to $H_{p}$ is an almost complex structure. We note again from (2.10) that $H_{p}^{\prime}$ is spanned by $\alpha$ and $\beta$ at $p$ at which $\|\alpha\| \neq 0$. Thus we have

$$
\operatorname{dim} H_{p} \geqq \operatorname{dim} S-2
$$

and the equality holds at $p$ at which we have $\|\alpha\|=0$.
The next proposition is a result of a direct computation.
Proposition 2.1. If $S$ is a totally geodesic submanifold of $M$, then the function $\varphi$ defined by (2.7) is constant and therefore $\operatorname{dim} H_{p}$ is constant on $S$.

We shall assume, from now on, that there is at least one point $p$ at which $\|\alpha\|$ does not vanish.

An assignment of $H_{p}$ to each $p$ of $S$ defines a distribution $D$, if $\operatorname{dim} H_{p}$ is constant. Let $X$ and $Y$ be any local vector fields which belong to $H_{p}$ in a sufficiently small neighborhood $\odot \mathcal{V}$. It is well known that the distribution $D$ is completely integrable if and only if $[X, Y]$ is also a local vector field belonging to $H_{p}$ in $\widetilde{V}$. This condition is equivalent to

$$
\begin{equation*}
\left.\langle[X, Y], \alpha\rangle_{S}=0 \quad \text { and } \quad\langle[X, Y], \beta\rangle_{S}=0 .{ }^{2}\right) \tag{2.11}
\end{equation*}
$$

The equations (2.11) can be written as

$$
\begin{align*}
& \left\langle\nabla_{X} \alpha, Y\right\rangle_{S}-\left\langle\nabla_{Y} \alpha, X\right\rangle_{S}=0,  \tag{2.12}\\
& \left\langle\nabla_{X} \beta, Y\right\rangle_{S}-\left\langle\nabla_{Y} \beta, X\right\rangle_{S}=0
\end{align*}
$$

by virtue of

$$
\langle X, \alpha\rangle_{S}=\langle X, \beta\rangle_{S}=\langle Y, \alpha\rangle_{S}=\langle Y, \beta\rangle_{S}=0,
$$

where $\nabla$ denotes the covariant differentiation along $S$ with respect to the connection induced on $S$ from the Riemannian connection of $M$. On the other hand we have, from (2.4) and (2.5),

$$
\begin{equation*}
\left\langle\nabla_{X} \alpha, Y\right\rangle_{S}=-\varphi \tilde{k}(X, Y)-\tilde{h}(A Y, X)+\tilde{l}(X) \tilde{\beta}(Y) \tag{2.13}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle\nabla_{X} \beta, Y\right\rangle_{S}=\varphi \tilde{h}(X, Y)-\tilde{k}(A Y, X)-\tilde{l}(X) \tilde{\alpha}(Y), \tag{2.14}
\end{equation*}
$$

where $\tilde{h}$ and $\tilde{k}$ are the second fundamental forms of $S$ and $\tilde{l}$ the third fundamental form of $S, X$ and $Y$ being arbitrary vector fields on $S$. Thus the equations (2.12) become

$$
\left\{\begin{array}{l}
\langle(A h+h A) X, Y\rangle_{S}=0,  \tag{2.15}\\
\langle(A k+k A) X, Y\rangle_{S}=0,
\end{array}\right.
$$

$h$ and $k$ being tensors on $S$ of type $(1,1)$ defined respectively by $\langle h X, Y\rangle_{S}$ $=\tilde{h}(X, Y)$ and $\langle k X, Y\rangle_{S}=\tilde{k}(X, Y)$. Thus we have

Proposition 2. 2. Suppose that $\operatorname{dim} H_{p}$ is constant on $S$. In order that the distribution $D: p \rightarrow H_{p}$ is completely integrable, it is necessary and sufficient that the equations (2.15) are valid for arbitrary vectors $X$ and $Y$ on $H_{p}$.

This proposition, together with Proposition 2.1, gives
Corollary 2.1. For a totally geodesic submanifold $S$ of $M$, the distribution $D: p \rightarrow H_{p}$ is always completely integrable.

In the case in which the distribution $D: p \rightarrow H_{p}$ is completely integrable, the integral manifold $H$ of the distribution $D$ is a minimal submanifold of $M$ (Theorem A). We denote submanifold maps by $\eta: H \rightarrow S$ and $\zeta: H \rightarrow M$ and their differentials by $\eta_{*}$ and $\zeta_{*}$ respectively. We shall use, in the neighborhood $U, C, D, \xi_{*} \alpha /\|\alpha\|$ and $\xi_{* \beta} \beta\|\beta\|$ as unit normal vector fields to $H$ and denote them by $C_{1}, C_{2}, C_{3}$ and $C_{4}$ respectively. Then we have

[^1]$$
\left\langle C_{x}, C_{y}\right\rangle_{M}=\delta_{x y} \quad(x, y=1,2,3,4) .
$$

Now we introduce van der Waerden-Bortolotti derivative along a submanifold of $M$. Let $M^{\prime}$ be a submanifold of $M$ and $\sigma$ a submanifold map from $M^{\prime}$ into $M$ whose differential is denoted by $\sigma_{*}$. We denote by $T_{s}^{r}(M)$ (resp. $\left.T_{s}^{r}\left(M^{\prime}\right)\right)$ the space of all tensor fields of type $(r, s)$ and let $T(M)=\sum_{r, s} T_{s}^{r}(M)$ (resp. $T\left(M^{\prime}\right)=\sum_{r, s} T_{s}^{r}\left(M^{\prime}\right)$ ). Given an element $X \in T_{0}^{1}\left(M^{\prime}\right)$, we define a derivation $\nabla_{X}^{\sigma}$ in the formal tensor product $T(M) \# T\left(M^{\prime}\right)$ by the following properties:
1)

$$
\nabla_{X}^{\sigma} V=\nabla_{\sigma, X} V \quad \text { for } \quad V \in T(M),
$$

where $\bar{V}$ denotes the covariant derivation with respect to an affine connection of $M$;
2)

$$
\nabla_{X}^{\sigma} W=\left(\text { the tangential part of } \nabla_{\sigma, X}\left(\sigma_{*} W\right)\right),
$$

for $W \in T\left(M^{\prime}\right)$ and
3)

$$
\nabla_{X}^{o}(V \# W)=\left(\nabla_{X}^{o} V\right) \# W+V \#\left(\nabla_{x}^{o} W\right),
$$

for $V \in T(M)$ and $W \epsilon T\left(M^{\prime}\right)$. Van der Waerden-Bortolotti derivative $\nabla^{\sigma}$ along $M^{\prime}$ is defined as the assignment: $\left(X, W^{*}\right) \rightarrow \nabla_{X}^{\sigma} W^{*}$ for $X \in T_{0}^{1}\left(M^{\prime}\right)$ and $W^{*} \in T(M) \# T\left(M^{\prime}\right)$. For detail, see Yano-Ishihara [2].

Van der Waerden-Bortolotti derivative $\nabla^{5}$ along $H$ as a submanifold of $M$ gives

$$
\begin{equation*}
\left\langle\nabla_{V}^{\xi} C_{x}, \zeta_{*} W\right\rangle_{M}=-\left\langle h^{(x)} V, W\right\rangle_{H} \tag{2.16}
\end{equation*}
$$

where $V$ and $W$ are tangent to $H$, each $h^{(x)}$ is a tensor on $H$ of type $(1,1)$ and $\langle,\rangle_{H}$ denotes the inner product on $H$ with respect to the metric induced from that of $M . \quad C_{1}$ and $C_{2}$ are respectively transformed under $F$ as follows:

$$
\begin{equation*}
F C_{1}=\varphi C_{2}-\|\alpha\| C_{3} \tag{2.17}
\end{equation*}
$$

and

$$
\begin{equation*}
F C_{2}=-\varphi C_{1}-\|\alpha\| C_{4} . \tag{2.18}
\end{equation*}
$$

Substituting (2.17) into (2.16) we have

$$
\left\langle F V_{V}^{\xi} C_{1}, \zeta_{*} W\right\rangle_{M}=-\varphi\left\langle h^{(2)} V, W\right\rangle_{H}+\|\alpha\|\left\langle\left\langle h^{(3)} V, W\right\rangle_{H},\right.
$$

because of the fact that $F$ is covariant constant. Since $\zeta_{*} W$ is tangent to the holomorphic submanifold $H$ so is also $F\left(\zeta_{*} W\right)$ and therefore we can put

$$
\begin{equation*}
F\left(\zeta_{*} W\right)=\zeta_{*}(f W), \tag{2.19}
\end{equation*}
$$

where $f$ is a tensor of type $(1,1)$ on $H$. We can easily show that

$$
f^{2}=-I,
$$

$I$ being the unit tensor. Thus we have

$$
\begin{equation*}
\left\langle h^{(1)} V, f W\right\rangle_{H}=-\varphi\left\langle h^{(2)} V, W\right\rangle_{H}+\|\alpha\|\left\langle h^{(3)} V, W\right\rangle_{H}, \tag{2.20}
\end{equation*}
$$

by virture of the relation

$$
\left\langle F \nabla_{V}^{\varsigma} C_{1}, \zeta_{*} W\right\rangle_{M}=-\left\langle V_{V}^{\xi} C_{1}, F\left(\zeta_{*} W\right)\right\rangle_{M} .
$$

A similar method gives

$$
\begin{equation*}
\left\langle h^{(2)} V, f W\right\rangle_{H}=\varphi\left\langle h^{(1)} V, W\right\rangle_{I I}+\|\alpha\|\left\langle h^{(4)} V, W\right\rangle_{H}, \tag{2.21}
\end{equation*}
$$

because of (2.18).
On the other hand, if we consider $H$ as a submanifold of $S$ and we choose $\alpha /\|\alpha\|$ and $\beta /\|\beta\|$ as fields of unit normals to $H$, then we have

$$
\begin{equation*}
\left\langle\nabla_{V}^{\eta}(\alpha /\|\alpha\|), \eta_{*} W\right\rangle_{S}=-\left\langle h^{\prime} V, W\right\rangle_{H}, \tag{2.22}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle\nabla_{V}^{\eta}(\beta /\|\beta\|), \eta_{*} W\right\rangle_{S}=-\left\langle k^{\prime} V, W\right\rangle_{H}, \tag{2.22}
\end{equation*}
$$

where $V$ and $W$ are tangent to $H$ and $\nabla^{\eta}$ denotes van der Waerden-Bortolotti covariant derivation along $H$ as a submanifold of $S . h^{\prime}$ and $k^{\prime}$ are the so-called second fundamental tensors of $H$ in $S$. We can easily verify, by the definition of van der Waerden-Bortolotti covariant derivation that

$$
\begin{equation*}
h^{\prime}=h^{(8)} \quad \text { and } \quad k^{\prime}=h^{(4)} \tag{2.24}
\end{equation*}
$$

if we take account of (2.16). By a similar argument we have

$$
\begin{equation*}
\left\langle h^{(1)} W, V\right\rangle_{H}=\left\langle h\left(\eta_{*} W\right), \eta_{*} V\right\rangle_{S} \tag{2.25}
\end{equation*}
$$

and

$$
\begin{equation*}
\left\langle h^{(2)} W, V\right\rangle_{I I}=\left\langle k\left(\eta_{*} W\right), \eta_{*} V\right\rangle_{S} \tag{2.26}
\end{equation*}
$$

where $V$ and $W$ are vector fields on $H$. Since $H$ is holomorphic submanıfold of $M$, we have

$$
\operatorname{Tr} h^{(x)}=0 \quad(x=1,2,3,4)
$$

and therefore

$$
\operatorname{Tr} h^{\prime}=0 \quad \text { and } \quad \operatorname{Tr} k^{\prime}=0 .
$$

These equations imply
Proposition 2.3. The integral manifold $H$ of the distribution $D: p \rightarrow H_{p}$ is a minimal submanifold of $S$.

We also have

$$
\operatorname{Tr} h=\left(\langle h \alpha, \alpha\rangle_{S}+\langle h \beta, \beta\rangle_{S}\right) /\|\alpha\|^{2}
$$

and

$$
\operatorname{Tr} k=\left(\langle k \alpha, \alpha\rangle_{S}+\langle k \beta, \beta\rangle_{S}\right) /\|\alpha\|^{2}
$$

by virtue of (2.25) and (2.26). Thus we have
Proposition 2.4. If $S$ is a minimal submanifold of $M$ and the integrability condition (2.15) of the distribution $D$ is satisfied, then we have

$$
\begin{equation*}
\langle h \alpha, \alpha\rangle_{S}+\langle h \beta, \beta\rangle_{S}=0 \tag{2.27}
\end{equation*}
$$

and

$$
\begin{equation*}
\langle k \alpha, \alpha\rangle_{S}+\langle k \beta, \beta\rangle_{S}=0 . \tag{2.28}
\end{equation*}
$$

We can write the equations (2.15) as

$$
\left\langle h^{(1)} X, f Y\right\rangle_{H}-\left\langle h^{(1)} Y, f X\right\rangle_{H}=0
$$

and

$$
\left\langle h^{(2)} X, f Y\right\rangle_{I I}-\left\langle h^{(2)} Y, f X\right\rangle_{I I}=0,
$$

if we take account of (2.19), (2.25) and (2.26).
A tensor $T$ of type $(1,1)$ is said to be hybrid with respect to $f$, if it satisfies

$$
f T+T f=0,
$$

where $f$ is a tensor of type (1,1). (See, e.g. Yano [3].)
Thus we have, taking account of (2.20),
Proposition 2.5. Under the integrability condition (2.15) of the distribution $D: p \rightarrow H_{p}$, each $h^{(x)}$ is hybrid tensor with respect to the almost complex structure $f$ on $H$ induced from the almost complex structure $F$ of $M$.

Let us define $\tilde{h}^{(x)}$ by

$$
\tilde{h}^{(x)}(V, W)=\left\langle h^{(x)} V, W\right\rangle_{H} \quad(x=1,2,3,4)
$$

for any pair of vectors $V$ and $W$ on $H$. Then we can see, from (2.25) and (2.26), that $\tilde{h}^{(1)}$ and $\tilde{h}^{(2)}$ are both symmetric bilinear form. As consequences of (2.15), we see that $\tilde{h}^{(3)}$ and $\tilde{h}^{(4)}$ become symmetric when the distribution $D$ is integrable. ${ }^{3)}$
3. We study, in this section, the integrability condition of the distribution $D^{\prime}$ which assigns $H_{p}^{\prime}$ to $p \in S$. If we use $\alpha$ and $\beta$ as a local basis of $D^{\prime}$ in a sufficiently small neighborhood of $p$, then the integrability condition of $D^{\prime}: p \rightarrow H_{p}^{\prime}$ is written as

$$
\begin{equation*}
\left\langle X_{\mu},[\alpha, \beta]\right\rangle_{S}=0, \tag{3.1}
\end{equation*}
$$

where $X_{\mu}$ is the local basis of the distribution $D: p \rightarrow H_{p}$. The equation (3.1) is written as

$$
\begin{equation*}
\left\langle X_{\mu},(h A-A h) \beta-(k A-A k) \alpha\right\rangle_{S}=0 \tag{3.2}
\end{equation*}
$$

3) If we define, in a sufficiently small neighborhood $\subset \geqslant$, van der Waerden-Bortoloth covariant derivative along a distribution $D$ and introduce tensors $L^{(x)}$ of type (1, 1) in a similar way as we did for $h^{(x)}$, i.e. by (normal part of $\left.\nabla_{X} Y\right)=\widetilde{L}^{(x)}(X, Y) C_{x}(x=1,2,3,1)$, for local vector fields $X, Y \in D$ then the integrability condition of the distribution $\underset{\sim}{D}$ is given by the hybridness of $L^{(1)}$ and $L^{(2)}$ or equivalently by the symmetry of $\widetilde{L}^{(3)}$ and $\widetilde{L}^{(4)}$.
or

$$
\begin{equation*}
\varphi\left(A_{\mu}+\bar{B}_{\mu}\right)+f_{\mu}{ }^{\nu}\left(B_{\nu}-\bar{A}_{\nu}\right)=0, \tag{3.3}
\end{equation*}
$$

because of $A \alpha=-\varphi \beta$ and $A \beta=\varphi \alpha$, where $A_{\mu}=\left\langle X_{\mu}, h \alpha\right\rangle_{S}, B_{\mu}=\left\langle X_{\mu}, h \beta\right\rangle_{S}, \bar{A}_{\mu}=\left\langle X_{\mu}, k \alpha\right\rangle_{S}$, $\bar{B}_{\mu}=\left\langle X_{\mu}, k \beta\right\rangle_{S}$ and $\left(f_{\mu}{ }^{\nu}\right)$ are components of the tensor $f$ on $H$. Thus we have

Proposition 3.1. Suppose that $\operatorname{dim} H_{p}^{\prime}=$ const. In order that the distribution $D^{\prime}: p \rightarrow H_{p}^{\prime}$ is completely integrable, it is necessary and sufficient that the equation (3.2) or (3.3) holds.

Corollary 3.1. If $S$ is a totally geodesic submanifold of $M$, then the distribution $D^{\prime}$ is completely integrable.

On the other hand, we have

$$
\left\langle X_{\mu}, \operatorname{grad} \varphi\right\rangle_{S}=-B_{\mu}+\bar{A}_{\mu}
$$

and thus we have
Corollary 3.2. For a submanifold $S$ on which $\varphi=0$, the distribution $D^{\prime}$ is completely integrable.

The number of equations (3.3) is $m-2$ and that of unknown variables $A_{\mu}, B_{\mu}$, $\bar{A}_{\mu}$ and $\bar{B}_{\mu}$ is $4(m-2)$. Therefore, it seems that a submanifold which satisfies the integrability condition (3.3) of the distribution $D^{\prime}$ is a very special one. We shall show, at the end of this section, an example of such submanifolds. In that example, the second fundamental tensors $h$ and $k$ which are considered as linear transformations on $T_{p}(S)$ leave invariant the holomorphic subspace $H_{p}$ of $T_{p}(S)$.

When the distribution $D^{\prime}: p \rightarrow H_{p}^{\prime}$ is integrable, we denote by $H^{\prime}$ the integral manifold of the distribution $D^{\prime}$ and by $\zeta^{\prime}$ a submanifold map $\zeta^{\prime}: H^{\prime} \rightarrow M$. We can choose $\xi_{*} X_{\lambda}, C, D$ as unit normal vector fields to $H^{\prime}$. By using van der WaerdenBortolotti covariant derivation $\nabla^{5^{\prime}}$ along $H^{\prime}$ we have

$$
\begin{aligned}
& \left\langle\nabla_{\dot{x}}^{\left.c_{X}^{\prime} C, \zeta_{*}^{\prime} Y\right\rangle_{M}=-\left\langle^{\prime} h^{(m-1)} X, Y\right\rangle_{H^{\prime}},}\right. \\
& \left\langle\nabla_{\dot{x}}^{\prime} D, \zeta_{*}^{\prime} Y\right\rangle_{M}=-\left\langle^{\prime} h^{(m)} X, Y\right\rangle_{H^{\prime}},
\end{aligned}
$$

and

$$
\left\langle\nabla_{X}^{v_{x}^{\prime}} \zeta_{*} X_{\lambda}, \zeta_{*}^{\prime} Y\right\rangle_{M}=-\left\langle^{\prime} h^{(\lambda)} X, Y\right\rangle_{H^{\prime}},
$$

where $X$ and $Y$ are arbitrary vector fields on $H^{\prime}$ and ${ }^{\prime} h^{(\lambda),} h^{(m-1)}$ and ${ }^{\prime} h^{(m)}$ are the second fundamental tensors of $H^{\prime}$ as a submanifold of $M$. Since we have chosen $\alpha$ and $\beta$ as a local basis of $H^{\prime}$, we have

$$
\begin{equation*}
\|\alpha\|^{2} \operatorname{Tr}^{\prime} h^{(m-1)}=\left(\langle h \alpha, \alpha\rangle_{S}+\langle h \beta, \beta\rangle_{S}\right) /\|\alpha\|^{2} \tag{3.4}
\end{equation*}
$$

and

$$
\begin{equation*}
\|\alpha\|^{2} \operatorname{Tr}^{\prime} h^{(m)}=\left(\langle k \alpha, \alpha\rangle_{S}+\langle k \beta, \beta\rangle_{S}\right) /\|\alpha\|^{2} . \tag{3.5}
\end{equation*}
$$

On the other hand, we have

$$
\|\alpha\|^{2} \operatorname{Tr}^{\prime} h^{(\lambda)}=\varphi\left(B_{\lambda}-\bar{A}_{\lambda}\right)-f_{\lambda}^{\mu}\left(A_{\mu}+\bar{B}_{\mu}\right),
$$

from which we obtain

$$
\operatorname{Tr}^{\prime} h^{(\lambda)}=-f_{\lambda^{\mu}}{ }^{\mu}\left(A_{\mu}+\bar{B}_{\mu}\right),
$$

by virtue of (3.3). On the other hand, we have, by a straightforward computation,

$$
\left\langle X_{\mu}, \operatorname{div} A\right\rangle_{S}=A_{\mu}+\bar{B}_{\mu},
$$

where $A$ is the tensor defined in $\S 2$. This proves
Proposition 3.2. A necessary condition that $H^{\prime}$ is a minimal submanifold of $M$ is

1) The second fundamental tensors $h$ and $k$ of $S$ satisfy (3.4) and (3.5) respectively and
2) $\operatorname{div} A \in H_{p}^{\prime}$.

Conversely, if we assume 1) and 2) mentioned above and we assume further $\operatorname{grad} \varphi \in H_{p}^{\prime}$ and $\operatorname{dim} H_{p}^{\prime}=$ const., then the distribution $D^{\prime}: p \rightarrow H_{p}^{\prime}$ is completely integrable and the integral manifold $H^{\prime}$ of the distribution $D^{\prime}$ is a minimal submanifold of $M$.

We shall now discuss a sufficient condition under which the distributions $D: p \rightarrow H_{p}$ and $D^{\prime}: p \rightarrow H_{p}^{\prime}$ are both integrable. We assume that the following equations are valid for any local vector field $X$ of the distribution $D$ :

$$
\left\{\begin{array}{l}
(A h+h A) X=0  \tag{3.6}\\
(A k+k A) X=0
\end{array}\right.
$$

The equation (2.15) shows that (3.6) is one of sufficient conditions under which the distribution $D$ is completely integrable. The next lemma is a result of a direct computation

Lemma 3.1. Under the condition (3.6), we have

$$
\begin{equation*}
A_{\mu}=B_{\mu}=\bar{A}_{\mu}=\bar{B}_{\mu}=0 \tag{3.7}
\end{equation*}
$$

and therefore $\operatorname{grad} \varphi$ and $\operatorname{div} A$ belong to $H_{p}^{\prime}$.
(3.7) proves

Proposition 3.3. If we assume (3.6), then the holomorphic subspace $H_{p}$ of $T_{p}(S)$ is left invariant under the linear transformations induced from the second fundamental tensors $h$ and $k$ of $S$.

Proposition 3.4. The distributions $D$ and $D^{\prime}$ are both completely integrable, if we assume $\operatorname{dim} H_{p}=$ const. and the equations (3.6).

Lemma 3.2. Let $S$ be a minimal submanifold of $M$. If the equations (3.6)
are valid for any $X \in H_{p}$, then we have

$$
\left\{\begin{array}{c}
A h+h A=0  \tag{3.8}\\
A k+k A=0
\end{array}\right.
$$

on $T_{p}(S)$.
Conversely, a submanifold $S$ whose second fundamental tensors $h$ and $k$ satisfy the equations (3.8), then $S$ is a minimal submanifold of $M$, if $\varphi$ does not vanish.

Proof.
Straightforward computations show that

$$
\begin{aligned}
\langle(A h+h A) \alpha, \alpha\rangle_{S} & =\langle A h \alpha, \alpha\rangle_{S}-\varphi\langle h \beta, \alpha\rangle_{S} \\
& =\varphi\langle h a, \beta\rangle_{S}-\varphi\langle h \beta, \alpha\rangle_{S}=0 ; \\
\langle(A h+h A) \alpha, \beta\rangle_{S} & =\langle A h \alpha, \beta\rangle_{S}-\varphi\langle h \beta, \beta\rangle_{S} \\
& =-\varphi\langle h \alpha, \alpha\rangle_{S}-\varphi\langle h \beta, \beta\rangle_{S} \\
& =0, \quad \text { by }(2.27) ; \\
\langle(A h+h A) \beta, \alpha\rangle_{S} & =-\langle(A h+h A) \alpha, \beta\rangle_{S}=0 ;
\end{aligned}
$$

and

$$
\langle(A h+h A) \beta, \beta\rangle_{S}=-\varphi\langle h \beta, \alpha\rangle_{S}+\varphi\langle h \alpha, \beta\rangle_{S}=0 .
$$

We have, from (3.6),

$$
\langle(A h+h A) \alpha, X\rangle_{S}=\langle(A h+h A) \beta, X\rangle_{S}=0
$$

for any $X \in H_{p}$. These equations give

$$
\begin{equation*}
(A h+h A) \alpha=(A h+h A) \beta=0 . \tag{3.9}
\end{equation*}
$$

Similar computations show that

$$
\begin{equation*}
(A k+k A) \alpha=(A k+k A) \beta=0 . \tag{3.10}
\end{equation*}
$$

The equations (3.8) follow from (3.6), (3.9) and (3.10).
The converse is now obvious by a straightforward computation. q.e.d.
Corollary 3.3. Let $S$ be a minimal submanifold of $M$. We assume that the equations (3.6) are valid and further the function $\varphi$ is constant. Then $\tilde{A}$ defined by (2.2) is harmonic form.

Proof. From the definition of $\tilde{A}$ and the equation

$$
\begin{aligned}
\left(\nabla_{X} \tilde{A}\right)(Y, Z)= & -\tilde{h}(X, Y) \tilde{\alpha}(Z)+\tilde{h}(X, Z) \tilde{\alpha}(Y) \\
& -\tilde{k}(X, Y) \tilde{\beta}(Z)+\tilde{k}(X, Z) \tilde{\beta}(Y),
\end{aligned}
$$

it is obvious that $\tilde{A}$ is skew-symmetric and closed. We can easily see that

$$
\langle\operatorname{div} A, \alpha\rangle_{S}=\langle h \alpha, \alpha\rangle_{S}+\langle k \beta, \alpha\rangle_{S}
$$

and

$$
\langle\operatorname{div} A, \beta\rangle_{S}=\langle h \alpha, \beta\rangle_{S}+\langle k \beta, \beta\rangle_{S} .
$$

On the other hand, we have

$$
\langle\operatorname{grad} \varphi, \alpha\rangle_{S}=\langle k \alpha, \alpha\rangle_{S}-\langle h \beta, \alpha\rangle_{S}
$$

and

$$
\langle\operatorname{grad} \varphi, \beta\rangle_{S}=\langle k \alpha, \beta\rangle_{S}-\langle h \beta, \beta\rangle_{S},
$$

from which we have

$$
\langle\operatorname{div} A, \alpha\rangle_{S}=\langle h \alpha, \alpha\rangle_{S}+\langle h \beta, \beta\rangle_{S}
$$

and
$\langle\operatorname{div} A, \beta\rangle_{S}=\langle k \alpha, \alpha\rangle_{S}+\langle k \beta, \beta\rangle_{S}$.
The right hand side of each equation above must be zero because of (2.27) and (2.28). Since $\operatorname{div} A$ belongs to $H_{p}^{\prime}$ (Lemma 3.1), we have $\operatorname{div} A=0$ which implies, together with $d \tilde{A}=0$, that $\tilde{A}$ is a harmonic form. q.e.d.

Summing up the results, we have
Theorem 3.1. Let $S$ be a minimal submanifold of a Kählerian manifold $M$ whose codimension is $2 .{ }^{4)}$ We assume that $\operatorname{dim} H_{p}=$ const. and the second fundamental tensors $h$ and $k$ of $S$ satisfy the condition (3.6). Then $S$ is locally decomposed into two submanifolds one of which is holomorphic in $M$ and the other is anti-holomorphic in $M$ both being minimal submanifolds of $S$ and at the same time of $M$. The dimension of the anti-holomorphic submanifold equals to the codimension of S .

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[^2]
[^0]:    Received June 19, 1967.

    1) We use here an identification of a differentiable manifold $S$ with $\xi(S)$, where $\xi$ is a differentiable immersion from $S$ into $M$, whose differential $\xi_{*}: T_{p}(S) \rightarrow T_{\xi(p)}(M)$ is injective. Manifolds, mappings and geometric objects considered in this paper are all assumed to be of differentiability class $C^{\infty}$.
[^1]:    2) Vectors $X$ and $Y$ on $H_{p}$ are regarded as vectors on $T_{p}(S)$ by the identification map.
[^2]:    4) We have assumed, throughout this paper, that codimension of $S$ is 2 , but we can also discuss in the same way the case in which the condimension of $S$ is even and smaller than or equals to the half of the dimension of $M$.
