## LECTURE 13

## THEORY OF THE LAST MULTIPLIER (CONT.) APPLICATIONS.

Now consider the most general system of equations (1):

$$
\begin{equation*}
\frac{d x}{X}=\frac{d x_{1}}{X_{1}}=\ldots=\frac{d x_{n}}{X_{n}} \tag{1}
\end{equation*}
$$

in which $X, X_{1}, \ldots, X_{n}$ are given functions of $x, x_{1}, \ldots, x_{n}$.
The system is integrated when one knows $n$ first integrals. Suppose that one knows only ( $n-$ 1) of them, and let them be:
( $\gamma$ )

From those $(n-1)$ integrals, we can infer $(n-1)$ of the variables $x_{i}$, for example $x_{2}, x_{3}, \ldots, x_{n}$, as functions of the other two, $x$ and $x_{1}$, and of the $\alpha$. That amounts to saying that the determinant $\Delta$ :

$$
\Delta=\left|\begin{array}{cccc}
\frac{\partial f_{2}}{\partial x_{2}} & \frac{\partial f_{2}}{\partial x_{3}} & \cdots & \frac{\partial f_{2}}{\partial x_{n}} \\
\frac{\partial f_{3}}{\partial x_{2}} & \frac{\partial f_{3}}{\partial x_{3}} & \cdots & \frac{\partial f_{3}}{\partial x_{n}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial f_{n}}{\partial x_{2}} & \frac{\partial f_{n}}{\partial x_{3}} & \cdots & \frac{\partial f_{n}}{\partial x_{n}}
\end{array}\right|
$$

is not identically zero.
More generally, if $F$ denotes an arbitrary function of $x_{1}, x_{2}, \ldots, x_{n}$ then we shall let $F^{\prime}$ represent what $F$ will become when we replace $x_{2}, x_{3}, \ldots, x_{n}$ with their values that we infer from the system $(\gamma)$. Having made that convention, it will be clear that $x$ and $x_{1}$ satisfy the equation:

$$
X_{1}^{\prime} d \alpha-X^{\prime} d x_{1}=0
$$

in which one leaves the $\alpha$ constant.
When one knows a solution $M$ of the equation:

$$
X \frac{\partial \log M}{\partial x}+X_{1} \frac{\partial \log M}{\partial x_{1}}+\cdots+X_{n} \frac{\partial \log M}{\partial x_{n}}+\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}=0,
$$

one will know an integrating factor to the equation (2). That factor is equal to $M^{\prime} / \Delta^{\prime}\left({ }^{1}\right)$.
In other words:

$$
\frac{M^{\prime}}{\Delta^{\prime}}\left(X_{1}^{\prime} d \alpha-X^{\prime} d x_{1}\right)=d F\left(x, x_{1}, \alpha_{2}, \alpha_{3}, \ldots, \alpha_{m}\right)
$$

The proof differs from the one that was given in the case of two equations only by the complexity of the calculations.

Suppose for the moment that one knows $n$ first integrals $f_{1}=\alpha_{1}, f_{2}=\alpha_{2}, \ldots, f_{n}=\alpha_{n}$. The equations:

$$
\begin{aligned}
& X \frac{\partial f}{\partial x}+X_{1} \frac{\partial f}{\partial x_{1}}+\cdots+X_{n} \frac{\partial f}{\partial x_{n}}=0, \\
& \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots \cdots, \\
& X \frac{\partial f_{n}}{\partial x}+X_{1} \frac{\partial f_{n}}{\partial x_{1}}+\cdots+X_{n} \frac{\partial f_{n}}{\partial x_{n}}=0
\end{aligned}
$$

will be equivalent to the following ones:

$$
\frac{A}{X}=\frac{A_{1}}{X_{1}}=\frac{A_{2}}{X_{2}}=\ldots=\frac{A_{n}}{X_{n}}=M_{1},
$$

when one sets:

$$
R=\left|\begin{array}{cccc}
\frac{\partial f_{1}}{\partial x} & \frac{\partial f_{1}}{\partial x_{1}} & \cdots & \frac{\partial f_{1}}{\partial x_{n}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial f_{n}}{\partial x} & \frac{\partial f_{n}}{\partial x_{1}} & \cdots & \frac{\partial f_{n}}{\partial x_{n}} \\
\alpha & \alpha_{1} & \cdots & \alpha_{n}
\end{array}\right|
$$

and

$$
A=\frac{\partial R}{\partial \alpha}, \quad A_{1}=\frac{\partial R}{\partial \alpha_{1}}, \quad \ldots, \quad A_{n}=\frac{\partial R}{\partial \alpha_{n}},
$$

We shall first show that the quantities $A_{i}$ verify the relation:

$$
\begin{equation*}
\frac{\partial A}{\partial x}+\frac{\partial A_{1}}{\partial x_{1}}+\cdots+\frac{\partial A_{n}}{\partial x_{n}} \equiv 0 . \tag{4}
\end{equation*}
$$

[^0]Indeed, observe that the left-hand side of (4) is a homogeneous linear function of the second derivatives of $f_{1}, f_{2}, \ldots, f_{n}$.

Now that function cannot contain a term in $\partial^{2} f_{k} / \partial x_{i}^{2}$, because the term in question would come from $\partial A_{i} / \partial x_{i}$, and $A_{i}=\partial R / \partial \alpha_{i}$ contains no derivative of $f$ with respect to the variable $x_{i}$.

It can no longer include a term in $\frac{\partial^{2} f_{k}}{\partial x_{i} \partial x_{j}}$, because that term would come from the sum $\frac{\partial A_{i}}{\partial x_{i}}+\frac{\partial A_{j}}{\partial x_{j}}:$ Set $\beta_{i}^{k}=\partial f_{k} / \partial x_{i}$, to abbreviate the notation. The term in $\frac{\partial^{2} f_{k}}{\partial x_{i} \partial x_{j}}$ that comes from $\frac{\partial A_{i}}{\partial x_{i}}$ will have the coefficient $\frac{\partial A_{i}}{\partial \beta_{j}^{k}}$ or $\frac{\partial^{2} R}{\partial \alpha_{i} \partial \beta_{j}^{k}}$. The two second-order minors of $R$ that are thus introduced are equal in absolute value, because they are obtained by suppressing the same rows and columns of $R$. Moreover, they have opposite signs because one can pass from one to the other by permuting two columns in $R$.

The sum of the terms in the left-hand side of (4) is therefore identically zero, and the relation (4) is then established. If one replaces $A, A_{1}, \ldots, A_{n}$ with $M_{1} X, M_{1} X_{1}, \ldots, M_{1} X_{n}$ in the relation then one will see that $M_{1}$ is an integral of the equation:

$$
\begin{equation*}
X \frac{\partial \log M}{\partial x}+X_{1} \frac{\partial \log M}{\partial x_{1}}+\cdots+X_{n} \frac{\partial \log M}{\partial x_{n}}+\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}=0 . \tag{3}
\end{equation*}
$$

Having said that, it will be easy to see that $M_{1}^{\prime} / \Delta^{\prime}$ is an integrating factor of equation (2). Since one has:

$$
A^{\prime}=M_{1}^{\prime} X^{\prime}, \quad A_{1}^{\prime}=M_{1}^{\prime} X_{1}^{\prime},
$$

it will suffice to prove that $1 / \Delta^{\prime}$ is an integrating factor of the equation:

$$
A_{1}^{\prime} d x-A^{\prime} d x_{1}=0
$$

To that end, express the variables $x_{1}, x_{2}, \ldots, x_{n}$ in $f_{1}\left(x, x_{1}, x_{2}, \ldots, x_{n}\right)$ as functions of the $(n-1)$ variables $\alpha_{2}, \alpha_{3}, \ldots, \alpha_{n}$ and of $x$ and $x_{1}$. We have:

$$
f_{1}\left(x, x_{1}, x_{2}, \ldots, x_{n}\right)=f_{1}^{\prime}\left(x, x_{1}, \alpha_{3}, \ldots, \alpha_{n}\right),
$$

and as a result:

$$
\begin{aligned}
& \frac{\partial f_{1}}{\partial x}=\frac{\partial f_{1}^{\prime}}{\partial x}+\sum_{i=2}^{n} \frac{\partial f_{1}^{\prime}}{\partial \alpha_{i}} \frac{\partial f_{i}}{\partial x} \\
& \frac{\partial f_{1}}{\partial x_{1}}=\frac{\partial f_{1}^{\prime}}{\partial x_{1}}+\sum_{i=2}^{n} \frac{\partial f_{1}^{\prime}}{\partial \alpha_{i}} \frac{\partial f_{i}}{\partial x_{1}}
\end{aligned}
$$

$$
\begin{aligned}
& \frac{\partial f_{1}}{\partial x_{2}}=\sum_{i=2}^{n} \frac{\partial f_{1}^{\prime}}{\partial \alpha_{i}} \frac{\partial f_{i}}{\partial x_{2}}, \\
& \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
& \frac{\partial f_{1}}{\partial x_{n}}=\sum_{i=2}^{n} \frac{\partial f_{1}^{\prime}}{\partial \alpha_{i}} \frac{\partial f_{i}}{\partial x_{n}} .
\end{aligned}
$$

Moreover, one knows that:

$$
\begin{aligned}
& A=\left|\begin{array}{cccc}
\frac{\partial f_{1}}{\partial x_{1}} & \frac{\partial f_{1}}{\partial x_{2}} & \cdots & \frac{\partial f_{1}}{\partial x_{n}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial f_{n}}{\partial x_{1}} & \frac{\partial f_{n}}{\partial x_{2}} & \cdots & \frac{\partial f_{n}}{\partial x_{n}}
\end{array}\right|, \\
& A_{1}=\left|\begin{array}{cccc}
\frac{\partial f_{1}}{\partial x} & \frac{\partial f_{1}}{\partial x_{2}} & \cdots & \frac{\partial f_{1}}{\partial x_{n}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial f_{n}}{\partial x} & \frac{\partial f_{n}}{\partial x_{2}} & \cdots & \frac{\partial f_{n}}{\partial x_{n}}
\end{array}\right| .
\end{aligned}
$$

If one replaces the $\partial f_{1} / \partial x_{i}$ with the values that are written above then one will see that upon supposing that $x_{2}, x_{3}, \ldots, x_{n}$ are expressed as functions of $\alpha_{2}, \alpha_{3}, \ldots, \alpha_{n}$, and $x, x_{1}$ everywhere, and the expressions $A$ and $A_{1}$ reduce to $\frac{\partial f_{1}^{\prime}}{\partial x_{1}} \Delta$ and $-\frac{\partial f_{1}^{\prime}}{\partial x_{1}} \Delta$, respectively, i.e., one will have:

$$
A^{\prime}=\frac{\partial f_{1}^{\prime}}{\partial x_{1}} \Delta, \quad A_{1}^{\prime}=-\frac{\partial f_{1}^{\prime}}{\partial x_{1}} \Delta
$$

As a result:

$$
\begin{equation*}
\frac{A_{1}^{\prime} d x-A^{\prime} d x_{1}}{\Delta^{\prime}}=-\left[\frac{\partial f_{1}^{\prime}}{\partial x} d x+\frac{\partial f_{1}^{\prime}}{\partial x_{1}} d x_{1}\right] \tag{5}
\end{equation*}
$$

If one then gives constant values to $\alpha_{2}, \alpha_{3}, \ldots, \alpha_{n}$ then the left-hand side of equation (5) will be the total differential of the function $f_{1}^{\prime}\left(x, x_{1}, \alpha_{3}, \ldots, \alpha_{n}\right)$, which proves that $M_{1}^{\prime} / \Delta^{\prime}$ is indeed an integrating factor of equation (2).

Finally, let $M$ be an arbitrary solution to (3). One soon sees that the ratio $M / M_{1}$ or $N$ satisfies the equation:

$$
X \frac{\partial N}{\partial x}+X_{1} \frac{\partial N}{\partial x_{1}}+\cdots+X_{n} \frac{\partial N}{\partial x_{n}}=0
$$

and as a result:

$$
N=\varphi\left(f_{1}, f_{2}, \ldots, f_{n}\right)
$$

Hence:

$$
M=M_{1} \varphi\left(f_{1}, f_{2}, \ldots, f_{n}\right) .
$$

If one replaces $M_{1}$ as a function of $M$ in equation (5) then it will become:

$$
\frac{M^{\prime}}{\Delta^{\prime}}\left(X_{1}^{\prime} d x-X^{\prime} d x_{1}\right)=\varphi\left(f_{1}, f_{2}, \ldots, f_{n}\right) d f_{1}=d F\left(x, x_{1}, \alpha_{2}, \ldots, \alpha_{n}\right) .
$$

$M^{\prime} / \Delta^{\prime}$ is then an integrating factor for equation (2). That is the theorem that we wish to prove.
When one knows an arbitrary solution $M$ of equation (3), it will then suffice to form ( $n-1$ ) distinct first integrals of (1) in order to achieve the integration by quadratures. The function $M$ is called a last multiplier of the system (1).

In certain cases, one can recognize a last multiplier immediately. That is what happens, in particular, when the expression:

$$
\delta=\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}=0
$$

is identically zero: One can then set $M=1$, and the final binomial will admit $1 / \Delta^{\prime}$ as an integrating factor.

It is fitting to observe, in that regard, that if one replaces $x_{1}, x_{2}, \ldots, x_{n}$ with functions of $x$ that satisfy the system (1) then that equation will become:

$$
\frac{d L M}{d x}+\frac{1}{X}\left(\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}\right) \equiv 0 .
$$

Conversely, if a function $U\left(x, x_{1}, \ldots, x_{n}\right)$ [when one replaces $x_{n}, x_{1}$ in it with an arbitrary system of in integrals of (1)] verifies the relation:

$$
\frac{d U}{d x}=\frac{1}{X}\left(\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}\right)
$$

then $M=e^{-U}$ will be a last multiplier of (1).
From that, if the expression:

$$
\frac{1}{X}\left(\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}\right)
$$

is a function of only $x$ then a last multiplier will be obtained by a quadrature. More generally, if one has:

$$
\frac{1}{X}\left(\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}\right) \equiv \psi\left(x_{i}\right)
$$

then $M=\exp \left[-\int \psi\left(x_{i}\right) d x_{i}\right]$ will be a multiplier.
Finally, we add that if $x, x_{1}, \ldots, x_{k}$ do not enter into any functions $X$ then any multiplier $M$ of (1) that is independent of $x, x_{1}, \ldots, x_{k}$ will also be a multiplier of the system:

$$
\frac{d x_{(k+1)}}{X_{k+1}}=\frac{d x_{(k+2)}}{X_{(k+2)}}=\ldots=\frac{d x_{n}}{X_{n}}
$$

Indeed, $M\left(x_{k+1}, x_{k+2}, \ldots, x_{n}\right)$ verifies the equation:

$$
X_{k+1} \frac{\partial \log M}{\partial x_{k+1}}+X_{k+2} \frac{\partial \log M}{\partial x_{k+2}}+\cdots+X_{n} \frac{\partial \log M}{\partial x_{n}}=0 .
$$

Remark in regard to the case in which one knows only $(n-k)$ first integrals. - In the foregoing, we assumed that we had formed $(n-1)$ first integrals. Knowing a last multiplier would then permit us to find a last integral by quadrature.

We now place ourselves in the case where we know only ( $n-k$ ) first integrals of equations (1):

$$
\begin{equation*}
\frac{d x}{X}=\frac{d x_{1}}{X_{1}}=\ldots=\frac{d x_{n}}{X_{n}} \tag{1}
\end{equation*}
$$

One can infer $(n-k)$ of the variables $x, x_{1}, \ldots, x_{n}$ (for example, $x_{n}, x_{n-1}, \ldots, x_{k+1}$ ) as functions of the $(k+1)$ other ones $x, x_{1}, \ldots, x_{k}$ and $(n-k)$ constants $\alpha$ from those $(n-k)$ integrals, say, $f_{n}=$ $\alpha_{n}, f_{(n-1)}=\alpha_{n-1}, \ldots, f_{(k+1)}=\alpha_{k+1}$. Let $(F)$ denote what a function $F\left(x, x_{1}, \ldots, x_{n}\right)$ will become after that substitution. The variables $x, x_{1}, \ldots, x_{n}$ satisfy the equations:

$$
\begin{equation*}
\frac{d x}{(X)}=\frac{d x_{1}}{\left(X_{1}\right)}=\ldots=\frac{d x_{k}}{\left(X_{k}\right)}, \tag{1'}
\end{equation*}
$$

which must now be integrated. Can knowing a multiplier $M$ of the system (1) serve to integrate the system ( $1^{\prime}$ )? We shall show that it is easy to deduce a multiplier for ( $1^{\prime}$ ) from a multiplier of (1).

From the foregoing, an arbitrary multiplier $M$ of (1) satisfies (if we preserve the notations on page 189) the relations:

$$
M X=\varphi A, \quad M X_{1}=\varphi A_{1}, \ldots, \quad M X_{n}=\varphi A_{n}
$$

in which $\varphi$ denotes an arbitrary function of $f_{1}, f_{2}, \ldots, f_{n}$. Conversely, any function $M\left(x, x_{1}, \ldots, x_{n}\right)$ that satisfies one of those relations, say, the relation:

$$
M X=\varphi\left(f_{1}, f_{2}, \ldots, f_{n}\right) A
$$

will be a multiplier of (1). Having recalled that, assume that one knows a first integral of (1) (say $f_{n}=\alpha_{n}$ ) and infer one of the variables (for example, $x_{n}$ ) from that integral as a function of the $x, x_{1}$, $\ldots, x_{n-1}$, and $\alpha_{n}$. Let $F^{\prime}$ and $\left(\frac{\partial F}{\partial x_{n}}\right)$ denote what the functions $F\left(x, x_{1}, \ldots, x_{n}\right)$ and $\partial F / \partial x_{n}$ will become when one replaces $x_{n}$ with that value. One will then have:

$$
f_{i}\left(x, x_{1}, \ldots, x_{n-1}, x_{n}\right)=f_{i}^{\prime}\left(x, x_{1}, \ldots, x_{n-1}, x_{n}\right),
$$

and as a result:

$$
\begin{aligned}
& \frac{\partial f_{i}}{\partial x_{j}}=\frac{\partial f_{i}^{\prime}}{\partial x_{j}}+\frac{\partial f_{i}^{\prime}}{\partial \alpha_{n}} \frac{\partial f_{n}}{\partial x_{j}} \quad\binom{i=1,2, \ldots,(n-1)}{j=1,2, \ldots,(n-1)} \\
& \frac{\partial f_{i}}{\partial x_{n}}=\frac{\partial f_{i}}{\partial \alpha_{n}} \frac{\partial f_{n}}{\partial x_{n}}
\end{aligned}
$$

If one replaces the $\frac{\partial f_{i}}{\partial x_{j}}, \frac{\partial f_{i}}{\partial x_{n}}$ in the determinant $A$ with those values then one will soon see that $A$ reduces to:

$$
\frac{\partial f_{n}}{\partial x_{n}}\left|\begin{array}{cccc}
\frac{\partial f_{1}^{\prime}}{\partial x_{1}} & \frac{\partial f_{1}^{\prime}}{\partial x_{2}} & \cdots & \frac{\partial f_{1}^{\prime}}{\partial x_{n-1}} \\
\frac{\partial f_{2}^{\prime}}{\partial x_{1}} & \frac{\partial f_{2}^{\prime}}{\partial x_{2}} & \cdots & \frac{\partial f_{2}^{\prime}}{\partial x_{n-1}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial f_{n-1}^{\prime}}{\partial x_{1}} & \frac{\partial f_{n-1}^{\prime}}{\partial x_{2}} & \cdots & \frac{\partial f_{n-1}^{\prime}}{\partial x_{n-1}}
\end{array}\right|=\frac{\partial f_{n}}{\partial x_{n}} \alpha .
$$

One can then write the equality as:

$$
M^{\prime} X^{\prime} \equiv \varphi\left(f_{1}^{\prime}, f_{2}^{\prime}, \ldots, f_{(n-1)}^{\prime}, \alpha_{n}\right)\left(\frac{\partial f_{n}}{\partial x_{n}}\right) \alpha
$$

On the other hand, the multipliers $\mu$ of the system:

$$
\begin{equation*}
\frac{d x}{X^{\prime}}=\frac{d x_{1}}{X_{1}^{\prime}}=\ldots=\frac{d x_{n-1}}{X_{(n-1)}^{\prime}} \tag{2}
\end{equation*}
$$

are given by the relation:

$$
\mu X^{\prime}=\psi\left(f_{1}^{\prime}, f_{2}^{\prime}, \ldots, f_{(n-1)}^{\prime}\right) \alpha
$$

in which $\psi$ is an arbitrary function of $f_{1}^{\prime}, f_{2}^{\prime}, \ldots, f_{(n-1)}^{\prime}$. Therefore, the function $\mu_{1}\left(x, x_{1}, \ldots, x_{n-1}\right.$, $\alpha_{n}$ ) that is defined by the equality:

$$
\mu_{1} X^{\prime}=\varphi\left(f_{1}^{\prime}, f_{2}^{\prime}, \ldots, f_{(n-1)}^{\prime}\right) \alpha
$$

will be a multiplier of (2). Since one has:

$$
\mu_{1}=\frac{M^{\prime}}{\left(\frac{\partial f_{n}}{\partial x_{n}}\right)}
$$

one will see that knowing a multiplier of (1) will imply that one knows a multiplier of (2).
It is clear that one can reason with the system (2) as one did with the system (1). If one knows a second first integral of $(1), f_{(n-1)}=\alpha_{(n-1)}$ then that integral will correspond to an integral of (2), $f_{(n-1)}^{\prime}=\alpha_{(n-1)}$. Upon inferring $x_{(n-1)}$ from the latter relation, one will form the system:

$$
\begin{equation*}
\frac{d \alpha}{X^{\prime \prime}}=\frac{d \alpha_{1}}{X_{1}^{\prime \prime}}=\ldots=\frac{d \alpha_{(n-2)}}{X_{(n-2)}^{\prime \prime}} \tag{3}
\end{equation*}
$$

The expression $\mu^{\prime} /\left(\frac{\partial f_{(n-1)}^{\prime}}{\partial x_{(n-1)}}\right)$ is a multiplier of (3). $\mu^{\prime}$ and $\left(\frac{\partial f_{(n-1)}^{\prime}}{\partial x_{(n-1)}}\right)$ denote what $\mu$ and $\frac{\partial f_{(n-1)}^{\prime}}{\partial x_{(n-1)}}$, resp., will become when one expresses $x_{n-1}$ in terms of $x, x_{1}, \ldots, x_{n-2}, \alpha_{n-1}, \alpha_{n}$.

One will arrive at that conclusion by pursuing the argument in the same manner: If one sets:

$$
\begin{aligned}
& f_{n}\left(x, x_{1}, \ldots, x_{n}\right)=\varphi_{n}\left(x, x_{1}, \ldots, x_{n}\right) \\
& f_{(n-1)}\left(x, x_{1}, \ldots, x_{n}\right)=\phi_{(n-1)}\left(x, x_{1}, \ldots, x_{n}\right) \\
& f_{(n-2)}\left(x, x_{1}, \ldots, x_{n}\right)=\phi_{(n-1)}, \\
& \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots
\end{aligned},
$$

then the expression:

$$
\frac{M}{\frac{\partial \varphi_{n}}{\partial x_{n}} \frac{\partial \varphi_{(n-1)}}{\partial x_{(n-1)}} \cdots \frac{\partial \varphi_{(k+1)}}{\partial x_{(k+1)}}}
$$

in which one replaces $x_{n}, x_{n-1}, \ldots, x_{k+1}$ as functions of $x, x_{1}, \ldots, x_{k}, \alpha_{k+1}, \alpha_{k+2}, \ldots, \alpha_{n}$, will be a last multiplier of the system (1').

In particular, if one sets $k=1$ then one will see that the product $\frac{\partial \varphi_{n}}{\partial x_{n}} \frac{\partial \varphi_{(n-1)}}{\partial x_{(n-1)}} \cdots \frac{\partial \varphi_{(k+1)}}{\partial x_{(k+1)}}$ coincides with the determinant $\Delta$ that was introduced above when one expresses $x_{n}, \ldots, x_{k+1}$ in terms of $x$, $\ldots, x_{k}, \alpha_{k+2}, \ldots, \alpha_{n}$.

The foregoing makes no assumptions about the values of the constants $\alpha$ and will persist $a$ fortiori when one gives particular values to those constants (or at least some of them). From that, let $\alpha_{n}^{0}$ be a certain value of the constant $\alpha_{n}$ (for example, $\alpha_{n}^{0}=0$ ). Assume that one knows a first integral of the system (2) for that value $\alpha_{n}$ :

$$
\varphi_{(n-1)}\left(x, x_{1}, \ldots, x_{n-1}\right)=\alpha_{n-1} .
$$

The expression $M / \frac{\partial \varphi_{n}}{\partial x_{n}} \frac{\partial \varphi_{(n-1)}}{\partial x_{(n-1)}}$ is once more a multiplier of the system (3) for $\alpha_{n}=\alpha_{n}^{0}$. Similarly, if one knows a first integral of the system (3) for $\alpha_{n}=\alpha_{n}^{0}, \alpha_{n-1}=\alpha_{n-1}^{0}$ :

$$
\varphi_{n-2}\left(x, x_{2}, \ldots, x_{n-2}\right)=\alpha_{n-2},
$$

and so on, then the expression (4) will again define a multiplier of ( $1^{\prime}$ ) for $\alpha_{n}=\alpha_{n}^{0}, \alpha_{n-1}=\alpha_{n-1}^{0}$, $\ldots, \alpha_{k+2}=\alpha_{k+2}^{0}$.

We shall soon have occasion to utilize that remark.
We shall now apply the theory of the last multiplier to the equations of dynamics. The first problem that naturally presents itself is the problem of the motion of a free point.

Application to the motion of a free point. - The equations of motion of a free point can be written:

$$
\begin{equation*}
d t=\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d z}{z^{\prime}}=\frac{m d x^{\prime}}{X}=\frac{m d y^{\prime}}{Y}=\frac{m d z^{\prime}}{Z}, \tag{1}
\end{equation*}
$$

in which $X, Y, Z$ are given functions of $x, y, z, x^{\prime}, y^{\prime}, z^{\prime}, t$.
In the general case, one must know six first integrals in order to integrate the system (1).
If $X, Y, Z$ do not depend upon $t$ then it will obviously suffice to know five independent functions of $t$, since the sixth one can be obtained by quadrature.

If $X, Y, Z$ do not depend upon $x^{\prime}, y^{\prime}, z^{\prime}$ then it will likewise suffice to know five first integrals (in which $t$ can appear). The expression $\delta=\frac{\partial X}{\partial x}+\frac{\partial X_{1}}{\partial x_{1}}+\cdots+\frac{\partial X_{n}}{\partial x_{n}}$ that was considered above will be identically zero here. $M=1$ is then a last multiplier of the system (1), and the sixth integral is obtained by quadrature.

Finally, if $X, Y, Z$ depend upon neither $x^{\prime}, y^{\prime}, z^{\prime}$ nor $t$ then from a remark that was made above, $M=1$ will also be a multiplier of the system:

$$
\begin{equation*}
\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d z}{z^{\prime}}=\frac{m d x^{\prime}}{X}=\frac{m d y^{\prime}}{Y}=\frac{m d z^{\prime}}{Z} \tag{1'}
\end{equation*}
$$

It will then suffice to know four first integrals of (1'), i.e., four integrals of (1) that are independent of $t$. The fifth integral of ( $1^{\prime}$ ) will then be given by quadratures, and $t$ is likewise obtained by a quadrature.

Hence, when a free point $M$ is subject to a force that depends upon only the position of the point, it will suffice to know four first integrals of motion (in which time does not figure) in order for the determination of the motion to be achieved by quadratures.

Suppose, for example, that the force is a central force that is a function of only the distance $r$ from the point $M$ to the center $O$ of the force. The theorem of moments provides three first integrals, and the vis viva theorem provides a fourth. Since $t$ does not enter into those integrals, the motion can be calculated by quadratures.

We verify that by applying the theory of the last multiplier to that particular case. Since the trajectories are planar, from the theorem of moments, we can take that plane to be the $x y$-plane. The equations of motion in that plane are (upon setting $m=1$ ):

$$
\begin{equation*}
d t=\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d x^{\prime}}{F \frac{x}{r}}=\frac{d y^{\prime}}{F \frac{y}{r}}, \tag{1}
\end{equation*}
$$

in which $F$ is a function of $r=\sqrt{x^{2}+y^{2}}$. Consider the system:

$$
\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d x^{\prime}}{F \frac{x}{r}}=\frac{d y^{\prime}}{F \frac{y}{r}}
$$

We know two integrals of that system:

$$
\begin{equation*}
f \equiv \frac{1}{2}\left(x^{\prime 2}+y^{\prime 2}\right)-U=\alpha, \text { in which } U=\int F(r) d r, \tag{a}
\end{equation*}
$$

$$
\begin{equation*}
\varphi \equiv x y^{\prime}-x^{\prime} y=\beta \tag{b}
\end{equation*}
$$

If one infers $x^{\prime}$ and $y^{\prime}$ from (a) and (b) and substitutes them in the equation:

$$
x y^{\prime}-x^{\prime} y=0
$$

then that equation will admit:

$$
\frac{1}{\frac{\partial f}{\partial x^{\prime}} \frac{\partial \varphi}{\partial y^{\prime}}-\frac{\partial f}{\partial y^{\prime}} \frac{\partial \varphi}{\partial x^{\prime}}} \equiv \frac{1}{x x^{\prime}+y y^{\prime}}
$$

for an integrating factor.
In order to confirm that, express $x x^{\prime}+y y^{\prime}$ in terms of $x$ and $y$. Upon multiplying (a) by $2\left(x^{2}+y^{2}\right)$ and subtracting the square of $(b)$, one will infer that:

$$
\left(x x^{\prime}+y y^{\prime}\right)^{2}=2 r^{2}(U+\alpha)-\beta^{2}=R(r) .
$$

One will then have:

$$
\begin{aligned}
& x x^{\prime}+y y^{\prime}=\sqrt{R(r)}, \\
& x y^{\prime}-x^{\prime} y=\beta .
\end{aligned}
$$

Hence:

$$
\begin{aligned}
& r^{2} x^{\prime}=-\beta x+x \sqrt{R(r)}, \\
& r^{2} y^{\prime}=+\beta y+y \sqrt{R(r)} .
\end{aligned}
$$

One will then have:
(c)

$$
\begin{aligned}
\frac{y^{\prime} d x-x^{\prime} d y}{x x^{\prime}+y y^{\prime}} & =\frac{1}{\sqrt{R(r)}} \frac{1}{r^{2}}[\beta(x d x+y d y)+(y d x-x d y) \sqrt{R(r)}] \\
& =\frac{\beta d r}{r \sqrt{R(r)}}+\frac{y d x-x d y}{x^{2}+y^{2}} .
\end{aligned}
$$

If one sets $\theta=\arctan y / x$ then one will see that the left-hand side of $(c)$ is the total differential of the function:

$$
\psi=\beta \int_{r_{0}}^{r} \frac{d r}{r \sqrt{R(r)}}+\theta .
$$

That third integral $\psi=\gamma$ permits one to calculate $y$ (and as a result $x^{\prime}, y^{\prime}$ ) as a function of $x$. As for $t$, it is given by the quadrature:

$$
d t=\frac{d x}{x^{\prime}} .
$$

Moreover, one can write:

$$
d t=\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{x d x+y d y}{x x^{\prime}+y y^{\prime}}=\frac{r d r}{\sqrt{R(r)}} .
$$

Therefore:

$$
t=\int_{r_{0}}^{r} \frac{r d r}{\sqrt{R(r)}}
$$

We thus verify the conclusions of the theory of the last multiplier in those particular cases.

## Application to the motion of material systems.

I. Motion of a point on a surface. - Let $M$ be a material point that moves without friction on the surface:

$$
\begin{equation*}
\varphi(x, y, z, t)=0 . \tag{s}
\end{equation*}
$$

If $X^{\prime}, Y^{\prime}, Z^{\prime}$ are the components of the active force then the equations of motion, in the first form that Lagrange gave, can be written (upon setting $m=1$ ):

$$
\begin{equation*}
d t=\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d z}{z^{\prime}}=\frac{d x^{\prime}}{X^{\prime}+\lambda \frac{\partial \varphi}{\partial x}}=\frac{d y^{\prime}}{Y^{\prime}+\lambda \frac{\partial \varphi}{\partial y}}=\frac{d z^{\prime}}{Z^{\prime}+\lambda \frac{\partial \varphi}{\partial z}} . \tag{1}
\end{equation*}
$$

We know that $\lambda$ can be expressed as a function of $x, y, z, x^{\prime}, y^{\prime}, z^{\prime}, t$. Upon differentiating equation $(s)$ twice with respect to $t$ and taking (1) into account, it will become:

$$
\begin{equation*}
\frac{\partial \varphi}{\partial x}\left[X^{\prime}+\lambda \frac{\partial \varphi}{\partial x}\right]+\frac{\partial \varphi}{\partial y}\left[Y^{\prime}+\lambda \frac{\partial \varphi}{\partial y}\right]+\frac{\partial \varphi}{\partial z}\left[Z^{\prime}+\lambda \frac{\partial \varphi}{\partial z}\right]+\left(\frac{\partial \varphi}{\partial t}+\frac{\partial \varphi}{\partial x} x^{\prime}+\frac{\partial \varphi}{\partial y} y^{\prime}+\frac{\partial \varphi}{\partial z} z^{\prime}\right)=0 . \tag{2}
\end{equation*}
$$

If one replaces $\lambda$ in (1) with its value that is deduced from (2) the equations (1) thus-obtained will admit the integral:

$$
\varphi(x, y, z, t)=\alpha t+\beta
$$

or, what amounts to the same thing, the first two integrals:

$$
\begin{align*}
& x^{\prime} \frac{\partial \varphi}{\partial x}+y^{\prime} \frac{\partial \varphi}{\partial y}+z^{\prime} \frac{\partial \varphi}{\partial z}+\frac{\partial \varphi}{\partial t}=\alpha,  \tag{a}\\
& \varphi-t\left(x^{\prime} \frac{\partial \varphi}{\partial x}+y^{\prime} \frac{\partial \varphi}{\partial y}+z^{\prime} \frac{\partial \varphi}{\partial z}+\frac{\partial \varphi}{\partial t}\right)=\beta
\end{align*}
$$

At least one of the derivatives $\frac{\partial \varphi}{\partial x}, \frac{\partial \varphi}{\partial y}, \frac{\partial \varphi}{\partial z}$ is not identically zero, say, $\frac{\partial \varphi}{\partial z}$. Infer $z^{\prime}$ from
(a) by setting $\alpha=0$ and then infer $z$ from (b), which will then reduce to $\varphi=\beta$ by setting $\beta=0$. Finally, substitute those values of $z$ and $z^{\prime}$ in the equations:

$$
\begin{equation*}
d t=\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d x^{\prime}}{X^{\prime}+\lambda \frac{\partial \varphi}{\partial x}}=\frac{d y^{\prime}}{Y^{\prime}+\lambda \frac{\partial \varphi}{\partial y}} \tag{3}
\end{equation*}
$$

We thus define four first-order equations for determining $x, y, x^{\prime}, y^{\prime}$ as functions of $t$, i.e., in order to determine the motion of the point on the surface $\varphi=0$. On the other hand, if $M(x, y, z$, $x^{\prime}, y^{\prime}, z^{\prime}, t$ ) is a multiplier of the system (1), in which $\lambda$ is defined by (2), then we know from a remark that was made above that the expression $M /\left(\frac{\partial \varphi}{\partial z}\right)^{2}$ will also be a multiplier of (3) after we have replaced $z$ and $z^{\prime}$ in it with their values that we infer from the equations:
(c)

$$
\left\{\begin{array}{l}
\varphi(x, y, z, t)=0 \\
\frac{\partial \varphi}{\partial x} x^{\prime}+\frac{\partial \varphi}{\partial y} y^{\prime}+\frac{\partial \varphi}{\partial z} z^{\prime}+\frac{\partial \varphi}{\partial t}=0 .
\end{array}\right.
$$

Having said that, I say that the system (1) admits the quantity:

$$
R=\left(\frac{\partial \varphi}{\partial x}\right)^{2}+\left(\frac{\partial \varphi}{\partial y}\right)^{2}+\left(\frac{\partial \varphi}{\partial z}\right)^{2}
$$

as a multiplier when the given force $\left(X^{\prime}, Y^{\prime}, Z^{\prime}\right)$ does not depend upon velocity.
The equation of the last multiplier here is:

$$
\frac{d}{d t}(\log M)+\frac{\partial \varphi}{\partial x} \frac{\partial \lambda}{\partial x^{\prime}}+\frac{\partial \varphi}{\partial y} \frac{\partial \lambda}{\partial y^{\prime}}+\frac{\partial \varphi}{\partial z} \frac{\partial \lambda}{\partial z^{\prime}}=0 .
$$

Now, if one differentiates the relation (8) with respect to $x^{\prime}, y^{\prime}, z^{\prime}$, in succession, while observing that $X^{\prime}, Y^{\prime}, Z^{\prime}$ do not depend upon those variables, then that will give:

$$
\frac{\partial \lambda}{\partial x}\left[\left(\frac{\partial \varphi}{\partial x}\right)^{2}+\left(\frac{\partial \varphi}{\partial y}\right)^{2}+\left(\frac{\partial \varphi}{\partial z}\right)^{2}\right]+2\left[\frac{\partial^{2} \varphi}{\partial x^{2}} x^{\prime}+\frac{\partial^{2} \varphi}{\partial x \partial y} y^{\prime}+\frac{\partial^{2} \varphi}{\partial x \partial z} z^{\prime}+\frac{\partial^{2} \varphi}{\partial x \partial t}\right]=0
$$

or rather:

$$
R \frac{\partial \lambda}{\partial x^{\prime}}+2 \frac{\partial}{\partial x}\left(\frac{d \varphi}{d t}\right)=0
$$

and similarly:

$$
\begin{aligned}
& R \frac{\partial \lambda}{\partial y^{\prime}}+2 \frac{\partial}{\partial y}\left(\frac{d \varphi}{d t}\right)=0, \\
& R \frac{\partial \lambda}{\partial z^{\prime}}+2 \frac{\partial}{\partial z}\left(\frac{d \varphi}{d t}\right)=0 .
\end{aligned}
$$

If one observes, as one had verified in the context of the Lagrange equations, that:

$$
\frac{\partial}{\partial x}\left(\frac{d \varphi}{d t}\right)=\frac{d}{d t}\left(\frac{\partial \varphi}{\partial x}\right)
$$

then one will deduce the following equality from the previous relations:

$$
\frac{\partial \varphi}{\partial x} \frac{\partial \lambda}{\partial x^{\prime}}+\frac{\partial \varphi}{\partial y} \frac{\partial \lambda}{\partial y^{\prime}}+\frac{\partial \varphi}{\partial z} \frac{\partial \lambda}{\partial z^{\prime}}+\frac{2}{R}\left[\frac{\partial \varphi}{\partial x} \frac{d}{d t}\left(\frac{\partial \varphi}{\partial x}\right)+\frac{\partial \varphi}{\partial y} \frac{d}{d t}\left(\frac{\partial \varphi}{\partial z}\right)+\frac{\partial \varphi}{\partial y} \frac{d}{d t}\left(\frac{\partial \varphi}{\partial z}\right)\right]=0
$$

or rather:

$$
\frac{\partial \varphi}{\partial x} \frac{\partial \lambda}{\partial x^{\prime}}+\frac{\partial \varphi}{\partial y} \frac{\partial \lambda}{\partial y^{\prime}}+\frac{\partial \varphi}{\partial z} \frac{\partial \lambda}{\partial z^{\prime}}+\frac{d R / d t}{R}=0
$$

That relation will be true if one supposes that $x, y, z, x^{\prime}, y^{\prime}, z^{\prime}$ are arbitrary functions of $t$ that satisfy equations (1). Consequently, $M=R$ will be a multiplier of (1).

With that, assume that one knows three first integrals of the motion of a point on the surface. Those integrals can always be put into the form:

$$
\left\{\begin{array}{l}
\psi_{1}\left(t, x, y, x^{\prime}, y^{\prime}\right)=\gamma_{1}  \tag{d}\\
\psi_{2}\left(t, x, y, x^{\prime}, y^{\prime}\right)=\gamma_{2} \\
\psi_{3}\left(t, x, y, x^{\prime}, y^{\prime}\right)=\gamma_{3}
\end{array}\right.
$$

One can infer $x^{\prime}, y^{\prime}, y$, for example, as functions of $x$ and $t$ from that relation and substitute them in the equation:

$$
x^{\prime} d t-d x=0
$$

That equation admits the quantity:

$$
\frac{R}{\left(\frac{\partial \varphi}{\partial z}\right)^{2} \Delta}
$$

as an integrating factor, in which $\Delta$ represents the determinant:

$$
\left|\begin{array}{lll}
\frac{\partial \psi_{1}}{\partial y} & \frac{\partial \psi_{1}}{\partial x^{\prime}} & \frac{\partial \psi_{1}}{\partial y^{\prime}} \\
\frac{\partial \psi_{2}}{\partial y} & \frac{\partial \psi_{2}}{\partial x^{\prime}} & \frac{\partial \psi_{2}}{\partial y^{\prime}} \\
\frac{\partial \psi_{3}}{\partial y} & \frac{\partial \psi_{3}}{\partial x^{\prime}} & \frac{\partial \psi_{3}}{\partial y^{\prime}}
\end{array}\right|
$$

and one supposes that $x^{\prime}, y^{\prime}, z^{\prime}, y$, and $z$ are expressed in terms of $x$ and $t$, according to (c) and (d).

When the surface and the active force do not depend upon time, $t$ will not appear in either $\varphi$ or equation (2), or in $R$ and $\lambda$, as a result. The multiplier $R /\left(\frac{\partial \varphi}{\partial z}\right)^{2}$ of the system (3) is also a multiplier of the system:

$$
\begin{equation*}
\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d x^{\prime}}{X^{\prime}+\lambda \frac{\partial \varphi}{\partial x}}=\frac{d y^{\prime}}{Y^{\prime}+\lambda \frac{\partial \varphi}{\partial y}} \tag{3'}
\end{equation*}
$$

into which $t$ will not enter. Consequently, in order to solve the problem by quadratures, it will suffice to know two first integrals of ( $3^{\prime}$ ).

We then arrive at the following conclusions: When a point moves without friction on a surface and is subject to a given force that does not depend upon the velocity of the point:

1. If the surface varies with time then in order for the motion to be determined by quadratures, it will suffice to know three first integrals of the motion.
2. If neither the surface nor the given force depends upon time then it will suffice to know two first integrals in which $t$ does not figure.
II. Motion of an arbitrary system. - Let $\Sigma$ be a system of $n$ material points that are subject to $p$ frictionless constraints. The coordinates $\left(x_{i}, y_{i}, z_{i}\right)$ of those points are restricted to verify the $p$ equations of constraint:
(A)

$$
\begin{aligned}
& \varphi\left(x_{1}, y_{1}, z_{1}, \ldots, x_{i}, y_{i}, z_{i}, \ldots, x_{n}, y_{n}, z_{n}, t\right)=0, \\
& \psi\left(x_{1}, y_{1}, z_{1}, \ldots, x_{i}, y_{i}, z_{i}, \ldots, x_{n}, y_{n}, z_{n}, t\right)=0 \text {, } \\
& \chi\left(x_{1}, y_{1}, z_{1}, \ldots, x_{i}, y_{i}, z_{i}, \ldots, x_{n}, y_{n}, z_{n}, t\right)=0,
\end{aligned}
$$

One can infer $p$ of the $3 n$ quantities $x_{i}, y_{i}, z_{i}$ from those $p$ equations, for example, $x_{n}, y_{n}, z_{n}, z_{n-1}$, $\ldots$ as functions of the $3 n-p$ other ones $x_{1}, y_{1}, z_{1}, x_{2}, \ldots$, and $t$. Suppose, for more clarity, that this solution has been performed and write the equations of constraint in the form:

$$
\begin{align*}
& \varphi \equiv z_{n}-\varphi_{1}\left(x_{1}, y_{1}, z_{1}, x_{2}, \ldots, t\right)=0 \\
& \psi \equiv y_{n}-\psi_{1}\left(x_{1}, y_{1}, z_{1}, x_{2}, \ldots, t\right)=0 \\
& \chi \equiv x_{n}-\chi_{1}\left(x_{1}, y_{1}, z_{1}, x_{2}, \ldots, t\right)=0
\end{align*}
$$

$\qquad$
$\qquad$
then the functions $\varphi_{1}, \psi_{1}, \chi_{1}, \ldots$ depend upon only the $(3 n-p)$ variables $x_{1}, y_{1}, z_{1}, x_{2}, \ldots$, and $t$.
If $X_{i}, Y_{i}, Z_{i}$ are the components of the active force that is exerted on the point $M_{1}$ then the motion of the system will be determined by the equations:

$$
\left\{\begin{array}{l}
m_{i} \frac{d^{2} x_{i}}{d t^{2}}=X_{i}^{\prime}+\lambda \frac{\partial \varphi}{\partial x_{i}}+\mu \frac{\partial \psi}{\partial x_{i}}+v \frac{\partial \chi}{\partial x_{i}}+\cdots=X_{i}, \\
m_{i} \frac{d^{2} y_{i}}{d t^{2}}=Y_{i}^{\prime}+\lambda \frac{\partial \varphi}{\partial y_{i}}+\mu \frac{\partial \psi}{\partial y_{i}}+v \frac{\partial \chi}{\partial y_{i}}+\cdots=Y_{i}, \quad(i=1,2, \ldots, n)  \tag{1}\\
m_{i} \frac{d^{2} z_{i}}{d t^{2}}=Z_{i}^{\prime}+\lambda \frac{\partial \varphi}{\partial z_{i}}+\mu \frac{\partial \psi}{\partial z_{i}}+v \frac{\partial \chi}{\partial z_{i}}+\cdots=Z_{i} .
\end{array}\right.
$$

One knows that one can express the coefficients $\lambda, \mu, \nu, \ldots$ as functions of $x_{i}, y_{i}, z_{i}, x_{i}^{\prime}, y_{i}^{\prime}, z_{i}^{\prime}$, and $t$. It will suffice to differentiate the equations of constraint twice with respect to $t$. Upon taking equations (1) into account, one will thus obtain $p$ equations such as the following:
(B)

$$
\left\{\begin{array}{c}
\sum_{i=1}^{n} \frac{1}{m_{i}}\left[\frac{\partial \varphi}{\partial x_{i}}\left(X_{i}^{\prime}+\lambda \frac{\partial \varphi}{\partial x_{i}}+\mu \frac{\partial \psi}{\partial x_{i}}+\cdots\right)+\frac{\partial \varphi}{\partial y_{i}}\left(Y_{i}^{\prime}+\lambda \frac{\partial \varphi}{\partial y_{i}}+\mu \frac{\partial \psi}{\partial y_{i}}+\cdots\right)+\frac{\partial \varphi}{\partial z_{i}}\left(Z_{i}^{\prime}+\lambda \frac{\partial \varphi}{\partial z_{i}}+\mu \frac{\partial \psi}{\partial z_{i}}+\cdots\right)\right] \\
+\left[\frac{\partial \varphi}{\partial t}+\sum_{i=1}^{n}\left(\frac{\partial \varphi}{\partial x_{i}} x_{i}^{\prime}+\frac{\partial \varphi}{\partial y_{i}} y_{i}^{\prime}+\frac{\partial \psi}{\partial y_{i}} z_{i}^{\prime}\right)\right]_{2}=0 .
\end{array}\right.
$$

As we have shown, those $p$ relations are soluble for $\lambda, \mu, v, \ldots$, and if one replaces $\lambda, \mu, v, \ldots$ with those values in equations (1) then the system thus-obtained, which is equivalent to the system:

$$
\begin{equation*}
d t=\frac{d x_{i}}{x_{i}^{\prime}}=\frac{d y_{i}}{y_{i}^{\prime}}=\frac{d z_{i}}{z_{i}^{\prime}}=\frac{m_{i} d x_{i}^{\prime}}{X_{i}}=\frac{m_{i} d y_{i}^{\prime}}{Y_{i}}=\frac{m_{i} d z_{i}^{\prime}}{Z_{i}} \tag{2}
\end{equation*}
$$

will admit the $2 p$ first integrals:
(a)

$$
\begin{aligned}
\varphi^{\prime}=\frac{\partial \varphi}{\partial t}+\sum_{i=1}^{n}\left(\frac{\partial \varphi}{\partial x_{i}} x_{i}^{\prime}+\frac{\partial \varphi}{\partial y_{i}} y_{i}^{\prime}+\frac{\partial \varphi}{\partial z_{i}} z_{i}^{\prime}\right) & =\alpha^{\prime}, \\
\varphi-t \varphi^{\prime} & =\alpha, \\
\psi^{\prime}=\frac{\partial \psi}{\partial t}+\sum_{i=1}^{n}\left(\frac{\partial \psi}{\partial x_{i}} x_{i}^{\prime}+\frac{\partial \psi}{\partial y_{i}} y_{i}^{\prime}+\frac{\partial \psi}{\partial z_{i}} z_{i}^{\prime}\right) & =\beta^{\prime}, \\
\psi-t \psi^{\prime} & =\beta, \quad \text { etc. }
\end{aligned}
$$

If the equations of constraint have been put into the form $\left(A^{\prime}\right)$ then the first integral (a) will contain only $z_{n}^{\prime}$, the second one will contain only $z_{n}$, the third one will contain only $y_{n}^{\prime}$, etc. One
then infers $z_{n}^{\prime}, z_{n}, y_{n}^{\prime}, \ldots$ from equations (a) while setting $\alpha^{\prime}=\alpha=\beta=\ldots=0$ and then substitute that in equations (2), in which we suppress the $2 p$ equations that include $d z_{n}^{\prime}, d z_{n}, d y_{n}^{\prime}$, etc. We thus form a system of (3) of $(6 n-2 p)$ first-order equations between $t$ and the $(6 n-2 p)$ variables $x_{1}, x_{1}^{\prime}, y_{1}, y_{1}^{\prime}$, etc. If $M$ is a multiplier of (2) then the expression $M / \frac{\partial \varphi^{\prime}}{\partial z_{n}^{\prime}} \frac{\partial \varphi}{\partial z_{n}} \frac{\partial \psi^{\prime}}{\partial y_{n}^{\prime}} \cdots$ (i.e., $M$ here) will also be a multiplier of the system (3), on the condition that we replace $z_{n}^{\prime}, z_{n}, y_{n}^{\prime}, \ldots$ with their values that we infer from $(a)$.

It is even clear, moreover, that if the relations (a) are not solved with respect to $p$ of the variables then one can again deduce a multiplier of (3) from $M$.

Finally, if the constraints and the given forces are independent of time then equations ( $3^{\prime}$ ) will be obtained by suppressing the first equation (in $d t$ ) in the system (3), which does not include $t$ because $t$ does not appear in either the $\lambda, \mu, v, \ldots$ or the $X_{i}^{\prime}, Y_{i}^{\prime}, Z_{i}^{\prime}$. Any multiplier of (3) that is independent of $t$ will then be a multiplier ( $3^{\prime}$ ).

Having said that, we shall show that we will always know a multiplier of equations (1) when the given forces do not depend upon velocities.

The equation for the multipliers of the system (1) can then be written:

$$
\begin{equation*}
\frac{d}{d t} \log M+\sum_{j=1}^{n}\left(\frac{\partial \lambda}{\partial x_{j}^{\prime}} \frac{\partial \varphi}{\partial x_{j}}+\frac{\partial \lambda}{\partial y_{j}^{\prime}} \frac{\partial \varphi}{\partial y_{j}}+\frac{\partial \lambda}{\partial z_{j}^{\prime}} \frac{\partial \varphi}{\partial z_{j}}\right)+\sum_{j=1}^{n}\left(\frac{\partial \mu}{\partial x_{j}^{\prime}} \frac{\partial \psi}{\partial x_{j}}+\frac{\partial \mu}{\partial y_{j}^{\prime}} \frac{\partial \psi}{\partial y_{j}}+\frac{\partial \mu}{\partial z_{j}^{\prime}} \frac{\partial \psi}{\partial z_{j}}\right)+\cdots=0 . \tag{4}
\end{equation*}
$$

One can replace $\frac{\partial \lambda}{\partial x_{j}^{\prime}}, \frac{\partial \lambda}{\partial y_{j}^{\prime}}, \ldots$ with their values that are deduced from equations $(B)$ in the latter equation. Now differentiate those equations with respect to $x_{j}^{\prime}, y_{j}^{\prime}, z_{j}^{\prime}$ upon setting:

$$
\begin{aligned}
& (\varphi, \varphi)=\sum_{i=1}^{n}\left[\left(\frac{\partial \varphi}{\partial x_{i}}\right)^{2}+\left(\frac{\partial \varphi}{\partial y_{i}}\right)^{2}+\left(\frac{\partial \varphi}{\partial z_{i}}\right)^{2}\right] \\
& (\varphi, \psi)=\sum_{i=1}^{n}\left(\frac{\partial \varphi}{\partial x_{i}} \frac{\partial \psi}{\partial x_{i}}+\frac{\partial \varphi}{\partial y_{i}} \frac{\partial \psi}{\partial y_{i}}+\frac{\partial \varphi}{\partial z_{i}} \frac{\partial \psi}{\partial z_{i}}\right), \quad \text { etc., }
\end{aligned}
$$

to abbreviate. That will give:

$$
\frac{\partial \lambda}{\partial x_{j}^{\prime}}(\varphi, \varphi)+\frac{\partial \lambda}{\partial x_{j}^{\prime}}(\varphi, \psi)+\frac{\partial \lambda}{\partial x_{j}^{\prime}}(\varphi, \chi)+\cdots+2 \frac{d}{d t}\left(\frac{\partial \varphi}{\partial x_{j}}\right)=0,
$$

or upon observing that $\frac{d}{d t}\left(\frac{\partial \varphi}{\partial x_{j}}\right)=\frac{\partial}{\partial x_{j}}\left(\frac{d \varphi}{d t}\right)$ and setting $\frac{d \varphi}{d t}=\varphi^{\prime}$ :

$$
\begin{equation*}
\frac{\partial \lambda}{\partial x_{j}^{\prime}}(\varphi, \varphi)+\frac{\partial \mu}{\partial x_{j}^{\prime}}(\varphi, \psi)+\frac{\partial v}{\partial x_{j}^{\prime}}(\varphi, \chi)+\cdots+2 \frac{\partial \varphi^{\prime}}{\partial x_{j}}=0 . \tag{5}
\end{equation*}
$$

Similarly :

$$
\frac{\partial \lambda}{\partial x_{j}^{\prime}}(\psi, \varphi)+\frac{\partial \mu}{\partial x_{j}^{\prime}}(\psi, \varphi)+\frac{\partial v}{\partial x_{j}^{\prime}}(\psi, \chi)+\cdots+2 \frac{\partial \psi^{\prime}}{\partial x_{j}}=0 .
$$

Consider the determinant:

$$
R=\left|\begin{array}{ccccc}
(\varphi, \varphi) & (\varphi, \psi) & (\varphi, \chi) & \cdots & \cdots \\
(\psi, \varphi) & (\psi, \psi) & (\psi, \chi) & \cdots & \cdots \\
(\chi, \varphi) & (\chi, \psi) & (\chi, \chi) & \cdots & \cdots \\
\vdots & \vdots & \vdots & \ddots & \\
\vdots & \vdots & \vdots & & \ddots
\end{array}\right|
$$

I say that this determinant satisfies equation (4), i.e., that it is a multiplier of (1). In order to verify that, solve the relation (5). One will find that:

$$
\frac{\partial \lambda}{\partial x_{j}^{\prime}}=-\frac{2}{R}\left[\frac{\partial \varphi^{\prime}}{\partial x_{j}} \cdot \frac{\partial R}{\partial(\varphi, \varphi)}+\frac{\partial \psi^{\prime}}{\partial x_{j}} \cdot \frac{\partial R}{\partial(\psi, \varphi)}+\frac{\partial \chi^{\prime}}{\partial x_{j}} \cdot \frac{\partial R}{\partial(\chi, \varphi)}+\cdots\right]
$$

Similarly:

$$
\begin{aligned}
& \frac{\partial \lambda}{\partial y_{j}^{\prime}}=-\frac{2}{R}\left[\frac{\partial \varphi^{\prime}}{\partial y_{j}} \cdot \frac{\partial R}{\partial(\varphi, \varphi)}+\frac{\partial \psi^{\prime}}{\partial y_{j}} \cdot \frac{\partial R}{\partial(\psi, \varphi)}+\frac{\partial \chi^{\prime}}{\partial y_{j}} \cdot \frac{\partial R}{\partial(\chi, \varphi)}+\cdots\right], \\
& \frac{\partial \lambda}{\partial z_{j}^{\prime}}=-\frac{2}{R}\left[\frac{\partial \varphi^{\prime}}{\partial z_{j}} \cdot \frac{\partial R}{\partial(\varphi, \varphi)}+\frac{\partial \psi^{\prime}}{\partial z_{j}} \cdot \frac{\partial R}{\partial(\psi, \varphi)}+\frac{\partial \chi^{\prime}}{\partial z_{j}} \cdot \frac{\partial R}{\partial(\chi, \varphi)}+\cdots\right] .
\end{aligned}
$$

As a result, one will get:

$$
\begin{aligned}
\sum_{j=1}^{n}\left(\frac{\partial \lambda}{\partial x_{j}^{\prime}} \frac{\partial \varphi}{\partial x_{j}}\right. & \left.+\frac{\partial \lambda}{\partial y_{j}^{\prime}} \frac{\partial \varphi}{\partial y_{j}}+\frac{\partial \lambda}{\partial z_{j}^{\prime}} \frac{\partial \varphi}{\partial z_{j}}\right)+\sum_{j=1}^{n}\left(\frac{\partial \mu}{\partial x_{j}^{\prime}} \frac{\partial \psi}{\partial x_{j}}+\frac{\partial \mu}{\partial y_{j}^{\prime}} \frac{\partial \psi}{\partial y_{j}}+\frac{\partial \mu}{\partial z_{j}^{\prime}} \frac{\partial \psi}{\partial z_{j}}\right)+\cdots \\
= & -\frac{2}{R}\left[\frac{\partial R}{\partial(\varphi, \varphi)} \sum\left(\frac{\partial \varphi^{\prime}}{\partial x_{j}} \frac{\partial \varphi}{\partial x_{j}}+\frac{\partial \varphi^{\prime}}{\partial y_{j}} \frac{\partial \varphi}{\partial y_{j}}+\frac{\partial \varphi^{\prime}}{\partial z_{j}} \frac{\partial \varphi}{\partial z_{j}}+\right)\right. \\
& \left.+\frac{\partial R}{\partial(\psi, \varphi)} \sum\left(\frac{\partial \psi^{\prime}}{\partial x_{j}} \frac{\partial \varphi}{\partial x_{j}}+\frac{\partial \psi^{\prime}}{\partial y_{j}} \frac{\partial \varphi}{\partial y_{j}}+\frac{\partial \psi^{\prime}}{\partial z_{j}} \frac{\partial \varphi}{\partial z_{j}}+\right)+\cdots\right]
\end{aligned}
$$

i.e.:

$$
\begin{aligned}
& =-\frac{1}{R}\left[\frac{\partial R}{\partial(\varphi, \varphi)} \frac{d(\varphi, \varphi)}{d t}+\frac{\partial R}{\partial(\psi, \psi)} \frac{d(\psi, \psi)}{d t}+\cdots+\frac{2 \partial R}{\partial(\varphi, \psi)} \frac{d(\varphi, \psi)}{d t}+\cdots\right] \\
& =-\frac{1}{R} \frac{d R}{d t}
\end{aligned}
$$

if one observes that $(\varphi, \psi)=(\psi, \varphi)$.
One can then put equation (4) into the form:

$$
\frac{d L M}{d t}-\frac{1}{R} \frac{d R}{d t}=0
$$

i.e., that $M=R$ is a multiplier of (1).

If the constraints do not depend upon time, any more than the given forces, then $t$ will not appear in $R$, which is also a multiplier of equation ( $3^{\prime}$ ).

One then has this theorem:

If the given forces that are exerted on a system of $n$ material points are subject to $p$ frictionless constraints and do not depend upon velocity then in order for the determination of the motion to be achieved by quadrature, it will suffice to know $(6 n-2 p-1)$ first integrals of the motion.

Moreover, when the constraints and the given forces are independent of time, it will suffice to consider $(6 n-2 p-2)$ first integrals into which $t$ does not enter.

The proof of that theorem will become much quicker if one applies the theory of the last multiplier to the canonical equations. That new proof can be extended to the theorem on continuous systems whose position depends upon a finite number of parameters, in addition, as we shall now prove.

## LECTURE 14

## APPLICATION OF THE THEORY OF THE LAST MULTIPLIER TO THE CANONICAL EQUATIONS.

Consider an arbitrary system of canonical equations:

$$
\begin{equation*}
d t=\frac{d q_{1}}{\frac{\partial K}{\partial p_{1}}}=\frac{d p_{1}}{\frac{\partial K}{\partial q_{1}}+Q_{1}}=\frac{d q_{2}}{\frac{\partial K}{\partial p_{2}}}=\ldots=\frac{d p_{k}}{\frac{\partial K}{\partial q_{k}}+Q_{k}} . \tag{1}
\end{equation*}
$$

The equations for the last multiplier $M$ relative to that system can be written as:

$$
\frac{d L M}{d t}+\frac{\partial Q_{1}}{\partial p_{1}}+\frac{\partial Q_{2}}{\partial p_{2}}+\cdots+\frac{\partial Q_{k}}{\partial p_{k}}=0,
$$

since $\frac{\partial^{2} K}{\partial q_{i} \partial p_{j}}-\frac{\partial^{2} K}{\partial p_{i} \partial q_{j}} \equiv 0$.
If the quantity $\sum \frac{\partial Q_{i}}{\partial p_{i}}$ is identically zero then $M=1$ will be a multiplier of equations (1). That will be true when the material system $\Sigma$ whose motion is determined by equations (1) is frictionless and the given forces do not depend upon velocity.

If one then knows $(2 k-1)$ distinct first integrals of (1) in this case then the last equation can be integrated by quadratures.

When neither the constraints nor the given forces depend upon time, moreover, $t$ will not figure in either $K$ or the $Q_{i}$. In order to achieve the integration by quadratures, it will suffice to know ( $2 k$ -2 ) first integrals in which $t$ does not enter.

Let us apply those generalities to the frictionless systems with constraints that are independent of time and whose position is defined by two parameters. If the given forces that are exerted upon such a system depend upon neither velocity nor time then if one is to calculate the motion by quadratures, it will suffice to know two integrals into which $t$ does not enter.

In particular, if the given forces admit a force function $U\left(q_{1}, q_{2}\right)$ then it will suffice to know a second integral that is distinct from that of vis viva and does not include $t$. An application of the theory of the last multiplier will then lead to some remarkable conclusions in this case that we shall now develop.

Write the canonical equations (while ignoring the first one):
(1')

$$
\frac{d q_{1}}{\frac{\partial H}{\partial p_{1}}}=\frac{d p_{1}}{\frac{\partial H}{\partial q_{1}}}=\frac{d q_{2}}{\frac{\partial H}{\partial p_{2}}}=\frac{d p_{2}}{\frac{\partial H}{\partial q_{2}}},
$$

in which:

$$
H=T-U .
$$

Let:

$$
\begin{equation*}
f\left(q_{1}, q_{2}, p_{1}, p_{2}\right)=\alpha \tag{2}
\end{equation*}
$$

be a first integral of $\left(1^{\prime}\right)$ that is distinct from the vis viva integral:

$$
\begin{equation*}
H\left(q_{1}, q_{2}, p_{1}, p_{2}\right)=h . \tag{3}
\end{equation*}
$$

Infer $p_{1}$ and $p_{2}$ as functions of $q_{1}, q_{2}$ from equations (2) and (3):

$$
p_{1}=\varphi_{1}\left(q_{1}, q_{2}, \alpha, h\right), \quad p_{2}=\varphi_{2}\left(q_{1}, q_{2}, \alpha, h\right) .
$$

If we replace $p_{1}$ and $p_{2}$ everywhere with those values then the equation:

$$
\frac{\partial H}{\partial p_{1}} d q_{2}-\frac{\partial H}{\partial p_{2}} d q_{1}=0
$$

will admit:

$$
\frac{1}{\delta}=\frac{1}{\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial p_{1}}}
$$

as an integrating factor. In other words, one has:

$$
\begin{equation*}
\frac{\frac{\partial H}{\partial p_{1}} d q_{2}-\frac{\partial H}{\partial p_{2}} d q_{1}}{\delta}=d \cdot F\left(q_{1}, q_{2}, \alpha, h\right) \tag{4}
\end{equation*}
$$

I say that the left-hand side of (4) coincides with:

$$
\frac{\partial p_{1}}{\partial \alpha} d q_{1}+\frac{\partial p_{2}}{\partial \alpha} d q_{2}
$$

Indeed, from the theory of functional determinants, one knows that:

$$
\delta=\frac{1}{\frac{\partial p_{1}}{\partial h} \frac{\partial p_{2}}{\partial \alpha}-\frac{\partial p_{1}}{\partial \alpha} \frac{\partial p_{2}}{\partial h}} .
$$

On the other hand, if one replaces $p_{1}, p_{2}$ with $\varphi_{1}, \varphi_{2}$ in (2) and (3) then those equations will be verified identically. Upon differentiating (3) with respect to $h$ and $\alpha$, one will find that:

$$
\begin{aligned}
& \frac{\partial H}{\partial p_{1}} \frac{\partial p_{1}}{\partial h}+\frac{\partial H}{\partial p_{2}} \frac{\partial p_{2}}{\partial h}=1, \\
& \frac{\partial H}{\partial p_{1}} \frac{\partial p_{1}}{\partial \alpha}+\frac{\partial H}{\partial p_{2}} \frac{\partial p_{2}}{\partial \alpha}=0 .
\end{aligned}
$$

Hence:

$$
\frac{\partial H}{\partial p_{1}}=\delta \frac{\partial p_{2}}{\partial \alpha}, \quad \frac{\partial H}{\partial p_{2}}=-\delta \frac{\partial p_{1}}{\partial \alpha} .
$$

The equation that remains to be integrated can then be written:

$$
\frac{\partial p_{1}}{\partial \alpha} d q_{1}+\frac{\partial p_{2}}{\partial \alpha} d q_{2}=0
$$

and one has assumed that its left-hand side is a total differential $d F\left(q_{1}, q_{2}, \alpha, h\right)$.
It is quite easy, moreover, to verify the conclusion to which we just arrived by direct calculation by showing that ( $p_{1} d q_{1}+p_{2} d q_{2}$ ) is an exact total differential.

In order for that to be true, it is necessary and sufficient that one should have:

$$
\frac{\partial \varphi_{1}}{\partial q_{2}} \equiv \frac{\partial \varphi_{2}}{\partial q_{1}} .
$$

Calculate $\partial p_{1} / \partial q_{2}$ and $\partial p_{2} / \partial q_{1}$ using equations (2) and (3). That will give:

$$
\begin{aligned}
& \frac{\partial H}{\partial q_{1}}+\frac{\partial H}{\partial p_{1}} \frac{\partial p_{1}}{\partial q_{1}}+\frac{\partial H}{\partial p_{2}} \frac{\partial p_{2}}{\partial q_{1}}=0 \\
& \frac{\partial f}{\partial q_{1}}+\frac{\partial f}{\partial p_{1}} \frac{\partial p_{1}}{\partial q_{1}}+\frac{\partial f}{\partial p_{2}} \frac{\partial p_{2}}{\partial q_{1}}=0 .
\end{aligned}
$$

Therefore, one infers that:

$$
\begin{equation*}
\left(\frac{\partial H}{\partial q_{1}} \frac{\partial f}{\partial p_{1}}-\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial q_{1}}\right)-\delta \frac{\partial p_{2}}{\partial q_{1}}=0 . \tag{5}
\end{equation*}
$$

One similarly finds that:

$$
\begin{equation*}
\left(\frac{\partial H}{\partial q_{2}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial q_{2}}\right)+\delta \frac{\partial p_{1}}{\partial q_{2}}=0 . \tag{6}
\end{equation*}
$$

Furthermore, the function $f$ satisfies the relation:

$$
\begin{equation*}
\frac{\partial H}{\partial q_{1}} \frac{\partial f}{\partial p_{1}}-\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial q_{1}}+\frac{\partial H}{\partial q_{2}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial q_{2}}=0 \tag{7}
\end{equation*}
$$

identically.
If we the add corresponding sides of (5) and (6) then that will give:

$$
\delta\left(\frac{\partial p_{2}}{\partial q_{1}}-\frac{\partial p_{1}}{\partial q_{2}}\right) \equiv 0 .
$$

In all of the foregoing, we have assumed that $\delta$ is not identically zero. Under that condition (which is always realized, as we will soon show), we see that the expression $p_{1} d q_{1}+p_{2} d q_{2}$ will be an exact total differential:

$$
\int p_{1} d q_{1}+p_{2} d q_{2}=W\left(q_{1}, q_{2}, \alpha, h\right)
$$

The integral of equation (4) is given by the equality $\partial W / \partial \alpha=\beta$.
Conversely, if a relation $f=\alpha$, combined with the equation $H=h$, determines functions $p_{1}, p_{2}$ of $\left(q_{1}, q_{2}, \alpha, h\right)$ such that $p_{1} d q_{1}+p_{2} d q_{2}$ is an exact differential then the function $f$ will verify equation (7) and will be a first integral of ( $1^{\prime}$ ), in addition.

From that, consider the partial differential equation:

$$
\begin{equation*}
H\left(q_{1}, q_{2}, \frac{\partial W}{\partial q_{1}}, \frac{\partial W}{\partial q_{2}}\right)=h \tag{8}
\end{equation*}
$$

and a second relation:

$$
f\left(q_{1}, q_{2}, \frac{\partial W}{\partial q_{1}}, \frac{\partial W}{\partial q_{2}}\right)=\alpha
$$

The necessary and sufficient condition for those two equations to admit a common integral $W$ ( $q_{1}$, $q_{2}, \alpha, h$ ) for each value of $\alpha$ and $h$ is that $f$ and $H$ should be coupled by equation (7); in other words, $f$ must be a first integral of $\left(1^{\prime}\right)$.

By definition, if one knows a first integral $f=\alpha$ of $\left(1^{\prime}\right)$ that is distinct from that of vis viva $H=$ $h$, and one can infer $p_{1}, p_{2}$ as functions of $q_{1}, q_{2}, \alpha, h$ using those two equations then the expression $p_{1} d q_{1}+p_{2} d q_{2}$ will be an exact differential:

$$
\int p_{1} d q_{1}+p_{2} d q_{2}=W\left(q_{1}, q_{2}, \alpha, h\right),
$$

and the motion of the system is determined by the equality:

$$
\frac{\partial W}{\partial \alpha}=\beta
$$

which defines $q_{2}$ as a function of $q_{1}$ and the three constants $h, \alpha, \beta$.
As for time $t$, it will be defined as a function of $q_{1}$ (for example, with the aid of a quadrature) when one expresses $q_{2}$ in terms of $q_{1}$. However, more symmetrically, it should be pointed out that $t$ satisfies the two equalities:

$$
d t=\frac{d q_{1}}{\frac{\partial H}{\partial p_{1}}}=\frac{d q_{2}}{\frac{\partial H}{\partial p_{2}}} .
$$

On the other hand, one has:

$$
\begin{aligned}
& \frac{\partial H}{\partial p_{1}} \frac{\partial p_{1}}{\partial h}+\frac{\partial H}{\partial p_{2}} \frac{\partial p_{2}}{\partial h}=1, \\
& \frac{\partial f}{\partial p_{1}} \frac{\partial p_{1}}{\partial h}+\frac{\partial f}{\partial p_{2}} \frac{\partial p_{2}}{\partial h}=0 .
\end{aligned}
$$

Hence:

$$
\frac{\partial p_{1}}{\partial h}=\frac{\partial f / \partial p_{2}}{\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial p_{1}}}, \quad \frac{\partial p_{2}}{\partial h}=\frac{-\partial f / \partial p_{1}}{\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial p_{1}}}
$$

One can then write:

$$
d t=\frac{\frac{\partial f}{\partial p_{2}} d q_{1}-\frac{\partial f}{\partial p_{1}} d q_{2}}{\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial p_{1}}}=\frac{\partial p_{1}}{\partial h} d q_{1}+\frac{\partial p_{2}}{\partial h} d q_{2},
$$

i.e.:

$$
t-t_{0}=\frac{\partial W}{\partial h}
$$

We will see in the next lecture that this theorem is a particular case of a theorem of Jacobi.

Remark. - All of the foregoing argument supposes essentially that we can solve the two equations $f=\alpha, H=h$ for $p_{1}, p_{2}$. We shall now show that this is always true. In the contrary case, the functional determinant $\delta=\left(\frac{\partial H}{\partial p_{1}} \frac{\partial f}{\partial p_{2}}-\frac{\partial H}{\partial p_{2}} \frac{\partial f}{\partial p_{1}}\right)$ will be identically zero, and there will exist a relation of the form:

$$
\Phi\left(H, f, q_{1}, q_{2}\right)=0
$$

between $f$ and $H$.
One can infer $f$ from that equation. Otherwise, $H$ would be a simple function of $q_{1}, q_{2}$, or rather $q_{2}$ would be expressible as a function of $q_{1}$ with no arbitrary constants. We then write the relation thus:

$$
f=\psi\left(H, q_{1}, q_{2}\right) .
$$

$\psi$ will define upon at least one of the variables $q_{1}, q_{2}$, since the integral $f=\alpha$ is distinct from that of vis viva, by hypothesis.

On the other hand, since $\psi=\alpha$ is an integral, the function $\psi$ must verify the equality:

$$
\frac{\partial \psi}{\partial H}(H, H)+\frac{\partial \psi}{\partial q_{1}} \frac{\partial H}{\partial p_{1}}+\frac{\partial \psi}{\partial q_{2}} \frac{\partial H}{\partial p_{2}} \equiv 0,
$$

or

$$
\frac{\partial \psi}{\partial q_{1}} \frac{\partial H}{\partial p_{1}}+\frac{\partial \psi}{\partial q_{2}} \frac{\partial H}{\partial p_{2}} \equiv 0
$$

or finally:

$$
\begin{equation*}
\frac{\frac{\partial \psi}{\partial q_{1}}}{\frac{\partial \psi}{\partial q_{2}}} \equiv-\frac{\frac{\partial H}{\partial p_{2}}}{\frac{\partial H}{\partial p_{1}}} \tag{9}
\end{equation*}
$$

$\frac{\partial H}{\partial p_{1}}$ and $\frac{\partial H}{\partial p_{2}}$ are homogeneous linear forms in $p_{1}, p_{2}$ whose determinant is not zero, as one knows.
The right-hand side of (9) is therefore a function of $p_{1} / p_{2}$, say $\chi\left(p_{1} / p_{2}, q_{1}, q_{2}\right)$. As for the lefthand side, it will coincide with $\chi$ only if it depends upon $p_{1}, p_{2}$, and as a result, upon $H$. One then infers from equation (9) that:

$$
H=F\left(\frac{p_{1}}{p_{2}}, q_{1}, q_{2}\right)
$$

which is an absurd equality, since $H$ has degree two in $p_{1}, p_{2}$. We thus arrive at the conclusion that one can always solve the distinct integrals $H=h, f=\alpha$ for $p_{1}, p_{2}$.

We shall now apply the preceding theorems to some particular examples.
I. Motion of a massive point that moves without friction on a paraboloid with a vertical axis. - Let us first apply the theory of the last multiplier to the equations of motion in the first form that Lagrange gave to them.

If the paraboloid is defined by the relation:

$$
\begin{equation*}
\frac{x^{2}}{\alpha}+\frac{y^{2}}{\beta}-2 z=0 \tag{A}
\end{equation*}
$$

then the equations of motion will be:
(B)

$$
\left\{\begin{array}{l}
x^{\prime \prime}=\lambda \frac{x}{\alpha} \\
y^{\prime \prime}=\lambda \frac{y}{\beta} \\
z^{\prime \prime}=-\lambda+g
\end{array}\right.
$$

Differentiate the relation $(A)$ twice. Upon taking $(B)$ into account, that will give:

$$
\lambda\left(\frac{x^{2}}{\alpha^{2}}+\frac{y^{2}}{\beta^{2}}+1\right)+\frac{x^{\prime 2}}{\alpha}+\frac{y^{\prime 2}}{\beta}-g=0 .
$$

If one substitutes that value of $\lambda$ in the first two equations $(B)$ and eliminates $z$ and $z^{\prime}$ using (A) then the two equations thus-obtained will admit the expression $\left(\frac{x^{2}}{\alpha^{2}}+\frac{y^{2}}{\beta^{2}}+1\right)$ as a last multiplier. Since $z$ and $z^{\prime}$ do not figure in $\lambda$, the elimination is found to take place because of that fact in its own right, and one immediately obtains:
( $B^{\prime}$ )

$$
\left\{\begin{array}{l}
\frac{x^{\prime \prime}}{\frac{x^{\prime 2}}{\alpha}+\frac{y^{\prime 2}}{\beta}-g}=-\frac{x}{\alpha} \frac{1}{\left(\frac{x^{2}}{\alpha^{2}}+\frac{y^{2}}{\beta^{2}}+1\right)} \\
\frac{y^{\prime \prime}}{\frac{x^{\prime 2}}{\alpha}+\frac{y^{\prime 2}}{\beta}-g}=-\frac{y}{\beta} \frac{1}{\left(\frac{x^{2}}{\alpha^{2}}+\frac{y^{2}}{\beta^{2}}+1\right)}
\end{array}\right.
$$

On the other hand, one knows one first integral of those equations, namely, the vis viva integral:
(a)

$$
2 T \equiv x^{\prime 2}\left(1+\frac{x^{2}}{\alpha^{2}}\right)+y^{\prime 2}\left(1+\frac{y^{2}}{\beta^{2}}\right)+\frac{2 x y}{\alpha \beta} x^{\prime} y^{\prime}=g\left(\frac{x^{2}}{\alpha}+\frac{y^{2}}{\beta}\right)+h .
$$

Now multiply the first equation in $\left(B^{\prime}\right)$ by $x^{\prime} / \alpha$, the second one by $y^{\prime} / \beta$, and add them. Upon integrating, that will give:

$$
\begin{equation*}
f=\left(\frac{x^{\prime 2}}{\alpha}+\frac{y^{\prime 2}}{\beta}-g\right)\left(1+\frac{x^{2}}{\alpha^{2}}+\frac{y^{2}}{\beta^{2}}\right)=C . \tag{b}
\end{equation*}
$$

Let $\delta$ denote the functional determinant of $T$ and $f$ with respect to the variables $x^{\prime}, y^{\prime}$. If we infer $x^{\prime}, y^{\prime}$ from equations $(a)$ and $(b)$ then the expression:

$$
\begin{equation*}
\left(\frac{x^{2}}{\alpha^{2}}+\frac{y^{2}}{\beta^{2}}+1\right) \frac{y^{\prime} d x-x^{\prime} d y}{\delta} \tag{c}
\end{equation*}
$$

will be an exact total differential $d \cdot F(x, y, h, C)$. The motion is defined by the equality $F=$ const.
If we let $\mu$ and $v$ denote the two values of the function:

$$
-\left(\alpha+\beta+\frac{x^{2}}{\alpha}+\frac{y^{2}}{\beta}\right) \pm \sqrt{\left(\alpha+\beta+\frac{x^{2}}{\alpha}+\frac{y^{2}}{\beta}\right)^{2}-4\left(\frac{\beta x^{2}}{\alpha}+\frac{\alpha y^{2}}{\beta}+\alpha \beta\right)}
$$

then a laborious calculation will verify that the expression $(c)$ is effectively equal to:

$$
\sqrt{\frac{\mu}{(\alpha+\mu)(\beta+\mu)\left(g \mu^{2}+h \mu+C\right)}} d \mu+\sqrt{\frac{v}{(\alpha+v)(\beta+v)\left(g v^{2}+h v+C\right)}} d v=d F
$$

However, we shall arrive at that result more easily in what follows by appealing to different coordinates.

Thus, here is a case in which Jacobi's theory will show that the motion can certainly be calculated by quadratures, although it would be quite difficult to show that directly, at least with the variables that are employed. Furthermore, one will arrive at the same result by appealing to the canonical equations. If one substitutes the variables $p_{1}$ and $p_{2}$ that are defined by the equalities:

$$
\begin{aligned}
& p_{1}=x^{\prime}\left(1+\frac{x^{2}}{\alpha^{2}}\right)+y^{\prime} \frac{x y}{\alpha \beta}, \\
& p_{2}=x^{\prime} \frac{x y}{\alpha \beta}+y^{\prime}\left(1+\frac{y^{2}}{\beta^{2}}\right)
\end{aligned}
$$

for the variables $x^{\prime}, y^{\prime}$ and in $(a)$ and $(b)$ then one will effortlessly see that the expression $\left(\frac{\partial p_{1}}{\partial ?} d x+\frac{\partial p_{2}}{\partial ?} d y\right)$ coincides with the expression (c).

We similarly propose to study the motion of a point $M$ that moves without friction on an ellipsoid and is attracted to the center of the ellipsoid in proportion to its distance from it.

If the ellipsoid is defined by the relation:

$$
\frac{x^{2}}{a^{2}}+\frac{y^{2}}{b^{2}}+\frac{z^{2}}{c^{2}}=1
$$

then the equations of motion can be written:

$$
\left\{\begin{array}{l}
x^{\prime \prime}=-k x+\lambda \frac{x}{a^{2}},  \tag{B}\\
y^{\prime \prime}=-k y+\lambda \frac{y}{b^{2}}, \\
z^{\prime \prime}=-k z+\lambda \frac{z}{c^{2}}
\end{array}\right.
$$

In addition to the vis viva integral, one defines a second integral in the following manner: Add corresponding sides of equations $(B)$ after multiplying them by $\frac{x^{\prime}}{a^{2}}, \frac{y^{\prime}}{b^{2}}, \frac{z^{\prime}}{c^{2}}$, respectively, and then by $\frac{x}{a^{2}}, \frac{y}{b^{2}}, \frac{z}{c^{2}}$, resp. Upon taking the relations:

$$
\frac{x^{\prime} x}{a^{2}}+\frac{y^{\prime} y}{b^{2}}+\frac{z^{\prime} z}{c^{2}}=0, \quad \frac{x^{\prime \prime} x}{a^{2}}+\frac{y^{\prime \prime} y}{b^{2}}+\frac{z^{\prime \prime} z}{c^{2}}+\frac{x^{\prime 2}}{a^{2}}+\frac{y^{\prime 2}}{b^{2}}+\frac{z^{\prime 2}}{c^{2}}=0
$$

into account, one will find that:

$$
\frac{x^{\prime} x^{\prime \prime}}{a^{2}}+\frac{y^{\prime} y^{\prime \prime}}{b^{2}}+\frac{z^{\prime} z^{\prime \prime}}{c^{2}}=\lambda\left[\frac{x x^{\prime}}{a^{4}}+\frac{y y^{\prime}}{b^{4}}+\frac{z z^{\prime}}{c^{4}}\right]
$$

and

$$
-\left(\frac{x^{\prime 2}}{a^{2}}+\frac{y^{\prime 2}}{b^{2}}+\frac{z^{\prime 2}}{c^{2}}\right)=-k+\lambda\left[\frac{x^{2}}{a^{4}}+\frac{y^{2}}{b^{4}}+\frac{z^{2}}{c^{4}}\right] .
$$

One will then deduce immediately that:

$$
\frac{-\frac{d}{d t}\left(\frac{x^{\prime 2}}{a^{2}}+\frac{y^{\prime 2}}{b^{2}}+\frac{z^{\prime 2}}{c^{2}}-k\right)}{\frac{x^{\prime 2}}{a^{2}}+\frac{y^{\prime 2}}{b^{2}}+\frac{z^{\prime 2}}{c^{2}}-k}=\frac{\frac{d}{d t}\left(\frac{x^{2}}{a^{4}}+\frac{y^{2}}{b^{4}}+\frac{z^{2}}{c^{4}}\right)}{\frac{x^{2}}{a^{4}}+\frac{y^{2}}{b^{4}}+\frac{z^{2}}{c^{4}}},
$$

and finally:

$$
\left(\frac{x^{2}}{a^{4}}+\frac{y^{2}}{b^{4}}+\frac{z^{2}}{c^{4}}\right)\left(\frac{x^{\prime 2}}{a^{2}}+\frac{y^{\prime 2}}{b^{2}}+\frac{z^{\prime 2}}{c^{2}}-k\right)=C .
$$

We are then certain that the problem can be achieved by quadratures. However, those quadratures can be carried out comfortably only with the aid of conveniently-chosen coordinates, which are the elliptic coordinates that we shall introduce later on.

## II. - Two points $M$ and $M_{1}$ are constrained to slide without friction on two helices:

$x=R \cos \theta, \quad y=R \sin \theta, \quad z=K \theta \quad$ and $\quad x_{1}=R_{1} \cos \theta_{1}, \quad y_{1}=R_{1} \sin \theta_{1}, \quad z_{1}=K \theta_{1}$.
The two points repel in proportion to the distance between them. Find the motion of the system.

We know how to form (see page 106) two first integrals of the motion, namely:

$$
2 T=m\left(R^{2}+K^{2}\right) \theta^{\prime 2}+m_{1}\left(R_{1}^{2}+K^{2}\right) \theta_{1}^{\prime 2}=-2 \mu R R_{1} \cos \left(\theta-\theta_{1}\right)+\mu K^{2}\left(\theta_{1}-\theta\right)^{2}+K
$$

and

$$
m\left(R^{2}+K^{2}\right) \theta^{\prime}+m_{1}\left(R_{1}^{2}+K^{2}\right) \theta_{1}^{\prime}=C
$$

Replace $\theta^{\prime}, \theta_{1}^{\prime}$ with the canonical variables:

$$
p=m\left(R^{2}+K^{2}\right) \theta^{\prime}, \quad p_{1}=m_{1}\left(R_{1}^{2}+K^{2}\right) \theta_{1}^{\prime} .
$$

Upon setting:

$$
\alpha=\frac{1}{m\left(R^{2}+K^{2}\right)}+\frac{1}{m_{1}\left(R_{1}^{2}+K^{2}\right)}, \quad \beta=m m_{1}\left(R^{2}+K^{2}\right)\left(R_{1}^{2}+K^{2}\right),
$$

and

$$
U=-2 \mu R R_{1} \cos \left(\theta-\theta_{1}\right)+\mu K^{2}\left(\theta_{1}-\theta\right)^{2}
$$

we will infer the following values of $p$ and $p_{1}$ from the two integrals:

$$
\begin{aligned}
& \alpha p_{1}=\frac{C}{m\left(R^{2}+K^{2}\right)}+\sqrt{\alpha(U+2 h)-\frac{C^{2}}{\beta}}, \\
& \alpha p=\frac{C}{m_{1}\left(R_{1}^{2}+K^{2}\right)}-\sqrt{\alpha(U+2 h)-\frac{C^{2}}{\beta}} .
\end{aligned}
$$

The expression ( $p d \theta+p_{1} d \theta_{1}$ ) is indeed an exact total differential $d W$. If one introduces the variable $\left(\theta_{1}-\theta\right)=\varphi$ then one will have:

$$
\alpha W=C\left[\frac{\theta}{m_{1}\left(R_{1}^{2}+K^{2}\right)}+\frac{\theta_{1}}{m\left(R^{2}+K^{2}\right)}\right]+\int_{\varphi_{0}}^{\varphi} \sqrt{\alpha(U+2 h)-\frac{C^{2}}{\beta}} d \varphi .
$$

The equalities $\partial W / \partial C=$ const., $t-t_{0}=\partial W / \partial h$, which determine the motion, are:

$$
\begin{gathered}
\frac{\theta}{m_{1}\left(R_{1}^{2}+K^{2}\right)}+\frac{\theta_{1}}{m\left(R^{2}+K^{2}\right)}+\frac{C}{\beta} \int_{\varphi_{0}}^{\varphi} \frac{d \varphi}{\sqrt{\alpha(U+2 h)-\frac{C^{2}}{\beta}}}=\text { const., } \\
t-t_{0}=\int_{\varphi_{0}}^{\varphi} \frac{d \varphi}{\sqrt{\alpha(U+2 h)-\frac{C^{2}}{\beta}}}
\end{gathered}
$$

here, which are equalities that can also be written:

$$
\begin{aligned}
& {\left[m\left(R^{2}+K^{2}\right)+m_{1}\left(R_{1}^{2}+K^{2}\right)\right] \theta=C t-m_{1}\left(R_{1}^{2}+K^{2}\right) \varphi+\text { const. },} \\
& {\left[m\left(R^{2}+K^{2}\right)+m_{1}\left(R_{1}^{2}+K^{2}\right)\right] \theta_{1}=C t+m\left(R^{2}+K^{2}\right) \varphi+\text { const. }}
\end{aligned}
$$

The problem depends upon only the single quadrature that gives $t$. We thus indeed recover the results that were already obtained on page 106.

## III. A point $M$ is attracted to the origin $O$ and released with zero initial velocity in the

 plane $x O y$. Find the motion of the point when the law of attraction has the form:$$
F=\mu r x^{m} y^{m} \quad\left(r=\sqrt{x^{2}+y^{2}}\right)
$$

The equations of motion are:

$$
\begin{aligned}
& x^{\prime \prime}=\mu x^{m+1} y^{m}, \\
& y^{\prime \prime}=\mu x^{m} y^{m+1}
\end{aligned}
$$

here.
Multiply the left-hand side by $y^{\prime}$, the right-hand side by $x^{\prime}$, and add them. That will give:

$$
d \cdot x^{\prime} y^{\prime}=\mu x^{m} y^{m}[x d y+y d x]
$$

or

$$
x^{\prime} y^{\prime}=\frac{\mu}{m+1} x^{m+1} y^{m+1}+a
$$

On the other hand, the area integral will give:

$$
x y^{\prime}-y x^{\prime}=C .
$$

One infers from this that:

$$
x^{\prime}=\frac{-C+\sqrt{R}}{2 y}, \quad y=\frac{C+\sqrt{R}}{2 x},
$$

when one sets:

$$
R=C^{2}+4 x y\left[\frac{\mu}{m+1} x^{m+1} y^{m+1}+\alpha\right] .
$$

From the theory of the last multiplier, the expression:

$$
\left(\frac{\partial x^{\prime}}{\partial C} \frac{\partial y^{\prime}}{\partial \alpha}-\frac{\partial x^{\prime}}{\partial \alpha} \frac{\partial y^{\prime}}{\partial C}\right)\left(y^{\prime} d x-x^{\prime} d y\right) \equiv \frac{y^{\prime} d x-x^{\prime} d y}{\sqrt{R}}
$$

is an exact total differential. One easily verifies that because one will have:

$$
\begin{gathered}
2 \frac{\left(y^{\prime} d x-x^{\prime} d y\right)}{\sqrt{R}}=\frac{d x}{x}-\frac{d y}{y}+\frac{C}{\sqrt{R}}\left(\frac{d x}{x}+\frac{d y}{y}\right) \\
=d \cdot L \frac{y}{x}+\frac{C d u}{u \sqrt{C^{2}+4 \cdot u\left(\frac{\mu}{m+1} u^{m+1}+\alpha\right)}}=d L \frac{y}{\alpha}+d F(u),
\end{gathered}
$$

upon setting $u=x y$. The motion is determined by the equality:

$$
y=\beta x e^{-F(x y)},
$$

to which one must append the relations:

$$
d t=\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{y d x+x d y}{x^{\prime} y+y^{\prime} x}=\frac{d u}{\sqrt{R(u)}}
$$

which is equivalent to the single relation:

$$
t=\int \frac{d u}{\sqrt{R(u)}}+\text { const. }
$$

Thus, here is an application of the theory to a case in which the vis viva theorem does not give an integral.
IV. Two massive points $M$ and $M_{1}$ of masses $m$ and $m_{1}$, resp., are constrained to slide without friction, one of them on the vertical $O z$ and the other on a cylinder of revolution around $O z$. The two points attract each other according to an arbitrary function of the distance. Motion of the system. -

Let $\theta$ be the angle that the plane $M O z$ makes with the $x O z$ plane. The position of the system depends upon three parameters $x, \theta$, and $z$. In order to solve the problem by quadratures, it will
suffice to know four first integrals in which $t$ does not appear, or rather, five integrals that depend upon $t$.

The theorem of moments relative to $O z$ gives:

$$
\theta=C t+C^{\prime}
$$

which is an equality that is equivalent to two integrals. On the other hand, from the theorem about the motion of the center of gravity, one will have:

$$
m z+m_{1} z_{1}=\frac{1}{2} g t^{2}+\alpha t+\beta
$$

If one combines those integrals with the vis viva integral:

$$
m_{1} z_{1}^{\prime 2}+m\left(R^{2} \theta^{\prime 2}+z^{\prime 2}\right)=2 g\left(m z+m_{1} z_{1}\right)+2 U\left(z-z_{1}\right)+\text { const } .
$$

then one will see that one knows five integrals in which $t$ appears. The motion is then calculated with the aid of quadratures.

Consider separately the four equations:

$$
\begin{equation*}
\frac{d z}{z^{\prime}}=\frac{d z_{1}}{z_{1}^{\prime}}=\frac{d z^{\prime}}{m g+\frac{\partial U}{\partial z}}=\frac{d z_{1}^{\prime}}{m_{1} g+\frac{\partial U}{\partial z_{1}}} \tag{A}
\end{equation*}
$$

Those equations are the equations of motion of a system whose vis viva is ( $m z^{\prime 2}+m_{1} z_{1}^{\prime 2}$ ) and is subjected to forces that admit the force function $\left[g\left(m z+m_{1} z_{1}\right)+U\right]$. In addition to the integral:

$$
\begin{equation*}
m z^{\prime 2}+m_{1} z_{1}^{\prime 2}=2 g\left(m z+m_{1} z_{1}\right)+2 U\left(z-z_{1}\right)+2 h, \tag{B}
\end{equation*}
$$

we know an integral of equations $(A)$, namely:

$$
\begin{equation*}
\frac{\left(m z^{\prime}+m_{1} z_{1}^{\prime}\right)^{2}}{m+m_{1}}=2 g\left(m z+m_{1} z_{1}\right)+a \tag{C}
\end{equation*}
$$

which is an integral that is deduced immediately from the foregoing. (It is the vis viva integral applied to the motion of the projection onto $O z$ of the center of gravity of the points $M, M_{1}$.)

On the other hand, if we remark that the canonical variables for equations $(A)$ are $p_{1}=m z^{\prime}, p_{2}$ $=m_{1} z_{1}^{\prime}$ then the theory that was developed above will show us that the expression:

$$
m z^{\prime} d z+m_{1} z_{1}^{\prime} d z_{1},
$$

in which we replace $z^{\prime}$ and $z_{1}^{\prime}$ with the values that we infer from $(R)$ and (C), must be an exact total differential $d W$. Indeed, a very simple calculation will give:

$$
\begin{aligned}
& m z^{\prime} d z+m_{1} z_{1}^{\prime} d z_{1} \\
& \begin{array}{c}
=\frac{1}{\sqrt{m+m_{1}}}\left[\sqrt{m m_{1}(2 U+2 h-\alpha)}\left(d z-d z_{1}\right)+\sqrt{2 g\left(m z+m_{1} z_{1}\right)+\alpha}\left(m d z+m_{1} d z_{1}\right)\right] \\
\quad=\frac{1}{\sqrt{m+m_{1}}}\left[\sqrt{m m_{1}[2 U(\eta)+2 h-\alpha]} d \eta+\sqrt{2 g \xi+\alpha} d \xi\right]=d W,
\end{array}
\end{aligned}
$$

if one sets $\left(z-z_{1}\right)=\eta$ and $m z+m_{1} z_{1}=\xi$.
The motion is then determined by the two equations:

$$
t=\frac{\partial W}{\partial h}+\text { const. }=\sqrt{\frac{m m_{1}}{m+m_{1}}} \int_{\eta_{1}}^{\eta} \frac{d \eta}{\sqrt{2 U(\eta)+2 h-\alpha}}+\text { const. }
$$

and

$$
\frac{d \xi}{\sqrt{2 g \xi+\alpha}}=\sqrt{\frac{m m_{1}}{2 U(\eta)+2 h-\alpha}} d \eta
$$

or rather

$$
\frac{1}{m+m_{1}}\left(\frac{d \xi}{d t}\right)^{2}=2 g \xi+\alpha
$$

The last equation is nothing but the integral ( $C$ ).
One will arrive at the result more quickly by immediately introducing the variables $\eta$ and $\xi$. Indeed, the vis viva of the $\operatorname{system}\left(M, M_{1}\right)$ is equal to:

$$
\frac{1}{m+m_{1}}\left(\eta^{\prime 2}+m m_{1} \xi^{\prime 2}\right)+m R^{2} \theta^{\prime 2}
$$

Since one has $\theta^{\prime}=\theta_{0}^{\prime}$ and $\frac{\xi^{\prime 2}}{m+m_{1}}=2 g \xi$, the vis viva theorem will imply the following equality:

$$
\frac{m m_{1}}{m+m_{1}} \eta^{\prime 2}=2 U+\text { const. }
$$

The motion is then determined by the three equalities:

$$
\theta=C t+C^{\prime}, \quad \frac{\xi}{m+m_{1}}=\frac{1}{2} g t^{2}+\alpha t+\beta
$$

and

$$
\sqrt{\frac{m+m_{1}}{m m_{1}}} d t=\frac{d \eta}{\sqrt{2 U(\eta)+2 h-\alpha}}
$$

It would be easy to apply the theory of the last multiplier to the examples that were treated in the previous chapters. However, this handful of exercises here will suffice to show the utility of that theory: In a great number of cases, it permits one to predict that the problem under study can be solved at one stroke by quadratures, even though the variables employed do not exhibit that fact clearly.

It is therefore appropriate to study the variables that give the simplest form to those quadratures.

## LECTURE 15

## JACOBI'S THEOREM. - LIOUVILLE'S THEOREM.

In the preceding lecture, we showed that if we know a second integral for a two-parameter mechanical problem, in addition to the vis viva integral, then the solution to the problem can be achieved by quadratures. That theorem is only a consequence of a general proposition by Jacobi that we shall now develop.

Let $S$ be a material system without friction whose constraints can depend upon time and whose position is defined by the parameters $q_{1}, q_{2}, \ldots, q_{k}$. If the given forces that are exerted on the system admit a force function $U\left(t, q_{1}, q_{2}, \ldots, q_{k}\right)$ then the canonical equations of motion of the system can be written:

$$
\left\{\begin{align*}
\frac{d q_{i}}{d t} & =\frac{\partial H}{\partial p_{i}}  \tag{1}\\
\frac{d p_{i}}{d t} & =-\frac{\partial H}{\partial q_{i}}
\end{align*} \quad(i=1,2, \ldots, k)\right.
$$

with

$$
H=K\left(t, q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}\right)-U\left(t, q_{1}, q_{2}, \ldots, q_{k}\right) .
$$

Replace the $p_{i}$ in $H$ with $\partial V / \partial q_{i}$ and consider the partial differential equation:

$$
\begin{equation*}
\frac{\partial V}{\partial t}+H\left(t, q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0 . \tag{2}
\end{equation*}
$$

It is a first-order equation that pertains to the function $V$ of the $(k+1)$ variables $t, q_{1}, q_{2}, \ldots, q_{k}$, and into which $V$ does not enter explicitly.

Jacobi showed that one can deduce the general integral of the canonical equations (1) from a complete integral of (2).

A complete integral $V$ of (2) is, by definition, an integral $V\left(t, q_{1}, q_{2}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{k}\right)$ that depends upon $k$ arbitrary constants $\alpha_{1}, \ldots, \alpha_{k}$ that permit one to attribute arbitrary values to the ( $k$ $+1)$ derivatives of $V$ for arbitrary values $t_{0}, q_{i}^{0}$ of $t$ and the $q_{i}$, resp., subject to only the condition that they must verify the relation (2).

Here, if the relation (2) is solved for the $\partial V / \partial q_{i}$ then the integral $V\left(t, q_{1}, q_{2}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{k}\right)$ will be a complete integral if one can select the $\alpha_{1}, \ldots, \alpha_{k}$ in such a manner as to give arbitrary
values to the derivatives $\left(\frac{\partial V}{\partial q_{1}}\right)_{0}, \ldots,\left(\frac{\partial V}{\partial q_{k}}\right)_{0} \cdot\left(\frac{\partial V}{\partial t}\right)_{0}$ will then take the value that is assigned to it by equation (2).

If one considers the equations:

$$
\left\{\begin{array}{l}
\frac{\partial V}{\partial q_{1}}=A_{1}  \tag{3}\\
\cdots \ldots \ldots \ldots \ldots \\
\frac{\partial V}{\partial q_{k}}=A_{k}
\end{array}\right.
$$

in which $A_{1}, \ldots, A_{k}$ are arbitrary constants, then one can solve that system (3) for the $\alpha_{1}, \ldots, \alpha_{k}$. Analytically, that amounts to saying that the functional determinant of the $k$ functions $\partial V / \partial q_{i}$ of $\alpha_{1}, \ldots, \alpha_{k}$ is not identically zero. That determinant is nothing but:

$$
\Delta=\left|\begin{array}{cccc}
\frac{\partial^{2} V}{\partial q_{1} \partial \alpha_{1}} & \frac{\partial^{2} V}{\partial q_{1} \partial \alpha_{2}} & \cdots & \frac{\partial^{2} V}{\partial q_{1} \partial \alpha_{k}} \\
\vdots & \vdots & \ddots & \vdots \\
\frac{\partial^{2} V}{\partial q_{k} \partial \alpha_{1}} & \frac{\partial^{2} V}{\partial q_{k} \partial \alpha_{2}} & \cdots & \frac{\partial^{2} V}{\partial q_{k} \partial \alpha_{k}}
\end{array}\right| .
$$

Hence, a complete integral is an integral that depends upon $k$ arbitrary constants such that $\Delta$ is non-zero. Having recalled that definition, I say that if one knows a complete integral to (2) then the canonical system (1) can be integrated from that fact alone.

Indeed, set:

$$
\left\{\begin{array}{l}
\frac{\partial V}{\partial \alpha_{1}}=\beta_{1}  \tag{4}\\
\ldots \ldots \ldots \ldots \ldots . . . . . \\
\frac{\partial V}{\partial \alpha_{k}}=\beta_{k}
\end{array}\right.
$$

$$
\left\{\begin{array}{l}
p_{1}=\frac{\partial V}{\partial q_{1}}  \tag{5}\\
\ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
p_{k}=\frac{\partial V}{\partial q_{k}}
\end{array}\right.
$$

The system (4-5), thus-formed, in which one considers the $\alpha$ and the $\beta$ to be $2 k$ arbitrary constants, will define the general integral $p_{1}(t), \ldots, p_{k}(t), q_{1}(t), \ldots, q_{k}(t)$ of equations (1).

First of all, the relations (4-5) determine the $p_{i}, q_{i}$ as functions of $t$ and $2 k$ distinct arbitrary constants.

Indeed, on the one hand, equations (4) can be solved for $q_{1}, q_{2}, \ldots, q_{k}$, because the functional determinant of those equations with respect to the $q_{1}, q_{2}, \ldots, q_{k}$ is nothing but $\Delta . q_{1}, q_{2}, \ldots, q_{k}$ are thus obtained as functions of $\alpha_{1}, \ldots, \alpha_{k}, \beta_{1}, \ldots, \beta_{k}$. If, on the other hand, one substitutes those values of $q_{i}$ in (5), then one will obtain $p_{1}, \ldots, p_{k}$ as functions of the same quantities.

On the other hand, those $2 k$ constants $\alpha_{i}, \beta_{i}$ are distinct. In other words, one can select them in such a way that one gives arbitrary values $q_{i}^{0}, p_{i}^{0}$ to $q_{i}, p_{i}$, resp., at $t=t_{0}$. Indeed, if one sets $t=t_{0}$, $q_{i}=q_{i}^{0}$ in equations (5) then one can select the $\alpha$ in such a fashion that one gives arbitrary values $p_{i}^{0}$ to the $p_{i}$. When the $\alpha$ are determined in that way, say $\alpha_{i}=\alpha_{i}^{0}$, the values of $\beta$ will be obtained immediately upon setting $t=t_{0}, q_{i}=q_{i}^{0}, \alpha_{i}=\alpha_{i}^{0}$ in (4).

From that, if one can prove that any system of functions $p_{i}(t), q_{i}(t)$ that verifies the relations (4-5) also verifies the equations (1), it will be clear that the relations (4-5) define the general integral of the canonical system.

In order to prove that this is true, calculate the derivatives $\frac{d p_{i}}{d t}, \frac{d q_{i}}{d t}$ using (4-5). That will give:



One can infer $d q_{i} / d t$ from (4') and $d p_{i} / d t$ from ( $5^{\prime}$ ). One must substitute those values in (1) and confirm that the relations thus-obtained are verified for any system $p_{i}(t), q_{i}(t)$ that satisfies the relations (4-5) [if $V$ is a complete integral of (2)]. Instead of doing that, it is legitimate for one to replace the $d q_{i} / d t, d p_{i} / d t$ in $\left(4^{\prime}\right)$ and ( $\left.5^{\prime}\right)$ with their values that are inferred from (1) and see if the conditions that are calculated in that way are consequences of equations (4-5).

First, make that substitution in (4'). One will find that:

$$
\begin{equation*}
\left\{\frac{\partial^{2} V}{\partial \alpha_{1} \partial t}+\frac{\partial^{2} V}{\partial \alpha_{1} \partial q_{1}} \frac{\partial H}{\partial p_{1}}+\frac{\partial^{2} V}{\partial \alpha_{1} \partial q_{2}} \frac{\partial H}{\partial p_{2}}+\cdots+\frac{\partial^{2} V}{\partial \alpha_{1} \partial q_{k}} \frac{\partial H}{\partial q_{k}}=0,\right. \tag{6}
\end{equation*}
$$

The equation that we just wrote down must be a consequence of the relations (4-5). By hypothesis, $V$ will satisfy equation (2) for any $t, q_{1}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{k} . V$ will then identically satisfy the equation that is obtained by differentiating (2) with respect to $\alpha_{1}$, i.e., the equation:

$$
\frac{\partial^{2} V}{\partial t \partial \alpha_{1}}+\frac{\partial H}{\partial\left(\frac{\partial V}{\partial q_{1}}\right)} \frac{\partial^{2} V}{\partial q_{1} \partial \alpha_{1}}+\frac{\partial H}{\partial\left(\frac{\partial V}{\partial q_{2}}\right)} \frac{\partial^{2} V}{\partial q_{2} \partial \alpha_{1}}+\cdots+\frac{\partial H}{\partial\left(\frac{\partial V}{\partial q_{k}}\right)} \frac{\partial^{2} V}{\partial q_{k} \partial \alpha_{1}} \equiv 0
$$

On the other hand, if one replaces the $p_{i}$ in the first equation in (6) with their values that one infers from (5) then it will coincide with the preceding identity. It will then be indeed a consequence of the relations (4-5). One will likewise verify the other equations (6).

That shows us that the first group of canonical equations:

$$
\frac{d q_{i}}{d t}=\frac{\partial H}{\partial p_{i}} \quad(i=1,2, \ldots, k)
$$

is verified by the solutions $p_{i}(t), q_{i}(t)$ to (4-5). If we now replace $\frac{d p_{i}}{d t}$ with $-\frac{\partial H}{\partial q_{i}}$, and $\frac{d q_{i}}{d t}$ with $\frac{\partial H}{\partial p_{i}}$ in equations (5') then that will give:

On the other hand, $V$ will satisfy the relation that is obtained by differentiating (2) with respect to $q_{1}$, which is a relation that can be written:

$$
\frac{\partial^{2} V}{\partial t \partial q_{1}}+\frac{\partial H}{\partial q_{1}}+\frac{\partial H}{\partial\left(\frac{\partial V}{\partial q_{1}}\right)} \frac{\partial^{2} V}{\partial q_{1}^{2}}+\frac{\partial H}{\partial\left(\frac{\partial V}{\partial q_{2}}\right)} \frac{\partial^{2} V}{\partial q_{2} \partial q_{1}}+\cdots+\frac{\partial H}{\partial\left(\frac{\partial V}{\partial q_{k}}\right)} \frac{\partial^{2} V}{\partial q_{k} \partial q_{1}}=0
$$

If one takes equations (5) into account then the first equation in (7) will be no different from that identity. One verifies the other equations in (7) similarly.

The proof is then complete, and one arrives at the following conclusion:
The integration of the system (1) can be achieved when one knows a complete integral of equation (2), and the equations that resolve the problem of mechanics are:

$$
\frac{\partial V}{\partial \alpha_{1}}=\beta_{1}
$$



$$
\frac{\partial V}{\partial \alpha_{k}}=\beta_{k}
$$

They determine $q_{1}, q_{2}, \ldots, q_{k}$ as functions of $t$ and the $2 k$ arbitrary constants $\alpha_{i}, \beta_{i}$.
We shall indicate a certain number of cases in which the Jacobi equation can be replaced with a simpler equation.

Case in which $H$ does not depend upon time. - Suppose that the constraints do not depend upon time, any more than the force function $U . t$ will not enter into $H$ then.

There will then exist complete integrals of equation (2) that have the form:

$$
V=-h t+W\left(q_{1}, q_{2}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{k-1}, h\right),
$$

in which $h$ is a constant.
Indeed, replace $V$ with $-h t+W\left(q_{1}, \ldots, q_{k}\right)$ in equation (2). $W$ must then satisfy the equation:

$$
H\left(q_{1}, \ldots, q_{k}, \frac{\partial W}{\partial q_{1}}, \ldots, \frac{\partial W}{\partial q_{k}}\right)=h .
$$

If $W\left(q_{1}, q_{2}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{k-1}, h\right)$ is a complete integral of (2) for each value of $h$ then one can choose the $\alpha_{1}, \ldots, \alpha_{k-1}$ in such a manner that one can give arbitrary values to $\frac{\partial W}{\partial q_{1}}, \ldots, \frac{\partial W}{\partial q_{k}}$, for example, for arbitrary $h$. It will follow from this that the function $V=-h t+W$ is a complete integral of equation (2). Indeed, if one wishes that $\frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}$ should take arbitrary values $p_{1}^{0}$, $\ldots, p_{k}^{0}$, resp., for $t_{0}, q_{1}^{0}, \ldots, q_{k}^{0}$ then one can begin by setting $h=H_{0}$, where $H_{0}$ is the value of $H$ that one will get when one replaces $q_{i}, \frac{\partial W}{\partial q_{i}}$ with $q_{i}^{0}, p_{i}^{0}$, resp., in it. One then chooses $\alpha_{1}, \ldots, \alpha_{k-1}$ in such a fashion that one will give the values $p_{1}^{0}, \ldots, p_{(k-1)}^{0}$ to $\frac{\partial W}{\partial q_{1}}, \ldots, \frac{\partial W}{\partial q_{(k-1)}}$, resp.

In this case, the motion is defined by the equalities:

$$
\frac{\partial V}{\partial h}=-t+\frac{\partial W}{\partial h}=\beta, \quad \text { or } \quad t=\frac{\partial W}{\partial h}+\beta
$$

and

$$
\frac{\partial W}{\partial \alpha_{i}}=\beta_{i} \quad(i=1,2, \ldots, k-1) .
$$

In particular, if $k=2$ then we will recover the theorem that was proved in the previous chapter. Indeed, suppose that one knows a complete integral $V=-h t+W\left(q_{1}, q_{2}, \alpha, h\right)$ of equation (2). The two equations:

$$
\frac{\partial W}{\partial q_{1}}=A_{1}, \quad \frac{\partial W}{\partial q_{2}}=A_{2}
$$

can be solved for $h$ and $\alpha$, so $W$ will satisfy the two equations:

$$
\begin{equation*}
H\left(q_{1}, q_{2}, \frac{\partial W}{\partial q_{1}}, \frac{\partial W}{\partial q_{2}}\right)=h, \quad f\left(q_{1}, q_{2}, \frac{\partial W}{\partial q_{1}}, \frac{\partial W}{\partial q_{2}}\right)=\alpha . \tag{8}
\end{equation*}
$$

In order for equations (8) to be compatible, as we know, it is necessary and sufficient that $f=$ $\alpha$ should be an integral of the canonical system (1). Therefore, any complete integral $W$ ( $q_{1}, q_{2}, \alpha$, $h$ ) of the equation $H=h$ can be regarded as a common integral to the two equations (8), in which $f=\alpha$ is a certain first integral of (1). On the other hand, if the function $W$ satisfies the two equations then we have proved that the motion of the system is determined by the equalities:

$$
\frac{\partial W}{\partial \alpha}=\beta, \quad t=\frac{\partial W}{\partial h}+\beta^{\prime} .
$$

Those are precisely the equalities to which Jacobi’s theorem will lead.
Knowing a first integral $f\left(q_{1}, q_{2}, p_{1}, p_{2}\right)=\alpha$ of (1) will then permit one to determine a complete integral of the Jacobi equation in this particular case with the aid of a simple quadrature of a total differential. In the next lecture, we will see what the analogue of that theorem is for equations (8) when the number of variables $k$ is arbitrary.

Case in which several parameters $q$ do not enter into $H$. - Suppose that $H$ has the form:

$$
H=H\left(t, q_{i+1}, q_{i+2}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}\right)
$$

One can look for an integral of the Jacobi equation that has the form:

$$
V=\alpha_{1} q_{1}+\alpha_{2} q_{2}+\ldots+\alpha_{i} q_{i}+W\left(t, q_{i+1}, \ldots, q_{k}\right)
$$

The function $W$ must satisfy the equation:

$$
\frac{\partial W}{\partial t}+H\left(t, q_{i+1}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{i}, \frac{\partial W}{\partial q_{i+1}}, \ldots, \frac{\partial W}{\partial q_{k}}\right)=0 .
$$

If one knows a complete integral $\left(t, q_{i+1}, q_{i+2}, \ldots, q_{k}, \alpha_{1}, \alpha_{2}, \ldots, \alpha_{i}, \alpha_{i+1}, \ldots \alpha_{k}\right)$ to the latter equation when one gives arbitrary values to the constants $\alpha_{1}, \ldots \alpha_{k}$ then the function $V=\alpha_{1} q_{1}+$ $\ldots+\alpha_{i} q_{i}+W$ will be a complete integral of the Jacobi equation, as one will see quite easily. The motion of the system is determined by the equations:

$$
\begin{aligned}
& q_{1}+\frac{\partial W}{\partial \alpha_{1}}=\beta_{1}, \quad \ldots, \quad q_{i}+\frac{\partial W}{\partial \alpha_{i}}=\beta_{i}, \\
& \frac{\partial W}{\partial \alpha_{i+1}}=\beta_{i+1}, \quad \ldots, \quad \frac{\partial W}{\partial \alpha_{k}}=\beta_{k} .
\end{aligned}
$$

When $H$ is independent of $t$ at the same time as $q_{1}, q_{2}, \ldots, q_{k}$, one sets:

$$
V=-h t+\alpha_{1} q_{1}+\alpha_{2} q_{2}+\ldots+\alpha_{i} q_{i}+W\left(q_{i+1}, \ldots, q_{k}\right)
$$

and one seeks a complete integral $W$ of the equation:

$$
H\left(q_{i+1}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{i}, \frac{\partial W}{\partial q_{i+1}}, \ldots, \frac{\partial W}{\partial q_{k}}\right)=h .
$$

The equalities that define the motion are then:

$$
\begin{aligned}
t=\frac{\partial W}{\partial h}+\gamma, \quad q_{1} & =\beta_{1}-\frac{\partial W}{\partial \alpha_{1}}, \quad \ldots, \quad q_{i}=\beta_{i}-\frac{\partial W}{\partial \alpha_{i}}, \\
\frac{\partial W}{\partial \alpha_{i+1}} & =\beta_{i+1}, \ldots, \quad \frac{\partial W}{\partial \alpha_{k-1}}=\beta_{k-1} .
\end{aligned}
$$

In particular, if $i=k-1$ then $W$ (which is a function of only the variable $q_{k}$ ) will be defined by a quadrature:

$$
H\left(q_{k}, \alpha_{1}, \ldots, \alpha_{k-1}, \frac{\partial W}{\partial q_{k}}\right)=h .
$$

That case often presents itself in problems with two parameters: $t$ and $q_{1}$ will then be given as functions of $q_{2}$ by the equations:

$$
t=\gamma+\frac{\partial W}{\partial h}\left(q_{2}, \alpha, h\right), \quad \quad q_{1}=\beta-\frac{\partial W}{\partial \alpha}\left(q_{2}, \alpha, h\right)
$$

In this particular case, the canonical equations admit the integral $p_{1}=\alpha$, and when that is combined with the equality $H=h$, that will determine the function:

$$
\int p_{1} d q_{1}+p_{2} d q_{2} \equiv \alpha q_{1}+W\left(q_{2}, \alpha, h\right)
$$

precisely.

Applications. - We shall now indicate some examples to which those remarks apply. We shall first study the motion of a material point $M$ under the action of a central force that is a function of only the distance $r=O M$ from the point $M$ to the center $O$ of the force.

If we refer the point to polar coordinates $r$ and $\theta$ in the plane of the trajectory then we will have:

$$
\begin{aligned}
& T=\frac{1}{2} m\left[r^{\prime 2}+r^{2} \theta^{\prime 2}\right]=\frac{1}{2 m}\left[p_{1}^{2}+\frac{p_{2}^{2}}{r^{2}}\right], \\
& H=\frac{1}{2 m}\left[p_{1}^{2}+\frac{p_{2}^{2}}{r^{2}}\right]-U(r) .
\end{aligned}
$$

The Jacobi equation admits an integral of the form $V=-h t+\alpha \theta+W(r, \alpha, h) . W$ is given by the quadrature:

$$
W=\int \frac{\sqrt{2 m r^{2}(U+h)-\alpha^{2}}}{r} d r
$$

while $t$ and $\theta$ are given by the equalities:

$$
t=\int \frac{m r d r}{\sqrt{2 m r^{2}(U+h)-\alpha^{2}}}+\text { const. }, \quad \theta=\int \frac{\alpha d r}{r \sqrt{2 m r^{2}(U+h)-\alpha^{2}}}+\text { const. }
$$

Those quadratures coincide with the ones that we obtained before using other methods. However, it would be appropriate to remark that those two quadratures are found to be performed by that fact in its own right if we know how to perform the single quadrature that gives $W$. Jacobi's method exhibits that fact quite clearly.

Similarly, let us calculate the motion of a point $M$ on a surface of revolution when the given force admits a force function $U(r)$, where $r$ denotes the distance from the point $M$ to the axis of revolution, and $\theta$ denotes the angle between the two planes $z O M$ and $x O z$. If $z=\varphi(r)$ is the equation of the surface then we will have:

$$
\begin{aligned}
& T=\frac{1}{2} m\left[r^{\prime 2}\left(1+\varphi^{\prime 2}\right)+r^{2} \theta^{\prime 2}\right]=\frac{1}{2 m}\left[\frac{p_{1}^{2}}{1+\varphi^{\prime 2}}+\frac{p_{2}^{2}}{r^{2}}\right] \\
& H=\frac{1}{2 m}\left[\frac{p_{1}^{2}}{1+\varphi^{\prime 2}}+\frac{p_{2}^{2}}{r^{2}}\right]-U(r)
\end{aligned}
$$

The functions $W(r, \alpha, h)$ and $t(r), \theta(r)$ are determined here by the equalities:

$$
\begin{gathered}
W=\int \frac{\sqrt{\left(1+\varphi^{\prime 2}\right)\left[2 m r^{2}(U+h)-\alpha^{2}\right]}}{r} d r, \\
t=\int \frac{m r \sqrt{1+\varphi^{\prime 2}} d r}{\sqrt{2 m r^{2}(U+h)-\alpha^{2}}}, \quad \theta=\int \frac{\alpha \sqrt{1+\varphi^{\prime 2}}}{r \sqrt{2 m r^{2}(U+h)-\alpha^{2}}} d r .
\end{gathered}
$$

It would be appropriate to make the same remark as before on the subject of those two integrals. Furthermore, that remark can be repeated in all of the analogous cases.

It would be easy to apply Jacobi's method to all of the examples that were treated before in which $H$ could be reduced to something that depended upon only $q_{k}$ and the variables $p_{i}$ by a convenient choice of variables. We shall confine ourselves here to recalling the example that was treated on page 116.

A massive, homogeneous, solid body of revolution is traversed along its axis by a needle to which it is subject and one of whose extremities slides without friction on one vertical $O z$, while the other slides on the horizontal plane $x O y$. Let us study its motion.

The vis viva $2 T$ of the system (see page 117) is equal to:

$$
\left(A+M d^{2}\right) \sin ^{2} \theta \psi^{\prime 2}+\left\{A+M\left[d^{2} \cos ^{2} \theta+(l-d)^{2} \sin ^{2} \theta\right]\right\} \theta^{\prime 2}+\left[C \cos \theta \psi^{\prime}+\varphi^{\prime}\right]^{2},
$$

and the force function $U$ is $M g(l-d) \cos \theta$. If one sets:

$$
p_{1}=\frac{\partial T}{\partial \theta^{\prime}}, \quad p_{2}=\frac{\partial T}{\partial \varphi^{\prime}}, \quad p_{3}=\frac{\partial T}{\partial \psi^{\prime}}
$$

then that will give:

$$
H=\frac{\left(p_{3}-C \cos \theta p_{2}\right)^{2}}{2\left(A+M d^{2}\right) \sin ^{2} \theta}+\frac{p_{1}^{2}}{2\left\{A+M\left[d^{2} \cos ^{2} \theta+(l-d)^{2} \sin ^{2} \theta\right]\right\}}+\frac{p_{2}^{2}}{2}-M g(l-d) \cos \theta .
$$

One can take the function $W$ to be the function that is defined by:

$$
W=\alpha \varphi+\alpha^{\prime} \psi+W_{1}(\theta),
$$

with
$W_{1}=\int d \theta \sqrt{\left\{A+M\left[d^{2} \cos ^{2} \theta+(l-d)^{2} \sin ^{2} \theta\right]\right\}\left\{2 M g(l-d) \cos \theta+2 h-d^{2}-\frac{\alpha^{\prime}-(\alpha \cos \theta)^{2}}{\left(A+M d^{2}\right) \sin ^{2} \theta}\right\}}$,
and the motion will then be determined by the equalities:

$$
t=\frac{\partial W_{1}}{\partial h}+\text { const., } \quad \varphi=\beta-\frac{\partial W_{1}}{\partial \alpha}, \quad \psi=\beta^{\prime}-\frac{\partial W_{1}}{\partial \alpha^{\prime}},
$$

in which $\alpha, \beta, \alpha^{\prime}, \beta^{\prime}, h$ are arbitrary constants. One thus recovers the equations on page 118. The three quadratures that give $t, \varphi$, and $\psi$ can be performed once one has performed the one that gives $W_{1}$.

Theorems of Liouville and Staeckel.

Now suppose that the vis viva of the system has the form:

$$
T=\frac{\varphi_{1}\left(q_{1}\right)+\varphi_{2}\left(q_{2}\right)+\cdots+\varphi_{k}\left(q_{k}\right)}{2}\left[A_{1}\left(q_{1}\right) q_{1}^{\prime 2}+A_{2}\left(q_{2}\right) q_{2}^{\prime 2}+\cdots+A_{k}\left(q_{k}\right) q_{k}^{\prime 2}\right]
$$

and that the forces are derived from the potential:

$$
U=\frac{\psi_{1}\left(q_{1}\right)+\psi_{2}\left(q_{2}\right)+\cdots+\psi_{k}\left(q_{k}\right)}{\varphi_{1}\left(q_{1}\right)+\varphi_{2}\left(q_{2}\right)+\cdots+\varphi_{k}\left(q_{k}\right)} .
$$

One has the expression for $H$ :

$$
H=\frac{1}{2\left(\varphi_{1}+\varphi_{2}+\cdots+\varphi_{k}\right)}\left[\frac{p_{1}^{2}}{A_{1}}+\frac{p_{2}^{2}}{A_{2}}+\cdots+\frac{p_{k}^{2}}{A_{k}}\right]-\frac{\psi_{1}+\psi_{2}+\cdots+\psi_{k}}{\varphi_{1}+\varphi_{2}+\cdots+\varphi_{k}} .
$$

Conversely, if $H$ has that form then $T$ and $U$ will have the indicated form (see pp. 165). The partial differential equation in $W$, namely:

$$
\sum_{i=1}^{k} \frac{1}{A_{i}}\left(\frac{d W}{d q_{i}}\right)^{2}=2 \sum_{i=1}^{k}\left(\psi_{i}+h \varphi_{i}\right)
$$

will then admit a complete integral of the form:

$$
W=W_{1}\left(q_{1}\right)+W_{1}\left(q_{1}\right)+\ldots+W_{k}\left(q_{k}\right) .
$$

Indeed, it suffices to take $W_{i}\left(q_{i}\right)$ to be the function that is given by the equality:

$$
W_{i}\left(q_{i}\right)=\int d q_{i} \sqrt{2 A_{i}\left(\psi_{i}+h \varphi_{i}+\alpha_{i}\right)},
$$

in which the $\alpha_{i}$ are constants that are constrained by the single condition that:

$$
\alpha_{1}+\alpha_{2}+\ldots+\alpha_{k}=0
$$

The function $W=\sum W_{i}$ will then satisfy equation ( $2^{\prime}$ ) identically.
If one sets $F_{i}=\sqrt{\psi_{i}+h \varphi_{i}+\alpha_{i}}$ then the motion will be defined by the equalities:

$$
\frac{\partial W_{i}}{\partial \alpha_{i}}-\frac{\partial W_{k}}{\partial \alpha_{k}}=\text { const. } \quad\left[i=1,2, \ldots,(k-1), \text { since } \alpha_{k}=-\left(\alpha_{1}+\ldots+\alpha_{k-1}\right)\right]
$$

and

$$
t+\text { const. }=\frac{\partial W}{\partial h}
$$

which are equalities that can be written:

$$
\int \sqrt{\frac{A_{1}}{F_{1}}} d q_{1}+\beta_{1}=\int \sqrt{\frac{A_{2}}{F_{2}}} d q_{2}+\beta_{2}=\ldots=\int \sqrt{\frac{A_{k}}{F_{k}}} d q_{k}+\beta_{k}
$$

and

$$
t \sqrt{2}+\text { const. }=\int \varphi_{1} \sqrt{\frac{A_{1}}{F_{1}}} d q_{1}+\varphi_{2} \sqrt{\frac{A_{2}}{F_{2}}} d q_{2}+\cdots+\varphi_{k} \sqrt{\frac{A_{k}}{F_{k}}} d q_{k}
$$

resp. One thus recovers the same equalities that one would get by starting from Lagrange's equations (see pages 120-123). However, on the one hand, those equalities will follow more directly from Jacobi's theorem. On the other hand, one sees that the $2 k$ integrals that define the motion can be performed if one knows how to perform the quadratures that define the $W_{i}$.

In another lecture, we will give some new applications of Liouville's theorem. I will now indicate a generalization of that theorem. That generalization is due to P. Staeckel, who extended it to an arbitrary number of parameters. In order to not complicate the notation, we shall adopt the case of three parameters.

Let $\Delta$ be the determinant:

$$
\Delta=\left|\begin{array}{lll}
\varphi_{1}\left(q_{1}\right) & \varphi_{2}\left(q_{2}\right) & \varphi_{3}\left(q_{3}\right) \\
\psi_{1}\left(q_{1}\right) & \psi_{2}\left(q_{2}\right) & \psi_{3}\left(q_{3}\right) \\
\chi_{1}\left(q_{1}\right) & \chi_{2}\left(q_{2}\right) & \chi_{3}\left(q_{3}\right)
\end{array}\right|,
$$

and let $\Phi_{1}, \Phi_{2}, \Phi_{3}$ be the minors of $\Delta$ with respect to the elements of the first row, while $\Psi_{1}, \ldots$ and $\Xi_{1}, \ldots$ are the minors relative to $\psi_{1}, \ldots, \chi_{1}, \ldots$, resp.

Suppose that the vis viva $2 T$ of a system and the force function $U$ have the forms:

$$
\begin{equation*}
2 T=\Delta\left(\frac{q_{1}^{\prime 2}}{\phi_{1}}+\frac{q_{2}^{\prime 2}}{\phi_{2}}+\frac{q_{3}^{\prime 2}}{\phi_{3}}\right), \quad U=\frac{f_{1}\left(q_{1}\right) \phi_{1}+f_{2}\left(q_{2}\right) \phi_{2}+f_{3}\left(q_{3}\right) \phi_{3}}{\Delta} . \tag{A}
\end{equation*}
$$

The Jacobi method will permit one to determine the motion by quadratures.
Indeed, we have:

$$
H=\frac{1}{2 \Delta}\left[p_{1}^{2} \Phi_{1}+p_{2}^{2} \Phi_{2}+p_{3}^{2} \Phi_{3}\right]-\frac{f_{1} \Phi_{1}+f_{2} \Phi_{2}+f_{3} \Phi_{3}}{\Delta}
$$

here, and equation ( $2^{\prime}$ ) for $W$ will be written:

$$
\left(\frac{\partial W}{\partial q_{1}}\right)^{2} \Phi_{1}+\left(\frac{\partial W}{\partial q_{2}}\right)^{2} \Phi_{2}+\left(\frac{\partial W}{\partial q_{3}}\right)^{2} \Phi_{3}=2\left(f_{1} \Phi_{1}+f_{2} \Phi_{2}+f_{3} \Phi_{3}+h \Delta\right)
$$

Let us see whether there exists a complete integral of the form:

$$
W=W_{1}\left(q_{1}\right)+W_{2}\left(q_{2}\right)+W_{3}\left(q_{3}\right) .
$$

If we observe that we have:

$$
\sum_{1,2,3} \varphi \Phi \equiv \Delta, \quad \sum_{1,2,3} \psi \Phi \equiv 0, \quad \sum_{1,2,3} \chi \Phi \equiv 0
$$

then when we set:

$$
\begin{aligned}
& \left(\frac{\partial W}{\partial q_{1}}\right)^{2}=2\left[f_{1}+h \varphi_{1}+\alpha \psi_{1}+\alpha^{\prime} \chi_{1}\right]=F_{1}\left(q_{1}\right), \\
& \left(\frac{\partial W}{\partial q_{2}}\right)^{2}=2\left[f_{2}+h \varphi_{2}+\alpha \psi_{2}+\alpha^{\prime} \chi_{2}\right]=F_{2}\left(q_{2}\right), \\
& \left(\frac{\partial W}{\partial q_{3}}\right)^{2}=2\left[f_{3}+h \varphi_{3}+\alpha \psi_{3}+\alpha^{\prime} \chi_{3}\right]=F_{3}\left(q_{3}\right),
\end{aligned}
$$

equation ( $2^{\prime}$ ) will be verified identically: $h, \alpha, \alpha^{\prime}$ denote constants.
The motion is defined by the equalities:

$$
\begin{aligned}
& \beta=\int \frac{\psi_{1} d q_{1}}{\sqrt{F_{1}\left(q_{1}\right)}}+\int \frac{\psi_{2} d q_{2}}{\sqrt{F_{2}\left(q_{2}\right)}}+\int \frac{\psi_{3} d q_{3}}{\sqrt{F_{3}\left(q_{3}\right)}}, \\
& \gamma=\int \frac{\chi_{1} d q_{1}}{\sqrt{F_{1}\left(q_{1}\right)}}+\int \frac{\chi_{2} d q_{2}}{\sqrt{F_{2}\left(q_{2}\right)}}+\int \frac{\chi_{3} d q_{3}}{\sqrt{F_{3}\left(q_{3}\right)}},
\end{aligned}
$$

and

$$
\sqrt{2} t+\text { cost. }=\int \frac{\varphi_{1} d q_{1}}{\sqrt{F_{1}\left(q_{1}\right)}}+\int \frac{\varphi_{2} d q_{2}}{\sqrt{F_{2}\left(q_{2}\right)}}+\int \frac{\varphi_{3} d q_{3}}{\sqrt{F_{3}\left(q_{3}\right)}}
$$

The preceding theorem includes Liouville's theorem as a special case, because if $T$ and $U$ have the form that Liouville's theorem requires then one will see directly that one can always put it into the form (A).

## LECTURE 16

## STUDY OF REAL TRAJECTORIES. EQUATIONS OF THE TRAJECTORIES WHEN THE FORCES ARE ZERO OR DERIVED FROM A POTENTIAL.

Number of arbitrary constants that the trajectories depend upon. - Consider a system $S$ in which neither the constraints nor the forces depend upon time. The motion of that system is determined by the Lagrange equations:

$$
\begin{equation*}
\frac{d}{d t}\left(\frac{\partial T}{\partial q_{i}^{\prime}}\right)-\frac{\partial T}{\partial q_{i}}=Q_{i}\left(q_{1}, q_{2}, \ldots, q_{k}, q_{1}^{\prime}, q_{2}^{\prime}, \ldots, q_{k}^{\prime}\right), \quad \frac{d q_{i}}{d t}=q_{i}^{\prime} \quad(i=1,2, \ldots, k) \tag{1}
\end{equation*}
$$

in which:

$$
2 T \equiv \sum_{i, j} q_{i}^{\prime} q_{j}^{\prime} A_{i j}\left(q_{1}, q_{2}, \ldots, q_{k}\right) \quad\left(A_{i j} \equiv A_{j i}\right)
$$

If $t_{0}, q_{2}^{0}, \ldots, q_{k}^{0}, q_{1}^{\prime 0}, \ldots, q_{k}^{\prime 0}$ are the values of $t, q_{2}, \ldots, q_{k}, q_{1}^{\prime}, \ldots, q_{k}^{\prime}$ for a given value of $q_{1}$ (say $q_{1}=0$ ) then those equations will define $q_{1}, q_{2}, \ldots, q_{k}$ as functions of $t-t_{0}$ and ( $2 k-1$ ) arbitrary constants $q_{2}^{0}, \ldots, q_{k}^{0}, q_{1}^{\prime 0}, \ldots, q_{k}^{\prime 0}$. As a result, $q_{2}, \ldots, q_{k}$ will be functions of $q_{1}$ that depend upon at most $2 k-1$ arbitrary constants $\left({ }^{2}\right)$. On the other hand, they depend upon at least $(2 k-2)$ constants, because $q_{2}, \ldots, q_{k}$, and $\frac{d q_{2}}{d q_{1}}\left(\right.$ or $\left.\frac{q_{2}^{\prime}}{q_{1}^{\prime}}\right), \ldots, \frac{d q_{k}}{d q_{1}}\left(\right.$ or $\left.\frac{q_{k}^{\prime}}{q_{1}^{\prime}}\right)$ can take arbitrary values for a given $q_{1}$. We shall see that in general the number $v$ of constants that the trajectories of the system depend upon is indeed equal to $2 k-1$ and find the conditions under which that number will reduce to $2 k-2$.

Solve equations (1) for the $q_{i}^{\prime \prime}$. If $\Delta$ denotes the discriminant of the quadratic form $T$, and $\alpha_{i}$ denotes what $\Delta$ will become when one replaces the elements of the $i^{\text {th }}$ column with $Q_{1}, Q_{2}, \ldots, Q_{k}$, and finally if $P_{i}$ represents quadratic form with respect to the $q_{i}^{\prime}$ then one can write:

$$
\begin{equation*}
q_{i}^{\prime \prime}=P_{i}+\frac{\alpha_{i}}{\Delta}=P_{i}+\beta_{i} \quad(i=1,2, \ldots, k) \tag{2}
\end{equation*}
$$

On the other hand, we have:

$$
q_{i}^{\prime \prime}=\frac{d^{2} q_{i}}{d q_{1}^{2}} q_{1}^{\prime 2}+\frac{d q_{i}}{d q_{1}} q_{1}^{\prime \prime}
$$

[^1]so upon replacing $q_{1}^{\prime \prime}$ and $q_{i}^{\prime \prime}$ with their values that we infer from (2):
\[

$$
\begin{equation*}
\frac{d^{2} q_{i}}{d q_{1}^{2}}=\frac{\left[P_{i}-\frac{d q_{i}}{d q_{1}} P_{1}\right]+\left[\beta_{i}-\frac{d q_{i}}{d q_{1}} \beta_{1}\right]}{q_{1}^{\prime 2}} \quad(i=2,3, \ldots, k) . \tag{3}
\end{equation*}
$$

\]

In what follows, let $q_{(i)}^{\prime}, q_{(i)}^{\prime \prime}, \ldots$ represent the derivatives $\frac{d q_{i}}{d q_{1}}, \frac{d^{2} q_{i}}{d q_{1}^{2}}, \ldots$, and let $\pi_{i}$ denote what $P_{i}$ will become when one replaces $q_{1}^{\prime}$ with $1, q_{2}^{\prime}$ with $q_{(2)}^{\prime}, \ldots, q_{k}^{\prime}$ with $q_{(k)}^{\prime}$.

Equation (3) will then become:

$$
\begin{equation*}
q_{(i)}^{\prime \prime}=\pi_{i}-q_{(i)}^{\prime} \pi_{1}+\frac{\beta_{i}-q_{(i)}^{\prime} \beta_{1}}{q_{1}^{\prime 2}} \quad(i=2,3, \ldots, k) \tag{4}
\end{equation*}
$$

Those $\beta$ are functions of $q_{1}, q_{2}, \ldots, q_{k}, q_{1}^{\prime}, q_{(2)}^{\prime}, \ldots, q_{(k)}^{\prime}$. There are two cases to be distinguished: If the right-hand side of equations (4) is independent of $q_{1}^{\prime}$ then the system (4) will form a system of $(k-1)$ second-order equations that involve the $k$ variables $q_{1}, q_{2}, \ldots, q_{k}$, and those equations will define $q_{2}, q_{3}, \ldots, q_{k}$ as functions of $q_{1}$ and $(2 k-2)$ arbitrary constants $q_{2}^{0}, \ldots, q_{k}^{0}, q_{(2)}^{\prime 0}, \ldots$, $q_{(k)}^{\prime 0}$. If, on the contrary, $q_{1}^{\prime}$ appears in at least one of equations (4) (say, the equation for $i=2$ ) then one can take:

$$
q_{2}, \ldots, q_{k}, q_{1}^{\prime}, q_{(2)}^{\prime}, \ldots, q_{(k)}^{\prime}
$$

arbitrarily for a given $q_{1}$ and choose $q_{1}^{\prime 0}$ in order to give $q_{(2)}^{\prime \prime 0}$ an arbitrary value, so the functions $q_{2}, q_{3}, \ldots, q_{k}$ of $q_{1}$ will then depend upon ( $2 k-2$ ) distinct constants.

In order for the right-hand sides of all of equations (4) to be independent of $q_{1}^{\prime}$, it is necessary and sufficient that the $(k-1)$ expressions $\frac{\beta_{i}-\left(q_{i}^{\prime} / q_{1}^{\prime}\right) \beta_{1}}{q_{1}^{\prime 2}}$ should be homogeneous of degree zero with respect to the $q_{i}^{\prime}$, in such a way that the result to which we will arrive can be stated thus:

The number $v$ of distinct constants upon which the trajectories of the system depend is equal to $(2 k-1)$ or $(2 k-2)$. In order for it to reduce to $(2 k-2)$, it is necessary and sufficient that the $(k$ -1 ) expressions ( $\beta_{i} q_{1}^{\prime}-\beta_{1} q_{i}^{\prime}$ ) should be homogeneous of degree three with respect to the $q_{i}^{\prime}$.

If one lets $a_{i j}$ denote the minor of $\Delta$ relative to $A_{i j}$ then one will have:

$$
\begin{equation*}
\beta_{i}=\frac{1}{\Delta}\left[a_{i 1} Q_{1}+a_{i 2} Q_{2}+\cdots+a_{i k} Q_{k}\right] . \tag{5}
\end{equation*}
$$

The $v$ will then be equal to ( $2 k-2$ ), in particular, whenever the forces $Q_{i}$ are homogeneous of degree two with respect to $q_{1}^{\prime}, q_{2}^{\prime}, \ldots, q_{k}^{\prime}$.

When the $Q_{i}$ do not depend upon the velocities, the same thing will be true for the $\beta_{i}$. In order for $v$ to be equal to $(2 k-2)$, it is necessary and sufficient that the expressions $\beta_{i} q_{1}^{\prime}-\beta_{1} q_{i}^{\prime}$ should be identically zero; in other words, that one should have:

$$
\begin{equation*}
\frac{\beta_{1}}{q_{1}^{\prime}}=\frac{\beta_{2}}{q_{2}^{\prime}}=\ldots=\frac{\beta_{k}}{q_{k}^{\prime}} . \tag{6}
\end{equation*}
$$

It is possible for the $\beta$ to not depend upon the $q_{i}^{\prime}$ only if one has:

$$
\beta_{1} \equiv \beta_{2} \equiv \ldots \equiv \beta_{k} \equiv 0
$$

and as a result $\left(^{3}\right)$ :

$$
Q_{1} \equiv Q_{2} \equiv \ldots \equiv Q_{k} \equiv 0 .
$$

The trajectories of a system for which the $Q_{i}$ are zero thus depend upon $(2 k-2)$ parameters. We shall give the name of geodesics for the $d s^{2}$ of $T$ to those trajectories upon setting:

$$
d s^{2}=\sum A_{i j} d q_{i} d q_{j}
$$

That is the case for a system without friction that is not subject to any given forces. When the forces $Q_{i}$, which are functions of only $q_{1}, q_{2}, \ldots, q_{k}$, are not all zero, the trajectories will always depend upon $(2 k-1)$ parameters.

A remarkable case is the one in which the forces depend upon velocities so the trajectories of the system will coincide with the geodesics of $d s^{2}$. In order for that to be true, it is necessary and sufficient that equations (6) should be verified; in other words, that one should have:

$$
\beta_{1}=\lambda q_{1}^{\prime}, \quad \beta_{2}=\lambda q_{2}^{\prime}, \quad \ldots, \quad \beta_{k}=\lambda q_{k}^{\prime}
$$

and as a result, from (5):

$$
\begin{equation*}
\lambda q_{i}^{\prime}=\frac{1}{\Delta}\left[a_{i 1} Q_{1}+a_{i 2} Q_{2}+\cdots+a_{i k} Q_{k}\right] \quad(i=1,2, \ldots, k) . \tag{7}
\end{equation*}
$$

However, on the other hand, let $p_{i}=\frac{\partial T}{\partial q_{i}^{\prime}}$, so one knows that one has:

$$
q_{i}^{\prime}=\frac{1}{\Delta}\left[a_{i 1} p_{1}+a_{i 2} p_{2}+\cdots+a_{i k} p_{k}\right] .
$$

[^2]Equations (7) will then be equivalent to the following ones:

$$
\begin{equation*}
Q_{1}=\lambda \frac{\partial T}{\partial q_{1}^{\prime}}, \quad Q_{2}=\lambda \frac{\partial T}{\partial q_{2}^{\prime}}, \quad \ldots, \quad Q_{k}=\lambda \frac{\partial T}{\partial q_{k}^{\prime}} \tag{8}
\end{equation*}
$$

In order for the trajectories of (1) to be the same when all of the forces $\left(Q_{i}\right)$ are zero, it will therefore be necessary and sufficient that $Q_{1}, Q_{2}, \ldots, Q_{k}$ should be proportional to $\frac{\partial T}{\partial q_{1}^{\prime}}, \frac{\partial T}{\partial q_{2}^{\prime}}, \ldots$, $\frac{\partial T}{\partial q_{k}^{\prime}}$.

Examples. - Let us interpret those conditions (8) in certain special cases. First of all, let $T$ be the vis viva of a free material point, so one has: $T=\frac{1}{2} m\left[x^{\prime 2}+y^{\prime 2}+z^{\prime 2}\right]$, and the conditions (8) express the idea that $X, Y, Z$ are proportional to $x^{\prime}, y^{\prime}, z^{\prime}$, i.e., the force that is exerted on that point will have the velocity for the line of action. When a free point is subject to a force that constantly points in the direction of its velocity (or in the opposite sense), it will describe a line, no matter what the initial conditions.

Now, if the system is composed of free material points $M_{i}$ then one will have: $T=$ $\sum \frac{1}{2} m_{i}\left[x_{i}^{\prime 2}+y_{i}^{\prime 2}+z_{i}^{\prime 2}\right]$, and the conditions (8) express the idea that the force ( $F_{i}$ ) that is exerted upon each point $M_{i}$ is directed along the velocity of $M_{i}$ (in the same or opposite sense), and that one should have $\frac{F_{i}}{m_{i} v_{i}}=\frac{F_{1}}{m_{1} v_{1}}(i=2,3, \ldots, k)$, moreover.

Finally, let us treat the case of a point $M$ that moves on a surface $S$. The conditions (8) then express the idea that the force that acts upon it is constantly in the normal plane to the surface and tangent to the trajectory. Indeed, let $x\left(q_{1}, q_{2}\right), y\left(q_{1}, q_{2}\right), z\left(q_{1}, q_{2}\right)$ be the Cartesian coordinates of a point on the surface. The total force that is exerted on $M$ is, as one knows, the resultant of the normal reaction and the given force $\left(F^{\prime}\right)$ [viz., the active force and the force of friction, if it exists]. Let $\left(F_{1}\right)$ or $\left(X_{1}, Y_{1}, Z_{1}\right)$ be the projection of the given force $\left(F^{\prime}\right)$ onto the tangent plane to $S$. One has the following relations between $Q_{1}, Q_{2}, X_{1}, Y_{1}, Z_{1}$ :

$$
\begin{equation*}
Q_{1}=X_{1} \frac{\partial x}{\partial q_{1}}+Y_{1} \frac{\partial y}{\partial q_{1}}+Z_{1} \frac{\partial z}{\partial q_{1}}, \quad Q_{2}=X_{1} \frac{\partial x}{\partial q_{2}}+Y_{1} \frac{\partial y}{\partial q_{2}}+Z_{1} \frac{\partial z}{\partial q_{2}} . \tag{9}
\end{equation*}
$$

Conversely, if ( $F_{1}$ ) denotes a segment that is tangent to $S$ whose projections satisfy equations (9) then it will denote the component of $\left(F^{\prime}\right)$ that is tangent to $S$. In order for that segment $\left(F_{1}\right)$ to be tangent to the trajectory, it is necessary that one must have:

$$
\frac{X_{1}}{\frac{\partial x}{\partial q_{1}} q_{1}^{\prime}+\frac{\partial x}{\partial q_{2}} q_{2}^{\prime}}=\frac{Y_{1}}{\frac{\partial y}{\partial q_{1}} q_{1}^{\prime}+\frac{\partial y}{\partial q_{2}} q_{2}^{\prime}}=\frac{Z_{1}}{\frac{\partial z}{\partial q_{1}} q_{1}^{\prime}+\frac{\partial z}{\partial q_{2}} q_{2}^{\prime}}
$$

and as a result, upon calling the common value of those ratios $\mu$ :

$$
Q_{1}=\mu\left\{q_{1}^{\prime}\left[\left(\frac{\partial x}{\partial q_{1}}\right)^{2}+\left(\frac{\partial y}{\partial q_{1}}\right)^{2}+\left(\frac{\partial z}{\partial q_{1}}\right)^{2}\right]+q_{1}^{\prime}\left(\frac{\partial x}{\partial q_{1}} \frac{\partial x}{\partial q_{2}}+\frac{\partial y}{\partial q_{1}} \frac{\partial y}{\partial q_{2}}+\frac{\partial z}{\partial q_{1}} \frac{\partial z}{\partial q_{2}}\right)\right\}=\frac{\mu}{m} \frac{\partial T}{\partial q_{1}^{\prime}}
$$

and similarly, $Q_{2}=\frac{\mu}{m} \frac{\partial T}{\partial q_{2}^{\prime}}$.
Conversely, if $Q_{1}$ and $Q_{2}$ have that form then the segment:

$$
X_{1}=\mu\left(\frac{\partial x}{\partial q_{1}} q_{1}^{\prime}+\frac{\partial x}{\partial q_{2}} q_{2}^{\prime}\right), \quad Y_{1}=\mu\left(\frac{\partial y}{\partial q_{1}} q_{1}^{\prime}+\frac{\partial y}{\partial q_{2}} q_{2}^{\prime}\right), \quad Z_{1}=\mu\left(\frac{\partial z}{\partial q_{1}} q_{1}^{\prime}+\frac{\partial z}{\partial q_{2}} q_{2}^{\prime}\right)
$$

will be a tangent segment to the trajectory, so to the surface that satisfies equations (9) and, as a result, it will represent the component of $\left(F^{\prime}\right)$ that is tangent to $S$. The conditions (8) then express the fact that $\left(F^{\prime}\right)$ is projected onto the tangent plane to $S$ along tangent to the trajectory. That will be the case, for example, with a point that moves without friction on a surface in a resisting medium with no other force acting upon it.

We shall now study the systems that are subject to forces that are independent of velocity exclusively.

Study of the trajectories in the case where the forces do not depend upon velocity. - From the foregoing, it would be suitable to subdivide this case into two other ones according to whether all of the forces are zero or not.
I. All of the coefficients $Q_{i}$ are zero. - The trajectories (viz., geodesics of $d s^{2}$ ) will then depend upon $(2 k-2)$ constants and are defined [see page 235] by $(k-1)$ equations of the form:

$$
q_{(i)}^{\prime \prime}=\pi_{i}-q_{(i)}^{\prime} \pi_{i} \quad(i=2,3, \ldots, k) .
$$

Once those equations are integrated, the motion of the system will be given by the equality:

$$
d t=\frac{d s}{\sqrt{h}}=\frac{d q_{1}}{\sqrt{h}} \cdot \sqrt{\sum A_{i j} q_{(i)}^{\prime} q_{(j)}^{\prime}},
$$

in which $q_{(i)}^{\prime}=d q_{i} / d q_{1}, q_{1}^{\prime}=1$, and $h$ represents an arbitrary constant for each geodesic. Therefore, if one expresses $q_{2}, q_{3}, \ldots, q_{k}$ as functions of $q_{1}$ and $(2 k-2)$ arbitrary constants $c_{1}, c_{2}$, $\ldots, c_{2 k-2}$ then one will have:

$$
t=c \int f\left[q_{1}, c_{1}, c_{2}, \ldots, c_{2 k-2}\right] d q_{1}+c^{\prime}
$$

in which $c$ represents any one of the constants. More generally, let:

$$
\varphi\left[q_{1}, q_{2}, \ldots, q_{k}, q_{(2)}^{\prime}, \ldots, q_{(k)}^{\prime}, c\right]=\alpha
$$

be a first integral of equations (4') that depends upon one arbitrary parameter $c$. One can write:

$$
d t=\varphi \sqrt{\sum A_{i j} q_{(i)}^{\prime} q_{(j)}^{\prime}} d q_{1},
$$

and that equality, when combined with equations (4'), define the same motion as (1). Conversely, if the equality:

$$
\begin{equation*}
d t=\psi\left[q_{1}, q_{2}, \ldots, q_{k}, q_{(2)}^{\prime}, \ldots, q_{(k)}^{\prime}, c\right] \sqrt{\sum A_{i j} q_{(i)}^{\prime} q_{(j)}^{\prime}} d q_{1} \tag{a}
\end{equation*}
$$

in which $\psi$ is an arbitrary function, is compatible with equations (1), in other words, if one combines it with equations (4') then it will define a motion on $S . \psi$ is a first integral of the geodesics.

Indeed, consider an arbitrary geodesic. From (1), one will have $\frac{d t}{d s}=\frac{1}{\sqrt{h}}$ all along that geodesic, and as a result, from $(a)$, one must also have: $\psi=1 / \sqrt{h}=$ const. Therefore, $\psi$ is a first integral of (4').

We will soon give an explicit form to equations (4').
For the moment, I shall insist upon only the fact that the same trajectory can correspond to an infinitude of distinct motions, and I intend the word "motion" to mean the same positions, but different velocities. Moreover, those motions are deduced from just one of them by multiplying all velocities by the same numerical constant.
II. Not all of the coefficients $Q_{i}\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ are zero. Equations of the trajectories. - We know that the trajectories depend upon $(2 k-1)$ parameters. In order to construct the differential equations that define those trajectories without the intermediary of $t$, observe that from (4) [see page 2], we first have:

$$
\begin{equation*}
\frac{q_{(i)}^{\prime \prime}+q_{(i)}^{\prime} \pi_{1}-\pi_{i}}{\beta_{i}-q_{(i)}^{\prime} \beta_{1}}=\frac{q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \pi_{1}-\pi_{2}}{\beta_{2}-q_{(2)}^{\prime} \beta_{1}} \quad(i=3,4, \ldots, k), \tag{10}
\end{equation*}
$$

in which the common value of those ratios is $\left(\frac{d t}{d q_{1}}\right)^{2}$.

On the other hand, if one differentiates the equality:

$$
q_{1}^{\prime 2}=\frac{\beta_{2}-q_{(2)}^{\prime} \beta_{1}}{q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \pi_{1}-\pi_{2}} \equiv \frac{\psi_{2}}{\chi_{2}}
$$

with respect to $q_{1}$ and points out that $\frac{d}{d q_{1}} \cdot q_{1}^{\prime 2}=2 q_{1}^{\prime} q_{1}^{\prime \prime} \frac{d t}{d q_{1}}=2 q_{1}^{\prime \prime}$, and replaces $q_{1}^{\prime \prime}$ with $\left(\pi_{1} q_{1}^{\prime 2}+\beta_{1}\right) \equiv \pi_{1}\left(\psi_{2} / \chi_{2}\right)+\beta_{1}$ then one will have:

$$
2 \pi_{1} \frac{\psi_{2}}{\chi_{2}}+2 \beta_{1}=\frac{d}{d q_{1}} \frac{\psi_{2}}{\chi_{2}},
$$

or rather:

$$
\begin{equation*}
\frac{\frac{d}{d q_{1}}\left[q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \pi_{1}-\pi_{2}\right]}{q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \pi_{1}-\pi_{2}}+\pi_{1}=\frac{\frac{d}{d q_{1}}\left(\beta_{2}-q_{(2)}^{\prime} \beta_{1}\right)}{\beta_{2}-q_{(2)}^{\prime} \beta_{1}} \cdot \frac{2 \beta_{1}\left[q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \pi_{1}-\pi_{2}\right]}{\beta_{2}-q_{(2)}^{\prime} \beta_{1}}, \tag{11}
\end{equation*}
$$

which is an equation of the form:

$$
q_{(2)}^{\prime \prime \prime}=\frac{-3 q_{(2)}^{\prime \prime 2}+q_{2}^{\prime \prime} M_{3}+M_{5}}{M_{0}-q_{(2)}^{\prime}},
$$

in which $M_{3}$ and $M_{5}$ denote polynomials of degrees three and five, resp., in $q_{(2)}^{\prime}, q_{(3)}^{\prime}, \ldots, q_{(k)}^{\prime}$, and $M_{0}$ is a function of the $q_{i}$.

Since equations (10) and (11) define the trajectories, it would be easy to give them a more symmetric form, but that is hardly important in the context of our objective.

Observe that in equations (10) and (11), the expressions $\pi_{1}, \ldots, \pi_{k}$ depend upon only the vis viva $T$. Only the coefficients $\beta_{i}$ vary with the forces $Q_{i}$. Furthermore, observe that the geodesics of $T$ are obtained by equating all of the numerators $\chi_{i}$ of the ratios (10) to zero.

A first consequence of those remarks is that the geodesics of $T$ belong to the trajectories no matter what the forces $Q_{i}\left(q_{1}, q_{2}, \ldots, q_{k}\right)$. Indeed, since those geodesics satisfy the equations $\chi_{i}=$ 0 , they will satisfy equations (10). On the other hand, equation (11) can be written:

$$
\frac{d \chi_{2}}{d q_{1}}=\chi_{2} A,
$$

and since one has $\chi_{2}=0, d \chi_{2} / d q_{1}=0$ for an arbitrary geodesic, it will also be verified. The geodesics then define a congruence of trajectories with $(2 k-2)$ parameters.

Determination of time. - When the system (10), (11) has been integrated, one will know $q_{2}$, $q_{3}, \ldots, q_{k}$ as functions of $q_{1}$ and $(2 k-1)$ arbitrary constants, namely, the initial values of $q_{2}^{\prime \prime}, \ldots$, $q_{k}^{\prime \prime}, q_{(2)}^{\prime 0}, q_{(3)}^{\prime 0}, \ldots, q_{(k)}^{\prime 0}, q_{(2)}^{\prime \prime 0}$ for $q_{1}^{0}$. One and only one trajectory will correspond to those initial values, when taken at random. All along that trajectory, the motion will be defined by any one of the equalities:

$$
d t=d q_{1} \sqrt{\frac{\chi_{i}}{\psi_{i}}} \equiv d q_{1} \sqrt{\frac{q_{(i)}^{\prime \prime}+q_{(i)}^{\prime} \pi_{1}-\pi_{i}}{\beta_{i}-q_{(i)}^{\prime} \beta_{1}}} .
$$

One sees that $t$ is calculated by a quadrature and that for a given trajectory, one will have:

$$
t= \pm f\left(q_{1}\right)+\text { const. }
$$

in which $f$ denotes a well-defined function of $q_{1}$ (which depends upon the trajectory).
The values $q_{1}^{0}, q_{2}^{0}, \ldots, q_{k}^{0}, q_{1}^{0}, \ldots, q_{k}^{\prime 0}$ correspond to one and only one system of values $q_{1}^{0}$, $q_{2}^{0}, \ldots, q_{k}^{0}, q_{(2)}^{\prime 0}, \ldots, q_{(k)}^{\prime 0}, q_{(2)}^{\prime 0}$, and that system will not change when one changes the signs of all the $q_{i}^{\prime 0}$. It follows from this that if each point of the material system occupies the same position with equal, but directly opposite, velocities at the instants $t_{1}$ and $t_{2}$ then one can pass from the first motion to the second one by changing $t$ into $t_{1}+t_{2}-t$.

Indeed, one has:

$$
t-t_{1}=+f\left(q_{1}\right)-f\left(q_{1}^{0}\right)
$$

for the first motion, and:

$$
t-t_{2}=-f\left(q_{2}\right)-f\left(q_{1}^{0}\right)
$$

for the second one, if $q_{1}^{0}$ is the common value of $q_{1}$ at the two instants $t_{1}$ and $t_{2}$.
One sees that a given trajectory can be traversed in only two distinct manners. For each position of the system along that trajectory, the velocities $q_{1}^{\prime}, q_{2}^{\prime}, \ldots, q_{k}^{\prime}$ will be determined, up to sign, by the equalities:

$$
q_{1}^{\prime}= \pm \sqrt{\frac{\psi_{i}}{\chi_{i}}}, q_{2}^{\prime}=q_{(2)}^{\prime} q_{1}^{\prime}, \ldots, q_{k}^{\prime}=q_{(k)}^{\prime} q_{1}^{\prime}
$$

In particular, if that trajectory is a geodesic then one will always have $q_{1}^{\prime}=\infty$, or if one prefers, $d t / d q_{1}=0$. Conversely, if one takes arbitrary values for $q_{1}^{0}, \ldots, q_{k}^{0}, q_{(2)}^{\prime 0}, \ldots, q_{(k)}^{\prime 0}, q_{(2)}^{\prime \prime 0}$ that annul $\chi_{2}$ then the trajectory that is defined by those values will be unique, and as a result, it will coincide with the geodesic that satisfies those initial conditions, since that geodesic is also a trajectory. All along that trajectory, $d t / d q_{1}$ will then be zero. One can further say that the congruence of geodesics is the $(2 k-2)$-parameter congruence that is obtained by subjecting the $(2 k-1)$ parameters of the trajectory to the condition that $1 / T_{0}=0$. In particular, if there exists a force function then each
value $h$ of the constant of the vis viva integral will correspond to ( $2 k-1$ ) parameters, so the congruence of geodesics will correspond to the value $h=\infty$.

Remarkable trajectories. - Nevertheless, there can exist exceptional trajectories that we call remarkable and for which the preceding conclusions break down: They are the trajectories that give the form $0 / 0$ to the ratios $\chi_{i} / \psi_{i}$; in other words, that they satisfy both of the equalities:

$$
q_{(i)}^{\prime \prime}+\pi_{1} q_{(i)}^{\prime}-\pi_{i}=0 ; \quad \beta_{i}-q_{(i)}^{\prime} \beta_{1}=0 \quad(i=2,3, \ldots, k)
$$

Those trajectories are then the geodesics that simultaneously satisfy all of the equations:

$$
\begin{equation*}
q_{(2)}^{\prime}=\frac{\beta_{2}}{\beta_{1}}, \ldots, q_{(k)}^{\prime}=\frac{\beta_{k}}{\beta_{1}} . \tag{12}
\end{equation*}
$$

In general, the system (12) and the equations of the geodesics have no common integral: In all cases, from (12), those common integrals cannot depend upon more than $(k-1)$ arbitrary constants.

There is an infinitude of possible motions on one of those remarkable trajectories. Indeed, we can replace the equations of motion with equations (4):

$$
\begin{equation*}
0=q_{(i)}^{\prime \prime}+q_{(i)}^{\prime} \pi_{1}-\pi_{i}+\frac{\beta_{1} q_{(i)}^{\prime}-\beta_{i}}{q_{1}^{\prime 2}} \equiv \chi_{i}-\frac{\psi_{i}}{q_{1}^{\prime 2}}, \tag{4}
\end{equation*}
$$

combined with one of the Lagrange equations, or if one prefers, with the vis viva equality:

$$
T=\int\left[Q_{1}+Q_{2} q_{(2)}^{\prime}+\cdots+Q_{k} q_{(k)}^{\prime}\right] d q_{1} .
$$

By hypothesis, the trajectory considered will satisfy the equations $\chi_{i}=0, \psi_{i}=0$, so it will satisfy the equations (4). The motion on that trajectory will then be defined by the single equality:

$$
\begin{equation*}
\frac{1}{2} q_{1}^{\prime 2} \sum A_{i j} q_{(i)}^{\prime} q_{(j)}^{\prime}=\int\left[Q_{1}+Q_{2} q_{(2)}^{\prime}+\cdots+Q_{k} q_{(k)}^{\prime}\right] d q_{1} \tag{13}
\end{equation*}
$$

or rather:

$$
d t=d q_{1} \sqrt{\frac{\varphi\left(q_{1}\right)}{2\left[f\left(q_{1}\right)+h\right]}},
$$

in which $\varphi$ and $f$ are two functions of $q_{1}$ that are determined by the trajectory considered.
Observe that one can always assume that the coefficient $\varphi$ of $\frac{1}{2} q_{1}^{\prime 2}$ in (13) will be non-zero if the trajectory corresponds to a real motion of the system. Indeed, one has $2 T=$ $\sum m_{j}\left(x_{j}^{\prime 2}+y_{j}^{\prime 2}+z_{j}^{\prime 2}\right)$, and it is always legitimate to take $\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ to be $k$ of the coordinates $x_{j}, y_{j}, z_{j}$ that are independent. Under those conditions, $2 T$ will be annulled under a real motion of
the system only if all of the $q^{\prime}$ are zero $\left({ }^{4}\right)$. Since one has $\varphi=2 T\left[q_{1}, q_{2}, \ldots, q_{k}, 1, q_{(2)}^{\prime}, \ldots, q_{(k)}^{\prime}\right], \varphi$ cannot be annulled for a real position of the system and real values of $q_{(2)}^{\prime}, q_{(3)}^{\prime}, \ldots, q_{(k)}^{\prime}$. Observe, moreover, that any real motion of $S$ will correspond to a real trajectory, but the converse is not necessarily true. It can happen that $q_{1}, q_{2}, \ldots, q_{k}$ are real, but certain coordinates $x_{j}, y_{j}, z_{j}$ of the points $M_{j}$ of $S$ are imaginary.

The motion of the system on a remarkable trajectory is that of a system with complete constraints. The material points of such a system will traverse only one trajectory, but they can traverse it in an infinitude of ways.

Suppose, for example, that one is dealing with a free material point that is referred to the rectangular coordinates $x, y, z$ and is subjected to the force $(F)$ or $X, Y, Z$. The remarkable trajectories are the lines $D$ (if they exist) such that all along $D$, the force has $D$ for its line of action. Notably, if $F$ is a force that issues from the origin $O$ then the remarkable trajectories will be composed of all of the lines that pass through $O$, which is a congruence with $(k-1)=2$ parameters.

Similarly, if $M$ is a point that moves on a surface $\Sigma$ then the exceptional trajectories are the geodesics of $\Sigma$ that are tangent at each point to the projection of the force onto the plane tangent to $\Sigma$.

## Study of real trajectories.

Consider a real trajectory that is not a remarkable trajectory. The motion along that trajectory is defined by any one of the equalities $\frac{d t}{d q_{1}}=\sqrt{\frac{\chi_{i}}{\psi_{i}}}$. It will then be real if $\chi_{i} / \psi_{i}$ is positive and imaginary if $\chi_{i} / \psi_{i}$ is negative. The real trajectories $(C)$ are then subdivided into two categories: The trajectories $\left(C^{\prime}\right)$ along which the motion is real and the trajectories $\left(C^{\prime \prime}\right)$ along which it is imaginary. However, that raises the question: Can part of the same analytic trajectory belong to the class $\left(C^{\prime}\right)$, while part of it belongs to the class $\left(C^{\prime \prime}\right)$ ? In order to answer that question concisely, it is necessary to discuss certain properties of motion in detail.

## Exhibiting some properties of motion:

In what follows, we shall suppose that the parameters $q_{i}$ have been chosen in such a fashion that any real position of the system $S$ corresponds to real values of $q_{i}$, and that $\Delta$, the discriminant of $T$, is not annulled for any real position of $S$. For any system of real values of $q_{i}$ that corresponds

[^3]to real positions of $S$, the functions $A_{i j}, Q_{i}$ can have several real determinations ( ${ }^{5}$ ). However, we assume that those determinations define just as many continuous functions of the real variables $q_{i}$ that possess continuous first and second derivatives, except perhaps for certain exceptional values that we shall call singular values of the values of $q_{i}$ that correspond to singular positions of the system $S$.

Now consider a system of real values of the $q_{i}$, namely, $q_{1}=\alpha_{1}, q_{2}=\alpha_{2}, \ldots, q_{k}=\alpha_{k}$, in the neighborhood of which the determination that is made for each of the functions $A_{i j}, Q_{i}$ is regular, by which I mean that they remain well-defined and continuous, along with their first and second derivatives. If we solve the Lagrange equations (1) that define the motion of $S$ for the $q_{i}^{\prime \prime}$ then we will get the equations:

$$
\begin{equation*}
\frac{d q_{i}}{d t}=q_{i}^{\prime}, \quad \frac{d q_{i}^{\prime}}{d t}=P_{i}+\beta_{i} \tag{2}
\end{equation*}
$$

whose right-hand sides are well-defined and continuous, along with their first derivatives, while the variables $q_{i}^{\prime}$ will remain finite and the $q_{i}$ will remain close to $a_{1}, a_{2}, \ldots, a_{k}$. There will then exist one and only one system of integrals of (2), namely, $q_{1}(t), \ldots, q_{k}(t), q_{1}^{\prime}(t), \ldots, q_{k}^{\prime}(t)$, such that for $t=t_{0}$, the $q_{i}$ will take the values $a_{i}$, and the $q_{i}^{\prime}$ will take the values $q_{i}^{\prime 0}$ that are given in advance.

Let us make that theorem more precise: One can find a number $\delta$ such that the functions $Q_{i}$ are regular for the values of $q_{i}$ that satisfy the equalities:

$$
\left|q_{i}-a_{i}\right| \leq \delta \quad(i=1,2, \ldots, k)
$$

If $L$ is an arbitrary number that we take to be greater than $\delta$ then the right-hand sides of equations (2) will remain less than a certain limit $M$ in absolute value, while the inequalities:

$$
\left|q_{i}-a_{i}\right| \leq \delta, \quad\left|q_{i}^{\prime}\right| \leq L \quad(i=1,2, \ldots, k)
$$

will be verified. From that, let $q_{i}^{0}, q_{i}^{\prime 0}$ be a system of values that satisfies the conditions:

$$
\left|q_{i}^{0}-\alpha_{i}\right| \leq \frac{\delta}{2}, \quad\left|q_{i}^{\prime 0}\right| \leq \frac{L}{2} \quad(i=1,2, \ldots, k) .
$$

For any system $q_{i}, \ldots, q_{i}^{\prime}$ such that one has:

[^4]$$
\left|q_{i}-q_{i}^{0}\right| \leq \frac{\delta}{2}, \quad\left|q_{i}^{\prime}-q_{i}^{\prime 0}\right| \leq \frac{\delta}{2},
$$
the right-hand sides of (2) will remain less than $M$ in modulus. From the fundamental theorem on differential equations $\left({ }^{6}\right)$, the system of integrals $q_{i}(t), q_{i}^{\prime}(t)$ that takes the values $q_{i}^{0}, q_{i}^{\prime 0}$ for $t=$ $t_{0}$ will be a system of functions of $t$ that is well-defined and continuous, at least in the interval from $t_{0}-\delta / 2 M$ to $t_{0}+\delta / 2 M$.

Let us study what happens for initial values $q_{i}^{\prime 0}$ that are very large in absolute value. Take one of the parameters $q$ to be an independent variable, which is chosen in such a way that its derivative $(d q / d t)$ is not less than any of the other derivatives ${q_{i}^{\prime 0}}^{0}$ in absolute value: Let $q_{1}$ be that parameter. If we set:

$$
\frac{1}{q_{1}^{\prime}} \equiv \frac{d t}{d q_{1}}=r_{1}
$$

then equations (2) will become:

$$
\begin{cases}\frac{d t}{d q_{1}}=r_{1}, & \frac{d r_{1}}{d q_{1}}=-r_{1} \Pi_{1}-r_{1}^{2} \beta_{1} \\ \frac{d r_{1}}{d q_{1}}=q_{(i)}^{\prime}, & \frac{d q_{(i)}^{\prime}}{d q_{1}}=\Pi_{1}-q_{(i)}^{\prime} \Pi_{1}+\left(\beta_{i}-q_{(i)}^{\prime} \beta_{1}\right) r_{1}^{2}\end{cases}
$$

Let $M_{1}$ be the maximum modulus of the right-hand sides of equations ( $2^{\prime}$ ) when one has, simultaneously:

$$
\left|q_{i}-a_{i}\right| \leq \delta, \quad r_{1} \leq L_{2}, \quad\left|q_{(i)}^{\prime}\right| \leq 1+\frac{\delta}{2} \quad(i=1,2, \ldots, k)
$$

(with the condition $L_{1}>\delta$ ). On the other hand, let the values $q_{i}^{0}, r_{1}^{0}, q_{(i)}^{\prime 0}$ satisfy the conditions:

$$
\left|q_{i}^{0}-\alpha_{i}\right| \leq \frac{\delta}{2}, \quad r_{1}^{0} \leq \frac{L_{1}}{2}, \quad\left|q_{i}^{\prime 0}\right| \leq 1 .
$$

There exists one and only one system of integrals of $\left(2^{\prime}\right) t\left(q_{1}\right), r_{1}\left(q_{1}\right), r_{i}\left(q_{1}\right), q_{(i)}^{\prime}\left(q_{1}\right)$ that takes the values $t_{0}, q_{i}^{0}, r_{1}^{0}, q_{(i)}^{\prime 0}$ for $q_{1}=q_{1}^{0}$, and those integrals will be continuous, at least in the interval from $q_{1}^{0}-\varepsilon_{1}$ to $q_{1}^{0}+\varepsilon_{1}$, when $\varepsilon_{1}$ denotes the quantity $\delta / 2 M_{1}$.

In particular, if $r_{1}^{0}=0$ then that system will have the form:

$$
t \equiv t_{0}, \quad r_{1} \equiv 0, \quad q_{i}=\varphi_{i}\left(q_{1}\right), \quad q_{(i)}^{\prime}=\varphi_{i}^{\prime}\left(q_{1}\right),
$$

[^5]as equations ( $2^{\prime}$ ) will show immediately $\left({ }^{7}\right)$.
That shows that $r_{1}$ will never be annulled in the interval $q_{1}^{0}-\varepsilon_{1}$ to $q_{1}^{0}+\varepsilon_{1}$, or at least it will not be identically zero. $t$ will then vary constantly in the same sense when $q_{1}$ grows from $q_{1}^{0}-\varepsilon_{1}$ to $q_{1}^{0}+\varepsilon_{1}$ and will pass from the value $t_{0}-\eta$ to the value $t_{0}+\eta^{\prime}$ or from the value $t_{0}+\eta^{\prime}$ to the value $t_{0}-\eta$. It follows from this that $q_{1}(t), q_{2}(t), \ldots, q_{k}(t)$ are functions of $t$ that are continuous, along with their first derivatives, in the interval from $t_{0}-\eta$ to $t_{0}+\eta^{\prime}$, and that $q_{1}$ passes from one of the values $q_{1}^{0} \pm \varepsilon_{1}$ to the other when $t$ varies in that interval.

The same argument can be repeated when the independent variable is a different parameter $q_{j}$. Equations (2) will correspond to analogous equations: Let $M_{j}$ be the maximum modulus of the right-hand sides of those equations when one has, at the same time:

$$
\left|q_{i}-a_{i}\right| \leq \delta, \quad\left|\frac{d t}{d q_{j}}\right| \leq L_{1}, \quad\left|\frac{d q_{i}}{d q_{j}}\right| \leq 1+\frac{\delta}{2} \quad(i=1,2, \ldots, k) .
$$

The quantity $\varepsilon_{1}=\delta / 2 M$ will correspond to the quantity $\varepsilon_{j}=\delta / 2 M_{j}$. I shall let $\varepsilon$ denote the smallest of the quantities $\varepsilon_{j}$.

Those propositions allow us to prove an important property of motion. Consider the system of integrals $q_{i}(t), q_{i}^{\prime}(t)$ of equations (2) such that for $t=t_{0}, q_{i}=q_{i}^{0}, q_{i}^{\prime}=q_{i}^{\prime 0}$, the values of $q_{i}^{0}$ are not singular values of $A_{i j}, Q_{i}$. When one makes $t$ increase when starting from $t_{0}$, several situations can present themselves: The motion might remain regular for any value of $t$. (By that, I mean that it remains finite and continuous, and the system $S$ does not pass through any singular position, moreover.) One or more of the parameters $q_{i}$ might become infinite or indeterminate when $t$ tends to a certain value $t_{1}$. Finally, the system $S$ might tend to a singular position. However, can it happen that when $t$ tends to $t_{1}$, the system $S$ will tend to a non-singular position and the velocities $q_{i}^{\prime}$ will become indeterminate or infinite? We shall see that this can never happen. More precisely, assume that when $t$ tends to $t_{1}$, the parameters $q_{1}, q_{2}, \ldots, q_{k}$ tend to the values $a_{1}, a_{2}, \ldots, a_{k}$, respectively, in whose neighborhood the determinations that were taken for the $A_{i j}, Q_{i}$ remain regular. Under those conditions, the $q_{i}^{\prime}$ will tend to finite limits, respectively, and the motion will remain regular outside of the instant $t_{1}$.

Indeed, there are two possibilities: Either the $q_{i}^{\prime}$ all tend to zero when $t$ tends to $t_{1}$ (which would prove the theorem) or the modulus of at least one of the $q_{i}^{\prime}$ is greater than a certain limit $\lambda$ for certain values of $t$ that are as close to $t_{1}$ as one desires. Then consider the number $\varepsilon$ that was introduced above and corresponds to the conditions:

$$
\left|q_{i}-a_{i}\right| \leq \delta, \quad\left|\frac{d t}{d q_{j}}\right| \leq \frac{1}{\lambda}, \quad\left|\frac{d q_{j}}{d q_{i}}\right| \leq 1+\frac{\delta}{2} \quad(i=1,2, \ldots, k) .
$$

[^6](1/ $\lambda$ replaces $L_{1}$ ). By hypothesis, we can find an instant $t^{\prime}$ that is sufficiently close to $t_{1}$ that when $t$ varies from $t^{\prime}$ to $t_{1}$, each variable $q_{i}$ will remain between $a_{i}-\alpha$ and $a_{i}+\alpha$, where $\alpha$ denotes an arbitrary number that is less than $\delta / 2$ and $\varepsilon / 2$. Now let $t_{0}$ be a value of $t$ that is found between $t^{\prime}$ and $t_{1}$ and for which the greatest of the moduli $\left|q_{i}^{\prime}\right|$, namely, $\left|q_{1}^{\prime}\right|$, exceeds $\lambda$. One has, for $t=t_{0}$ :
$$
\left|q_{i}^{0}-a_{i}\right|<\frac{\delta}{2}, \quad\left|\frac{d t}{d q_{1}}\right|_{0}<\frac{1}{\lambda},\left|\frac{d q_{i}}{d q_{1}}\right|_{0} \leq 1 \quad(i=1,2, \ldots, k)
$$

When $q_{1}$ varies from $q_{i}^{0}-\varepsilon$ to $q_{i}^{0}+\varepsilon$, $t$ will vary from $t_{0}-\eta$ to $t_{0}+\eta^{\prime}$ (or from $t_{0}+\eta^{\prime}$ to $t_{0}-\eta$ ). I say that $t_{1}$ is found between $t_{0}$ and $t_{0}+\eta^{\prime}$. In other words, when $t$ varies from $t_{0}$ to $t_{0}+\eta^{\prime}$, and as a result between $t^{\prime}$ and $t_{1}, q_{1}$ will vary between $q_{1}^{0}$ and $q_{1}^{0} \pm \varepsilon$. However, between $t^{\prime}$ and $t_{1}$, one has: $\left|q_{1}-a_{i}\right| \leq \alpha \leq \varepsilon / 2$, so there will be two values of $q_{1}$ in that interval that cannot differ by $\varepsilon$. Therefore, the instant $t_{1}$ is found between $t_{0}$ and $t_{0}+\eta^{\prime}$, and since the functions $q_{i}(t), q_{i}^{\prime}(t)$ are continuous in that interval, the motion will remain regular at the instant $t_{1}$ and beyond. Q.E.D.

We can then state the following theorem:

## Theorem:

When the system $S$ tends to a non-singular position as tends to $t_{1}$, its velocities will tend to $a$ limit, and the motion can be continued regularly beyond $t_{1}$.

If all of the $q_{i}^{\prime}$ are non-zero for $t=t_{1}$ (say $q_{1}^{\prime} \neq 0$ ) then the ratios $q_{i}^{\prime} / q_{1}^{\prime}$ will have well-defined values. The same thing will be true when the $q_{i}^{\prime}$ 's are all annulled for $t=t_{1}$. Indeed, one has:

$$
q_{i}^{\prime \prime}=\frac{\left(t-t_{1}\right)^{2}}{1 \cdot 2}\left[\beta_{i}\left(a_{1}, a_{2}, \ldots, a_{k}\right)+\delta_{i}\right]=\frac{\left(t-t_{1}\right)^{2}}{2}\left(\beta_{i}^{0}+\delta_{i}\right) \quad(i=1,2, \ldots, k)
$$

in that case, in which the $\delta_{i}$ tend to zero with $t-t_{1}$. Furthermore, none of the $\beta_{i}^{0}$ are zero. In other words, the unique system of integrals that satisfies the initial conditions $q_{i}=a_{i}, q_{i}^{\prime}=0\left(\right.$ for $\left.t=t_{1}\right)$ will be the system $q_{i}(t) \equiv a_{i}, q_{i}^{\prime}(t) \equiv 0\left({ }^{8}\right)$. Therefore, let $\beta_{1}^{0} \neq 0$. The ratios $\frac{q_{i}^{\prime}}{q_{1}^{\prime}}=\frac{d q_{i}}{d q_{1}}$ take the values $\frac{\beta_{i}^{0}}{\beta_{1}^{0}}$ for $t=t_{1}$. That shows that as $t$ tends to $t_{1}$, the system $S$ cannot tend to a (regular) equilibrium position with a vis viva that tends to zero.

[^7]I must add that the integrals $q_{i}(t)$ are even functions of $t-t_{1}$ then. In other words, if one sets $t$ $-t_{1}=\tau=\theta^{2}$ then one will have:

$$
q_{i}=\alpha_{i}+\beta_{i} \theta+c_{i} \theta^{2}+\ldots \quad(i=1,2, \ldots, k)
$$

Indeed, if one changes $\tau$ into $-\tau$ in the integrals $q_{i}(\tau)$ then one will again get a system of integrals. When the first system satisfies the initial conditions $q_{i}^{0}, q_{i}^{0}$ for $t=t_{1}$, the second one will satisfy the conditions $q_{i}^{0},-q_{i}^{\prime 0}$. Since the $q_{i}^{\prime 0}$ are zero here, the initial conditions will remain the same, and the two systems of integrals will then coincide: $q_{i}(\tau) \equiv q_{i}(-\tau)$. When $t$ goes beyond the instant $t_{1}$, the system $S$ will reverse: At the instant $t_{1}+\tau$, it will pass through the same position that it passed through at the instant $t_{1}-\tau$, but the velocities will have changed sense.

Now suppose that the system $S$ tends to a non-singular position $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ as $t$ increases indefinitely. I say that all of the velocities necessarily tend to zero. Indeed, assume that this is not true and repeat the argument that was made above while keeping the same notation: By hypothesis, for any value of $t$ that is greater than a certain limit $t^{\prime}$, one will have: $\left|q_{i}-a_{i}\right| \leq \alpha \leq \varepsilon / 2(i=1,2$, $\ldots, k)$. However, on the other hand, there exist values $t_{0}$ of $t$ that are greater than $t^{\prime}$, and are such that at least one of the parameters, say $q_{1}$, varies by $\varepsilon$ when $t$ varies from $t_{0}$ to $t_{0}+\eta^{\prime}$. There is then a contradiction. The $q_{i}^{\prime}$ tend to zero with $1 / t$.

Moreover, the position $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ is a position of equilibrium of the system $S$. Indeed, make the change of variables $t=1 / \theta$. One will have:

$$
\frac{d q_{i}}{d t}=q_{i}^{\prime}=-\theta^{2} \frac{d q_{i}}{d \theta}, \quad \frac{d q_{i}^{\prime}}{d t}=-\theta^{2} \frac{d q_{i}^{\prime}}{d \theta}
$$

and as a result:

$$
\frac{d q_{i}}{d \theta}=-\frac{q_{i}^{\prime}}{\theta^{2}}, \quad \frac{d q_{i}^{\prime}}{d \theta}=\frac{-\beta_{i}+\Pi_{i}}{\theta^{2}}
$$

By hypothesis, when $\theta$ tends to zero, the $q_{1}, \ldots, q_{k}$ will tend to $a_{1}, \ldots, a_{k}$, and the $q_{i}^{\prime}$ will tend to zero.

It follows from this that the $\beta_{i}\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ must be zero because if $\beta_{i}\left(a_{1}, a_{2}, \ldots, a_{k}\right)=\beta_{i}^{0}$ then one will have:

$$
\frac{d q_{i}^{\prime}}{d \theta}=\frac{-\beta_{i}^{0}+\delta_{i}^{\prime}}{\theta^{2}}
$$

$\delta_{i}^{\prime}$ will tend to zero with $\theta$, and $q_{i}^{\prime}$ will increase indefinitely when $\theta$ tends to zero. It is then necessary that $\beta_{1}^{0}, \beta_{2}^{0}, \ldots, \beta_{k}^{0}$ must be zero.

We know, moreover, that conversely when $t$ increases, the system $S$ cannot tend to a (regular) equilibrium position with a vis viva that is annulled without $t$ increasing beyond any limit.

It would be fitting to observe here that the $q_{i}^{\prime}$ tend to zero without the ratios $\frac{q_{i}^{\prime}}{q_{1}^{\prime}}=\frac{d q_{i}}{d q_{1}}$ necessarily having a limit. In order to convince ourselves of that, consider the motion in a plane of a point $(x, y)$ that is subject to the force $X=2 y, Y=-2 x$. The equations of motion:

$$
x^{\prime \prime}=2 y, \quad y^{\prime \prime}=-2 x
$$

admit the family of integrals:

$$
x=e^{-t}[\alpha \cos t-\beta \sin t], \quad y=e^{-t}[\alpha \sin t+\beta \cos t],
$$

in which $\alpha, \beta$ are two arbitrary real constants. When $t$ increases indefinitely, $x$ and $y$ will tend to zero, as well as $x^{\prime}, y^{\prime}$, but the ratio $\frac{d y}{d x}=\frac{\alpha \tan t+\beta}{\alpha-\beta \tan t}$ will not tend to any limit.

In the foregoing, we supposed that $t$ is increasing. However, all of the conclusions will obviously persist if one makes $t$ decrease, since it is legitimate to change $t$ into $-t$.

Properties of real trajectories. Return to the study of trajectories. - To abbreviate the language, let us agree to regard $q_{1}, q_{2}, \ldots, q_{k}$ as the $k$ rectangular coordinates of a point $M$ in the $k$ dimensional space $E_{k}$. The trajectories $q_{i}=\varphi_{i}\left(q_{1}\right)[i=2, \ldots, k]$ will be curves $C$ in that space, and the differentials $d q_{1}, d q_{2}, \ldots, d q_{k}$ will define the direction of the tangent at a point on one such curve.

Finally, set:

$$
d \sigma=\sqrt{d q_{1}^{2}+d q_{2}^{2}+\cdots+d q_{k}^{2}},
$$

in which the arc-length $\sigma$ denotes the length of the segment of the curve $C$ that is found between two points $M$ and $M^{\prime}$, and extend the integral $\int \sqrt{d q_{1}^{2}+d q_{2}^{2}+\cdots+d q_{k}^{2}}$ along the curve $M M^{\prime}$ (all of the elements being positive). When a curve $(C)$ is regular (by that, I mean that it admits a continuous tangent at each point), one can suppose that $q_{1}, q_{2}, \ldots, q_{k}$ are expressed as functions of arc-length $\sigma$, which is measured by starting from a fixed point $M_{0}$ and proceeding positively in one sense and negatively in the other. Each value of $\sigma$ will then correspond to a well-defined point ( $q_{1}$, $q_{2}, \ldots, q_{k}$ ) on (C).

We shall consider only the domain in the real space $E_{k}$ in which the points $\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ correspond to real positions of the system, and we shall study the trajectories that belong to that domain exclusively. (What we say will apply to other real trajectories, moreover.) We say singular points in the space $E_{k}$ to mean the points at which the functions $A_{i j}, Q_{i}$ cease to be regular. Under the most unfavorable hypothesis, those points will form a $(k-1)$-dimensional surface in $E_{k}$, namely, the surface $\psi\left(q_{1}, q_{2}, \ldots, q_{k}\right)=0$. In particular, that surface constitutes the boundary of the domain $E_{k}$ when that space is not considered to be enveloped by all of space.

Finally, we shall say equilibrium points $N$ to mean the points $q_{1}, q_{2}, \ldots, q_{k}$ that correspond to an equilibrium position of $S$, i.e., where $\beta_{1}, \beta_{2}, \ldots, \beta_{k}$ are annulled. In the most general case, those points will be isolated points.

True trajectories. Conjugate trajectories. Mixed trajectories. - If one changes $t$ into it in the equations of motion (1) then those equations will not be altered, except for the fact that the $Q_{i}$ will all change sign. The trajectories that are defined by the system (1) and the system that is obtained by changing $Q_{i}$ into $-Q_{i}$ will then coincide.

That new motion will be called the motion that is conjugate to the true motion. That motion is the motion of the system $S$ when one changes the senses of all the given forces without changing their directions or magnitudes. If the true motion is imaginary along a real trajectory $C^{\prime \prime}$ then the conjugate motion along that same trajectory will be real, since $\left(\frac{d t}{d q_{1}}\right)^{2}$, which is negative under the first motion, will change sign when one changes $t$ into $i t$. We give the name of true trajectories to the real arcs $\left(C^{\prime}\right)$ of the trajectory along which the true motion is real and the name of conjugate trajectories to the arcs $\left(C^{\prime \prime}\right)$ along which the conjugate motion is real. When one changes $t$ into $i t$, the two classes of trajectories $\left(C^{\prime}\right)$ and $\left(C^{\prime \prime}\right)$ will permute.

Having said that, give the system $S$ the real initial conditions $q_{1}^{0}, q_{2}^{0}, \ldots, q_{k}^{0}, q_{1}^{\prime 0}, q_{2}^{\prime 0}, \ldots, q_{k}^{\prime 0}$ at the instant $t_{0}$ and measure the arc-length $\sigma$ of the trajectory $C$ by starting from the initial point $M_{0}$ and proceeding in the sense that makes $\sigma$ begin by increasing with $t$. $\sigma$ will continue to increase with $t$ as long as $t$ does not attain a value $t_{1}$ for which the motion ceases to be regular or a value $t_{1}$ for which $d \sigma / d t$, and as a result, all of $q_{i}^{\prime}$, are annulled.

We adopt the first hypothesis to begin with: When $t$ tends to $t_{1}$, either the point $\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ does not tend to any point $M$ at a finite distance in the space $E_{k}$ (in which case, $\sigma$ would increase indefinitely) or the point ( $q_{1}, q_{2}, \ldots, q_{k}$ ) tends to a singular point $N$ in $E_{k}$. From the foregoing, no other case would be possible.

Under the second hypothesis, in which all of the $q_{i}^{\prime}$ are annulled when $t$ tends to $t_{1}, \sigma$ will increase up to a certain limit $\sigma_{1}$, then decrease and take on the same value at $t_{1}+\alpha$ that it had at $t_{1}$ $-\alpha$. The point ( $q_{1}, q_{2}, \ldots, q_{k}$ ) moves backward in its trajectory.

Indeed, the point $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$, or $M_{1}$, which corresponds to the value $\sigma_{1}$ of $\sigma$, cannot be an equilibrium point (otherwise, $t_{1}$ would be infinite), and the trajectory will be defined by equations of the form:

$$
q_{i}=a_{i}+b_{i} \theta+c_{i} \theta^{2}+\cdots \quad(i=1,2, \ldots, k)
$$

in the neighborhood of $M_{1}$, in which $\theta=\left(t-t_{1}\right)^{2}$, and none of the $b_{i}$ are zero. That shows us that $(C)$ can be extended up to the point $M_{1}$ while always admitting a continuous tangent. Furthermore, since one has:

$$
\frac{d \sigma}{d t}=\left(t-t_{1}\right)\left[-\sqrt{b_{1}^{2}+b_{2}^{2}+\cdots+b_{k}^{2}}-\varepsilon\right]=\left(t-t_{1}\right)(B+\varepsilon)
$$

(in which $\varepsilon$ tends to zero with $t-t_{1}$ ), $\sigma$ will pass through a maximum $\sigma_{1}$ for $t=t_{1}$. The point $M_{1}$ will be called a point of regression (point d'arret) of the trajectory ( $C$ ). The conjugate motion will be real along the segment $M_{1} M$ of $(C)$, which follows the segment $M_{0} M_{1}$. Finally, if the trajectory $(C)$ is not a remarkable trajectory then $\left(\frac{d \sigma}{d t}\right)^{2}$ will have a well-defined value $\varphi(\sigma)$ at each point of $M_{0} M$ : The equality $\left(\sigma-\sigma_{1}\right)=\left(t-t_{1}\right)^{2}\left(B / 2+\varepsilon^{\prime}\right)$ and its consequence $\left(\frac{d \sigma}{d t}\right)^{2}=$ $\left(\sigma-\sigma_{1}\right)\left(2 B+\varepsilon^{\prime}\right)$ prove that $\left(\frac{d \sigma}{d t}\right)^{2}$ will remain a continuous function of $\sigma$ (but with its sign changed) when one crosses a point of regression while varying $M$ from $M_{0}$ to $M$ along ( $C$ ).

We give the name of mixed trajectories to those trajectories $\Gamma$ that possess at least one point of regression $M_{1}$. They define a family that depends upon $k$ arbitrary constants, for example, the coordinates $a_{1}, a_{2}, \ldots, a_{k}$ of a point of regression. Indeed, take one of the derivatives $q_{i}^{\prime}\left(\right.$ say, $\left.q_{1}^{\prime}\right)$ to be the independent variable and study $q_{1}, q_{2}, \ldots, q_{k}, q_{2}^{\prime}, \ldots, q_{k}^{\prime}$ as functions of $q_{1}^{\prime}$. If one considers all of the trajectories then one can give the values $q_{1}^{0}, q_{2}^{0}, \ldots, q_{k}^{0}, q_{2}^{\prime 0}, \ldots, q_{k}^{\prime 0}$, for $q_{1}^{0}$ $=0$ arbitrarily. In order to get the mixed trajectories, one sets $q_{2}^{\prime 0}=0, \ldots, q_{k}^{\prime 0}=0$. The congruence of trajectories $\Gamma$ then depends upon $k$ arbitrary constants: An infinitude of trajectories that depend upon one parameter then passes through a given point $\left(q_{1}, q_{2}, \ldots, q_{k}\right)$. It can nonetheless happen that those $k$ parameters $a_{1}, a_{2}, \ldots, a_{k}$ are not distinct: In order for that to happen, it is necessary that an arbitrary mixed trajectory $\Gamma$ should correspond to an infinitude of values of the constants $a_{1}, a_{2}$, $\ldots, a_{k}$ such that one can take at least one of them arbitrarily.

It will then be necessary that all of the points of a segment of $\Gamma$ must be points of regression, and a result, that an infinitude of motions will be possible on $\Gamma$. In other words, the mixed trajectories must be remarkable trajectories. On the other hand, since at least one mixed trajectory will pass through an arbitrary point $M_{0}$ or $\left(q_{1}^{0}, q_{2}^{0}, \ldots, q_{k}^{0}\right)$ (namely, the one that admits $M_{0}$ as a point of regression), the congruence $\Gamma$ will depend upon at least $(k-1)$ distinct constants: We then arrive at the following conclusion: The congruence of mixed trajectories ( $\Gamma$ ) depends upon $k$ distinct constants, except in the case where there exists a $(k-1)$-parameter congruence of remarkable trajectories $\left({ }^{9}\right)$, in which case, that congruence will coincide with the congruence $(\Gamma)$.

Moreover, a remarkable trajectory $(\gamma)$ must be regarded as a mixed trajectory, in the sense that an arbitrary point $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ of the trajectory $\gamma$ must be a point of regression for one of the motions of $S$ along $\gamma$. Indeed, all of those motions are defined by an equality of the form (see page 242):

$$
\begin{equation*}
T=f\left(q_{1}\right)+k, \tag{a}
\end{equation*}
$$

[^8]in which $h$ is an arbitrary constant. Let $q_{1}=a_{1}$ for $t=t_{0}$, and let $h=-f\left(a_{1}\right)$ : At the point $\left(a_{1}, a_{2}\right.$, $\ldots, a_{k}$ ) of ( $\gamma$ ), the vis viva, and as a result $q_{1}^{\prime}, q_{2}^{\prime}, \ldots, q_{k}^{\prime}$, will be annulled, and the system $S$ will reverse along $(\gamma)$ when $t$ passes from $t_{0}-\varepsilon$ to $t_{0}+\varepsilon$ for the motion in question. However, there can exist motions that always take place in the same sense along $(\gamma)$ : In other words, it can happen that $f\left(q_{1}\right)+h$ is never annulled for other values of $h$, as we will soon verify in an example.

Now suppose that as $t$ increases indefinitely, the motion remains regular and the vis viva is not annulled: There two possible cases: Either the point $\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ does not tend to any point at a finite distance in the space $E_{k}$ (in which case, $\sigma$ would increase indefinitely with $t$ ) or $\left(q_{1}, q_{2}, \ldots\right.$, $\left.q_{k}\right)$ will tend to a point $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ that will then be an equilibrium point $N^{\prime}$. The trajectory does not necessarily have a tangent at that point and cannot be continued analytically beyond it.

From the foregoing, we can state the following conclusions:

Let $M_{0} M$ be a continuous fragment of the same real trajectory $(C)$ that does not pass through either a singular point $N$ of $E_{k}$ or an equilibrium point $N^{\prime}$ : The curve ( $C$ ) admits a continuous tangent along the $\operatorname{arc} M_{0} M$, and the total length of that axis is a certain finite number $\sigma$.

Furthermore, if $(C)$ is not a mixed trajectory (which is the general case in which the mixed trajectories depend upon only $k$ constants) then the arc $M_{0} M$ will always be traversed in the same sense during a finite time, whether the motion is a true or a conjugate one. That will also be true when the trajectory $(C)$ is mixed if it possesses no point of regression between $M_{0}$ and $M$. The entire arc $M M^{\prime}$ will then belong to either the class $\left(C^{\prime}\right)$ or the class $\left(C^{\prime \prime}\right)$.

If $(C)$ is a mixed trajectory (without being a remarkable trajectory) then, in general, it can possess only one point of regression $M_{1}$. When that point $M_{1}$ belongs to the arc $M_{0} M$, that arc will decompose into two parts $M_{0} M_{1}$ and $M_{1} M$, both of which are traversed twice in opposite senses (in a finite length of time), one of which will be the true motion, and the other of which will be the conjugate one. However, it can happen that there exist several points of regression $M_{1}, M_{2}, \ldots$ between $M_{0}$ and $M\left({ }^{10}\right)$, but there is always just a finite number of them. Indeed, suppose that the function $\left(\frac{d \sigma}{d t}\right)^{2}=\varphi(\sigma)$ admits an infinitude of zeroes between $M_{0}$ and $M$ that correspond to values (increasing, for example) $\sigma_{1}, \sigma_{2}, \ldots, \sigma_{n}, \ldots$ of $\sigma . \sigma_{n}$ remains less than $\sigma^{\prime}$ (viz., the length of $M_{0} M$, while tending to $\sigma^{\prime}$ as $n$ increases indefinitely. The function $\varphi(\sigma)$ is continuous, so it will be annulled when $\sigma=\sigma^{\prime}$, and one will have (when the corresponding point $M^{\prime}$ is not an equilibrium point):

$$
\varphi(\sigma)=\left(\sigma-\sigma^{\prime}\right)(2 B+\varepsilon) \quad(B \neq 0)
$$

which shows that $\varphi(\sigma)$ will admit no other zeroes besides $\sigma^{\prime}$ in the neighborhood of $\sigma^{\prime}$. The hypothesis is therefore absurd. From that, one can always decompose the arc $M_{0} M$ into a finite number of segments that either belong to the class $\left(C^{\prime}\right)$ or to the class $\left(C^{\prime \prime}\right)$ entirely.

[^9]If $(C)$ is a remarkable trajectory then each point of $M_{0} M$ is a point of regression for a one of the corresponding motions, but one can always take $h$ in the equality:

$$
T=f\left(q_{1}\right)+h=F(\sigma)+h
$$

to be sufficiently large that $M_{0} M$ is traversed in the same sense in its entirety. Furthermore, none of the motions that take place along $C$ can admit two points of regression between $M_{0}$ and $M$. Indeed, assume that $T$ is annulled at $M_{1}$ and $M_{2}$ : From the equality $(\alpha), F^{\prime}(\sigma)$ (which is continuous between $M_{0}$ and $M$ ) will be annulled between $M_{1}$ and $M_{2}$, and therefore between $M_{0}$ and $M$. Let $\sigma^{\prime}$ be the first zero of $F^{\prime}(\sigma)$ that one encounters upon starting from a point $\mu$ of $M_{0} M$ where $F^{\prime}(\sigma)$ is not zero and proceeding towards $M$ (or towards $M_{0}$ ). I say that the point $M^{\prime}$ or $\sigma^{\prime}$ of $(C)$ is an equilibrium point $N^{\prime}$. In order to see that, it will suffice to consider the motion along (C) that is defined by the equality:

$$
T=F(\sigma)-F\left(\sigma^{\prime}\right)=\left(\sigma-\sigma^{\prime}\right)^{2} F_{1}(\sigma),
$$

since $F_{1}(\sigma)$ is annulled at most once between $\mu$ and $M^{\prime}$ (say, at $M_{1}$ ), and the sign of $F_{1}(\sigma)$ is constant along a finite arc $M_{1} M^{\prime}$. Upon supposing that it is positive (which is legitimate, since otherwise one could change $t$ into $i t$ ), one will have:

$$
d t=\frac{d \sigma}{\sigma-\sigma^{\prime}} G(\sigma)
$$

in which $G$ remains greater than a certain positive number when $\sigma$ varies between $\sigma^{\prime}-\alpha$ and $\sigma^{\prime}$. $t$ will then increase indefinitely when $\sigma$ tends to $\sigma^{\prime}$, or rather $\sigma$ will tend to $\sigma^{\prime}$ when $t$ increases indefinitely. That will be possible only if $M^{\prime}$ is an equilibrium point $\left({ }^{11}\right)$. It will follow from this
$\left({ }^{11}\right)$ It is easy to verify that conclusion as follows: One has:

$$
F^{\prime}(\sigma)=Q_{1} \frac{d q_{1}}{d \sigma}+\cdots+Q_{k} \frac{d q_{k}}{d \sigma},
$$

but on the other hand, since the trajectory $(C)$ is remarkable, one knows that $d q_{1} / \beta_{1}=d q_{i} / \beta_{i}$. Therefore:

$$
F^{\prime}(\sigma)=\frac{Q_{1} \beta_{1}+\cdots+Q_{k} \beta_{k}}{\sqrt{\beta_{1}^{2}+\cdots+\beta_{k}^{2}}},
$$

or rather [since $Q_{i}=\sum_{j=1}^{k} A_{i j} B_{j}$, from equations (6) on page 236]:

$$
F^{\prime}(\sigma)=\frac{T\left(q_{1}, q_{2}, \ldots, q_{k}, \beta_{1}, \beta_{2}, \ldots, \beta_{k}\right)}{\sqrt{\beta_{1}^{2}+\cdots+\beta_{k}^{2}}},
$$

which is an expression that be zero only if all of the $\beta_{i}$ are zero. When all of the points of $(C)$ are equilibrium points, the geodesic ( $C$ ) will be traversed with a constant vis viva (that is arbitrary, moreover), since there are no given forces.
that for $h$ sufficiently large ( $h>h_{1}$ ), the segment $M_{0} M$ will be traversed in its entirety under the corresponding motion. For $h$ less than a certain limit $h_{2}$, the same segment will be traversed in its entirety under the conjugate motion. For $h$ between $h_{1}$ and $h_{2}$, the segment will decompose into two continuous segments, one of which is traversed under the true motion, while the other one is traversed under the conjugate motion.

Branches of singular trajectories. - Now assume that the segment considered $M_{0} M$ includes equilibrium points $N^{\prime}$ (but not singular points). Let $N^{\prime}$ be the first equilibrium point that one encounters along $M_{0} M$ upon starting from $M_{0}$. Now study the segment $M_{0} N^{\prime}$, while first supposing that $(C)$ is not a remarkable trajectory.

Under that hypothesis, $\left(\frac{d \sigma}{d t}\right)^{2}$ will be a function of $\sigma . \varphi(\sigma)$ will be continuous as long as $M$ remains between $M_{0}$ and $N^{\prime}$, but what will happen that when one makes $M$ tend to $N^{\prime}$ ? I say that $\varphi(\sigma)$ will tend to a limit. First of all, if $\varphi(\sigma)$ tends to zero then the proposition will have been proved. If that is not true then one can find points $M$ along $(C)$ that are close as one wishes to $N^{\prime}$ $\left({ }^{12}\right)$ and are such that $f(\sigma)$ (and as a result, at least one of the $q_{1}^{\prime 2}$ ) has an absolute value at $M$ that is greater than a certain fixed number $\lambda^{2}$. Upon repeating the argument on page 246) identically (either with the true motion or the conjugate motion), one will see that the system $S$ must reach the position $N^{\prime}$ in a finite length of time with a well-defined, finite vis viva, and as a result, the motion, like the trajectory $(C)$, can be regularly extended beyond $N^{\prime}$. In summary, $f(\sigma)$ will tend to a finite limit $f\left(N^{\prime}\right)$ when one makes $M$ tend to $N^{\prime}$ along the $\operatorname{arc} M_{0} N^{\prime}$, and if that limit $f\left(N^{\prime}\right)$ is not zero then the point $N^{\prime}$ will be an ordinary point of $(C)$.

If, on the contrary, $f\left(N^{\prime}\right)=0$ then we say that the arc $M_{0} N^{\prime}$ is a singular branch: There are two cases to be distinguished according to whether $f(\sigma)$ is or is not annulled an infinitude of times between $M_{0}$ and $N^{\prime}$. We first consider the latter case.

1. There exist only a finite number of zeroes of $f(\sigma)$ between $M_{0}$ and $N^{\prime}$. It will then suffice to consider the segment $M_{1} N^{\prime}$ that is adjacent to $N^{\prime}$ and in which $f(\sigma)$ keeps a constant sign, and it is legitimate to suppose that it is positive. When $t$ increases and the system is placed between $M_{1}$ and $M$ with a positive value of $d \sigma / d t$, it will tend to $N^{\prime}$, and cannot attain the equilibrium position in a finite length of time. The system will then tend to $N^{\prime}$ along $(C)$ when $t$ increases indefinitely. It can happen that the real curve ( $C$ ) has no tangent at the point $N^{\prime}$ and that it cannot be prolonged beyond $N^{\prime}$. As for the arc-length $\sigma$, or $M_{0} N$, it can tend to a finite limit $\sigma^{\prime}$ or increase indefinitely when $M$ tends to $N^{\prime}$. The equations $x^{\prime \prime}=2 y, y^{\prime \prime}=-2 x$ that were cited above (page 249) offer us an example of this first case.

[^10]2. $f(\sigma)$ admits an infinitude of zeroes $M_{1}, M_{2}, \ldots, M_{n}, \ldots M_{0}$ and $N^{\prime}$. The points of regression $M_{1}, M_{2}, \ldots, M_{n}, \ldots$ tend to $N^{\prime}$ when $n$ increases indefinitely. The segment $M_{0} M$ will then decompose into an infinitude of segments $M_{1} M_{2}, M_{2} M_{3}, \ldots$ that tend to $N^{\prime}$ and correspond to as many periodic motions (which alternate between true and conjugate ones) whose amplitudes will tend to zero. As an example of that case (which is clearly an exception), we cite the equations:
\[

\left\{$$
\begin{array}{l}
x^{\prime \prime}=\frac{1}{2} x y\left[\frac{7}{4}\left(x^{2}+y^{2}\right)^{2}-3\right]+\frac{1}{4}\left(x^{2}+y^{2}\right)\left[5 y^{2}-x^{2}\right]  \tag{A}\\
y^{\prime \prime}=\frac{7}{8} y^{2}\left(x^{2}+y^{2}\right)^{2}-\frac{3}{2} x y\left(x^{2}+y^{2}\right)+\frac{1}{2} x^{2}-y^{2}
\end{array}
$$\right.
\]

One of the trajectories that are defined by $(A)$ is the following one: $x=\frac{\cos \theta}{\theta^{1 / 2}}, y=\frac{\sin \theta}{\theta^{1 / 2}}$, in which once more $\theta=1 / r^{2}$, in terms of polar coordinates. That trajectory passes through the origin (an asymptotic point), which is an equilibrium point, and the corresponding motion will be defined by the equality:

$$
d t=\frac{d \theta \theta^{1 / 4}}{\sqrt{\sin \theta}}=\frac{-2}{t^{3}} \frac{d r}{\sqrt{r \sin 1 / r^{2}}} .
$$

Along each arc of the curve $2 n \pi<\theta<(2 n+1) \pi$, the motion will be real and periodic. Along the $\operatorname{arcs}(2 n+1) \pi<\theta<2(n+1) \pi$, it will be the conjugate motion that is periodic.

When $M$ tends to $N^{\prime}, \sigma$ will tend to a limit $\sigma^{\prime}$ or increase indefinitely, so the real curve ( $C$ ) might or might not be analytically continued beyond $N^{\prime}$ according the situation. Finally, it can happen that it includes an infinitude of equilibrium positions that form a sequence.

If one considers the equality $d t=d \sigma / \sqrt{\varphi(\sigma)}$ then $\varphi(\sigma)$ will change sign at each of the zeroes $\sigma_{1}, \sigma_{2}, \ldots$, of $\varphi(\sigma)$, and when $M$ tends to $N^{\prime}, t$ will satisfy the relation:

$$
t-t_{0}=\int_{\sigma_{1}}^{\sigma_{1}} \frac{d \sigma}{\sqrt{\varphi(\sigma)}}+i \int_{\sigma_{1}}^{\sigma_{2}} \frac{d \sigma}{\sqrt{-\varphi(\sigma)}}+\int_{\sigma_{2}}^{\sigma_{3}} \frac{d \sigma}{\sqrt{\varphi(\sigma)}}+\cdots
$$

It is always legitimate to choose the positive value of the radical in each case, so the value of $t-t_{0}$ will then have the form:

$$
t-t_{0}=\left(\alpha_{1}+\alpha_{3}+\alpha_{5}+\ldots\right)+i\left(\alpha_{2}+\alpha_{4}+\ldots\right),
$$

in which all of the $\alpha$ are positive. It will then follow from this that $\left|t-t_{0}\right|$ will increase indefinitely when $M$ tends to $N^{\prime}$. In other words, $t$ will tend to a limit $A+i A^{\prime}$, and if one replaces $t$ with $(A+$ $\left.i A^{\prime}\right)+t$ then one will see that when $t$ tends to zero (for imaginary values, but that is unimportant), the system will tend to a regular equilibrium position, while the vis viva tends to zero, which is impossible.

One will then obtain all of the singular branches of the trajectories that pass through an equilibrium point $N^{\prime}$ (or $\alpha_{1}, \alpha_{2}, \ldots, \alpha_{k}$ ) by searching for all of the integrals $q_{i}(\theta), q_{i}^{\prime}(\theta)$ of the system:

$$
\frac{d q_{i}}{d \theta}=\frac{-q_{i}^{\prime}}{\theta^{2}}, \quad \frac{d q_{i}^{\prime}}{d \theta}=\frac{-\beta_{i}+\Pi_{i}}{\theta^{2}},
$$

and for real values, the $q_{i}$ will tend to $\alpha_{i}$ and the $q_{i}^{\prime}$ will tend to zero when $\theta$ tends to zero according to a certain law.

It remains for us to discuss the case that we have left aside in which the trajectory $M_{0} M^{\prime}$ is remarkable. First observe that any point $M_{1}$ of a remarkable trajectory $M_{0} M_{1}$ or $(\gamma)$ is a regular point on that trajectory. Indeed, we can place the point $M$ at a point of $M_{0} M_{1}$ that is as close to $M_{1}$ as we desire and has a velocity tangent to $(\gamma)$ that is as large as we desire. The argument on page 246 then shows that $M$ will go past $M_{1}$ for a regular motion. The trajectory $(\gamma)$ will then continue regularly beyond $M_{1}$.

Having said that, let $N^{\prime}$ be an equilibrium point that is situated on $(\gamma)$ and consider the equality:

$$
T=F(\sigma)+h .
$$

$F\left(\sigma^{\prime}\right)$ will have a finite value for $\sigma=\sigma^{\prime}$, say $F\left(\sigma^{\prime}\right)=0$. If $h$ is zero then the point $M$ will tend to $N^{\prime}$ when $t$ increases indefinitely (under either the true motion or the conjugate motion). For the other values of $h, N^{\prime}$ will be an ordinary point of the motion, which can then present two points of regression that include the point $N^{\prime}$. If $\sigma^{\prime}$ is only a double zero of $F(\sigma)$ then one will have:

$$
T=h+h+\left(\sigma-\sigma^{\prime}\right)^{2}[A+\varepsilon]
$$

in the neighborhood of $\sigma^{\prime}$. When the number $A$ is negative, the motions that correspond to small positive values of $h$ will be periodic around $N^{\prime}$. If $A$ is positive then the same remark will apply to the conjugate motion.

Furthermore, the trajectory $(\gamma)$ can include an infinite number of equilibrium points $N^{\prime}$ that form a sequence. That can be seen in the example of the two equations:

$$
2 x^{\prime \prime}=x^{3}\left[5 x \sin \frac{1}{x}-\cos \frac{1}{x}\right], \quad y^{\prime \prime}=0
$$

which admit the remarkable trajectories $y=y_{0}$, along which, the motion will be defined by the relation:

$$
\left(\frac{d x}{d t}\right)^{2}=x^{5} \sin \frac{1}{x}+h
$$

All of the roots $x_{i}$ of the equality $\tan 1 / x=x / 5$ correspond to equilibrium points $N^{\prime}$, and the roots $x_{i}$ have $x=0$ for a limit.

In that example, the force $X$ is continuous, along with its first derivatives $\partial X / \partial x, \partial X / \partial y \equiv 0$, in the neighborhood of $x=0$. However, it is appropriate to remark that the singularity in question will not present itself when the coefficients $A_{i j}, Q_{i}$ are holomorphic functions of $q_{1}, q_{2}, \ldots, q_{k}$ in the domain of $a_{1}, a_{2}, \ldots, a_{k}$ or $N^{\prime}$. Indeed, the point $N^{\prime}$ will then be a regular analytic point of $(\gamma)$, and the variables $q_{2}, \ldots, q_{k}$ will be holomorphic functions of $\sigma$ when $\sigma$ is in the neighborhood of $\sigma^{\prime}$, and the point $F^{\prime}(\sigma)=Q_{1} \frac{d q_{1}}{d \sigma}+Q_{2} \frac{d q_{2}}{d \sigma}+\cdots+Q_{k} \frac{d q_{k}}{d \sigma}$ will then be holomorphic in the neighborhood of $\sigma^{\prime}$, and the point will necessarily be an isolated zero of $F^{\prime}(\sigma)$.

Finally, observe that the $\beta_{i}$ can be zero all along ( $\gamma$ ). The vis viva of motion along ( $\gamma$ ) will then be constant. In order for that to be true, it is necessary and sufficient that a geodesic of $T$ should be a locus of equilibrium points $N^{\prime}$.

Remark concerning the case in which the forces $Q_{i}$ are derived from a potential $U\left(q_{1}, q_{2}\right.$, $\ldots, q_{k}$ ). - If there exists a force function $U$, and if the equilibrium points $N^{\prime}$ are isolated, moreover, then the trajectories $(C)$ that pass through an equilibrium point will be necessarily exceptional.

First of all, the trajectories $(C)$ for which $N^{\prime}$ is an ordinary point $\left(T \neq 0\right.$ at $\left.N^{\prime}\right)$ depend upon only $k$ parameters. As for the ones for which $N$ is a singular point ( $T=0$ at $N^{\prime}$ ), they satisfy the condition: $U\left(a_{1}, a_{2}, \ldots, a_{k}\right)+h=\sigma$. Those trajectories then correspond to particular values of the vis viva constant, and as a result, they cannot depend upon more than $2 k-2$ parameters.

Conclusion: All of the preceding discussion can then be summarized by:
I. - When the forces $Q_{i}\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ are not all zero, any arbitrary trajectory $(C)$ can be traversed in only two distinct manners $\left({ }^{13}\right)$, and the second motion is deduced from the first one by changing the sign of $t$.

There can nonetheless be exceptions for certain trajectories $(\gamma)$ that are called remarkable, which correspond to an infinitude of motions (which are pairwise in the opposite sense) that depend upon an arbitrary constant. Such trajectories (which do not exist, in general) are necessarily geodesics of $T$ and depend upon more than $(k-1)$ parameters.
II. - Every continuous arc of a real trajectory ( $C$ ) that does not pass through either a singular point $N$ or an equilibrium point $N^{\prime}$ of $E_{k}$ is regular (i.e., it admits a continuous tangent and curvature at each point).

The only points that can be singular points or extremities of a trajectory are the points $N$ and $N^{\prime}$ then. Any continuous arc ( $C$ ) of a real trajectory that does not pass through either a point $N$ or

[^11]a point $N^{\prime}$ will then be traversed in the same sense in its entirety, whether the motion is a true real motion or its conjugate. In the former case, we say that $(C)$ is a true trajectory $\left(C^{\prime}\right)$, and in the latter case that it is a conjugate trajectory $\left(C^{\prime \prime}\right)$.

There nonetheless exists an exceptional class of trajectories $(\Gamma)$ that we call mixed trajectories, and on them there will be finite segments that include no points like $N$ or $N^{\prime}$, so they will be composed, in part, of arcs of type $\left(C^{\prime}\right)$ and in part of arcs of type $\left(C^{\prime \prime}\right)$. Any point $M$ that separates two adjacent arcs $\left(C^{\prime}\right)$ and $\left(C^{\prime \prime}\right)$ along $(\Gamma)$ is a point of regression, where the real motion changes sense. The remarkable trajectories must be regarded as mixed trajectories. The congruence of mixed trajectories will then depend upon exactly $k$ parameters, except in the particular case where the congruence of remarkable trajectories attains its maximum number $(k-1)$ of parameters, in which case, the two congruences will coincide.
III. - If an arbitrary point $M$ in the space $E_{k}$ is neither a singular point $N$ nor an equilibrium point $N^{\prime}$ then an infinitude of regular trajectories will pass through it in a neighborhood of $M$ that will depend upon $k$ parameters, which are the values of $q_{1}^{\prime}, q_{2}^{\prime}, \ldots, q_{k}^{\prime}$ at $M$, and no other trajectories will pass through $M$. Among those trajectories, there are ones of type ( $C^{\prime}$ ) and ones of type $\left(C^{\prime \prime}\right)$ in the neighborhood of $M$. Nevertheless, there is one and only one of them that admits the point $M$ as a point of regression, namely, the mixed trajectory $\left(\Gamma_{1}\right)$, which corresponds to a zero vis viva at $M$. That trajectory $\left(\Gamma_{1}\right)$ belongs to the one-parameter congruence of mixed trajectories $(\Gamma)$ that pass through $M$. Nonetheless, in the case where the congruence of remarkable trajectories depends upon $(k-1)$ parameters, and as a result agrees with the congruence of trajectories $(\Gamma)$, the trajectory $\left(\Gamma_{1}\right)$ will be the only mixed trajectory that passes through $M$.

If the point $M$ in question is an equilibrium point $N^{\prime}$ then once more an infinitude of trajectories that depend upon $k$ parameters will pass through that point that are regular in the neighborhood $N^{\prime}$. The trajectory $\Gamma_{1}$ will then reduce to the point $N^{\prime}\left(q_{1} \equiv a_{1}, \ldots, q_{k} \equiv a_{k}\right)$. However, the most important fact is that other trajectories ( $C$ ) can exist besides those trajectories that also pass through $N^{\prime}$ and have arbitrary singularities at $N^{\prime}$, and are such that the segment $M N^{\prime}$ (which is sufficiently small) is never traversed by a real motion (whether true or conjugate) in a finite time.

All of the trajectories $(C)$ are obtained by searching for the integrals $q_{i}(\theta), q_{i}^{\prime}(\theta)$ of the system:

$$
\frac{d q_{i}}{d \theta}=-\frac{q_{i}^{\prime}}{\theta^{2}}, \quad \frac{d q_{i}^{\prime}}{d \theta}=-\frac{\beta_{i}+\Pi_{i}}{\theta^{2}},
$$

such that the $q_{i}$ tend to $a_{1}, \ldots, a_{k}$, and the $q_{i}^{\prime}$ tend to zero for real values when one makes $\theta$ tend to zero according to a certain law.

When a trajectory $(C)$ that is not a remarkable trajectory passes through the point $N^{\prime}$ without coinciding with one of the singular branches $(C)$, the point $N^{\prime}$ will be an ordinary point of $(C)$, and it will be crossed by a regular motion (whether true or conjugate). If the trajectory ( $C$ )
coincides with one those singular branches $N^{\prime} M$ then either the system will tend to $N^{\prime}$ along $M N^{\prime}$ (with a vis viva that tends to zero) when $t$ increases indefinitely or $M N^{\prime}$ will decompose into an infinitude of arcs (which will tend to $N^{\prime}$ ) that correspond to just as many periodic motions (true or conjugate). When a remarkable trajectory ( $\gamma$ ) passes through $N^{\prime}, N^{\prime}$ will always be an ordinary point of $(\gamma)$, and there will always exist an infinitude of periodic motions on ( $\gamma$ ) for which $M$ oscillates around $N^{\prime}$ (which are true or conjugate motions) and motions for which $M$ tends to either $N^{\prime}$ or some other equilibrium point $N^{\prime \prime}$ that is as close to $N^{\prime}$ as one desires as $t$ increases indefinitely.

Finally, when the forces are derived from a potential and there exist only isolated equilibrium positions, a trajectory that is taken at random will not include any singular arcs $(C)$. It will then follow from this that (except for certain exceptional trajectories) any continuous arc of the real trajectory that includes no singular points $N$ of $E_{k}$ will be traversed completely in the same sense under either a true motion or its conjugate.

Remark. - In all of the foregoing discussion, we have overlooked the case in which the trajectory passes through a singular point $N$ of $E_{k}$. When the system $S$ tends to a singular position $N$ (as $t$ tends to $t_{1}$ ), it might happen that the $q^{\prime}$ do not tend to any limit. In general, it can also happen that even when the $q^{\prime}$ do have a limit, knowing the velocities at the position $N$ will be insufficient for one to determine the ultimate motion. There would then be no reason to pursue the analytical study of motion any further.

By definition, when $t$ starts from $t_{0}$ and increases, it can happen that the system $S$ goes to infinity when $t$ tends to $t_{1}$, or it might not tend to any limit point, or finally, it might tend to a singular position $N$. However, as long as that is not the case, we have seen that the velocities will remain well-defined for each value of $t$, and the preceding discussion will be valid.

Before applying the preceding considerations to some examples, we shall make a few more observations on the subject of similitude in mechanics.

## On similitude in mechanics. Conjugate motions.

When one changes $T$ into $C T$ in a system:

$$
\begin{equation*}
\frac{d}{d t}\left(\frac{\partial T}{\partial q_{i}^{\prime}}\right)-\frac{\partial T}{\partial q_{i}}=Q_{i}\left(q_{1}, q_{2}, \ldots, q_{k}\right), \quad \frac{d q_{i}^{\prime}}{d t}=q_{i}^{\prime} \quad(i=1,2, \ldots, k) \tag{1}
\end{equation*}
$$

and changes the $Q_{i}$ into $c Q_{i}$, in which $C$ and $c$ denote two constants, if one wishes to pass from the first motion to the second then it will suffice to change $t$ into $\sqrt{\frac{c}{C}} t_{1}$.

Indeed, write the new equations:

$$
\begin{equation*}
\frac{d}{d t_{1}} \sum_{j} A_{i j} \frac{d q_{j}}{d t_{1}}-\frac{1}{2} \sum_{j, l} \frac{\partial A_{j l}}{\partial q_{i}} \frac{d q_{j}}{d t_{1}} \frac{d q_{l}}{d t_{1}}=\frac{c}{C} Q_{i} . \tag{1'}
\end{equation*}
$$

If one sets $t_{1}=\sqrt{\frac{C}{c}} t$ then one will have $\frac{d q_{i}}{d t_{1}}=\sqrt{\frac{c}{C}} \frac{d q_{i}}{d t}, d t_{1}=\sqrt{\frac{C}{c}} d t$, and equations (1') will be transformed into the system (1).

It follows from this that equations (1) and (1') will define the same trajectories. However, the motion on those trajectories will be the same only if $C / c$ is equal to unity. If $c / C$ is positive then the true real trajectories will be the same for the two systems. If $c / C$ is negative then the true trajectories of the same system will be the conjugate trajectories of the second one, and conversely. In particular, if $C=1$ and $c=-1$ then one can pass from the first motion to the second one by changing $t$ into $i t$. One will recover the result that was pointed out before, namely, that the motion that is conjugate to the true motion will be the motion of the system when one changes the sign of $Q_{i}$, i.e., when one changes the senses of all given forces without changing their direction or magnitudes.

From that, consider a system $S$ of material points $x_{i}, y_{i}, z_{i}$ that are subject to certain constraints and a homothetic system $\Sigma$ or $\xi_{i}=\lambda x_{i}, \eta_{i}=\lambda y_{i}, \zeta_{i}=\lambda z_{i}$ that is subject to corresponding constraints. The vis viva of the new system will be $\lambda^{2} T$, if $T$ is that of the first one. On the other hand, subject the new system to some given forces $\left(\Phi_{i}\right)$ or $\Xi_{i}, \mathrm{H}_{i}, \mathrm{Z}_{i}$ that are homothetic to the forces $\left(F_{i}\right)$ or $X_{i}, Y_{i}, Z_{i}$ that are exerted on the first one: $\Xi_{i}=\mu X_{i}, \mathrm{H}_{i}=\mu Y_{i}, Z_{i}=\mu Z_{i}$. If $Q_{j}^{\prime}$ denotes the coefficient $Q_{j}$ relative to the second system then one will have:

$$
Q_{j}^{\prime}=\sum \Xi \frac{\partial \xi}{\partial q_{j}}+\mathrm{H} \frac{\partial \eta}{\partial q_{j}}+\mathrm{Z} \frac{\partial \zeta}{\partial q_{j}}=\lambda \mu \sum X \frac{\partial x}{\partial q_{j}}+Y \frac{\partial y}{\partial q_{j}}+Z \frac{\partial z}{\partial q_{j}}=\lambda \mu Q_{j}
$$

It follows from this that the relations between the $q_{i}$ will be the same for the two motions and that one passes from the first motion to the second one by changing $t$ into $\sqrt{\frac{\mu}{\lambda}} t$. The trajectories of second system $\Sigma$ will then be homothetic to the trajectories of $S$ and will be deduced from the formulas $\xi_{i}=\lambda x_{i}, \eta_{i}=\lambda y_{i}, \zeta_{i}=\lambda z_{i}$. The motion of $\Sigma$ is deduced from the motion of $S$ with the aid of the preceding formula and changing $t$ into $\sqrt{\frac{\mu}{\lambda}} t$, moreover. If the homothetic correspondence between $S$ and $\Sigma$ and the correspondence between the $(F)$ and the $(\Phi)$ have the same (viz., $\lambda \mu>0$ ) then the true trajectories of $\Sigma$ will be the transforms of the true trajectories of $S$. Otherwise $(\lambda \mu<0)$, they will be the transforms of the conjugate trajectories of $S$. If $\mu=\lambda$, i.e., if one transforms the points $S$ and the forces $F$ together, then the new motion can be deduced from the first one by the formulas $\xi_{i}=\lambda x_{i}, \eta_{i}=\lambda y_{i}, \zeta_{i}=\lambda z_{i}$. If $\mu=-\lambda$ then the new motion can be deduced from the motion that is conjugate to the first one by using the same formulas.

One sees that, by definition, if two similar systems are subject to similar forces then the trajectories will also be similar. That principle of similitude in mechanics was introduced for the first time by Bertrand.

We have said that, on the one hand, the equations of motion will not change when one changes $t$ into $-t$. In other words, when a system $S$ without friction that has constraints that are independent of time is subject to forces that depend upon neither the velocities nor time, the motions of that system will be reversible.

I shall add that Appell has inferred an interpretation of imaginary time in mechanics from a consideration of conjugate motions that applies to many interesting special cases. We shall confine ourselves to citing the case of the simple pendulum as one example.


Let O be a circle in the vertical plane $x O y$, and let $O y$ be the direction of gravity. The position $M$ of a point that moves without friction on the circle is determined by the angle $\theta=M O y$. If one releases the point $M$ with zero velocity from $M_{0}(\theta=\alpha)$ then the motion will be defined by the equality:

$$
\sqrt{\frac{g}{t}} d t=\frac{d \theta}{2 \sqrt{\sin ^{2} \frac{\alpha}{2}-\sin ^{2} \frac{\theta}{2}}}
$$

or rather, upon setting $\sin \theta / 2=u \sin \alpha / 2$ :
( $\alpha$ )

$$
\sqrt{\frac{g}{l}} t=\int_{1}^{u} \frac{d u}{-\sqrt{\left(1-u^{2}\right)\left(1-k^{2} u^{2}\right)}} \quad\left(k^{2}=\sin ^{2} \frac{\alpha}{2}\right)
$$

i.e., $u=\operatorname{sn}\left(t \sqrt{\frac{g}{l}}+\right.$ const. $)$. The function sn admits periods that are expressed by $4 k$ and $2 i k^{\prime}$ with:

$$
k=\int_{0}^{1} \frac{d u}{+\sqrt{\left(1-u^{2}\right)\left(1-k^{2} u^{2}\right)}}, \quad k^{\prime}=\int_{1}^{1 / k} \frac{d u}{+\sqrt{\left(1-u^{2}\right)\left(1-k^{2} u^{2}\right)}} .
$$

The time that it takes for $M$ to go from $M_{0}$ to $A$ is equal to $\sqrt{\frac{l}{g}} k$. If one now changes the sense of gravity (without changing its direction or magnitude) then the equation of the new motion will be obtained by changing $t$ into $i t$, and one will have:

$$
\sqrt{\frac{g}{l}} t=\int_{1}^{u} \frac{d u}{+\sqrt{\left(1-u^{2}\right)\left(k^{2} u^{2}-1\right)}}
$$

$\theta$ will then increase from $\alpha$ to $\pi$, and the time that it takes for $M$ to go from $M_{0}$ (or $u=1$ ) to the point $A^{\prime}$ (or $u=\frac{1}{\sin \alpha / 2}=\frac{1}{k}$ ) will be equal to $\sqrt{\frac{l}{g}} k^{\prime}$.

One can further say that if $4 k$ and $2 i k^{\prime}$ are periods of the function sn that correspond to the modulus $k^{2}=\sin ^{2} \alpha / 2$ then the time that it takes for the moving point that is released without velocity at $M_{0}(\theta=\alpha)$ or at $M_{0}^{\prime}(\theta=\pi-\alpha)$ to arrive at $A$ will be equal to $\sqrt{\frac{l}{g}} k$ under the first hypothesis and $\sqrt{\frac{l}{g}} k^{\prime}$ under the second one.

Remark. - We just saw that if we replace the forces $Q_{i}$ with the forces $Q_{i}^{\prime}=c Q_{i}(c$ being a constant) then the trajectories will not be modified. It is appropriate to point out that those forces $c Q_{i}$ are the only forces $Q_{i}^{\prime}\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ that will generate the same trajectories when they are substituted for the forces $Q_{i}$ in the system (1).

Indeed, write the differential equations of the trajectories (see page 239) as:

$$
\begin{equation*}
\frac{q_{(i)}^{\prime \prime}+q_{(i)}^{\prime} \Pi_{1}-\Pi_{i}}{\beta_{i}-q_{(i)}^{\prime} \beta_{1}}=\frac{q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \Pi_{1}-\Pi_{i}}{\beta_{2}-q_{(2)}^{\prime} \beta_{1}} \quad(i=3,4, \ldots, k) \tag{10}
\end{equation*}
$$

and

$$
\begin{equation*}
\frac{\frac{d}{d q_{1}}\left[q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \Pi_{1}-\Pi_{2}\right]}{q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \Pi_{1}-\Pi_{2}}+2 \Pi_{1}=\frac{\frac{d}{d q_{1}}\left[\beta_{2}-q_{(2)}^{\prime} \beta_{1}\right]-2 \beta_{1}\left[q_{(2)}^{\prime \prime}+q_{(2)}^{\prime} \Pi_{1}-\Pi_{2}\right]}{\beta_{2}-q_{(2)}^{\prime} \beta_{1}} . \tag{11}
\end{equation*}
$$

The differential system will have the form:

$$
\frac{d^{2} q_{i}}{d q_{1}^{2}}=\frac{d^{2} q_{2}}{d q_{i}^{2}} \frac{\beta_{i}-\beta_{1} \frac{d q_{i}}{d q_{1}}}{\beta_{2}-\beta_{1} \frac{d q_{2}}{d q_{1}}}+L_{i} \quad(i=3,4, \ldots, k)
$$

and

$$
\begin{equation*}
\frac{d^{3} q_{i}}{d q_{1}^{3}}=\chi_{2} \frac{d}{d q_{1}} \log \beta_{1}+L_{2}^{\prime} \tag{11'}
\end{equation*}
$$

in which $L_{i}$ includes only the first derivatives, and $L_{2}^{\prime}$ is defined with the aid of the coefficients of $T$ and the ratios $\beta_{i} / \beta_{1}$.

Having said that, assume that one replaces the forces $Q_{i}$ with some other forces $Q_{i}^{\prime}\left(q_{1}, q_{1}, \ldots\right.$, $\left.q_{k}\right)$. The coefficients $\beta_{i}\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ will become $\beta_{i}^{\prime}\left(q_{1}, q_{2}, \ldots, q_{k}\right)$. In order for equations (10) and (11) remain unaltered by that substitution, it is first necessary that one must have:

$$
\frac{\beta_{i}-\beta_{1} \frac{d q_{i}}{d q_{1}}}{\beta_{2}-\beta_{1} \frac{d q_{2}}{d q_{1}}} \equiv \frac{\beta_{i}^{\prime}-\beta_{1}^{\prime} \frac{d q_{i}}{d q_{1}}}{\beta_{2}^{\prime}-\beta_{1}^{\prime} \frac{d q_{2}}{d q_{1}}}
$$

identically, and as a result:

$$
\begin{equation*}
\frac{\beta_{1}^{\prime}}{\beta_{1}} \equiv \frac{\beta_{2}^{\prime}}{\beta_{2}} \equiv \ldots \equiv \frac{\beta_{k}^{\prime}}{\beta_{k}} . \tag{12}
\end{equation*}
$$

It will then follow [from $\left(11^{\prime}\right)$ ] that:

$$
\frac{d}{d q_{1}} \log \beta_{1} \equiv \frac{d}{d q_{1}} \log \beta_{1}^{\prime}
$$

or rather that:

$$
\frac{\partial}{\partial q_{1}} \log \beta_{1} \equiv \frac{\partial}{\partial q_{1}} \log \beta_{1}^{\prime}, \quad \ldots, \quad \frac{\partial}{\partial q_{k}} \log \beta_{1} \equiv \frac{\partial}{\partial q_{k}} \log \beta_{1}^{\prime},
$$

which will imply the consequence:

$$
\beta_{1}^{\prime} \equiv c \beta_{1}
$$

so, from (12):

$$
\beta_{1}^{\prime}=c \beta_{1}, \quad \beta_{2}^{\prime}=c \beta_{2}, \quad \ldots, \quad \beta_{k}^{\prime}=c \beta_{k},
$$

and one immediately deduces the equalities:

$$
Q_{1}^{\prime}=c Q_{1}, \quad Q_{2}^{\prime}=c Q_{2}, \quad \ldots, \quad Q_{k}^{\prime}=c Q_{k}
$$

Q.E.D.

More generally, if one replaces $T$ with $C T$ and the $Q_{i}$ with $Q_{i}^{\prime}\left(q_{1}, q_{1}, \ldots, q_{k}\right)$, where $C$ denotes a constant, then in order for the trajectories to the remain the same, is necessary and sufficient for one to have $Q_{i}^{\prime} \equiv C Q_{i}$.

## Applications of the preceding generalities to some examples.

Consider a free material point $(x, y, z)$ that is subjected to a force $F$ whose projections $X, Y, Z$ are analytic functions of $x, y, z$ that are holomorphic for all real values of those variables.

The geodesics here are the lines in space: The remarkable trajectories (if they exist) will then be lines $D$, and if they form a congruence then it will depend upon at most two parameters. In what
case will those trajectory lines depend upon precisely two parameters? In order to account for that case, it is sufficient to observe that at each point $(x, y, z)$ of the line $D$, the line of action of $(F)$ must coincide with $D$. If one passes a line $(D)$ through any point $(x, y, z)$ then the set of all lines of action of the force $(F)$ must coincide with the congruence of lines $D$. Conversely, if the lines of action of $(F)$, which generally form a complex, form a congruence then every line of that congruence will be a remarkable trajectory. The only case in which the remarkable trajectories depend upon two parameters is then the case in which the complex of forces $(F)$ reduced to a congruence. In particular, when there exists a force function $U(x, y, z)$, the forces $(F)$ will form a complex, unless the level surfaces $U=$ const. are parallel. The normals to those surfaces will then form a congruence.

As for the singular trajectories, if $(a, b, c)$ is an equilibrium point then they will be obtained in any case by searching for all of the integrals of the system:

$$
\frac{d x}{x^{\prime}}=\frac{d y}{y^{\prime}}=\frac{d z}{z^{\prime}}=\frac{d x^{\prime}}{X}=\frac{d y^{\prime}}{Y}=\frac{d z^{\prime}}{Z}
$$

that satisfy the initial conditions: $x=a, y=b, z=c, x^{\prime}=y^{\prime}=z^{\prime}=0$.
Any continuous arc of the trajectory is necessarily regular, except perhaps at an equilibrium point ( $a, b, c$ ), and can be extended indefinitely in a regular fashion ( ${ }^{14}$ ), as long as one does not encounter an equilibrium point. However, the singular branches that pass through an equilibrium point $(a, b, c)$ can present arbitrary singularities at that point, and in particular, they can terminate there.

Now assume that the forces are derived from a potential $U$ and that the equilibrium points are isolated, moreover, i.e., that the three derivatives of $U$ are annulled simultaneously only at isolated points. A trajectory that is taken at random will not include any singular branch. Any continuous arc of such a thing will be traversed in its entirety in the same sense by a regular motion, whether true or conjugate. There will be an exception only for some special trajectories, namely, the mixed trajectories and the singular trajectories. The former depend upon three parameters, unless the level surfaces are parallel, in which case the mixed trajectories will coincide with the congruence of remarkable trajectories. As far as singular trajectories such as $M N^{\prime}$ are concerned ( $N^{\prime}$ being an equilibrium point), we know only that the arc $M N^{\prime}$ is never traversed by the moving point in a finite time by a real motion (whether true or conjugate): The moving point tends to $N^{\prime}$ when $t$ increases indefinitely, or rather, $M N^{\prime}$ decomposes into an infinitude of arcs that correspond to periodic motions.

When the level surfaces are parallel, any equilibrium point $N^{\prime}$ will be a multiple point of the level surface that contains it, and an infinitude of normals will pass through that point that form a cone, which are just as many remarkable trajectories $D$. One can define initial conditions for the moving point such that it will tend to $N^{\prime}$ on $D$ when $t$ increases indefinitely (for a true or conjugate motion, as the case may be).

[^12]I shall add some observations in regard to the periodic motions. For an arbitrary true motion to be periodic, it is necessary and sufficient that an arbitrary true trajectory should be a closed curve. Indeed, take a random trajectory: At any point $M$ of that trajectory (that is not a double point), the velocity of the moving point will have a well-defined value. On the other hand, the trajectory will be traversed in its entirety in the same sense by a regular motion in a finite time interval $t_{1}$, and after the time $t_{1}$, the moving point will return to the starting point $M$ with the same velocity and point in the same sense, and the same motion will begin again.

Similarly, in order for any true motion whose initial conditions are subject to a certain inequality to be periodic, it is necessary and sufficient that all of the true trajectories (whose $2 k-$ 1 parameters satisfy a certain inequality) should be closed curves.

If one desires that the particular motions should be periodic then there are several cases to be distinguished.

In order for a true trajectory that is neither mixed nor remarkable to correspond to a periodic motion, it is necessary and sufficient that it should be closed and contain no singular branch.

In order for a mixed trajectory to correspond to a periodic motion, it is necessary and sufficient that it should present at least two points of regression that do not include any singular branch, and between which the true motion will be real (and not the conjugate motion).

In order for a remarkable trajectory (which is then a line $D$ ) to correspond to a periodic motion, it is necessary and sufficient that the line $D$ should pass through an equilibrium point $N^{\prime}$, where the value of the function $U$ along the line possesses a maximum. In that case, there exists an infinitude of periodic motions in which the moving point will oscillate about $N^{\prime}$ along $D$ with a constant amplitude of oscillation that is as small as one desires.

## Examples:

1. As one particular application, let us study the motion of a gravitating point.

The forces form a congruence here, namely, the congruence of vertical lines. Those lines will be both remarkable trajectories and mixed trajectories for the point. There is no equilibrium position, so there will be no singular branches.

If one defines the direction and sense of gravity then the trajectories will be parabolas:

$$
\begin{aligned}
& y=\lambda x+\mu, \\
& z=\alpha x^{2}+\beta x+\gamma,
\end{aligned}
$$

in which $\alpha, \beta, \gamma, \lambda, \mu$ denote arbitrary constants. Those trajectories are true for $\alpha>0$ and conjugate for $\alpha<0$. The two classes will permute when one changes the sense of gravity. Any arc of the parabola will be traversed in the same sense in a finite time under either the true motion or the conjugate motion.

As for the motion along a vertical, it will necessarily present one and only one point of regression (the culmination point) when $t$ varies from $-\infty$ to $+\infty$.

There are no periodic motions.
2. Now consider a material point (of mass 1) that is attracted to the origin in proportion to the distance: $F=-k^{2} r$.

The forces once more form a congruence, namely, the congruence of lines that issue from the origin. Those lines $D$ are both remarkable and mixed trajectories of the moving point. The origin $O$ is an equilibrium point, and it is the only one.

The real trajectories are the real conics that have the origin for their centers. The true trajectories are the ellipses, and the conjugate trajectories are hyperbolas.

All of the true motions are periodic. Along a line $D$, the function $U$ has a maximum at the origin $O$. All of the motions along $D$ will present two points of regression that are equidistant from $O$. As one knows, the oscillations are tautochronous.

If one considers a point that is repelled by the origin in proportion to the distance then the real trajectories will not change, the true trajectories will be hyperbolas, and the conjugate trajectories will be ellipses. Along a line $D$, the motion will present one and only one point of regression between $t=-\infty$ and $t=+\infty$ or none of them, according to whether the absolute value $V_{0}$ of the initial velocity is less than or greater than $+k r_{0}$, respectively, where $r_{0}$ is the initial distance from the origin. If $V_{0}=k r_{0}$ then the moving point will tend to $O$ when $t$ tends to $+\infty$ (or to $-\infty$ ). No motion is periodic.

In the latter case, $x, y, z$ are rational functions of $e^{k t}$, and as a result, they will admit the imaginary period $2 i \pi / k$. That period corresponds to the real period $2 \pi / k$ of the conjugate motion, and indeed when it is attractive, $x, y, z$ will be rational functions of $\tan k t / 2$
3. Finally, let us study the motion of a material point of mass 1 that is subject to the force $X$ $=2 y, Y=-2 x, Z=0$.

The forces form a complex here. The $z$-axis is a locus of equilibrium points, and since it is also a geodesic, it will be a trajectory that the moving point can traverse with a constant and arbitrary velocity.

Do there exist other remarkable trajectories? The trajectories must be lines and verify the equations:

$$
\frac{d x}{y}=\frac{d y}{-x}=\frac{d z}{0} .
$$

The line $x=0, y=0$ is the only line that satisfies those conditions.
If one observes that the equations of motion can be written:

$$
\frac{d^{2}(x+i y)}{d t^{2}}+2 i(x+i y)=0, \quad z^{\prime \prime}=0
$$

then one will see that the motion is defined by the equalities:

$$
x+i y=(\alpha+i \beta) e^{(1-i) t}+(\gamma+i \delta) e^{(-1+i) t}, \quad z=\lambda t+\mu,
$$

or rather:
(C)

$$
\left\{\begin{array}{l}
x=e^{t}(\alpha \cos t+\beta \sin t)+e^{-t}(\gamma \cos t-\delta \sin t), \\
y=e^{t}(\beta \cos t-\alpha \sin t)+e^{-t}(\delta \cos t+\gamma \sin t), \\
z=\lambda t+\mu,
\end{array}\right.
$$

in which $\alpha, \beta, \gamma, \delta, \lambda, \mu$ are arbitrary constants. Since $x$ is always annulled a certain number of times (along with $y$ ) when $t$ varies from $-\infty$ to $+\infty$, one can suppose that $x$ is equal to zero for $t=$ 0 , which amounts to setting $\alpha=-\gamma$ in the expression for $x, y$.

The true real trajectories are obtained by giving real values to $\alpha, \beta, \gamma, \delta, \lambda$, and $t$, while the conjugate real trajectories are obtained by changing $t$ into $i t^{\prime}$ in the expression for $x+i y . x, y, z$ will then be expressed in the following form:

$$
\left\{\begin{array}{l}
x=e^{t^{\prime}}\left(\alpha^{\prime} \cos t^{\prime}-\beta^{\prime} \sin t^{\prime}\right)+e^{-t^{\prime}}\left(\gamma^{\prime} \cos t^{\prime}+\delta^{\prime} \sin t^{\prime}\right), \\
y=e^{t^{\prime}}\left(\beta^{\prime} \cos t^{\prime}+\alpha^{\prime} \sin t^{\prime}\right)+e^{-t^{\prime}}\left(\delta^{\prime} \cos t^{\prime}-\gamma^{\prime} \sin t^{\prime}\right), \\
z=\lambda^{\prime} t^{\prime}+\mu^{\prime},
\end{array}\right.
$$

and one will see that if one sets $\alpha^{\prime}=\alpha, \beta^{\prime}=-\beta, \gamma^{\prime}=\gamma, \delta^{\prime}=-\delta$, and if one changes $y$ into $-y$ then one will recover equations $(C)$. The conjugate trajectories are then symmetric to the true trajectories with respect to the $z x$-plane (as well as with respect to the $z y$-plane). Furthermore, the symmetric curve of a true trajectory with respect to an arbitrary point of $O z$ is again a true trajectory. The set of true trajectories also admits symmetry planes that are easy to see.

Among the real trajectories (which depend upon five constants), the mixed trajectories form a congruence that depends upon three constants. In order for a trajectory to mixed, it is necessary and sufficient that it should present a point (that is not an equilibrium point) where $x^{\prime}, y^{\prime}, z^{\prime}$ are simultaneously annulled. It will first be necessary then that $\lambda$ should be zero, so as a result, that $x^{\prime}+i y^{\prime}$ will be annulled, for a real value of $t$ that one can always suppose to be zero. In order for $x^{\prime}+i y^{\prime}$ to be annulled with $t$, it is necessary and sufficient that $\alpha+i \beta=\gamma+i \delta$, or that $\alpha=\gamma, \beta=$ $\delta$. The mixed trajectories will then be given by the equalities:

$$
\left\{\begin{array}{l}
x=\alpha \cos t\left(e^{t}+e^{-t}\right)+\beta \sin t\left(e^{t}-e^{-t}\right),  \tag{C}\\
y=\alpha \sin t\left(e^{-t}-e^{t}\right)+\beta \cos t\left(e^{t}+e^{-t}\right), \\
z=\mu .
\end{array}\right.
$$

They present one and only one point of regression, $x=\alpha, y=\beta, z=\mu . x, y, z$ will take those same values again for values of $t$ that are equal and of opposite sign. However, if one sets $t=i t^{\prime}$ then one will get the other real portion of the mixed trajectory:

$$
\left\{\begin{array}{l}
x=\alpha \cos t^{\prime}\left(e^{t^{\prime}}+e^{-t^{\prime}}\right)-\beta \sin t^{\prime}\left(e^{t^{\prime \prime}}-e^{-t^{\prime}}\right), \\
y=\alpha \sin t^{\prime}\left(e^{-t^{\prime}}-e^{t^{\prime}}\right)+\beta \cos t^{\prime}\left(e^{t^{\prime}}+e^{-t^{\prime}}\right), \\
z=\mu .
\end{array}\right.
$$

If one sets $\theta=t^{2}=-t^{\prime 2}$ then the functions $x(\theta), y(\theta)$, and $z=\mu$ that are defined by $(C)$ or $\left(C^{\prime}\right)$ will be the same analytic functions of $\theta$, and when $\theta$ varies from $-\infty$ to $+\infty$, the point $(x, y$, $z$ ) will traverse all of the real mixed trajectory. That trajectory decomposes into two parts, one of which is true and the other of which is conjugate, that are separated by the point $\theta=0$. The curve that is symmetric to a mixed trajectory with respect to the $z y$-plane (or the $z x$-plane) is also a mixed trajectory, but the symmetric part of a true segment will be a conjugate segment, and conversely.

Let us now consider the singular trajectories. Those trajectories cannot include an infinite number of points of regression, since any one of the trajectories will possess at most one point of regression. The true singular branches are then obtained by increasing (or decreasing) $t$ indefinitely through real values and determining whether the moving point will tend to a limiting position $N^{\prime}$. That point will necessarily be an equilibrium point, and therefore, a point on $O z$. In order for $x$ and $y$ to tend to zero as $t$ decreases indefinitely, it is necessary and sufficient that $\gamma$ and $\delta$ should be zero. In order for $|z|$ to not increase indefinitely, $\lambda$ must be zero. Since, on the other hand, one can suppose that $\alpha=-\gamma$, so one will have $\alpha=0$ here, one will have, by definition, the trajectories:

$$
x=\beta e^{t} \sin t, y=\beta e^{t} \cos t, z=\mu,
$$

which depend upon two arbitrary constants. One will get the same trajectories by making $t$ increase indefinitely (it will suffice to change $t$ into $-t$ ). The conjugate singular branches are obtained by changing $t$ into $i t^{\prime}$, which will give:

$$
x^{\prime}=-\beta^{\prime} e^{t^{\prime}} \sin t^{\prime}, \quad y^{\prime}=\beta^{\prime} e^{t^{\prime}} \cos t^{\prime}, \quad z^{\prime}=\mu .
$$

These are the symmetric images with respect to the $y z$-plane of the true singular branches. If one eliminates $t$ then one will find that $r=C e^{-\theta}, z=\mu$ for true branches and $r=C e^{\theta}, z=\mu$ for the conjugate ones, in which $r$ and $\theta$ denote the polar coordinates of a point in the $x y$-plane. Those curves admit the origin as an asymptotic point. When the moving point describes one of those singular trajectories, if it is launched, for example, in the sense of the point $N^{\prime}$ then it will tend to that point without ever reaching it when $t$ increases indefinitely $\left({ }^{15}\right)$.

Any true trajectory that does not belong to either the mixed trajectories (a three-parameter congruence) or the singular trajectories (a two-parameter congruence) will be traversed in its entirety in the same sense when $t$ increases from $-\infty$ to $+\infty$.

[^13]
## Equations of the trajectories when the forces are zero or derive from a potential.

When the forces are zero, we saw (pp. 238) that the trajectories will depend upon $(2 k-2)$ constants, and we have indicated the means of forming the differential equations of those trajectories. We shall now give an explicit form for those differential equations

Let:

$$
2 T \equiv \sum A_{i j} q_{i}^{\prime} q_{j}^{\prime} \quad\left(A_{i j}=A_{j i}\right)
$$

be the vis viva of the system, and consider the Lagrange equations:

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial T}{\partial q_{i}^{\prime}}-\frac{\partial T}{\partial q_{i}}=0 \quad(i=1,2, \ldots, k) \tag{1}
\end{equation*}
$$

which define the motion in the absence of forces.
As one knows, those equations will imply the consequence: $T=h$. If one lets $T_{1}$ denote what $T$ will become when one replaces $q_{1}^{\prime}$ with $1, q_{2}^{\prime}$ with $d q_{2} / d q_{1}=q_{(2)}^{\prime}, \ldots$, and $q_{k}^{\prime}$ with $d q_{k} / d q_{1}=$ $q_{(k)}^{\prime}$ in it then the first integral can be written:

$$
q_{1}^{\prime 2} T_{1}=h
$$

or rather:

$$
\begin{equation*}
d t=d q_{1} \sqrt{\frac{T_{1}}{h}} \tag{2}
\end{equation*}
$$

It is legitimate to replace equations (1) with the last $(k-1)$ of those equations:

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial T}{\partial q_{i}^{\prime}}-\frac{\partial T}{\partial q_{i}}=0 \quad(i=2,3, \ldots, k) \tag{1'}
\end{equation*}
$$

combined with equation (2). If one now replaces $d t$ with $d q_{1} \sqrt{\frac{T_{1}}{h}}$ everywhere in (1) then one will define the differential equations of the trajectories. Upon remarking that $\frac{\partial T}{\partial q_{i}^{\prime}}$ is linear with respect to $q_{i}^{\prime}$, and is equal to $q_{1}^{\prime} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}$ or $\frac{\sqrt{h}}{T_{1}} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}$ as a result, and that $\frac{\partial T}{\partial q_{i}}$ is likewise equal to $\frac{h}{T_{1}} \frac{\partial T_{1}}{\partial q_{i}}$, one will see that equations (1), thus-transformed, will become:

$$
\frac{\sqrt{h}}{T_{1}} \frac{d}{d q_{1}}\left(\frac{\sqrt{h}}{T_{1}} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}\right)-\frac{h}{T_{1}} \frac{\partial T_{1}}{\partial q_{i}}=0 \quad(i=2,3, \ldots, k)
$$

or rather:

$$
\begin{equation*}
\frac{d}{d q_{1}}\left(\frac{1}{T_{1}} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}\right)-\frac{1}{\sqrt{T_{1}}} \frac{\partial T_{1}}{\partial q_{i}}=0 \quad(i=2,3, \ldots, k) \tag{3}
\end{equation*}
$$

If we now set:

$$
f=\sqrt{T_{1}}=\sqrt{A_{11}+2 A_{12} q_{(2)}^{\prime}+\cdots+2 A_{1 k} q_{(k)}^{\prime}+\sum_{i>1, j>1} A_{i j} q_{(i)}^{\prime} q_{(j)}^{\prime}}
$$

then the system (3) will be written:

$$
\begin{equation*}
\frac{d}{d q_{1}} \frac{\partial f}{\partial q_{(i)}^{\prime}}-\frac{\partial f}{\partial q_{i}}=0 \quad(i=2,3, \ldots, k) \tag{4}
\end{equation*}
$$

For example, if $k=2$ then the geodesics of the surface whose $d s^{2}$ is $A_{11} d q_{1}^{2}+2 A_{12} d q_{1} d q_{2}+$ $A_{22} d q_{2}^{2}$ will be given by the second-order equation:

$$
\frac{d}{d q_{1}} \frac{\partial f}{\partial q_{(2)}^{\prime}}-\frac{\partial f}{\partial q_{2}}=0
$$

in which $f=\sqrt{A_{11}+2 A_{12} q_{(2)}^{\prime}+A_{22} q_{(2)}^{\prime 2}}$ and $q_{(2)}^{\prime}=d q_{2} / d q_{1}$.
Equations (4) are the ones that one finds by annulling the variation of the integral $\int f d q_{1}$, when it is taken along an arbitrary curve $q_{i}=\varphi_{i}\left(q_{1}\right)$ whose extremities are fixed.

One indeed verifies in that manner that the geodesics depend upon only $(2 k-2)$ constants, and one forms the $(k-1)$ second-order differential equations that define them explicitly. It would be easy to prove directly that those equations, which are linear with respect to the $q_{(i)}^{\prime \prime}$, can be solved for those variables, but that would result from what was said at the beginning of this chapter.

Principle of least action. - The procedure that we just employed applies just as well to the case in which the forces are not zero, but are derived from a potential $U$. Indeed, under that hypothesis, write down the Lagrange equations:

$$
\frac{d}{d t} \frac{\partial T}{\partial q_{i}^{\prime}}-\frac{\partial T}{\partial q_{i}}=\frac{\partial U}{\partial q_{i}} \quad(i=1,2, \ldots, k)
$$

Those equations imply the consequence that $T-U=h$, so it will be legitimate for us to replace them with $(k-1)$ of them (the last ones, for example) with:

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial T}{\partial q_{i}^{\prime}}-\frac{\partial T}{\partial q_{i}}=\frac{\partial U}{\partial q_{i}} \quad(i=2,3, \ldots, k) \tag{1'}
\end{equation*}
$$

combined with the equation:

$$
\begin{equation*}
d t=d q_{1} \sqrt{\frac{T_{1}}{U+h}}, \tag{2}
\end{equation*}
$$

in which the notation is the same as it was recently.
If one expresses $d t$ in terms of $d q_{1}$ as in (2) everywhere in equations ( $1^{\prime}$ ) then since one has:

$$
\frac{\partial T}{\partial q_{i}^{\prime}}=\sqrt{\frac{U+h}{T_{1}}} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}, \quad \frac{\partial T}{\partial q_{i}}=\frac{U+h}{T_{1}} \frac{\partial T_{1}}{\partial q_{i}}
$$

here, it will become:

$$
\begin{equation*}
\sqrt{\frac{U+h}{T_{1}}} \frac{d}{d q_{1}}\left(\sqrt{\frac{U+h}{T_{1}}} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}\right)-\frac{U+h}{T_{1}} \frac{\partial T_{1}}{\partial q_{i}}=\frac{\partial U}{\partial q_{i}} \quad(i=2,3, \ldots, k) \tag{3}
\end{equation*}
$$

and those equations are the differential equations of the trajectories that correspond to the value $h$ of the constant of the vis viva integral. Those trajectories define a ( $2 k-2$ )-parameter congruence.

Equations (3) can also be written:

$$
\begin{equation*}
\frac{d}{d q_{1}}\left(\sqrt{\frac{U+h}{T_{1}}} \frac{\partial T_{1}}{\partial q_{(i)}^{\prime}}\right)-\sqrt{\frac{U+h}{T_{1}}} \frac{\partial T_{1}}{\partial q_{i}}-\sqrt{\frac{T_{1}}{U+h}} \frac{\partial U}{\partial q_{i}}=0 \quad(i=2,3, \ldots, k) \tag{4}
\end{equation*}
$$

In that form, one sees that if one sets:

$$
f=\sqrt{(U+h) T_{1}}
$$

then they will coincide with the equations:

$$
\begin{equation*}
\frac{d}{d q_{1}} \frac{\partial f}{\partial q_{(i)}^{\prime}}-\frac{\partial f}{\partial q_{i}}=0, \quad q_{(i)}^{\prime}=\frac{d q_{i}}{d q_{1}} \tag{5}
\end{equation*}
$$

$$
(i=2,3, \ldots, k)
$$

Let us examine this last equation in more detail. Consider an arbitrary curve $C$ in $k$-dimensional space: $q_{2}=\varphi_{2}\left(q_{1}\right), \ldots, q_{k}=\varphi_{k}\left(q_{1}\right)$ that is constrained to have both of its extremities fixed at $a_{1}$, $a_{2}, \ldots, a_{k}$ and $b_{1}, b_{2}, \ldots, b_{k}$. If one expresses the idea that the integral $\int_{C} f d q_{1}$ has zero variation when one passes from one particular curve $C_{1}$ to a neighboring curve $C$ then one will find that the curve $C_{1}$ must satisfy the differential equations (5) precisely. Moreover, one proves that if the fixed extremities $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$, or $M_{1}$, and $\left(b_{1}, b_{2}, \ldots, b_{k}\right)$, or $M_{2}$, are sufficiently close then the integral that is taken along $C$ will be a minimum $\left({ }^{16}\right)$. One can then say that the integral $\int_{M_{1}}^{M_{2}} f d q_{1}$ is minimal for the natural trajectory of the system $\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ that corresponds to the value $h$ of the vis viva

[^14]constant and passes through the points $M_{1}$ and $M_{2}$. That is what the principle of least action consists of.

Hamilton's principle. - When the forces are derived from a potential and one leaves the Lagrange equations in their original form:

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial T}{\partial q_{i}^{\prime}}-\frac{\partial(T+U)}{\partial q_{i}}=0 \quad(i=1,2, \ldots, k) \tag{1}
\end{equation*}
$$

one can attach them to the calculus of variations in another way.
Indeed, let $q_{1}, q_{2}, \ldots, q_{k}$ be $k$ arbitrary functions of $t$ that are subject to the condition that $q_{1}, q_{2}$, $\ldots, q_{k}$ have fixed the values $\left(a_{1}, a_{2}, \ldots, a_{k}\right)$ for $t=t_{0}$ and $\left(b_{1}, b_{2}, \ldots, b_{k}\right)$ for $t=t_{1}$. Now consider the integral $\int_{t_{0}}^{t_{1}}(T+U) d t$. If one expresses the idea that the variation of that integral is zero when one passes from a particular system of functions $q_{i}=\varphi_{i}(t)$ to a neighboring system then one will find that those functions $q_{i}=\varphi_{i}(t)$ must verify equations (1) precisely. Those equations then express the idea that the variation of the integral $\int_{t_{0}}^{t_{1}}(T+U) d t$ must be zero for the motion $q_{i}=$ $\varphi_{i}(t)$ of the material system, which is a motion that is defined by the conditions that the system must occupy two well-defined positions ( $a_{1}, a_{2}, \ldots, a_{k}$ ) and ( $b_{1}, b_{2}, \ldots, b_{k}$ ) for $t=t_{0}$ and $t=t_{1}$. That is what Hamilton's principle consists of.

Remarks on the principle of least action. - When the forces are derived from a potential, the principle of least action reduces the search for trajectories to the search for geodesics of an $d s^{2}$ that depends upon the arbitrary constant $h$. Indeed, if one sets:

$$
T^{\prime}=(U+h) T \quad \text { or } \quad d s^{\prime 2}=(U+h) d s^{2}=(U+h) \sum A_{i j} d q_{i} d q_{j}
$$

then the geodesics of $d s^{\prime 2}$ will be given by the equations:

$$
\begin{equation*}
\frac{d}{d q_{1}} \frac{\partial f}{\partial q_{(i)}^{\prime}}-\frac{\partial f}{\partial q_{i}}=0, \quad \frac{d q_{i}}{d q_{1}}=q_{(i)}^{\prime} \quad(i=2,3, \ldots, k), \tag{5}
\end{equation*}
$$

in which $f=d s^{\prime} / d q_{1}=\sqrt{(U+h) T_{1}}$.
If one agrees to call any $(2 k-2)$-parameter congruence of trajectories that corresponds to a well-defined value $a$ of the constant $h$ of the vis viva integral a natural congruence then one will see that the natural congruence $h=a$ will coincide with the geodesics of $(U+a) d s^{2}$, or if one
prefers, with the geodesics of $(\lambda U+\mu) d s^{2}$, where $\mu / \lambda=a$. The geodesics of $d s^{2}$ define a natural congruence that will then correspond to $a=\infty$, as we know already.

The study of the trajectories when there exists a force function then comes down to integrating the system (5), i.e., to an explicit system of $(k-1)$ second-order equations that depend upon one arbitrary parameter $h$. It is important to remark that those equations are of the same type that the calculus of variations provides, and as a result, as we will soon see, they will possess some important properties that facilitate their integration. In particular, we know a last multiplier.

Once the trajectories are known, the motion will be defined with the help of just one quadrature. What sort of relations exist between the motions that are defined, on the one hand, by the first system of Lagrange equations:

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial T}{\partial q_{i}^{\prime}}-\frac{\partial T}{\partial q_{i}}=\frac{\partial U}{\partial q_{i}}, \quad \frac{d q_{i}}{d t}=q_{i}^{\prime} \quad(i=1,2, \ldots, k), \tag{1}
\end{equation*}
$$

and on the other hand, by the system with no forces:

$$
\begin{equation*}
\frac{d}{d t} \frac{\partial T^{\prime}}{\partial q_{i}^{\prime}}-\frac{\partial T^{\prime}}{\partial q_{i}}=0, \quad \frac{d q_{i}}{d t}=q_{i}^{\prime} \quad(i=1,2, \ldots, k) \tag{1'}
\end{equation*}
$$

in which $T^{\prime}=(U+h) T$ ?
We know that the trajectories are the same for the two systems. Let us consider one of those trajectories. The motion along that trajectory, according to (1), is defined by the equality:

$$
d t^{2}=\frac{d s^{2}}{U+h}
$$

and from (1'), by:

$$
d t_{1}^{2}=\alpha(U+h) d s^{2},
$$

in which $\alpha$ denotes an arbitrary constant. One then passes from the first motion to the second one by changing $d t^{2}$ into $\frac{d t_{1}^{2}}{\alpha(U+h)^{2}}$, in which $\alpha$ is an arbitrary constant.

## Darboux's transformation.

We just said that when the forces are derived from a potential, the trajectories of a system will coincide for each value of $h$ with the geodesics of $(U+h) d s^{2}$. From that, consider a system whose $d s^{2}$ is equal to $(\alpha U+\beta) d s^{2}$, and is subject to forces whose potential is $\frac{\gamma U+\delta}{\alpha U+\beta}$. The trajectories of that new system will coincide with the geodesics of $d \sigma^{2}$, when we set:

$$
d \sigma^{2}=\left[\gamma U+\delta+h_{1}(\alpha U+\beta)\right] d s^{2},
$$

in which $h_{1}$ denotes the constant of the new vis viva integral. If one establishes the following relation between $h$ and $h_{1}$ :

$$
h=\frac{\delta+\beta h_{1}}{\gamma+\alpha h_{1}}
$$

then the geodesics of $d s^{2}$ and $(U+h) d s^{2}$ (which differ by only a constant factor) will coincide. Hence, one concludes that:

If one replaces $T$ with $(\alpha U+\beta) T$ and $U$ with $U^{\prime}=\frac{\gamma U+\delta}{\alpha U+\beta}$ in a system $[T, U]$ of Lagrange equations then the trajectories will not change. Every natural congruence $h=a$ that is composed of the former trajectories is a natural congruence $h_{1}=a_{1}$ that is composed of the new trajectories. The value of $h_{1}$ that corresponds to one value of $h$ is given by the equality:

$$
h_{1}=\frac{\delta-\gamma h_{1}}{\alpha h-\beta} .
$$

That transformation was pointed out for the first time by Darboux.
The motions that are defined by each of the two systems will not be the same on the same trajectory. Indeed, if we represent time by $t$ in the former motion and by $t_{1}$ in the latter one then we will have the two equalities:

$$
d s^{2}=(U+h) d t^{2}=\left(U+\frac{\delta+\beta h_{1}}{\gamma+\alpha h_{1}}\right) d t^{2}
$$

and

$$
(\alpha U+\beta) d s^{2}=\left(\frac{\gamma U+\delta}{\alpha U+\beta}+h_{1}\right) d t_{1}^{2} .
$$

Hence, upon eliminating $h_{1}$, we will infer that:

$$
\begin{equation*}
(\alpha \delta-\beta \gamma) d t_{1}^{2}=(\alpha U+\beta)^{2}\left[\alpha d s^{2}-(\alpha U+\beta) d t^{2}\right] \tag{6}
\end{equation*}
$$

That is the relation that exists between $d t$ and $d t_{1}$.
Since it is legitimate to add a constant to the force function, one can always suppose that $U^{\prime}$ has the form $\delta / \alpha U$ (in which $\alpha$ is non-zero). The relation (6) will then become (with $\beta=0$ ):

$$
\left(\frac{d t_{1}}{d t}\right)^{2}=\frac{\alpha^{2}}{\delta} U^{2}\left[\frac{d s^{2}}{d t^{2}}-U\right]=\frac{\alpha}{\delta} U^{2} h
$$

or rather:

$$
\left(\frac{d t}{d t_{1}}\right)^{2}=\frac{1}{\alpha U^{2}}\left[\alpha U \frac{d s^{2}}{d t_{1}^{2}}-\frac{\delta}{\alpha U}\right]=\frac{h_{1}}{\alpha U^{2}} .
$$

Those equalities show us that the expressions $\frac{1}{U} \frac{d t_{1}}{d t}$ and $U \frac{d t}{d t_{1}}$ are the first integrals of the two respective systems, namely, the two vis viva integrals.

Any first integral of the first system corresponds to an integral of the second one that is obtained by replacing $d t$ as a function of $d t_{1}$ using (6). The two systems are reciprocal, moreover, i.e., conversely, the first one is deduced from the second one by a Darboux transformation. All of the transforms of the transformed system coincide with the transforms of the original system.

## Properties of differential equations that are produced by the calculus of variations.

Let $x_{1}, x_{2}, \ldots, x_{n}$ be $n$ arbitrary functions of one variable $x$, and let $x_{1}^{\prime}, x_{2}^{\prime}, \ldots, x_{n}^{\prime}$ be their derivatives. Consider the integral $\int_{a}^{b} f\left(x, x_{1}, \ldots, x_{n}, x_{1}^{\prime}, \ldots, x_{n}^{\prime}\right) d x$, in which the $x_{i}$ are subject to the single condition that they must take given values at $x=a$ and $x=b$. $f$ represents a given arbitrary function on which one makes the single hypothesis that its Hessian $\Delta$ relative to the $n$ variables $x_{i}^{\prime}$ is not identically zero. If one annuls the variation of the integral in question then one will find, as we recalled, that the functions $x_{i}(x)$ must verify the differential system:

$$
\begin{equation*}
\frac{d}{d x} \frac{\partial f}{\partial x_{i}^{\prime}}-\frac{\partial f}{\partial x_{i}}=0, \quad \frac{d x_{i}}{d x}=x_{i}^{\prime} \quad(i=1,2, \ldots, n) \tag{A}
\end{equation*}
$$

It is a system of second-order equations that can be solved for the $x_{i}^{\prime \prime}$ that enter into them linearly because the determinant of the coefficients of the $x_{i}^{\prime \prime}$ is nothing but the Hessian $\Delta$ of $f$. Such a system ( $A$ ) enjoys some properties that are analogous to those of a Lagrange system for which $U$ exists. More generally, any system:

$$
\begin{equation*}
\frac{d}{d x} \frac{\partial f}{\partial x_{i}^{\prime}}-\frac{\partial f}{\partial x_{i}}=X_{i}\left(x_{1}, x_{2}, \ldots, x_{n}\right), \quad \frac{d x_{i}}{d x}=x_{i}^{\prime} \quad(i=1,2, \ldots, n) \tag{B}
\end{equation*}
$$

enjoys properties that are analogous to those of a Lagrange system in which forces do not depend upon velocities.

First of all, one can reduce an arbitrary system $(B)$ to the canonical form. It suffices to set:

$$
p_{i}=\frac{\partial f}{\partial x_{i}^{\prime}} \quad(i=1,2, \ldots, n)
$$

and

$$
K=p_{1} x_{1}^{\prime}+p_{2} x_{2}^{\prime}+\cdots+p_{n} x_{n}^{\prime}-f .
$$

If one replaces the $x_{i}$ as functions of the $p_{i}$ in $F$ then one will find that:

$$
\begin{equation*}
x_{i}^{\prime}=\frac{\partial K}{\partial p_{i}}, \quad-\frac{\partial K}{\partial p_{i}}=\frac{\partial f}{\partial x_{i}} \tag{a}
\end{equation*}
$$

In order to see that, it will suffice to repeat the argument that was made for $f=T$ (page 160), which made no assumption about the form of $T$.

Equations $(B)$ can then be replaced with the following ones:

$$
\begin{equation*}
\frac{d p_{i}}{d x}=-\frac{\partial K}{\partial x_{i}}+X_{i}, \quad \frac{d x_{i}}{d x}=\frac{\partial K}{\partial p_{i}} \quad(i=1,2, \ldots, n) . \tag{C}
\end{equation*}
$$

If the $X_{i}$ are the partial derivatives of a function $U\left(x, x_{1}, x_{2}, \ldots, x_{n}\right)$ then one will again have, upon setting $H=K-U$ :

$$
\frac{d p_{i}}{d x}=-\frac{\partial H}{\partial x_{i}}, \quad \frac{d x_{i}}{d x}=\frac{\partial H}{\partial p_{i}} \quad(i=1,2, \ldots, n) .
$$

Any first integral of $(C)$, say $\varphi\left(x, x_{1}, x_{2}, \ldots, x_{n}, p_{1}, p_{2}, \ldots, p_{n}\right)=$ const., must verify the condition:

$$
\begin{equation*}
\frac{\partial \varphi}{\partial x}+(\varphi, K)+\sum_{i=1}^{n} \frac{\partial \varphi}{\partial p_{i}} X_{i} \equiv 0 \tag{b}
\end{equation*}
$$

(see page 172), in which the symbol $(\varphi, K)$ represents $\sum_{i=1}^{n}\left(\frac{\partial \varphi}{\partial x_{i}} \frac{\partial K}{\partial p_{i}}-\frac{\partial \varphi}{\partial p_{i}} \frac{\partial K}{\partial x_{i}}\right)$, as always.
In particular, if $x$ does not enter into $f$, nor as a consequence $\left({ }^{17}\right)$ into $K$, then the independent first integrals of $x$ are characterized by the condition:

$$
\frac{\partial \varphi}{\partial x}+(\varphi, K)+\sum_{i=1}^{n} \frac{\partial \varphi}{\partial p_{i}} X_{i}=0
$$

When the $X_{i}$ are derived from a function $U\left(x, x_{1}, x_{2}, \ldots, x_{n}\right)$, equation (b) can be written:

$$
\frac{\partial \varphi}{\partial x}+(\varphi, H)=0 .
$$

[^15]If $f$ and $U$ do not depend upon $x$ then the independent first integrals of $x$ will simply verify the condition that:

$$
(\varphi, H)=0 .
$$

In the latter case, $H \equiv K-U=$ const. is obviously an integral. That integral will reduce to $K=$ const. if $X_{1} \equiv X_{2} \equiv \ldots \equiv X_{n} \equiv 0$.

Finally, equations $(C)$ admit unity as a last multiplier (see pp. 238). It will then suffice to know $(2 n-1)$ first integrals of $(C)$ if one is to achieve the integration by quadratures. If $x$ does not appear in $(C)$ then it will suffice to know $(2 n-2)$ independent integrals of $x$.

As for the original system $(B)$, it is easy to see that it admits the Hessian $\Delta$ of $f$ as a last multiplier. Indeed, admit that one knows $(2 n-1)$ first integrals of $(C)$, namely:

$$
\varphi_{i}\left(x, x_{1}, \ldots, x_{n}, p_{1}, p_{2}, \ldots, p_{n}\right)=c_{j} \quad(j=1,2, \ldots, 2 n-1) .
$$

If one infers $p_{1}, p_{2}, \ldots, p_{n}, x, x_{1}, \ldots, x_{n-2}$ as functions of $x_{n-1}$ and $x_{n}$, for example, then the expression:

$$
\frac{1}{\delta}\left(\frac{\partial K}{\partial p_{n}} d x_{n-1}-\frac{\partial K}{\partial p_{n-1}} d x_{n}\right) \equiv \frac{1}{\delta}\left[x_{n}^{\prime} d x_{n-1}-x_{n-1}^{\prime} d x_{n}\right]
$$

will be an exact differential, in which $\delta$ denotes the functional determinant:

$$
\frac{D\left(\varphi_{1}, \varphi_{2}, \varphi_{3}, \ldots, \varphi_{2 n-1}\right)}{D\left(x, x_{1}, \ldots, x_{n-2}, p_{1}, p_{2}, \ldots, p_{n}\right)} .
$$

However, if one supposes that the integrals are expressed with the aid of the $x_{i}^{\prime}$ then if one lets $\delta_{1}$ denote the determinant:

$$
\frac{D\left(\varphi_{1}, \varphi_{2}, \varphi_{3}, \ldots, \varphi_{2 n-1}\right)}{D\left(x, x_{1}, \ldots, x_{n-2}, x_{1}^{\prime}, x_{2}^{\prime}, \ldots, x_{n}^{\prime}\right)},
$$

one will have:

$$
\delta_{1}=\frac{D\left(\varphi_{1}, \varphi_{2}, \varphi_{3}, \ldots, \varphi_{2 n-1}\right)}{D\left(x, x_{1}, \ldots, x_{n-2}, p_{1}, p_{2}, \ldots, p_{n}\right)} \times \frac{D\left(\varphi_{1}, \varphi_{2}, \varphi_{3}, \ldots, \varphi_{2 n-1}\right)}{D\left(x, x_{1}, \ldots, x_{n-2}, x_{1}^{\prime}, x_{2}^{\prime}, \ldots, x_{n}^{\prime}\right)}=\frac{D\left(p_{1}, p_{2}, \ldots, p_{n}\right)}{D\left(x_{1}^{\prime}, x_{2}^{\prime}, \ldots, x_{n}^{\prime}\right)}=\delta \Delta .
$$

Hence, the expression:

$$
\frac{\Delta}{\delta_{1}}\left(x_{n}^{\prime} d x_{n-1}-x_{n-1}^{\prime} d x_{n}\right)
$$

will be an exact differential when one takes into account the $(2 n-1)$ integrals:

$$
\varphi_{j}\left(x, x_{1}, \ldots, x_{n}, x_{1}^{\prime}, x_{2}^{\prime}, \ldots, x_{n}^{\prime}\right)=\text { const. }
$$

of the system (3). The system then admits $\Delta$ as a last multiplier.
In particular, if $x$ does not enter explicitly into either $f$ or the $X_{i}$ then $\Delta$ will be a multiplier of the system:

$$
\frac{d x_{1}}{x_{1}^{\prime}}=\frac{d x_{2}}{x_{2}^{\prime}}=\ldots=\frac{d x_{n}}{x_{n}^{\prime}}=\frac{d \cdot \frac{\partial f}{\partial x_{1}^{\prime}}}{\frac{\partial f}{\partial x_{1}}+X_{1}}=\ldots=\frac{d \cdot \frac{\partial f}{\partial x_{n}^{\prime}}}{\frac{\partial f}{\partial x_{n}}+X_{n}}
$$

Finally, the integration of the canonical system $\left(C^{\prime}\right)$ can be reduced to the study of a complete integral of the partial differential equation:

$$
\frac{\partial V}{\partial x}+H\left(x, x_{1}, x_{2}, \ldots, x_{n}, \frac{\partial V}{\partial x_{1}}, \ldots, \frac{\partial V}{\partial x_{n}}\right)=0
$$

As we said in Lecture Fifteen (pp. 222-226), that makes no assumption about the form of $H$.

## Case in which $f$ is homogeneous with respect to the $x_{i}^{\prime}$.

In particular, suppose that $f$ is homogeneous with respect to the $x_{i}^{\prime}$. The degree of homogeneity $\mu$ (which can be arbitrary, real, or imaginary, moreover) must not be equal to zero or unity. Otherwise, the Euler identities would show that the Hessian $\Delta$ would be identically zero.

If one calculates the canonical function $K$ in that case then one will have:

$$
K=\sum_{i=1}^{n} x_{i}^{\prime} \frac{\partial f}{\partial x_{i}^{\prime}}-f=\mu f-f=(\mu-1) f
$$

When $\mu$ is equal to $2, K$ will coincide with $f$ on the condition that one must replace the $x_{i}^{\prime}$ with aid of the $p_{i}$ in $f$.

Furthermore, if $x$ does not enter into either $f$ or the $X_{i}$ then the system $(B)$ will imply the consequence:

$$
(\mu-1) f=\int X_{1} d x_{1}+X_{2} d x_{2}+\cdots+X_{n} d x_{n}
$$

and thus, the integral:

$$
(\mu-1) f=U+h
$$

when the $X_{i}$ are the derivatives of a function $U\left(x_{1}, x_{2}, \ldots, x_{n}\right)$.
We shall exclusively adopt the hypothesis that $x$ does not enter into $(B)$, so the relations between $x_{1}, x_{2}, \ldots, x_{n}$ will then depend upon at most $(2 n-1)$ constants. If $f$ is homogeneous then
the number of constants will be precisely $2 n-1$, as long as none of the $X_{i}$ are identically zero, in which case that number will drop down to $2 n-2$.

In order to see that, it will suffice to argue as on page 234. If one observes that the coefficients of the $x_{i}^{\prime \prime}$ in equations $(B)$ are homogeneous and of degree $\mu-2$ with respect to the $x_{i}^{\prime}$ and that the other first-order terms are homogeneous of degree $\mu$ then one will see that when equations ( $B$ ) are solved for the $x_{i}^{\prime \prime}$, they will have the form:

$$
x_{i}^{\prime \prime}=x_{1}^{\prime 2} \Pi_{i}+x_{1}^{\prime 2-\mu} \beta_{i} \quad(i=1,2, \ldots, n),
$$

in which the $\Pi_{i}, \beta_{i}$ are homogeneous and of degree zero with respect to the $x_{i}^{\prime}$. If one calculates $d^{2} x_{i} / d x_{1}^{2}$ using (B) then one will find that:

$$
\begin{equation*}
\frac{d^{2} x_{i}}{d x_{1}^{2}}=\Pi_{i} \cdot \frac{d x_{i}}{d x_{1}} \Pi_{1}+\frac{\beta_{i}-\frac{d x_{i}}{d x_{1}} \beta_{1}}{x_{1}^{\prime \mu}} \quad(i=1,2, \ldots, n) \tag{C}
\end{equation*}
$$

It follows from this, as it does for the equations of mechanics, that the relations between $x_{1}, x_{2}$, $\ldots, x_{n}$ will depend upon $2 n-1$ constants, unless the ratios:

$$
\frac{\beta_{i}-\frac{d x_{i}}{d x_{1}} \beta_{1}}{x_{1}^{\prime \mu}}
$$

are not independent of the $x_{1}^{\prime}$, which is possible only if those ratios are identically zero, which will imply that:

$$
\frac{\beta_{1}}{x_{1}^{\prime}} \equiv \frac{\beta_{2}}{x_{2}^{\prime}} \equiv \ldots \equiv \frac{\beta_{n}}{x_{n}^{\prime}} \equiv \lambda,
$$

and therefore, one will have the following values for the $X_{i}$ [see pp. 235, eq. (5)]:

$$
X_{i}=\lambda\left[\frac{\partial^{2} f}{\partial x_{i}^{\prime} \partial x_{1}^{\prime}} x_{1}^{\prime}+\frac{\partial^{2} f}{\partial x_{i}^{\prime} \partial x_{2}^{\prime}} x_{2}^{\prime}+\cdots+\frac{\partial^{2} f}{\partial x_{i}^{\prime} \partial x_{n}^{\prime}} x_{n}^{\prime}\right]=(\mu-1) \lambda \frac{\partial f}{\partial x_{i}^{\prime}} \quad(i=1,2, \ldots, n) .
$$

Those equalities demand that $\lambda$ and the $X_{i}$ must be identically zero. In other words, the derivatives $\partial f / \partial x_{i}^{\prime}$, when considered to be functions of the $x_{i}^{\prime}$, will differ by only a constant factor, and their functional determinant $\Delta$ will be identically zero.

By definition, when the functions $X_{i}$ are not all zero, the relations between $x_{1}, x_{2}, \ldots, x_{n}$ will depend upon $(2 n-1)$ constants. One can form the differential equations that define the relations in the same way as in the special case where $f$ is a quadratic form $T$ (see pp. 239). When those
equations are integrated, $x$ can be calculated as a function of $x_{1}$ by a simple quadrature by using the relation ( $C$ ) that gives $d x / d x_{1}$.

On the contrary, if all of the functions $X_{i}$ are zero then the relations between $x_{1}, x_{2}, \ldots, x_{n}$ will depend upon $(2 n-2)$ constants, and they will be defined by the system:

$$
\begin{equation*}
\frac{d^{2} x_{i}}{d x_{i}^{2}}=\Pi_{i}-\frac{d x_{i}}{d x_{1}} \Pi_{1} \quad(i=2,3, \ldots, n) \tag{D}
\end{equation*}
$$

$x$ will then be determined (once that system is integrated) by the first integral $f=h$, which can be written:

$$
\left(\frac{d x^{\prime}}{d x_{1}}\right)^{\mu}=\frac{f\left(x_{1}, x_{2}, \ldots, x_{n}, 1, \frac{d x_{2}}{d x_{1}}, \ldots, \frac{d x_{n}}{d x_{1}}\right)}{h}
$$

Finally, one can give an explicit form to the system $(D)$, which extends to the case in which there exist functions $X_{i}$ that are derivatives of $U\left(x_{1}, x_{2}, \ldots, x_{n}\right)$. Indeed, replace the system $(B)$ with the last $(n-1)$ equations of that system:

$$
\begin{equation*}
\frac{d}{d x} \frac{\partial f}{\partial x_{i}^{\prime}}-\frac{\partial f}{\partial x_{i}}=\frac{\partial U}{\partial x_{i}}, \quad \frac{d x_{i}}{d x}=x_{i}^{\prime} \quad(i=2,3, \ldots, n) \tag{E}
\end{equation*}
$$

combined with the first integral:

$$
(\mu-1) f-U=h,
$$

which can be written:

$$
\begin{equation*}
d x=d x_{1} \sqrt[\mu]{\frac{(\mu-1) f_{1}}{U+h}}, \tag{F}
\end{equation*}
$$

when one lets $f_{1}$ denote what $f$ will become when one replaces $x_{1}^{\prime}$ with $1, x_{2}^{\prime}$ with $d x_{2} / d x_{1}=x_{(2)}^{\prime}$, $\ldots$, and $x_{n}^{\prime}$ with $d x_{n} / d x_{1}=x_{(n)}^{\prime}$.

If one expresses $d x$ as a function of $d x_{1}$ using $(F)$ everywhere in equations $(E)$ then since one has:

$$
\frac{\partial f}{\partial x_{i}^{\prime}}=x_{1}^{\prime \mu-1} \frac{\partial f_{1}}{\partial x_{(i)}^{\prime}}, \quad \frac{\partial f}{\partial x_{i}}=x_{1}^{\prime \mu} \frac{\partial f}{\partial x_{i}^{\prime}},
$$

they will become (see page 270):

$$
\left[\frac{U+h}{(\mu-1) f_{1}}\right]^{1 / \mu} \times \frac{d}{d x_{1}}\left[\left(\frac{U+h}{(\mu-1) f_{1}}\right)^{1-1 / \mu} \frac{\partial f_{1}}{\partial x_{(i)}^{\prime}}\right]-\frac{U+h}{(\mu-1) f_{1}} \times \frac{\partial f_{1}}{\partial x_{i}}=\frac{\partial U}{\partial x_{i}} \quad(i=2,3, \ldots, n),
$$

which can be written:

$$
\frac{d}{d x_{1}}\left[\left(\frac{U+h}{f_{1}}\right)^{1-1 / \mu} \frac{\partial f_{1}}{\partial x_{(i)}^{\prime}}\right]-\left(\frac{U+h}{f_{1}}\right)^{1-1 / \mu} \frac{\partial f_{1}}{\partial x_{i}}=(\mu-1)\left(\frac{f_{1}}{U+h}\right)^{1 / \mu} \frac{\partial U}{\partial x_{i}},
$$

or rather, upon setting $P=(U+h)^{(1-1 / \mu)} f_{1}^{1 / \mu}$ :
(G)

$$
\frac{d}{d x_{1}} \frac{\partial P}{\partial x_{(i)}^{\prime}}-\frac{\partial P}{\partial x_{i}^{\prime}}=0 \quad(i=2,3, \ldots, n)
$$

In particular, if $X_{1} \equiv X_{2} \equiv \ldots \equiv X_{n} \equiv 0$ then equations ( $S$ ), which define the relations between $x_{1}, x_{2}, \ldots, x_{n}$, can be put into the form $(G)$, in which $P \equiv f_{1}^{1 / \mu}$. Let us say the geodesics of $d s^{\mu} \equiv$ $f\left(x_{1}, \ldots, x_{n}, d x_{1}, \ldots, d x_{n}\right)$ to mean those relations between the $x_{i}$ that depend upon ( $2 n-2$ ) constants. On the other hand, let us say the trajectories of a system $(B)$ that is defined by $\left\{f \equiv d s^{\mu} / d x^{\mu}, U\right\}$ to mean the relations between the $x_{1}, x_{2}, \ldots, x_{n}$ that imply that system. For each value $h$ of the constant in the integral $(\mu-1) f-U=h$, the trajectories will define a ( $2 n-2$ )-parameter congruence that will coincide with the geodesics of $(U+h)^{(\mu-1)} d s^{\mu}$. That is a generalization of the principle of least action.

The trajectories will not be changed if one replaces $f$ with $(a U+b)^{\mu-1} f$ and $U$ with $\frac{\gamma U+\delta}{\alpha U+\beta}$ in a system $(B)$ that is defined by $\{f, U\}$, in which $f$ is homogeneous and of degree $\mu$ with respect to the $x_{i}^{\prime}$.

One sees that the most important properties of the equations of mechanics extend to the more general equations $(B)$, in which $f$ is an arbitrary homogeneous function of the $x_{i}^{\prime}$, rather than a quadratic form $T$.

## LECTURE 17

## PROPERTIES OF FIRST INTEGRALS WHEN THE FORCES ARE DERIVED FROM A POTENTIAL. POISSON PARENTHESES. RETURN TO THE JACOBI EQUATION.

Previously, we established some properties of the first integrals of dynamics. When there exists a force function, we can complete those properties with a remarkable proposition that is due to Poisson and whose importance was shown by Jacobi.

In the case that we shall address, the motion is defined by the canonical system:

$$
\left\{\begin{align*}
\frac{d q_{i}}{d t} & =\frac{\partial H}{\partial p_{i}}  \tag{1}\\
\frac{d p_{i}}{d t} & =-\frac{\partial H}{\partial q_{i}}
\end{align*}\right.
$$

in which $H$ is a function of the $p_{i}, q_{i}$, and also $t$ when the constraints or forces depend upon time.
In order for the equality:

$$
f\left(t, q_{1}, q_{2}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}\right)=\alpha
$$

to be a first integral of the system (1), it is necessary and sufficient that $f$ should be an integral of the first-order differential equation:

$$
\begin{equation*}
\frac{\partial f}{\partial t}+(f, H)=0 \tag{2}
\end{equation*}
$$

in which the parentheses $(f, H)$ represent the sum:

$$
\sum_{i=1}^{n}\left(\frac{\partial f}{\partial q_{i}} \frac{\partial H}{\partial q_{i}}-\frac{\partial f}{\partial p_{i}} \frac{\partial H}{\partial q_{i}}\right)
$$

with Poisson's notation.
The integration of the system (1) is equivalent to the search for $2 k$ distinct particular integrals to equation (2).

Jacobi's theory of the last multiplier teaches us that it will suffice to know $(2 k-1)$ integrals of (2) in order for the last one to be obtained by quadrature.

In the case where $t$ does not enter into $H$, the first integrals that are independent of $t$ will verify the equation:

$$
(f, H)=0 .
$$

It will then suffice to know $(2 k-3)$ integrals of the equation (3) that are distinct from the integral $H=h$ in order for the integration of (1) to be achieved by quadratures.

The theorem of Poisson that we shall address is then stated:

If $f_{1}$ and $f_{2}$ are two first integrals of the system (1) then the expression $\left(f_{1}, f_{2}\right)$ will again be an integral of (1).

However, it is appropriate to observe immediately that the integral ( $f_{1}, f_{2}$ ) cannot be distinct from $f_{1}$ and $f_{2}$, and in particular, it can reduce to an absolute constant.

In order to prove the theorem, we shall begin by establishing some properties of the Poisson parentheses.

First of all, the following identities will result from definition itself of the symbol $(f, \varphi)$ :

$$
\begin{aligned}
& (f, \varphi) \equiv-(\varphi, f) \\
& (f, f) \equiv 0 \\
& (f, c) \equiv 0, \text { where } c \text { denotes a constant. }
\end{aligned}
$$

In the second place, let $f=F\left(x_{1}, x_{2}, \ldots, x_{n}\right), \varphi=\Phi\left(x_{1}, x_{2}, \ldots, x_{n}\right)$, in which the $x_{i}$ represent arbitrary functions of $q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}$. One has identically:

$$
\begin{equation*}
(f, \varphi) \equiv \sum_{i, j}\left(x_{i}, x_{j}\right)\left[\frac{\partial F}{\partial x_{i}} \frac{\partial \varphi}{\partial x_{j}}-\frac{\partial F}{\partial x_{j}} \frac{\partial \varphi}{\partial x_{i}}\right] . \tag{4}
\end{equation*}
$$

[The number of terms on the right-hand side is $n(n-1) / 2$.]
In order to verify the identity, it will suffice to express the derivatives of $f$ and $\varphi$ with respect to the $p, q$ in $(f, \varphi)$ as functions of the derivatives of $F, \Phi$ with respect to the $x_{i}$ and the derivatives of the $x_{i}$ with respect to the $p, q$ and then look for the coefficient of the product $\frac{\partial x_{i}}{\partial p_{r}} \frac{\partial x_{j}}{\partial q_{r}}$. That coefficient will coincide with its analogue in the right-hand side.

In particular, let $f=F\left(x_{1}, x_{2}, \ldots, x_{n-1}\right), \varphi=x_{n}$. One will have:

$$
(f, \varphi)=\left(x_{1}, \varphi\right) \frac{\partial F}{\partial x_{1}}+\left(x_{2}, \varphi\right) \frac{\partial F}{\partial x_{2}}+\cdots+\left(x_{n-1}, \varphi\right) \frac{\partial F}{\partial x_{n-1}} .
$$

Notably, if $f=u+v$ then one will get:

$$
(u+v, \varphi)=(u, \varphi)+(v, \varphi) .
$$

If $f=u v$ then one will get:

$$
(u v, \varphi)=v(u, \varphi)+u(v, \varphi) .
$$

However, a lemma that is much more important than that one is:

## Lemma:

Let $A, B, C$ be three functions of $p_{1}, p_{2}, \ldots, p_{k}, q_{1}, q_{2}, \ldots, q_{k}($ which can contain $t)\left({ }^{18}\right)$. If one sets:

$$
(B, C)=A^{\prime}, \quad(C, A)=B^{\prime}, \quad(A, B)=C^{\prime}
$$

then one will have:

$$
\left(A, A^{\prime}\right)+\left(B, B^{\prime}\right)+\left(C, C^{\prime}\right)=0
$$

identically.
One must then verify that the sum:

$$
S=(A,(B, C))+(B,(C, A))+(C,(A, B))
$$

is identically zero.
$(B, C)$ is an expression that is linear and homogeneous with respect to the first derivatives of $B$ and $C$. As a result, $(A,(B, C))$ will be an expression that is linear and homogeneous with respect to the second derivatives of $B$ and $C$. The sum $S$ will then be a linear, homogeneous expression in the second derivatives of $A, B, C$, and everything will come down to showing that the coefficient of any of those derivatives is identically zero.

Now, one has:

$$
\begin{gathered}
A^{\prime}=(B, C)=\sum_{i=1}^{k}\left(\frac{\partial B}{\partial q_{i}} \frac{\partial C}{\partial p_{i}}-\frac{\partial B}{\partial p_{i}} \frac{\partial C}{\partial q_{i}}\right) \\
\left(A, A^{\prime}\right)=\sum_{j=1}^{k}\left[\frac{\partial A}{\partial q_{j}} \sum_{i=1}^{k}\left(\frac{\partial B}{\partial q_{i}} \frac{\partial^{2} C}{\partial p_{i} \partial p_{j}}-\frac{\partial B}{\partial p_{i}} \frac{\partial^{2} C}{\partial q_{i} \partial p_{j}}\right)-\frac{\partial A}{\partial p_{j}} \sum_{i=1}^{k}\left(\frac{\partial B}{\partial q_{i}} \frac{\partial^{2} C}{\partial p_{i} \partial q_{j}}-\frac{\partial B}{\partial p_{i}} \frac{\partial^{2} C}{\partial q_{i} \partial q_{j}}\right)\right]+\cdots
\end{gathered}
$$

The ellipses ... represent the terms that do not contain the second derivatives of $C$.
The coefficient of $\frac{\partial^{2} C}{\partial p_{i} \partial p_{j}}(i, j \leq k)$ in $\left(A, A^{\prime}\right)$ is:

$$
\left(\frac{\partial A}{\partial q_{j}} \frac{\partial B}{\partial q_{i}}+\frac{\partial A}{\partial q_{i}} \frac{\partial B}{\partial q_{j}}\right)
$$

However, $\left(B, B^{\prime}\right) \equiv(B,(C, A)) \equiv-(B,(A, C))$ also contains a term in $\frac{\partial^{2} C}{\partial p_{j} \partial p_{i}}$, and since $(B$, $(A, C))$ differs from $\left(A, A^{\prime}\right)$ by only the permutation of $A$ and $B$, the coefficient of that term will

[^16]be the previous one (up to sign), or one will have to change $B$ into $A$ and $A$ into $B$. Since that coefficient will not change under that permutation, the two terms in $\frac{\partial^{2} C}{\partial p_{i} \partial p_{j}}$ will cancel in $S$.

Similarly, $\left(A, A^{\prime}\right)$ will provide terms:

$$
-\frac{\partial^{2} C}{\partial p_{j} \partial q_{i}}\left(\frac{\partial A}{\partial q_{j}} \frac{\partial B}{\partial p_{i}}+\frac{\partial A}{\partial p_{i}} \frac{\partial B}{\partial q_{j}}\right)+\frac{\partial^{2} C}{\partial q_{j} \partial q_{i}}\left(\frac{\partial A}{\partial p_{j}} \frac{\partial B}{\partial p_{i}}+\frac{\partial A}{\partial p_{i}} \frac{\partial B}{\partial p_{j}}\right)
$$

that do not change under the permutation of $A$ and $B$, and which will, as a result, cancel with the corresponding terms in $\left(B, B^{\prime}\right)$ or $(B,(A, C))$.

It results from this that the sum $S$ is identically zero.

## Poisson's theorem:

Let:

$$
f_{1}\left(t, q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}\right)=\alpha_{1}, \quad f_{2}\left(t, q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}\right)=\alpha_{2}
$$

be two first integrals of the canonical system. The expression:

$$
\left(f_{1}, f_{2}\right) \equiv \sum_{i=1}^{k}\left(\frac{\partial f_{1}}{\partial q_{i}} \frac{\partial f_{2}}{\partial p_{i}}-\frac{\partial f_{1}}{\partial p_{i}} \frac{\partial f_{2}}{\partial q_{j}}\right)
$$

is also an integral of that system.

If time $t$ does not enter into either $H$ or $f_{1}, f_{2}$ then the proof will be immediate. By hypothesis, one will have:

$$
\left(f_{1}, H\right)=0, \quad\left(f_{2}, H\right)=0 .
$$

Thus, one also has:

$$
\left(f_{1},\left(f_{2}, H\right)\right)=0, \quad\left(f_{2},\left(f_{1}, H\right)\right)=0
$$

However, when the preceding lemma is applied to the function $f_{1}, f_{2}, H$, that will give:

$$
\left(H,\left(f_{1}, f_{2}\right)\right)+\left(f_{1},\left(f_{2}, H\right)\right)+\left(f_{2},\left(f_{1}, H\right)\right)=0,
$$

and as a result:

$$
\left(H,\left(f_{1}, f_{2}\right)\right)=0,
$$

and the equality:

$$
\left(f_{1}, f_{2}\right)=\alpha
$$

will give an integral of (1) that is independent of $t$.
Now suppose that $t$ enters into $H, f_{1}, f_{2}$. By hypothesis, one has:

$$
\frac{\partial f_{1}}{\partial t}+\left(H, f_{1}\right)=0, \quad \frac{\partial f_{1}}{\partial t}+\left(H, f_{1}\right)=0
$$

and everything comes down to showing that one also has:

$$
\begin{equation*}
\frac{\partial\left(f_{1}, f_{2}\right)}{\partial t}+\left(H,\left(f_{1}, f_{2}\right)\right)=0 . \tag{5}
\end{equation*}
$$

Now, from the lemma, one can write:

$$
\left(H,\left(f_{1}, f_{2}\right)\right)+\left(f_{1},\left(f_{2}, H\right)\right)+\left(f_{2},\left(f_{1}, H\right)\right)=0
$$

or rather:

$$
\begin{equation*}
\left(H,\left(f_{1}, f_{2}\right)\right)+\left(f_{1}, \frac{\partial f_{2}}{\partial t}\right)-\left(f_{2}, \frac{\partial f_{1}}{\partial t}\right)=0 . \tag{6}
\end{equation*}
$$

However, from the definition of $\left(f_{1}, f_{2}\right)$, one concludes that:

$$
\frac{\partial\left(f_{1}, f_{2}\right)}{\partial t}=\sum_{i=1}^{k}\left[\frac{\partial f_{2}}{\partial p_{i}} \frac{\partial}{\partial q_{i}}\left(\frac{\partial f_{1}}{\partial t}\right)-\frac{\partial f_{2}}{\partial q_{i}} \frac{\partial}{\partial p_{i}}\left(\frac{\partial f_{1}}{\partial t}\right)\right]+\sum_{i=1}^{k}\left[\frac{\partial f_{1}}{\partial q_{i}} \frac{\partial}{\partial p_{i}}\left(\frac{\partial f_{2}}{\partial t}\right)-\frac{\partial f_{1}}{\partial p_{i}} \frac{\partial}{\partial q_{i}}\left(\frac{\partial f_{2}}{\partial t}\right)\right],
$$

i.e.:

$$
\frac{\partial\left(f_{1}, f_{2}\right)}{\partial t}=\left(\frac{\partial f_{1}}{\partial t}, f_{2}\right)+\left(f_{1}, \frac{\partial f_{2}}{\partial t}\right)
$$

Therefore, if one takes (6) into account then the equality (5) to be proved will have been established.

Remark. - From that theorem, it would seem that it would suffice to know two first integrals, $f_{1}$ and $f_{2}$, in general, for one to succeed in integrating the problem algebraically. Indeed, one can form the integral $f_{3}=\left(f_{1}, f_{2}\right)$ from $f_{1}$ and $f_{2}$, and then the integrals $f_{4}=\left(f_{1}, f_{3}\right), f_{5}=\left(f_{2}, f_{3}\right)$ from $f_{1}, f_{2}$, $f_{3}$, and so on, until one has obtained $2 k$ integrals. In reality, those integrals are not distinct in most applications, and upon repeating the Poisson process, one will come to integrals that were used already or that one can recognize immediately.

In order for us to account for that, let us see what will happen if we apply Poisson's theorem to the integrals that are provided by the theorem of the motion of the center of gravity and the theorem of areas in the case of a system of $n$ free points that are subject to only internal forces.

Let $M_{i}$ or $\left(x_{i}, y_{i}, z_{i}\right)$ be one of those points, and let $m_{i}$ be its mass:

$$
T=\sum_{i=1}^{n} \frac{1}{2} m_{i}\left(x_{i}^{\prime 2}+y_{i}^{\prime 2}+z_{i}^{\prime 2}\right) .
$$

Take the canonical variables to be:

$$
p_{i}=m_{i} x_{i}^{\prime}, \quad q_{i}=m_{i} y_{i}^{\prime}, \quad r_{i}=m_{i} z_{i}^{\prime} \quad(i=1,2, \ldots, n),
$$

so:

$$
H=\frac{1}{2} \sum_{i=1}^{n} \frac{1}{m_{i}}\left(p_{i}^{2}+q_{i}^{2}+r_{i}^{2}\right)-U\left(\ldots, x_{i}, y_{i}, z_{i}, \ldots\right)
$$

The theorem of the motion of the center of gravity gives the integrals:

$$
\sum p_{i}=\alpha, \quad \sum q_{i}=\beta, \quad \sum r_{i}=\gamma
$$

The theorem of areas gives the integrals:

$$
\begin{aligned}
& \sum\left(y_{i} r_{i}-z_{i} q_{i}\right)=A, \\
& \sum\left(z_{i} p_{i}-x_{i} r_{i}\right)=B, \\
& \sum\left(x_{i} q_{i}-y_{i} p_{i}\right)=C .
\end{aligned}
$$

Those integrals will be valid whenever the force function satisfies the conditions:

$$
\begin{aligned}
& \sum \frac{\partial U}{\partial x_{i}}=0, \quad \sum\left(y_{i} \frac{\partial U}{\partial z_{i}}-z_{i} \frac{\partial U}{\partial y_{i}}\right)=0 \\
& \sum \frac{\partial U}{\partial y_{i}}=0, \quad \sum\left(z_{i} \frac{\partial U}{\partial x_{i}}-x_{i} \frac{\partial U}{\partial z_{i}}\right)=0 \\
& \sum \frac{\partial U}{\partial z_{i}}=0, \quad \sum\left(x_{i} \frac{\partial U}{\partial y_{i}}-y_{i} \frac{\partial U}{\partial x_{i}}\right)=0
\end{aligned}
$$

Does Poisson's theorem allow us to deduce any new integrals from those six integrals? More generally, one has ( $f_{1}=$ const. and $f_{2}=$ const. denote two of those six integrals):

$$
\left(f_{1}, f_{2}\right)=\sum_{i=1}^{n}\left(\frac{\partial f_{1}}{\partial x_{i}} \frac{\partial f_{2}}{\partial p_{i}}-\frac{\partial f_{1}}{\partial p_{i}} \frac{\partial f_{2}}{\partial x_{i}}\right)+\left(\frac{\partial f_{1}}{\partial y_{i}} \frac{\partial f_{2}}{\partial q_{i}}-\frac{\partial f_{1}}{\partial q_{i}} \frac{\partial f_{2}}{\partial y_{i}}\right)+\left(\frac{\partial f_{1}}{\partial z_{i}} \frac{\partial f_{2}}{\partial r_{i}}-\frac{\partial f_{1}}{\partial r_{i}} \frac{\partial f_{2}}{\partial z_{i}}\right)
$$

here.

1. Make the associations:

$$
f_{1}=\sum p_{i}, \quad f_{2}=\sum q_{i}
$$

That will obviously give:

$$
\left(f_{1}, f_{2}\right) \equiv 0 .
$$

Poisson's theorem then leads to an identity.
2. Make the associations:

$$
f_{1}=\sum p_{i}, \quad f_{2}=\sum\left(y_{i} r_{i}-z_{i} q_{i}\right)
$$

That will again give:

$$
\left(f_{1}, f_{2}\right) \equiv 0
$$

i.e., an identity.
3. Make the associations:

$$
f_{1}=\sum p_{i}, \quad f_{2}=\sum\left(z_{i} p_{i}-x_{i} r_{i}\right)
$$

That will give:

$$
\left(f_{1}, f_{2}\right)=-\sum r_{i} .
$$

Poisson's theorem then gives back the third integral of the motion of the center of gravity.
4. Make the associations:

$$
f_{1}=\sum\left(y_{i} r_{i}-z_{i} q_{i}\right), \quad f_{2}=\sum\left(z_{i} p_{i}-x_{i} r_{i}\right)
$$

That will give:

$$
\left(f_{1}, f_{2}\right)=\sum\left(q_{i} x_{i}-p_{i} y_{i}\right)
$$

which gives back the third integral of area.
Thus, an arbitrary combination of two of the six known integrals will lead us to only one of those integrals or to an identity.

One can pose this question: Upon forming two combinations:

$$
F_{1}\left(f_{1}, f_{2}, \ldots, f_{6}\right)=\alpha_{1}, \quad F_{2}\left(f_{1}, f_{2}, \ldots, f_{6}\right)=\alpha_{2}
$$

will one obtain a new integral: $\left(F_{1}, F_{2}\right)=$ const.?
The answer is no: One will arrive at a combination of known integrals. That is a general property.

Let $f_{1}=\alpha_{1}, f_{2}=\alpha_{2}, \ldots, f_{r}=\alpha_{r}$ be $r$ first integrals such that any combination $\left(f_{i}, f_{j}\right)$ is a function of $f_{1}, f_{2}, \ldots, f_{r}(i, j \leq r)$. The combination $\left(F_{1}, F_{2}\right)$, where $F_{1}$ and $F_{2}$ are two arbitrary functions of $f_{1}, f_{2}, \ldots, f_{r}$, will also be a function of $f_{1}, f_{2}, \ldots, f_{r}$. (According to Lie, one then says that the $r$ integrals $f_{1}, \ldots, f_{r}$ form a group.)

That results immediately from the equality (4) on page 283:

$$
\begin{equation*}
\left(F_{1}, F_{2}\right) \equiv \sum_{i, j}\left(f_{i}, f_{j}\right)\left[\frac{\partial F_{1}}{\partial f_{i}} \frac{\partial F_{2}}{\partial f_{j}}-\frac{\partial F_{1}}{\partial f_{j}} \frac{\partial F_{2}}{\partial f_{i}}\right] \tag{4}
\end{equation*}
$$

Path to follow in order to apply Poisson's theorem. - Suppose that one knows $r$ distinct first integrals $f_{1}=\alpha_{1}, f_{2}=\alpha_{2}, \ldots, f_{r}=\alpha_{r}$. One then forms all combinations $\varphi=\left(f_{i}, f_{j}\right)$. If $\varphi$ is a function
of $f_{1}, f_{2}, \ldots, f_{r}$ (for any $i, j \leq r$ ) then Poisson's theorem will give nothing when $\varphi$ is identically zero or constant, in particular. Otherwise, among the $r(r-1) / 2$ combinations $\varphi$, there will exist $s$ of them, say, $f_{r+1}, \ldots, f_{r+s}$ [ $s$ can be equal to $r(r-1) / 2$ ], such that the $r+s$ integrals $f_{1}, f_{2}, \ldots, f_{r}, f_{r+1}$, $\ldots, f_{r+s}$ are distinct and all of the other combinations $\varphi$ are functions of $f_{1}, \ldots, f_{r+s}$. One then proceeds with the new system of integrals $f_{1}, \ldots, f_{r+s}$ as one did with the first one, and so on, until one arrives at a system of distinct integrals $f_{1}, f_{2}, \ldots, f_{l}(l>r)$ such that all of the combinations $\left(f_{i}\right.$, $f_{j}$ ) are functions of $f_{1}, f_{2}, \ldots, f_{l}$. That will necessarily be true after a finite number of such operations, because each operation will increase the number of integrals $f_{i}$, and that number cannot exceed $2 k$. One then forms, by definition, a system of distinct integrals $f_{1}, f_{2}, \ldots, f_{l}(l \leq 2 k)$ such that no integral $\left(f_{i}, f_{j}\right)$ is distinct from the first ones, and as a result, they will form a group.

One has thus exhausted all of the consequences of Poisson's theorem, because the combinations $\varphi$ that were overlooked along the way are all functions of $f_{1}, f_{2}, \ldots, f_{l}$, and as a result, the combinations $\left(f_{i}, \varphi\right)$, as well, from the remark that was made above. If $l=2 k-1$ then the problem will be solved by quadratures.

For example, assume that one knows two integrals $f_{1}, f_{2}$. One then forms the combination $\varphi=$ $\left(f_{1}, f_{2}\right)$. Two cases are possible according to whether $\varphi$ is or is not a function of $f_{1}, f_{2}$. In the former case, Poisson's theorem will give nothing new $\left({ }^{19}\right)$, while in the latter case, $\varphi$ will be a new integral $f_{3}$. One forms $\left(f_{1}, f_{3}\right)$ and $\left(f_{2}, f_{3}\right)$. If one supposes that those two combinations are functions of $f_{1}$, $f_{2}$ then Poisson's theorem will not permit one to add any integral to $f_{1}, f_{2}, f_{3}\left({ }^{20}\right)$.

In the case where $H$ is independent of $t, H=h$ will be one integral. Let $f_{1}=\alpha_{1}$ be a second integral that is independent of $t$, so one will have $\left(f_{1}, H\right)=0$ identically. There is then good reason to apply Poisson's theorem only if one knows at least two integrals $f_{1}, f_{2}$ (that are independent of time) that are distinct from that of vis viva. Poisson's theorem cannot provide more than ( $2 k-2$ ) such integrals (that are distinct from the vis viva integral). If it provides $2 k-3$ of them then the integration can be achieved by quadratures.

If $f_{1}$ is an integral that depends upon $t$ then the equality:

$$
\frac{\partial f_{1}}{\partial t}+\left(f_{1}, H\right)=0
$$

will show that $\partial f_{1} / \partial t$ will again be an integral, and as a result, so will all of the derivatives of $f_{1}$ with respect to $t$.

In the applications, Poisson's theorem will rarely give results that one might not have expected from the outset. However, that theorem still plays an important role, as we shall see, in the

[^17]integration of first-order partial differential equations, which is a problem that one will come to in the integration of a canonical system, as one knows.

## Return to Jacobi's method.

We said that the canonical system $\left({ }^{21}\right)$ :

$$
\left\{\begin{align*}
\frac{d q_{i}}{d t} & =\frac{\partial H}{\partial p_{i}},  \tag{1}\\
\frac{d q_{i}}{d t} & =-\frac{\partial H}{\partial q_{i}}
\end{align*} \quad(i=1,2, \ldots, k)\right.
$$

is integrated when one knows a complete integral of the partial differential equation:

$$
\begin{equation*}
\frac{\partial V}{\partial t}+H\left(t, q_{1}, q_{2}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \frac{\partial V}{\partial q_{2}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0 \tag{A}
\end{equation*}
$$

If $V\left(t, q_{1}, q_{2}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}, a_{1}, a_{2}, \ldots, a_{k}\right)$ denotes that complete integral then it will suffice to set:

$$
\begin{equation*}
\frac{\partial V}{\partial a_{i}}=b_{i}, \quad p_{i}=\frac{\partial V}{\partial q_{i}} \quad(i=1,2, \ldots, k), \tag{B}
\end{equation*}
$$

in which the $a_{i}, b_{i}$ denote constants, in order for the motion to be well-defined.
We have indicated (pp. 229-231) some important cases in which the substitution of the Jacobi equation for the canonical system is advantageous. However, more generally, is there any benefit to replacing the integration of the canonical system (1) with the search for a complete integral of equation $(A)$ ?

Theoretically, the two problems are equivalent. Indeed, once the system (1) has been integrated, it is easy to form a complete integral of $(A)$. Let $p_{i}=a_{i}$ and $q_{i}=b_{i}$ for $t=t_{0}$ and let:

$$
\begin{gather*}
p_{i}=\varphi_{i}\left(q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{k}, b_{1}, \ldots, b_{k}\right), \quad q_{i}=\psi_{i}\left(q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{k}, b_{1}, \ldots, b_{k}\right)  \tag{C}\\
(i=1,2, \ldots, k)
\end{gather*}
$$

[^18]be the general integral (1). Let $\chi(t)$ be the function of $t$ that is defined by the expression $\left({ }^{22}\right)$ :
$$
p_{1} \frac{\partial H}{\partial p_{1}}+p_{2} \frac{\partial H}{\partial p_{2}}+\cdots+p_{k} \frac{\partial H}{\partial p_{k}}-H
$$
when one replaces the $p_{i}, q_{i}$ as functions of $t$ using ( $C$ ).
If one sets $u(t)=\int_{t_{0}}^{t} \chi(t) d t$ and:
$$
v=a_{1} b_{1}+\ldots+a_{k} b_{k}+u(t)
$$
then the function $v$ will be a function of $t$ and the constants $a_{1}, \ldots, a_{k}, b_{1}, \ldots, b_{k}$. Infer the $b_{1}, \ldots, b_{k}$ as functions of $t, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{k}$ using equations $(C)$ and substitute those values in $v$. The function $V\left(t, q_{1}, q_{2}, \ldots, q_{k}, a_{1}, a_{2}, \ldots, a_{k}\right)$ thus-obtained is a complete integral of $(A)$. I shall confine myself to merely stating that theorem, whose proof is quite simple, but it would not be useful for our purposes to have it.

In reality, the question that one should pose is the following one: "Is there any advantage to studying the Jacobi equation directly, rather than the canonical system?" In order for us to account for that, we shall compare the various methods with the aid of which one can integrate the system (1), on the one hand, and the system (2), on the other.

When one starts from the canonical system, it will be necessary to integrate the ordinary differential equations (1) directly, or what amounts to the same thing, to find $2 k$ distinct integrals $f$ of the linear homogeneous partial differential equation:

$$
\begin{equation*}
\frac{\partial f}{\partial t}+(f, H)=0 \tag{2}
\end{equation*}
$$

We do not need to return to the study of those integrals. We know that, in the general case, all of the difficulty involved with that comes down to forming ( $2 k-1$ ) integrals of (2), and in the case where $t$ does not enter into $H$, one forms ( $2 k-3$ ) integrals other than $H$ of the equation $(f, H)=0$.

How can one determine a complete integral when one starts from the Jacobi equation? From a previous remark, a first method consists of integrating the corresponding canonical system (1) $\left({ }^{23}\right)$. In the case that we are concerned with, where $V$ does not enter explicitly into the equation, Cauchy's method of characteristics will lead to the same system of differential equations (1). If one wishes to integrate $(A)$ by one or the other of those two methods then there will then be no reason for replacing the equation $(A)$ with the system (1). However, Jacobi has indicated a method for integrating $(A)$ that is completely different, but has remained the same with several reprises,

[^19]and which Mayer and S. Lie have brought to a state of perfection. We shall briefly discuss that method, while referring to the ample presentation in the well-known book by Goursat on firstorder partial differential equations.

The basic principle of the method is the following:
Let $V\left(t, q_{1}, q_{2}, \ldots, q_{k}, a_{1}, a_{2}, \ldots, a_{k}\right)$ be a complete integral of equation (A). The equalities:

$$
p_{i}=\frac{\partial V}{\partial q_{i}} \quad(i=1,2, \ldots, k)
$$

can always be solved for $k$ constants (from the definition of the complete integral). If one performs that solution then one will form $k$ first integrals of the system (1).

$$
\begin{equation*}
f i\left(t, q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}\right)=a_{i} \quad(i=1,2, \ldots, k) \tag{D}
\end{equation*}
$$

and equations $(D)$ are verified identically when one replaces the $p_{i}$ with the derivatives $\partial V / \partial q_{i}$ of the complete integral. Since there always exists an infinitude of complete integrals of $(A)$, one will see that there is always an infinitude of ways of forming $k$ first integrals $f_{i}=a_{i}$ of the system (1), which can be solved for the $p_{i}$, and in such a way that the expression:

$$
\begin{equation*}
-H d t+p_{1} d q_{1}+\ldots+p_{k} d q_{k} \tag{E}
\end{equation*}
$$

is an exact total differential $d V$ when one replaces the $p_{i}$ as functions of $t, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{k}$. It is clear that once those $k$ integrals $f_{i}$ have been formed, the function $V\left(t, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{k}\right)$ that is given by the integration of that total differential will be a complete integral of $(A)$, and that as a result, the system (1) will be integrated.

Jacobi's new method will then replace the search for $2 k$ distinct integrals $f$ of equation (2) with the search for only $k$ integrals, but those $k$ integrals must be such that the expression $(E)$ is an exact total differential. We shall indicate the path to follow in order to effectively determine one such system of $k$ integrals, and we verify that this determination is always possible by virtue of that fact itself, which we knew beforehand.

When $t$ does not enter into $H$, one must determine $(k-1)$ integrals $f_{i}=a_{i}$ that are independent of $t$, and when they are combined with the integral $H=h$ and solved for $p_{i}$, that will make the expression $p_{1} d q_{1}+p_{2} d q_{2}+\ldots+p_{k} d q_{k}$ into an exact differential $d W$. The function $-h t+W\left(q_{1}\right.$, $\left.\ldots, q_{k}, a_{1}, \ldots, a_{k-1}, h\right)$ is then a complete integral of $(A)$.

In the latter case, when one starts from the canonical system, in order to arrive at quadrature, it will be necessary to determine $(2 k-3)$ integrals $f$ that are distinct from each other and from $H$, but those $(2 k-3)$ integrals are arbitrary.

There is one case in which the two methods are found to coincide, which is the one in which $t$ does not enter into $H$, and the number of parameters $k$ is equal to 2 . One will then have: $2 k-3=k$ $-1=1$. On the other hand, we said (see pp. 208-211 and pp. 226-227) that any first integral $f\left(q_{1}\right.$, $\left.q_{2}, p_{1}, p_{2}\right)=a$, when combined with $H=h$, will make ( $p_{1} d q_{1}+p_{2} d q_{2}$ ) into an exact differential. The Jacobi method, like that of the last multiplier, will therefore require only quadratures when
one knows an arbitrary integral $f$ of $(f, H)=0$. We have seen that the calculations to which those two methods lead are the same.

However, when $k$ exceeds 2 , or when $t$ enters into $H$, the same thing will no longer be true. In order to present a general form for the Jacobi method, we shall not make $t$ play a special role, and we shall consider an arbitrary equation $(A)$ :

$$
F\left(q_{1}, q_{2}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \frac{\partial V}{\partial q_{2}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0
$$

in which the number of variables $q_{i}$ is arbitrary and $V$ does not enter explicitly.

## Presentation of the method of Jacobi and Mayer.

The theorem upon which the method is based is this:
Let $k$ relations $\left({ }^{24}\right)$ :

$$
f_{i}\left(q_{1}, q_{2}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}\right)=0 \quad(i=1,2, \ldots, k)
$$

be soluble for $p_{1}, p_{2}, \ldots, p_{k}$ and such that the functional determinant $\frac{D\left(f_{1}, f_{2}, \ldots, f_{k}\right)}{D\left(p_{1}, p_{2}, \ldots, p_{k}\right)}$ is not annulled identically when one replaces the $p_{i}$ with their values $\varphi_{i}\left(q_{1}, \ldots, q_{k}\right)$ that one infers from $(\alpha)$. In order for the expression $p_{1} d q_{1}+p_{2} d q_{2}+\ldots, p_{k} d q_{k}=\sum \varphi_{i} d q_{i}$ to be an exact differential, it is necessary and sufficient that all of the equalities:

$$
\left(f_{i}, f_{j}\right)=0 \quad(i, j \leq k)
$$

(in which the $p_{i}, q_{i}$ are independent variables) must be consequences of the system ( $\alpha$ ).

The theorem supposes only that the functional determinant $\frac{D\left(f_{1}, f_{2}, \ldots, f_{k}\right)}{D\left(p_{1}, p_{2}, \ldots, p_{k}\right)}$ is not annulled identically when one replaces the $p_{i}$ as functions of the $q_{i}$ using $(\alpha)$.

In order for the expression $\sum p_{i} d q_{i}$ to be an exact total differential, it is necessary and sufficient that the functions $p_{i}=\varphi_{i}\left(q_{1}, \ldots, q_{k}\right)$ satisfy the $k(k-1) / 2$ conditions:

[^20]$$
\frac{\partial p_{i}}{\partial q_{j}} \equiv \frac{\partial p_{j}}{\partial q_{i}} \quad(i, j=1,2, \ldots, k)
$$

We shall show that those conditions are equivalent to the stated condition by appealing to the following lemma:

## Lemma:

If two simultaneous partial differential equations:

$$
f_{1}\left(q_{1}, q_{2}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}\right)=0, \quad f_{2}\left(q_{1}, q_{2}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}\right)=0
$$

in which $p_{i}=\partial V / \partial q_{i}$, admit a common integral $V\left(q_{1}, q_{2}, \ldots, q_{k}\right)$ then that integral will also verify the equation:

$$
\left(f_{1}, f_{2}\right)=0 .
$$

Indeed, let $p_{1}, \ldots, p_{k}$ be arbitrary functions of $q_{1}, \ldots, q_{k}$ that verify equations $(\beta)$. One has:

$$
\left.\begin{array}{l}
\frac{\partial f_{1}}{\partial q_{i}}+\sum_{j=1}^{k} \frac{\partial f_{1}}{\partial p_{j}} \frac{\partial p_{j}}{\partial q_{i}}=0 \\
\frac{\partial f_{2}}{\partial q_{i}}+\sum_{j=1}^{k} \frac{\partial f_{2}}{\partial p_{j}} \frac{\partial p_{j}}{\partial q_{i}}=0
\end{array}\right\} \quad(i=1,2, \ldots, k)
$$

Multiply the first one by $\partial f_{2} / \partial p_{i}$, the second of them by $\partial f_{1} / \partial p_{i}$, and add corresponding sides. In the result, one then gives the values $1,2, \ldots, k$ to $i$ in succession and takes the sum of the equations thus-obtained:

$$
\left(f_{1}, f_{2}\right)+\sum_{j=1}^{k} \sum_{i=1}^{k}\left(\frac{\partial f_{2}}{\partial p_{i}} \frac{\partial f_{1}}{\partial p_{j}}-\frac{\partial f_{2}}{\partial p_{j}} \frac{\partial f_{1}}{\partial p_{i}}\right) \frac{\partial p_{j}}{\partial q_{i}}=0
$$

which can be further written:

$$
\left(f_{1}, f_{2}\right)+\sum_{j=1}^{k} \sum_{i=1}^{k} \frac{\partial f_{2}}{\partial p_{i}} \frac{\partial f_{1}}{\partial p_{j}}\left(\frac{\partial p_{j}}{\partial q_{i}}-\frac{\partial p_{i}}{\partial q_{j}}\right)=0 .
$$

If $p_{1}, \ldots, p_{k}$ are the derivatives $\partial V / \partial q_{1}, \ldots, \partial V / \partial q_{k}$ of a function $V$ then the latter equality will become:

$$
\left(f_{1}, f_{2}\right)=0,
$$

which proves the lemma in question.
An immediate consequence of that lemma is that any integral $V\left(q_{1}, \ldots, q_{k}\right)$ that is common to equations ( $\alpha$ ) will also verify the equations:

$$
\left(f_{i}, f_{j}\right)=0 \quad(i, j=1,2, \ldots, k)
$$

If one then infers the $p_{i}$ as functions of the $q_{i}$ from $(\alpha)$ then since one has $p_{i}=\partial V / \partial q_{i}$, the conditions $\left(f_{i}, f_{j}\right)=0$ must be verified identically when one replaces the $p_{i}$ with those values. In other words, the equalities $\left(f_{i}, f_{j}\right)=0$ must be consequences of the system $(\alpha)$.

The stated conditions then appear to be necessary. It remains for us to prove that they are sufficient.

In order to do that, recall that the system of functions $p_{i}=\varphi_{i}\left(q_{1}, \ldots, q_{k}\right)$ that is defined by an arbitrary system $(\alpha)$ verifies the equations:

$$
\left(f_{r}, f_{s}\right)+\sum_{j=1}^{k} \sum_{i=1}^{k} \frac{\partial f_{r}}{\partial p_{i}} \frac{\partial f_{s}}{\partial p_{j}}\left(\frac{\partial p_{j}}{\partial q_{i}}-\frac{\partial p_{i}}{\partial q_{j}}\right)=0
$$

Replace the variables $p_{1}, \ldots, p_{k}$ with $\varphi_{1}, \ldots, \varphi_{k}$ in $(\gamma)$. By hypothesis, the expressions $\left(f_{r}, f_{s}\right)$ are then annulled identically and equations ( $\gamma$ ) will then become:

$$
\sum_{j=1}^{k} \sum_{i=1}^{k} \frac{\partial f_{r}}{\partial p_{i}} \frac{\partial f_{s}}{\partial p_{j}}\left(\frac{\partial p_{j}}{\partial q_{i}}-\frac{\partial p_{i}}{\partial q_{j}}\right) \equiv \sum_{j=1}^{k} \frac{\partial f_{r}}{\partial p_{j}} \sum_{i=1}^{k} \frac{\partial f_{s}}{\partial p_{i}}\left(\frac{\partial p_{i}}{\partial q_{j}}-\frac{\partial p_{j}}{\partial q_{i}}\right) \equiv \sum_{j=1}^{k} \frac{\partial f_{r}}{\partial p_{j}} \xi_{j}=0
$$

Since the determinant $\frac{D\left(f_{1}, \ldots, f_{k}\right)}{D\left(p_{1}, \ldots, p_{k}\right)}$ is not zero identically when one replaces the $p_{i}$ with $\varphi_{i}$ ( $q_{1}$, $\ldots, q_{k}$ ), the $k$ homogeneous linear equalities in the $\xi_{j}$ that one obtains by giving the values $1,2, \ldots$, $k$ to $r$ in the latter equality will imply the consequences:

$$
\xi_{1}=0, \quad \xi_{2}=0, \quad \ldots, \quad \xi_{k}=0
$$

i.e.:

$$
\sum_{i=1}^{k} \frac{\partial f_{s}}{\partial p_{i}}\left(\frac{\partial p_{i}}{\partial q_{j}}-\frac{\partial p_{j}}{\partial q_{i}}\right)=0 \quad(s=1,2, \ldots, k)
$$

which are equalities that imply the consequences:

$$
\frac{\partial p_{i}}{\partial q_{j}}-\frac{\partial p_{j}}{\partial q_{i}}=0 \quad(i=1,2, \ldots, k)
$$

for the same reason, and that is true for any $j(j=1,2, \ldots, k)$. The expression $\sum p_{i} d q_{i}$ is then indeed an exact total differential. The stated conditions are then sufficient. The theorem is thus proved completely.

Remark on the preceding theorem. - If the equalities $(\alpha)$ have the form:
$\left(\alpha^{\prime}\right)$

$$
F_{i}\left(p_{1}, p_{2}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)-a_{i}=0
$$

(in which the $a_{i}$ denote arbitrary constants, and the determinant $\frac{D\left(F_{1}, \ldots, F_{k}\right)}{D\left(p_{1}, \ldots, p_{k}\right)}$ is not identically zero) then in order for the expression $\sum p_{i} d q_{i}$ to be an exact differential, it is necessary and sufficient that the brackets $\left(f_{i}, f_{j}\right)$ should be identically zero.

More generally, suppose that the $f_{i}$ in the system $(\alpha)$ depend upon $n$ constants $a$, and that equations $(\alpha)$ are solved for $l$ of the variables $p_{i}, m$ of the variables $q_{i}$, and $n$ of the constants $a(l+$ $m+n=k)$ : The conditions $\left(f_{i}, f_{j}\right)$ cannot be consequences of the system $(\alpha)$ without being verified identically.

Before applying the theorem to the integration of the Jacobi equation, we shall establish some properties of systems of homogeneous linear partial differential equations.

## Complete linear systems.

One knows that $j$ functions $f_{1}, f_{2}, \ldots, f_{j}$ of the $n$ arbitrary variables $\left(x_{1}, x_{2}, \ldots, x_{n}\right)(j<n)$ are called independent or distinct when there exists no identity relation $G\left(f_{1}, f_{2}, \ldots, f_{j}\right) \equiv 0$ between those functions. In order for $f_{1}, f_{2}, \ldots, f_{j}$ to not be distinct, it is necessary and sufficient that all of the functional determinants $\frac{D\left(f_{1}, f_{2}, \ldots, f_{k}\right)}{D\left(x_{h}, x_{l}, \ldots, x_{r}\right)}$ should be identically zero.

An arbitrary linear equation:

$$
X_{1}\left(x_{1}, x_{2}, \ldots, x_{n}\right) \frac{\partial f}{\partial x_{1}}+X_{2}\left(x_{1}, x_{2}, \ldots, x_{n}\right) \frac{\partial f}{\partial x_{2}}+\cdots+X_{n}\left(x_{1}, x_{2}, \ldots, x_{n}\right) \frac{\partial f}{\partial x_{n}}=0
$$

admits $(n-1)$ independent integrals $f_{1}, f_{2}, \ldots, f_{n-1}$, and its general integral has the form $f=\varphi\left(f_{1}, f_{2}\right.$, $\left.\ldots, f_{n-1}\right)$.

In particular, let $F$ be a given function of the $2 k$ variables:

$$
\left(q_{1}, \ldots, q_{k}, p_{1}, p_{2}, \ldots, p_{k}\right) .
$$

The equation:

$$
(F, f)=\sum\left(\frac{\partial F}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}-\frac{\partial F}{\partial p_{i}} \frac{\partial f}{\partial q_{i}}\right)=0
$$

admits $(2 k-1)$ independent integrals.
Having recalled that, let us prove this more general proposition:

## Theorem:

Let $m$ distinct equations ( $m \leq k$ ) be given:

$$
\left(f_{1}, f\right)=0, \quad\left(f_{2}, f\right)=0, \quad \ldots, \quad\left(f_{m}, f\right)=0
$$

in which $f_{1}, f_{2}, \ldots, f_{m}$ are given functions of the $q_{i}, p_{i}$ that verify the $m(m-1) / 2$ relations:

$$
\left(f_{i}, f_{j}\right)=0 \quad(i, j=1,2, \ldots, m)
$$

The system ( $\delta$ ) admits $(2 k-m$ ) distinct first integrals.
[Jacobi gave the name of complete system to such a system $\left({ }^{25}\right)$.]
Equations ( $\delta$ ), which are linear and homogeneous in the $\frac{\partial f}{\partial q_{i}}, \frac{\partial f}{\partial p_{i}}$, will be distinct unless all of the determinants $\frac{D\left(f_{1}, f_{2}, \ldots, f_{m}\right)}{D\left(x_{h}, x_{l}, \ldots, x_{m}\right)}$ (in which $x_{1}, x_{2}, \ldots, x_{m}$ are any $m$ of the variables $q_{i}, p_{j}$ ) are not identically zero. To say that equations $(\gamma)$ are distinct is to say that the $m$ functions $f_{1}, f_{2}, \ldots, f_{m}$ are distinct.

The number $\mu$ of distinct integrals cannot exceed $2 k-m$ then. In order to see that, let $\varphi_{1}, \varphi_{2}$, $\ldots, \varphi_{2 k-m+1}$ be $(2 k-m+1)$ independent integrals such that, as a result, at least one of the determinants $\frac{D\left(\varphi_{1}, \varphi_{2}, \ldots, \varphi_{2 k-m+1}\right)}{D\left(\ldots, p_{i}, \ldots, q_{j}, \ldots\right)}$ is non-zero. The function $\psi=\sum a_{i} \varphi_{i}$ (in which the $a_{i}$ are constants) is once more an integral of ( $\delta$ ), and one can arrange the constants $a_{i}$ in such a fashion that among the derivatives $\frac{\partial \psi}{\partial q_{i}}, \frac{\partial \psi}{\partial p_{i}},(2 k-m+1)$ of them take arbitrary values of $p_{i}, q_{i}$ : That is absurd is the equations $(\delta)$ are distinct. The theorem expresses the idea that $\mu$ is equal to $2 k-m$ precisely.

The theorem is true for $m=1$. Assume that it is true for $m$ and prove that it is true for $m+1$. Let $\left(\delta^{\prime}\right)$ be the new system that is obtained adding the equation:

$$
\left(f_{m+1}, f\right)=0
$$

to the system $(\delta)$, in which $f_{m+1}$ denotes a given function that verifies the conditions:

$$
\left(f_{1}, f_{m+1}\right)=0, \quad\left(f_{2}, f_{m+1}\right)=0, \quad \ldots, \quad\left(f_{m}, f_{m+1}\right)=0,
$$

by hypothesis, i.e., it is an integral of the first system ( $\delta$. Furthermore, let $f_{1}, f_{2}, \ldots, f_{m+1}, \varphi_{1}, \varphi_{2}$, $\ldots, \varphi_{2 k-2 m-1}$ be a system of $(2 k-m)$ distinct integrals of $(\delta)$. Any function $f=F\left(\ldots, f_{i}, \ldots, \varphi_{j}, \ldots\right)$

[^21]will also verify the system $(\delta)$, and any integral of $(\delta)$ will have that form. Let us seek to determine the function $F$ in such a fashion that the equation:
$$
0=\left(f_{m+1}, f\right) \equiv \sum_{i=1}^{m+1}\left(f_{m+1}, f_{i}\right) \frac{\partial F}{\partial f_{i}}+\sum_{j=1}^{2 k-2 m-1}\left(f_{m+1}, \varphi_{j}\right) \frac{\partial F}{\partial \varphi_{j}}
$$
is also verified. On the one hand, all of the combinations $\left(f_{m+1}, f_{i}\right)$ are identically zero. On the other hand, from Poisson's theorem, the combinations $\left(f_{m+1}, \varphi_{j}\right)$ are integrals of each equation ( $\delta$ ), and therefore of the system $(\delta)$, which then gives the identity:
$$
\left(f_{m+1}, \varphi_{j}\right) \equiv g_{j}\left(f_{2}, \ldots, f_{m}, \varphi_{1}, \varphi_{2}, \ldots, \varphi_{2 k-2 m-1}\right) \quad(j=1,2, \ldots, 2 k-2 m-1)
$$

It is then necessary and sufficient that the function $F\left(\ldots, f_{i}, \ldots, \varphi_{j}, \ldots\right)$ should satisfy the relation:

$$
\begin{equation*}
g_{1} \frac{\partial F}{\partial \varphi_{1}}+g_{2} \frac{\partial F}{\partial \varphi_{2}}+\cdots+g_{2 k-2 m-1} \frac{\partial F}{\partial \varphi_{2 k-2 m-1}}=0, \tag{d}
\end{equation*}
$$

in which the $g$ are functions of $\ldots, f_{i}, \ldots, \varphi_{j}, \ldots$ Those coefficients are not identically zero, moreover, since otherwise the system $\left(\delta^{\prime}\right)$ would admit all of the integrals of $(\delta)$, i.e., $(2 k-m)$ distinct integrals, and equations ( $\delta^{\prime}$ ) would not be distinct. Since equation ( $d$ ) admits $2 k-m-1$ distinct integrals, the same thing will be true for the system $\left(\delta^{\prime}\right)$. Q.E.D.

If $k$ is equal to $m$ then the $m$ distinct integrals of ( $\delta$ ) will be $f_{1}, f_{2}, \ldots, f_{m}$. There can exist no more than $k$ distinct functions $f_{i}$ that satisfy the conditions that $\left(f_{i}, f_{j}\right)=0(i, j=1,2, \ldots, k)$.

Remark. - When at least one of the functional determinants of $f_{1}, f_{2}, \ldots, f_{m}$ relative to $m$ of the variables $p_{1}, p_{2}, \ldots, p_{k}(m<k)$ is non-zero, say $\frac{D\left(f_{1}, f_{2}, \ldots, f_{m}\right)}{D\left(p_{1}, p_{2}, \ldots, p_{m}\right)} \neq 0$, one can always find integrals $f$ of the system $(d)$ such that the determinant $\frac{D\left(f_{1}, f_{2}, \ldots, f_{m}, f\right)}{D\left(p_{1}, p_{2}, \ldots, p_{m}, p_{m+1}\right)}$ is non-zero.

In other words, the equality:

$$
\frac{D\left(f_{1}, f_{2}, \ldots, f_{m}, f\right)}{D\left(p_{1}, p_{2}, \ldots, p_{m}, p_{m+1}\right)} \equiv A_{1} \frac{\partial f}{\partial p_{1}}+A_{2} \frac{\partial f}{\partial p_{2}}+\cdots+A_{m+1} \frac{\partial f}{\partial p_{m+1}}=0 \quad\left(A_{m+1} \neq 0\right)
$$

will be a consequence of the system ( $\delta$ ), i.e., one can eliminate the $\partial f$ / $\partial q_{i}$ from the $m$ equations $(\delta)$, which is absurd, since the determinant of the coefficients of $\frac{\partial f}{\partial q_{1}}, \frac{\partial f}{\partial q_{2}}, \ldots, \frac{\partial f}{\partial q_{m}}$ in that system is the determinant $\frac{D\left(f_{1}, f_{2}, \ldots, f_{m}\right)}{D\left(p_{1}, p_{2}, \ldots, p_{m}\right)} \equiv A_{m+1}$, which is non-zero.

From that, let $m$ functions $f_{1}, f_{2}, \ldots, f_{m}(m \leq k)$ be given whose functional determinant relative to $m$ of the variables $p_{i}$ is not identically zero: If those functions verify the conditions:

$$
\left(f_{i}, f_{j}\right) \equiv 0 \quad(i, j=1,2, \ldots, m)
$$

then it will always be possible to find $(k-m)$ other functions $f_{m+1}, \ldots, f_{k}$ that satisfy the conditions:

$$
\left(f_{i}, f_{j}\right) \equiv 0 \quad(i, j=1,2, \ldots, k)
$$

and are such that the system:

$$
f_{i}=a_{i} \quad(i=1,2, \ldots, k)
$$

is soluble for the $p_{i}$. The expression $\sum p_{i} d q_{i}$ will then be an exact differential.

An important theorem results from that remark, which can be stated as follows:

## Theorem:

Let $m$ functions $f_{1}, f_{2}, \ldots, f_{m}(m \leq k)$ be given that satisfy the conditions $\left(f_{i}, f_{j}\right)=0(i, j=1,2, \ldots$, m) identically. If the equalities:

$$
f_{i}=0 \quad(i=1,2, \ldots, m)
$$

are soluble with respect to $m$ of the variables $p_{i}$, namely:

$$
\begin{gathered}
p_{1}-\varphi_{1}\left(p_{m+1}, p_{m+2}, \ldots, p_{k}, q_{1}, q_{2}, \ldots, q_{k}\right)=0 \\
\ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
p_{m}-\varphi_{m}\left(p_{m+1}, p_{m+2}, \ldots, p_{k}, q_{1}, q_{2}, \ldots, q_{k}\right)=0
\end{gathered}
$$

then the functions $F_{i} \equiv p_{i}-\varphi_{i}$ will further satisfy the conditions $\left(F_{i}, F_{j}\right)=0$ identically, with the only proviso being that the functional determinant:

$$
\Delta=\frac{D\left(f_{1}, f_{2}, \ldots, f_{m}\right)}{D\left(p_{1}, p_{2}, \ldots, p_{m}\right)}
$$

should not be annulled identically when one replaces $p_{1}, p_{2}, \ldots, p_{m}$ with $\varphi_{1}, \varphi_{2}, \ldots, \varphi_{m}$ in it.

When the latter restriction is satisfied, the functions $\varphi_{1}, \ldots, \varphi_{m}$ will have continuous first derivatives, except for the exceptional values of $p_{m+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}$ that annul $\Delta\left(p_{m+1}, \ldots, p_{k}\right.$, $q_{1}, \ldots, q_{k}$ ). Therefore, take a system of those variables $(\Delta \neq 0)$ at random and write down the equalities:

$$
f_{i}=a_{i} \quad(i=1,2, \ldots, m)
$$

If one solves those equalities for $p_{1}, p_{2}, \ldots, p_{m}$, namely:

$$
G_{i} \equiv p_{i}-P_{i}\left(p_{m+1}, p_{m+2}, \ldots, p_{k}, a_{1}, \ldots, a_{k}\right)=0
$$

then the function $P_{i}$ and its first derivatives will tend to $\varphi_{i}$ and its first derivatives when $a_{1}, a_{2}, \ldots$, $a_{k}$ tend to zero. The expressions $\left(G_{i}, G_{j}\right)$ tend to $\left(F_{i}, F_{j}\right)$ when $a_{1}, a_{2}, \ldots, a_{k}$ tend to zero. As a result, in order to prove the theorem, it will suffice to prove that the brackets $\left(G_{i}, G_{j}\right)$ are identically zero.

Now, one can combine the equalities:

$$
f_{i}=a_{i} \quad(i=1,2, \ldots, m)
$$

with $(k-m)$ other ones:

$$
f_{m+1}=a_{m+1}, \quad f_{m+2}=a_{m+2}, \quad \ldots, \quad f_{k}=a_{k},
$$

such that all of the equalities $\left(f_{i}, f_{j}\right)=0$ will be verified $(i, j=1,2, \ldots, k)$ and the system $f_{j}=a_{j}(j=$ $1,2, \ldots, k)$ will be soluble for the $p_{1}, \ldots, p_{k}$. The expression $\sum p_{i} d q_{i}$ will then be an exact differential. Now solve the system $f_{j}=a_{j}$ for the $p_{1}, p_{2}, \ldots, p_{m}, a_{m+1}, \ldots, a_{k}$.

That will give:

$$
G_{1} \equiv p_{1}-P_{1} \quad\left(p_{m+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{m}\right)=0,
$$

$$
G_{m} \equiv p_{m}-P_{m} \quad\left(p_{m+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{m}\right)=0
$$

$$
G_{m+1} \equiv a_{m+1}-P_{m+1}\left(p_{m+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{m}\right)=0
$$

$$
G_{k} \equiv a_{k}-P_{k} \quad\left(p_{m+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}, a_{1}, \ldots, a_{m}\right)=0 .
$$

From a previous remark, all of the conditions $\left(G_{i}, G_{j}\right)=0$ must be verified identically in order for $\sum p_{i} d q_{i}$ to be an exact differential. Thus, one will have:

$$
\left(G_{i}, G_{j}\right)=0 \quad(i, j=1,2, \ldots, m),
$$

in particular. Q.E.D. $\left({ }^{26}\right)$
We are now in a position to present the integration methods of Jacobi and Mayer.

[^22]
## Jacobi's method.

The first form that Jacobi gave to his method was the following one:
Let $f_{0}\left(p_{1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=0$ be a first-order partial differential equation that one can always suppose to be solved for one of the derivatives ( $p_{1}$, for example) or with respect to a constant when $f_{0}$ depends upon constants. In order to form a complete integral of that equation, it will suffice to determine $(k-1)$ first integrals $f_{1}, \ldots, f_{k-1}$ of the equation:

$$
\begin{equation*}
\left(f_{0}, f\right)=0 \tag{1}
\end{equation*}
$$

such that all of the conditions $\left(f_{i}, f_{j}\right) \equiv 0$ are verified identically, and the determinant $\frac{D\left(f_{0}, \ldots, f_{k-1}\right)}{D\left(p_{1}, \ldots, p_{k}\right)}$ is non-zero, moreover.

One begins by determining an integral $f_{1}$ of equation (1), which one combines with $f_{0}$, and then looks for a common integral of the system:

$$
\begin{equation*}
\left(f_{0}, f\right)=0, \quad\left(f_{1}, f\right)=0 \tag{2}
\end{equation*}
$$

which is a system that admits $(2 k-2)$ distinct integrals. In order to do that, one first determines an integral $f_{2}$ of equation (1) and then forms the combination $\left(f_{1}, f_{2}\right)$, which is again an integral of (1). If that parenthesis is identically zero then $f_{2}$ will be the desired integral. Otherwise, one sets $\left(f_{1}, f_{2}\right)$ $=f_{3}$, then forms ( $f_{1}, f_{3}$ ), and so on, until one arrives at an integral $f_{h}$ such that $\left(f_{1}, f_{h}\right)$ is a function of $f_{0}, f_{1}, \ldots, f_{h}$. One then sets $f=\varphi\left(f_{0}, f_{1}, \ldots, f_{h}\right)$, and seeks to determine $\varphi$ in such a fashion that the parenthesis $\left(f_{1}, \varphi\right)$ is zero. Now one has:

$$
\begin{aligned}
\left(f_{1}, \varphi\right) & \equiv \frac{\partial \varphi}{\partial f_{2}}\left(f_{1}, f_{2}\right)+\frac{\partial \varphi}{\partial f_{3}}\left(f_{1}, f_{3}\right)+\cdots+\frac{\partial \varphi}{\partial f_{h}}\left(f_{1}, f_{h}\right) \\
& \equiv \frac{\partial \varphi}{\partial f_{2}} f_{3}+\frac{\partial \varphi}{\partial f_{3}} f_{4}+\cdots+\frac{\partial \varphi}{\partial f_{h}} F\left(f_{0}, f_{1}, f_{2}, \ldots, f_{h}\right)=0 .
\end{aligned}
$$

One sees that $\varphi$ must satisfy a linear equation that is equivalent to a system of $h-2$ first-order ordinary differential equations (upon regarding $f_{0}, f_{1}$ as parameters).

There is nonetheless one exceptional case, namely, that of $h=2$, i.e., the case in which one has $\left(f_{r}, f_{2}\right) \equiv F\left(f_{0}, f_{1}, f_{2}\right)$, where $F$ is not identically zero. The equation in $\varphi$ will then become:

$$
\frac{\partial \varphi}{\partial f_{2}} \cdot F\left(f_{0}, f_{1}, f_{2}\right)=0
$$

and will give $f=\varphi\left(f_{0}, f_{1}\right)$, i.e., an integral that is distinct from $f_{0}, f_{1}$. In that case, it will be necessary to calculate a third integral $f_{3}$ of equation (1). If the number $h$ that relates to $f_{h}$ is not equal to 2 then
one can form an integral of (2) with $f_{3}$ by using the procedure that was just indicated. Otherwise, one sets $f=\varphi\left(f_{0}, f_{1}, f_{2}, f_{3}\right)$ and determines $\varphi$ in such a fashion that $\left(f_{1}, \varphi\right)$ is zero: That determination requires only quadratures. Indeed, $\varphi$ must satisfy the equation:

$$
\frac{\partial \varphi}{\partial f_{2}} \cdot F\left(f_{0}, f_{1}, f_{2}\right)+\frac{\partial \varphi}{\partial f_{3}} \cdot F^{\prime}\left(f_{0}, f_{1}, f_{2}\right)=0,
$$

and as a result, it will suffice to take $\varphi$ to be the function:

$$
\varphi=\int \frac{d f_{2}}{F\left(f_{0}, f_{1}, f_{2}\right)}-\int \frac{d f_{3}}{F^{\prime}\left(f_{0}, f_{1}, f_{2}\right)}
$$

in which $f_{0}$ and $f_{1}$ are regarded as parameters.
Hence, assume that one has calculated an integral of the system (2), which is an integral that we shall represent by $f_{2}$, and look for an integral of the system:

$$
\begin{equation*}
\left(f_{0}, f\right)=0, \quad\left(f_{1}, f\right)=0, \quad\left(f_{2}, f\right)=0 \tag{3}
\end{equation*}
$$

In order to do that, one first calculates a new integral, namely, $f_{3}$, of the system (2), and one forms the combination $\left(f_{2}, f_{3}\right)$, which is again an integral of (2), from Poisson's theorem. If that parenthesis is identically zero then $f_{3}$ is the desired integral. Otherwise, one sets $\left(f_{2}, f_{3}\right)=f_{4}$, then forms ( $f_{2}, f_{4}$ ) $=f_{5}$, and so on, until one arrives at an integral $f_{h}$ such that $\left(f_{2}, f_{h}\right)$ is a function $F$ of $\left(f_{0}, f_{1}, \ldots, f_{h}\right)$. One can then determine $\varphi\left(f_{0}, f_{1}, f_{2}, \ldots, f_{h}\right)$ in such a fashion that $\varphi$ is an integral of the system (3). It will suffice that $\varphi$ must verify the equation:

$$
\frac{\partial \varphi}{\partial f_{3}} f_{4}+\frac{\partial \varphi}{\partial f_{4}} f_{5}+\cdots+\frac{\partial \varphi}{\partial f_{h}} F\left(f_{1}, f_{2}, \ldots, f_{h}\right)=0
$$

There is nonetheless one exceptional case, namely, the one in which $f_{h}$ coincides with $f_{3}$ without $F$ being identically zero. It will then be necessary to determine a new integral of (2), namely, $f_{4}$. If one does not have $\left(f_{2}, f_{4}\right) \equiv F^{\prime}\left(f_{0}, f_{1}, f_{2}, f_{4}\right)$ then the preceding procedure will permit one to deduce an integral of (3) from $f_{4}$. On the contrary, if that identity is satisfied then one sets $f=\varphi\left(f_{0}, f_{1}, f_{2}\right.$, $f_{3}, f_{4}$ ), and the function:

$$
\varphi=\int \frac{d f_{2}}{F\left(f_{0}, f_{1}, f_{2}\right)}-\int \frac{d f_{3}}{F^{\prime}\left(f_{0}, f_{1}, f_{2}\right)}
$$

will be an integral of (3).
It is clear that the argument that was just made for the systems (2) and (3) can be repeated until one arrives at a system of $k$ equations. More precisely, assume that one has learned how to find an arbitrary integral of a system:

$$
\begin{equation*}
\left(f_{0}, f\right)=0, \quad\left(f_{1}, f\right)=0, \quad \ldots, \quad\left(f_{i-1}, f\right)=0 \quad(i \leq k-1), \tag{4}
\end{equation*}
$$

in which all of the conditions $\left(f_{i}, f_{j}\right) \equiv 0$ are verified $(i, j=0,1,2, \ldots, i-1)$.
Let $f_{i}$ be an integral of the system that is distinct from $f_{0}, f_{1}, \ldots, f_{i-1}$, and consider the equations:

$$
\begin{equation*}
\left(f_{0}, f\right)=0, \quad\left(f_{1}, f\right)=0, \quad \ldots, \quad\left(f_{i-1}, f\right)=0, \quad\left(f_{i}, f\right)=0 . \tag{5}
\end{equation*}
$$

If $i$ is less than $k-1$ then one can always determine an integral of (5) that is distinct from $f_{0}, f_{1}$, $\ldots, f_{i}$ in the following manner: One determines a new integral $f_{i+1}$ of (4) and then forms $\left(f_{i}, f_{i+1}\right)=$ $f_{i+2}$, etc., until one has obtained an integral $f_{h}$ such that $\left(f_{i}, f_{h}\right) \equiv F\left(f_{0}, f_{1}, \ldots, f_{i}, f_{i+1}, \ldots, f_{h}\right)$.

If one then sets $f=\varphi\left(f_{0}, f_{1}, f_{2}, \ldots, f_{i}, \ldots, f_{h}\right)$, one can determine $\varphi$ in such a fashion that $\varphi$ is an integral of (5) that is distinct from $f_{0}, f_{1}, \ldots, f_{i}$. That will be impossible only when one has $\left(f_{1}, f_{i+1}\right)$ $\equiv F\left(f_{0}, f_{1}, \ldots, f_{i}, f_{i+1}\right)$ unless $F$ is identically zero. It would then be necessary to introduce a new integral of the system (4).

When one follows that method, one will thus arrive at a system of $k$ equations:

$$
\begin{equation*}
\left(f_{0}, f\right)=0, \quad\left(f_{1}, f\right)=0, \quad \ldots, \quad\left(f_{k-1}, f\right)=0 \tag{6}
\end{equation*}
$$

such that all of the conditions $\left(f_{i}, f_{j}\right) \equiv 0$ are verified.
If one supposes, in addition, that condition is fulfilled $\left({ }^{27}\right)$ that the integrals $f_{0}, f_{1}, \ldots, f_{k-1}$ are distinct when they are considered to be functions of only the $p_{i}$ then the system (6) will define (up to an additive constant) a complete integral of the given equation $f_{0}=0$. Indeed, it will suffice to write that:

$$
\begin{equation*}
f_{0}=0, \quad f_{1}=C_{1}, \quad f_{2}=C_{2}, \quad \ldots, \quad f_{k-1}=C_{k-1} . \tag{7}
\end{equation*}
$$

The expression $p_{1} d q_{1}+p_{2} d q_{2}+\ldots+p_{k} d q_{k}$, in which one replaces the $p_{i}$ with their values that one infers from (6), will be the exact differential of a function:

$$
V\left(q_{1}, \ldots, q_{k}, C_{1}, C_{2}, \ldots, C_{k-1}\right)
$$

which is a complete integral of $f_{0}=0$.
One sees that in order to apply the preceding method, it is necessary to know at least ( $m-1$ ) distinct integrals of the equation $\left(f_{0}, f\right)=0$, in addition to the integral $f_{0}$, and to calculate the integrals of the intermediate equations in $\varphi$.

Observe that those equations in $\varphi$ bear upon a number of variables that is even larger, and as a result, they are equivalent to an ordinary differential system whose order is even higher than the one for which Poisson's theorem is most advantageous. Now, Jacobi's method does not take into account the simplification that can result from knowing new integrals of (1). That gap was filled by Sophus Lie. The method of Sophus Lie is the most perfect of all of them. It requires only the minimum number of integrations that is necessary. It completes and extends the theory of the last multiplier (in the case where the forces derive from a potential). However, presenting that method

[^23]would require too much development. We shall confine ourselves to indicating the improvements that Jacobi himself had made to his method, along with the ones that Mayer introduced. For the results of Sophus Lie, we shall refer to the works of that author, as well as the book by Goursat that was cited before. (Chapters X, XI, and XII).

## Method of Jacobi and Mayer.

Path of the calculations. - The second path that Jacobi indicated is the following one: Solve the given partial differential equation for one of its derivatives ( $p_{1}$, for example), and let:

$$
\begin{equation*}
f_{0} \equiv p_{1}-F_{1}\left(p_{2}, p_{3}, \ldots, p_{h}, q_{1}, q_{2}, \ldots, q_{h}\right)=0 \tag{1}
\end{equation*}
$$

be that equation. Form the equation:

$$
\begin{equation*}
0=\left(f_{0}, f\right) \equiv \frac{\partial f}{\partial q_{1}}+\frac{\partial F_{1}}{\partial q_{1}} \frac{\partial f}{\partial p_{1}}+\sum_{i=2}^{k} \frac{\partial F_{1}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}-\frac{\partial F_{1}}{\partial p_{i}} \frac{\partial f}{\partial q_{i}} . \tag{2}
\end{equation*}
$$

Since $p_{1}$ does not enter into (2) explicitly, calculate an integral $f$ of (2) that is independent of $p_{1}$. In other words, look for an integral $f\left(p_{2}, \ldots, p_{h}, q_{1}, \ldots, q_{h}\right)$ of the equation that is obtained by suppressing the term in $\partial f / \partial p_{1}$ from (2), which is an integral that we suppose depends upon at least one of the variables $p_{i}$, say, $p_{2}$. Only the exceptional integrals of (2) cannot satisfy that condition. We then write down the equalities:

$$
p_{1}-F_{1}\left(p_{2}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=0, \quad f_{1}\left(p_{2}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=C_{1}
$$

and solve them for $p_{1}, p_{2}$ : That will give:

$$
\left\{\begin{array}{l}
f_{0}^{\prime} \equiv p_{1}-F_{1}^{\prime}\left(p_{3}, p_{4}, \ldots, p_{k}, q_{1}, q_{2}, \ldots, q_{k}\right)=0  \tag{1'}\\
f_{1}^{\prime} \equiv p_{2}-F_{2}^{\prime}\left(p_{3}, p_{4}, \ldots, p_{k}, q_{1}, q_{2}, \ldots, q_{k}\right)=0
\end{array}\right.
$$

and from Theorem III of this lecture, the condition $\left(f_{0}^{\prime}, f_{1}^{\prime}\right) \equiv 0$ is verified, since the condition $\left(f_{0}\right.$, $\left.f_{1}\right) \equiv 0$ is itself verified $\left({ }^{28}\right)$.

Now form the equations:

[^24]\[

\left\{$$
\begin{array}{l}
0=-\left(f_{0}^{\prime}, f\right) \equiv \frac{\partial f}{\partial q_{1}}+\frac{\partial F_{1}}{\partial q_{1}} \frac{\partial f}{\partial p_{1}}+\frac{\partial F_{1}}{\partial q_{2}} \frac{\partial f}{\partial p_{2}}+\sum_{i=3}^{k} \frac{\partial F_{1}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}-\frac{\partial F_{1}}{\partial p_{i}} \frac{\partial f}{\partial q_{i}}  \tag{2}\\
0=-\left(f_{1}^{\prime}, f\right) \equiv \frac{\partial f}{\partial q_{2}}+\frac{\partial F_{2}}{\partial q_{1}} \frac{\partial f}{\partial p_{1}}+\frac{\partial F_{2}}{\partial q_{2}} \frac{\partial f}{\partial p_{2}}+\sum_{i=3}^{k} \frac{\partial F_{2}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}-\frac{\partial F_{2}}{\partial p_{i}} \frac{\partial f}{\partial q_{i}}
\end{array}
$$\right.
\]

We know that this system, into which $p_{1}, p_{2}$ do not enter explicitly, admits $(2 k-2)$ distinct integrals. Let us assume that we have determined an integral $f\left(p_{3}, p_{4}, \ldots, p_{h}, q_{1}, \ldots, q_{h}\right)$ that depends upon $p_{3}$ explicitly (we shall soon see the means for doing that). Write down the equalities:

$$
f_{0}^{\prime}=0, \quad f_{1}^{\prime}=0, \quad f\left(p_{3}, p_{4}, \ldots, p_{h}, q_{1}, \ldots, q_{h}\right)=C_{2} .
$$

Solve the last equality for $p_{3}$ and substitute that in the first one. We thus form a system:

$$
\left\{\begin{array}{l}
f_{0}^{\prime \prime} \equiv p_{1}-F_{1}^{\prime \prime}\left(p_{4}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=0,  \tag{1"}\\
f_{1}^{\prime \prime} \equiv p_{2}-F_{2}^{\prime \prime}\left(p_{4}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=0, \\
f_{2}^{\prime \prime} \equiv p_{3}-F_{3}^{\prime \prime}\left(p_{4}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=0
\end{array}\right.
$$

that satisfies the conditions:

$$
\left(f_{0}^{\prime \prime}, f_{1}^{\prime \prime}\right) \equiv 0, \quad\left(f_{0}^{\prime \prime}, f_{2}^{\prime \prime}\right) \equiv 0, \quad\left(f_{1}^{\prime \prime}, f_{2}^{\prime \prime}\right) \equiv 0
$$

from a theorem that was recalled before.
One then seeks an integral $f\left(p_{4}, \ldots, p_{h}, q_{1}, \ldots, q_{h}\right)$ of the system:

$$
\begin{equation*}
\left(f_{0}^{\prime \prime}, f\right) \equiv 0, \quad\left(f_{1}^{\prime \prime}, f\right) \equiv 0, \quad\left(f_{2}^{\prime \prime}, f\right) \equiv 0 \tag{2"}
\end{equation*}
$$

One solves the equality $f=C_{4}$ for $p_{4}$, and so on, until one succeeds in having $k$ equations that are solved for $p_{1}, \ldots, p_{h}$, namely:

$$
\begin{equation*}
\varphi \equiv p_{i}-P_{i}\left(p_{4}, \ldots, p_{h}, q_{1}, \ldots, q_{h}\right)=0 \quad(i=1,2, \ldots, k), \tag{A}
\end{equation*}
$$

which satisfies all of the conditions $\left(\varphi_{i}, \varphi_{j}\right) \equiv 0$.
Since the $P_{i}$ depend upon $(k-1)$ constants $C_{1}, C_{2}, \ldots, C_{k-1}$ (which are, conversely, expressed as functions of the $p_{2}, \ldots, p_{h}, q_{1}, \ldots, q_{h}$, from the preceding), the equalities $(A)$ will define the differential $\sum p_{i} d q_{i}$ of a complete integral $V\left(q_{1}, \ldots, q_{k}\right)$ of the given equation (1).

Integrations that must be performed. - All of the difficulty then comes down to determining an integral $f\left(p_{r+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)$ of each intermediate system:
(r)

$$
\begin{aligned}
& 0=-\left(f_{1}, f\right) \equiv \frac{\partial f}{\partial q_{1}}+\sum_{i=1}^{r} \frac{\partial F_{1}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}+\sum_{j=r+1}^{k}\left(\frac{\partial F_{1}}{\partial q_{j}} \frac{\partial f}{\partial p_{j}}-\frac{\partial F_{1}}{\partial p_{j}} \frac{\partial f}{\partial q_{j}}\right), \\
& 0=-\left(f_{2}, f\right) \equiv \frac{\partial f}{\partial q_{2}}+\sum_{i=1}^{r} \frac{\partial F_{2}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}+\sum_{j=r+1}^{k}\left(\frac{\partial F_{2}}{\partial q_{j}} \frac{\partial f}{\partial p_{j}}-\frac{\partial F_{2}}{\partial p_{j}} \frac{\partial f}{\partial q_{j}}\right) \text {, } \\
& 0=-\left(f_{r-1}, f\right) \equiv \frac{\partial f}{\partial q_{r}}+\sum_{i=1}^{r} \frac{\partial F_{r}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}+\sum_{j=r+1}^{k}\left(\frac{\partial F_{r}}{\partial q_{j}} \frac{\partial f}{\partial p_{j}}-\frac{\partial F_{r}}{\partial p_{j}} \frac{\partial f}{\partial q_{j}}\right),
\end{aligned}
$$

in which the functions $\left({ }^{29}\right) f_{l} \equiv p_{l}-F_{l}\left(p_{r+1}, \ldots, p_{r}, q_{1}, \ldots, q_{k}\right)$ satisfy all of the conditions $\left(f_{l}, f_{m}\right) \equiv$ $0(l, m=0,1,2, \ldots, r-1)$.

Observe, first of all, that each system ( $r$ ) (where $r$ is equal to at most $k-1$ ) admits $2(k-r)$ distinct integrals in which $p_{1}, p_{2}, \ldots, p_{r}$ do not appear. In other words, the system:
$\left(r^{\prime}\right)$

$$
\left\{\begin{array}{l}
\frac{\partial f}{\partial q_{1}}+\sum_{j=r+1}^{k} \frac{\partial F_{1}}{\partial q_{j}} \frac{\partial f}{\partial p_{j}}-\frac{\partial F_{1}}{\partial p_{j}} \frac{\partial f}{\partial q_{j}}=0, \\
\ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
\frac{\partial f}{\partial q_{r}}+\sum_{j=r+1}^{k} \frac{\partial F_{r}}{\partial q_{j}} \frac{\partial f}{\partial p_{j}}-\frac{\partial F_{r}}{\partial p_{j}} \frac{\partial f}{\partial q_{j}}=0
\end{array}\right.
$$

admits $2(k-r)$ distinct integrals $f\left(p_{r+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)$. In order to see that rigorously, it will suffice to repeat the argument on page 297. First of all, the first equation in $(r)$ admits $2 k-r-1$ distinct integrals of the form $f\left(p_{r+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)$. Assume that the first $l$ equations in $(r)(f<$ $r$ ) possess $2 k-r-l$ distinct common integrals of the same form and prove that an analogous theorem is true for $l+1$. If one sets $f=\Phi\left(\varphi_{1}, \ldots, \varphi_{2 k-r-l}\right)$ then one will see that, as on the cited page, in order to verify the $(l+1)^{\text {th }}$ equation in $(r), F$ must satisfy the condition:

$$
G_{1} \frac{\partial \phi}{\partial \varphi_{1}}+\cdots+G_{(2 k-r-l)} \frac{\partial \phi}{\partial \varphi_{(2 k-r-l)}}=0,
$$

in which $G_{i} \equiv\left[f_{l}, \varphi_{i}\right] \equiv g_{i}\left(\varphi_{1}, \ldots, \varphi_{2 k-r-l}\right)$. Furthermore, those coefficients $G_{i}$ are not zero. In other words, the $(j+1)^{\text {th }}$ equation in $(r)$ will be a consequence of the first $j$, which is impossible (see pp. 291).

One thus arrives at the conclusion that the system $\left(r^{\prime}\right)$ admits $2(k-r)$ distinct integrals, namely, $\varphi_{i}\left(q_{1}, \ldots, q_{k}, p_{r+1}, p_{r+2}, \ldots, p_{k}\right)[i=1,2, \ldots, 2(k-r)]$.

[^25]Let me add that the general integral $f=F\left(\varphi_{1}, \varphi_{2}, \ldots, \varphi_{2(k-r)}\right)$ of the system $\left(r^{\prime}\right)$ depends upon $p_{r+1}$ explicitly. Indeed, at least one of the functional determinants of $\varphi_{1}, \varphi_{2}, \ldots, \varphi_{2(k-r)}$ with respect to $2(k-r)$ of the variables $p_{j}, q_{j}$ is non-zero. Now, if the determinant:

$$
\frac{D\left(\varphi_{1}, \ldots, \varphi_{2(k-r)}\right)}{D\left(p_{r+1}, \ldots, p_{k}, q_{r+1}, \ldots, q_{k}\right)}
$$

is zero then all of the other ones will be so from equations $\left(r^{\prime}\right)$. It then follows from this that not all of the derivatives $\partial \varphi_{i} / \partial p_{r+1}$ can be zero, and that at least one of the $\varphi_{i}$ must depend upon $p_{r+1}$, and therefore $f$.

Finally, let $q_{1}^{0}, \ldots, q_{k}^{0}, p_{r+1}^{0}, \ldots, p_{k}^{0}$ be a system of numerical values for the $q_{i}, p_{i}$ that are taken at random, and in whose neighborhood the functions $f_{i}, F_{j}$ will be holomorphic as a result. Any integral of $\left(r^{\prime}\right)$ that is holomorphic in the neighborhood of those values will be obtained by taking $f$ to be a function $\Phi\left(\varphi_{1}, \varphi_{2}, \ldots, \varphi_{2(k-r)}\right)$ that is holomorphic in the neighborhood of the corresponding values of the $\varphi_{i}$. I say that one can choose $\Phi$ in such a fashion that for $q_{1}^{0}, \ldots, q_{r}^{0}$, $f$ coincides with a given function $A\left(q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)$ that is holomorphic in the neighborhood of $q_{r+1}^{0}, \ldots, q_{k}^{0}, p_{r+1}^{0}, \ldots, p_{k}^{0}$. Indeed, set:

$$
\varphi_{i}\left(q_{r+1}^{0}, \ldots, q_{k}^{0}, p_{r+1}^{0}, \ldots, p_{k}^{0}\right)=\xi_{i} \quad[i=1,2, \ldots, 2(k-r)]
$$

and solve those equalities for $p_{r+j}, q_{r+j}$ (which is possible). That will give (in which $q_{1}^{0}, \ldots, q_{r}^{0}$ are numbers):

$$
q_{r+j}=u_{j}\left(\xi_{1}, \ldots, \xi_{2(k-r)}\right), \quad p_{r+j}=v_{j}\left(\xi_{1}, \ldots, \xi_{2(k-r)}\right) \quad[j=1,2, \ldots,(k-r)] .
$$

If one substitutes those values in $A\left(\ldots, q_{r+j}, \ldots, p_{r+j}, \ldots\right)$ then one will obtain a function $\alpha\left(\xi_{1}, \ldots\right.$, $\left.\xi_{2(k-r)}\right)$ that must coincide with $\Phi\left(\xi_{1}, \ldots, \xi_{2(k-r)}\right)$ :

$$
\Phi\left(\xi_{1}, \ldots, \xi_{2(k-r)}\right) \equiv \alpha\left(\xi_{1}, \ldots, \xi_{2(k-r)}\right),
$$

which will determine $F$ unambiguously. There will then exist one and only one integral of the $\operatorname{system}(r)$ that coincides with $A\left(q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)$ for $q_{1}^{0}, \ldots, q_{r}^{0}$.

In particular, there exist $2(k-r)$ distinct integrals that will reduce to $q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}$, respectively, for $q_{1}^{0}, \ldots, q_{r}^{0}$.

Having said that, we shall now indicate Jacobi's integration procedure, and then that of Mayer.

## Jacobi's integration procedure.

In order to find an integral $f\left(p_{r+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)$ of the system $\left(r^{\prime}\right)$, Jacobi employed a procedure that is analogous to the one that was discussed in the context of the first method (see pp. 300). One first determines an integral $f=\varphi_{1}$ of the first equation $\left(r^{\prime}\right)$, namely:

$$
\begin{equation*}
\frac{\partial f}{\partial q_{1}}+\sum_{i=r+1}^{k} \frac{\partial F_{1}}{\partial q_{i}} \frac{\partial f}{\partial p_{i}}-\frac{\partial F_{1}}{\partial p_{i}} \frac{\partial f}{\partial q_{i}}=0 \tag{s}
\end{equation*}
$$

Since one can regard $q_{1}, \ldots, q_{k}$ as parameters in that equation, equation (1) will involve $2(k-r)+$ 1 variables, and the search for one integral $f$ will be equivalent to the search for a first integral of a system of $2(k-r)$ first-order equations.

One then deduces a common integral to the first two equations in $\left(r^{\prime}\right)$ from the integral $f=\varphi_{1}$ of $(s)$. To that effect, one forms $\left(f_{1}, \varphi_{1}\right)=\varphi_{2},\left(f_{1}, \varphi_{2}\right)=\varphi_{3}, \ldots$, which are just as many integrals of (1), until one arrives at an integral $\varphi_{h}$ such that $\left(f_{1}, \varphi_{h}\right)$ is expressed as a function of $\varphi_{1}, \ldots, \varphi_{h}$, and $q_{2}, \ldots, q_{r}$. One then sets $f=\Phi\left(q_{2}, \ldots, q_{r}, \varphi_{1}, \ldots, \varphi_{h}\right)$, and one determines $\Phi$ in such a fashion that the equality:

$$
0=\left(f_{1}, f\right) \equiv-\frac{\partial \Phi}{\partial q_{2}}+\frac{\partial \Phi}{\partial \varphi_{1}} \varphi_{2}+\cdots+\frac{\partial \Phi}{\partial \varphi_{h}} F\left(q_{2}, \ldots, q_{r}, \varphi_{1}, \ldots, \varphi_{h}\right)
$$

is verified. Since $h$ is equal to at most $2(k-r)$, one will see that it suffices to determine a first integral $\Phi$ of a system of $h$ first-order equations [ $h \leq 2(k-r)]\left({ }^{30}\right)$.

It is clear that the argument can be pursued further. Assume that one has determined an integral that is common to the first $l$ equations in $\left(r^{\prime}\right)$, say, $f=\psi_{1}\left(p_{r+1}, \ldots, p_{h}, q_{2}, \ldots, q_{r}\right)$. In order to deduce an integral of the first $(l+1)$ equations $\left(r^{\prime}\right)$, one forms $\left(f_{l}, \psi_{1}\right)=\psi_{2},\left(f_{l}, \psi_{2}\right)=\psi_{3}$, etc., until one arrives at an integral $\psi_{h}$ such that $\left(f_{l}, \psi_{h}\right)$ is expressed as a function of $\psi_{1}, \ldots, \psi_{h}$, and the variables $q_{l+1}, \ldots, q_{r}, \ldots$, which one can regard as parameters in the first $l$ equations of $(r)$ or $\left(r^{\prime}\right)$. One can then set:

$$
f=\Psi\left(q_{l+1}, \ldots, q_{r}, \psi_{1}, \ldots, \psi_{h}\right) .
$$

$\left({ }^{30}\right)$ In order for one to be able to deduce an integral $f=\varphi$ that also depends upon $p_{r+1}$, it will suffice that the integral $\varphi_{1}$ that one starts from should depend upon $p_{r+1}$ (which is always possible). Indeed, $\frac{\partial f}{\partial p_{r+1}}=$ $\frac{\partial \Phi}{\partial \varphi_{1}} \frac{\partial \varphi_{1}}{\partial p_{r+1}}+\cdots+\frac{\partial \Phi}{\partial \varphi_{h}} \frac{\partial \varphi_{h}}{\partial p_{r+1}}$, and on the other hand, one can give the values of the derivatives $\frac{\partial \Phi}{\partial \varphi_{1}}, \ldots, \frac{\partial \Phi}{\partial \varphi_{h}}$ of an integral $\Phi$ of the auxiliary equation $(\alpha)$ arbitrarily for $q_{1}^{0}, \ldots, q_{r}^{0}, \varphi_{1}^{0}, \ldots, \varphi_{h}^{0}$. Therefore, if $\frac{\partial \varphi_{1}}{\partial p_{r+1}}$ is not zero then $\frac{\partial f}{\partial p_{r+1}}$ will not be zero for any arbitrary integral $\Phi$.

In order for $f$, which is a common integral of the first $l$ equations in $(r)$, to satisfy the $(l+1)^{\text {th }}$ one, it is necessary and sufficient that one must have:

$$
0=\left(f_{l}, f\right) \equiv \frac{\partial \Psi}{\partial q_{l+1}}+\frac{\partial \Psi}{\partial \psi_{1}} \psi_{2}+\cdots+\frac{\partial \Psi}{\partial \psi_{h}} F\left(q_{l+1}, \ldots, q_{r}, \psi_{1}, \ldots, \psi_{h}\right) .
$$

Since $h$ is equal to at most $2 k-r-l-(r-l) \equiv 2(k-r)$, one will then once more have to find a first integral of an ordinary differential system of order $2(k-r)$. Let me add that the integral $f$ corresponds to an arbitrary integral $\Psi$ of $(\beta)$ that will depend upon $p_{r+1}$ if $\psi_{1}$ depends upon it, as was shown in the footnote above.

One then sees that the exceptional case that presents itself in the first method has thus been eliminated: From an arbitrary integral (that depends upon $p_{r+1}$ ) of the first equation ( $s$ ), one will deduce an integral $f_{r}$ of the system $\left(r^{\prime}\right)$ that also depends upon $p_{r+1}$, is independent of $p_{1}, \ldots, p_{r}$, and is distinct from $f_{0}, \ldots, f_{r-1}$, as a result. One can solve the system:

$$
f_{0} \equiv p_{1}-F_{1}=0, \quad f_{1} \equiv p_{2}-F_{2}=0, \quad \ldots, \quad f_{r} \equiv c_{r}
$$

for $p_{1}, \ldots, p_{r}, p_{r+1}$.
By definition, in order to find an integral $f\left(p_{r+1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)$ of the system ( $r$ ), it will be necessary to find an integral of the first equation in $\left(r^{\prime}\right)$, which is equivalent to a differential system of order $2(k-r)$, and then an integral of each of the other $(r-1)$ successive equations $(\beta), \ldots,(\lambda)$, $\ldots$, which are each equivalent to a differential system of order equal to at most $2(k-r)$.

In order to find a complete integral of the original given equation, it is necessary to find an integral of each of the $(k-1)$ successive systems $(r)$ that correspond to the values of $r: r=1, r=$ $2, \ldots, r=k-1$.

## Mayer's integration procedure.

The Mayer procedure presents a great advantage over the preceding one: It permits one to obtain an integral of the system ( $r$ ) by calculating only a first integral of a single system of differential equations of order $2(k-r)$.

We know that the system $\left(r^{\prime}\right)$ :
( $r^{\prime}$ )
(in which the parentheses extend over the variable $\left.p_{r+1}, q_{r+1}, \ldots, p_{k}, q_{k}\right)$ admits an integral $f\left(q_{1}, \ldots\right.$, $\left.q_{r}, q_{r+1}, \ldots, q_{h}, p_{r+1}, \ldots, p_{k}\right)$, namely, $f=\varphi_{1}$, which reduces an integral $A\left(q_{r+1}, \ldots, q_{h}, p_{r+1}, \ldots, p_{k}\right)$ for $q_{1}^{0}, q_{2}^{0}, \ldots, q_{r}^{0}$. If the $q_{1}^{0}, \ldots, q_{r}^{0}$ are taken at random then the functions $F_{1}, \ldots, F_{r}$ will be holomorphic for arbitrary values of:

$$
q_{r+j}, p_{r+j}, \quad \text { namely, } \quad q_{r+1}^{0}, \ldots, q_{h}^{0}, p_{r+1}^{0}, \ldots, p_{h}^{0}
$$

(as well as $A$, by hypothesis).
Having done that, make the change of variables:

$$
q_{1}=q_{1}^{0}+u_{1}, \quad q_{2}=q_{2}^{0}+u_{1} u_{2}, \ldots, \quad q_{r}=q_{r}^{0}+u_{1} u_{r}
$$

and the equalities:

$$
\begin{aligned}
& \frac{\partial f}{\partial u_{1}}=\frac{\partial f}{\partial q_{1}}+u_{2} \frac{\partial f}{\partial q_{2}}+\cdots+u_{r} \frac{\partial f}{\partial q_{r}} \\
& \frac{\partial f}{\partial u_{2}} \equiv u_{1} \frac{\partial f}{\partial q_{2}}, \ldots, \frac{\partial f}{\partial u_{r}} \equiv u_{1} \frac{\partial f}{\partial q_{r}}
\end{aligned}
$$

will show that the function $f$ of the new variables will satisfy the system:
( $R^{\prime}$ )

$$
\begin{aligned}
& \frac{\partial f}{\partial u_{1}}+\left(F_{1}^{\prime}, f\right)+u_{2}\left(F_{2}^{\prime}, f\right)+\cdots+u_{r}\left(F_{r}^{\prime}, f\right)=0 \\
& \frac{\partial f}{\partial u_{2}}+u_{1}\left(F_{2}^{\prime}, f\right)=0 \\
& \cdots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \ldots \\
& \frac{\partial f}{\partial u_{r}}+u_{1}\left(F_{r}^{\prime}, f\right)=0
\end{aligned}
$$

in which $F_{i}^{\prime}$ denotes what $F_{i}$ will become after one changes the variables, and the parentheses extend over the variables $p_{r+1}, q_{r+1}, \ldots, p_{k}, q_{k}$, as always.

The integral $f=\varphi_{1}$ of $\left(r^{\prime}\right)$ will become a function:

$$
\Psi_{1}\left(u_{1}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)
$$

that will reduce to $A\left(q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)$ for $u_{1}=0\left(u_{1}, \ldots, u_{r}\right.$ have arbitrary values). However, the functions $F_{i}^{\prime}$ are holomorphic in the domain:

$$
u_{1}=0, u_{2}^{0}, \ldots, u_{r}^{0}, q_{r+1}^{0}, \ldots, q_{k}^{0}, p_{r+1}^{0}, \ldots, p_{k}^{0},
$$

so the first equation in ( $R^{\prime}$ ) (in which $u_{2}, \ldots, u_{r}$ are parameters that can take on arbitrary values) will admit one and only one integral that reduces to $A\left(q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)$ for $u_{1}=0$. That integral will then coincide with $\psi_{1}$.

It follows from this that in order to integrate the system $\left(r^{\prime}\right)$, it suffices to integrate the first equation in $\left(R^{\prime}\right)$. Indeed, it will suffice to determine the $2(k-r)$ integrals of the latter equation that reduce to $q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}$, respectively, for $u_{1}=0$. Upon reverting to the old variables, one will have $2(k-r)$ distinct integrals of $\left(r^{\prime}\right)$.

That first equation in ( $R^{\prime}$ ) can be written:

$$
\begin{equation*}
\frac{\partial f}{\partial u_{1}}+\lambda_{1} \frac{\partial f}{\partial p_{r+1}}+\cdots+\lambda_{k-r} \frac{\partial f}{\partial p_{k}}-\mu_{1} \frac{\partial f}{\partial q_{r+1}}-\cdots-\mu_{(k-r)} \frac{\partial f}{\partial q_{k}}=0, \tag{S}
\end{equation*}
$$

in which the $\lambda_{i}$ are functions of $\mu_{1}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}$ in which one regards $u_{2}, \ldots, u_{r}$ as parameters. In order to integrate that equation, it is necessary for one to find $2(k-r)$ distinct integrals of the system:

$$
\frac{d u_{1}}{1}=\frac{d p_{r+1}}{\lambda_{1}}=\ldots=\frac{d p_{k}}{\lambda_{k-r}}=-\frac{d q_{r+1}}{\mu_{1}}=-\frac{d q_{k}}{\mu_{k-r}} .
$$

Assume that one has effectively calculated $2(k-r)$ distinct integrals of $(S)$, say, $\chi_{1}, \ldots, \chi_{2(k-r)}$. How does one deduce the $2(k-r)$ integrals $\psi_{1}, \ldots, \psi_{2(k-r)}$ that will reduce to $\ldots, q_{r+j}, \ldots, p_{r+j}, \ldots$ for $u_{1}=0$ for any $u_{2}, \ldots, u_{r}$ ? In order to do that, observe that any integral of $(S)$ such as:

$$
\chi\left(u_{1}, u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)
$$

can be expressed as a function of $\psi_{1}, \ldots, \psi_{2(k-r)}$, and $u_{2}, \ldots, u_{r}$ :

$$
\chi\left(u_{1}, u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right) \equiv F\left(u_{2}, \ldots, u_{r}, \psi_{1}, \ldots, \psi_{2(k-r)}\right),
$$

but for $u_{1}=0, \psi_{1}$ will coincide with $q_{r+1}$, etc., $\psi_{2(k-r)}$ with $p_{k}$, so one will have:

$$
\chi\left(0, u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right) \equiv F\left(u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right) .
$$

One can then write the equalities:

$$
\chi_{i}=\chi_{i}\left(0, u_{2}, \ldots, u_{r}, \psi_{1}, \ldots, \psi_{2(k-r)}\right) \quad(i=1,2, \ldots, 2(k-r)),
$$

and upon solving them $\left({ }^{31}\right)$ for $\psi_{1}, \ldots, \psi_{2(k-r)}$, one will obtain the desired integrals of $(S)$. In order to obtain $2(k-\mathrm{r})$ distinct integrals of $\left(r^{\prime}\right)$, it will suffice to replace $u_{2}, \ldots, u_{r}$ with functions of $q_{1}$, $\ldots, q_{r}$ in it.

The integration of the system $\left(r^{\prime}\right)$ is then reduced entirely to the integration of the single equation $(S)$, which is equivalent to a system of $2(k-r)$ first-order differential equations. However, we need only one particular integral of $\left(r^{\prime}\right)$. Mayer's method permits one to calculate such an integral when one knows just one integral of $(S)$.

Indeed, let $\chi\left(u_{1}, u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)$ be an integral of equation $(S)$ that we naturally assume is not simply a function of $u_{2}, \ldots, u_{r}$. For $u_{1}=0$, that integral will depend upon at least one of the variables $q_{r+1}, \ldots, p_{k}$, say $q_{r+1}$. Since it reduces (for $\left.u_{1}=0\right)$ to the form $\varpi\left(u_{2}\right.$, $\left.\ldots, u_{r}\right)$, it will coincide identically with the integral $\varpi\left(u_{2}, \ldots, u_{r}\right)$ of $(S)$. On the other hand, we have the equality:

$$
\begin{equation*}
\chi\left(u_{1}, u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right) \equiv \chi\left(0, u_{2}, \ldots, u_{r}, \psi_{1}, \psi_{2}, \ldots, \psi_{2(k-r)}\right) . \tag{l}
\end{equation*}
$$

By hypothesis, the $\psi_{i}$ satisfy not only the equation $(S)$, but also all of equations $\left(R^{\prime}\right)$, since they will be integrals of $\left(r^{\prime}\right)$ when one reverts to the original variables. The problem is then to define one of the integrals $\psi_{i}$ in terms of the integral $\chi$.

First of all, if the function $\chi$ is independent of $u_{2}, \ldots, u_{r}$ for $u_{1}=0$ then the integral $\chi$ will be the desired integral. If that is not true then one solves equation $(l)$ for $\psi_{1}$ :

$$
\psi_{1}=g_{1}\left(u_{1}, u_{2}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}, \psi_{1}, \psi_{2}, \ldots, \psi_{2(k-r)}\right) .
$$

If $g_{1}$ is independent of $\psi_{1}, \ldots, \psi_{2(k-r)}$ then $g_{1}$ will define the desired integral $\left({ }^{32}\right)$. Otherwise, express the idea that $\psi_{1}$ satisfies all of equations $\left(r^{\prime}\right)$, while taking into account the fact that $\psi_{2}$, $\ldots, \psi_{2(k-r)}$ are also integrals of that system. Each equation in $\left(r^{\prime}\right)$ will correspond to a relation of the form:

$$
\frac{\partial g_{1}}{\partial u_{1}}+v_{r+1} \frac{\partial g_{1}}{\partial p_{r+1}}+\cdots+v_{k-r} \frac{\partial g_{1}}{\partial p_{k}}-\varpi_{r+1} \frac{\partial g_{1}}{\partial q_{r+1}}-\cdots-\varpi_{k} \frac{\partial g_{1}}{\partial q_{k}}=0,
$$

in which the $v, \varpi$ are functions of $u_{1}, \ldots, u_{r}, q_{r+j}, \ldots, p_{r+j}, \ldots$ If all of the left-hand sides of those $r$ equations are identically zero (no matter what the $u, p, q, \psi$ ) then the function $g_{1}$, in which one

[^26]gives arbitrary numerical values to $\psi_{2}, \ldots, \psi_{2(k-r)}$ (which is function that cannot reduce to a constant, from the foregoing), will be the desired integral. Otherwise, at least one of those lefthand sides will be non-zero and contain the variables $\psi_{2}, \ldots, \psi_{2(k-r)}$, and one can solve the corresponding equation for $\psi_{2}$ (for example). Otherwise, there will exist a relation between the independent variables $u, p, q$. Therefore, perform that solution:
$$
\psi_{2}=g_{2}\left(u_{1}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}, \psi_{1}, \psi_{2}, \ldots, \psi_{2(k-r)}\right)
$$
and argue with $g_{2}$ as one did with $g_{1}$. Upon continuing in that way, either one will come to a function $g_{l}[l<2(k-r)]$ that satisfies the relations:
$$
\frac{\partial g_{l}}{\partial u_{1}}+v_{r+1} \frac{\partial g_{l}}{\partial p_{r+1}}+\cdots-\varpi_{r+1} \frac{\partial g_{l}}{\partial q_{r+1}}-\cdots=0
$$
identically, and that function, in which one gives arbitrary numerical values to $\psi_{l+1}, \ldots, \psi_{2(k-r)}$, will be the desired integral $\psi$, or one will arrive at an equality:
$$
\psi_{2(k-r)}=g_{2(k-r)}\left(u_{1}, \ldots, u_{r}, q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)
$$
that gives the desired integral $\psi$. Observe that in the latter case, the preceding relation $\psi_{2(k-r)-1}=$ $g_{2(k-r)-1}$ will give $\psi_{2(k-r)-1}$ as a function of $\psi_{2(k-r)}$ and independent variables, so as a function of those latter variables. Upon reassembling the series of relations $\psi_{l}=g_{l}$, one will thus deduce all of the integrals $\psi_{1}, \ldots, \psi_{2(k-r)}$ from the integral $\chi$.

However, the essential point is the following one:
In order to obtain an integral of the system $\left(r^{\prime}\right)$, it will suffice to form an integral of equation (3), which is equivalent to a system of $2(k-r)$ first-order equations.

Conclusion. - By definition, in order to integrate the given partial differential equation:

$$
f_{0} \equiv p_{1}-F_{1}\left(p_{1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)=0
$$

when one employs the Mayer procedure, one must successively calculate a first integral of an ordinary differential system of order $2(k-1)$, then a system of order $2(k-2)$, etc., and finally, of a system of order 2. At each intermediate integration, the order of the auxiliary differential system will diminish by two units.

On the contrary, when one seeks to integrate the canonical system:

$$
\frac{d q_{1}}{\frac{\partial f_{0}}{\partial p_{1}}} \equiv \ldots=\frac{d q_{k}}{\frac{\partial f_{0}}{\partial p_{k}}}=-\frac{d p_{1}}{\frac{\partial f_{0}}{\partial q_{1}}}=\ldots=-\frac{d p_{k}}{\frac{\partial f_{0}}{\partial q_{k}}}
$$

directly (where $f_{0} \equiv$ const. is one first integral), the search for one integral (that is distinct from $f_{0}$ ) will once more be a differential operation of order equal to $2 k-2$. Once that integral has been calculated, the search for a new integral will come down to the search for an integral of a system of $2 k-3$ first-order equations, and so on. The determination of an intermediate integral will permit one to lower the order of the differential system by one unit each time. One will then be led to successively look for a particular integral of a differential system of order $2 k-2$, then a system of order $2 k-3$, etc., a system of order 3, and finally a system of order 2 . The calculation is then achieved by quadratures, from the theory of the last multiplier.

That comparison will suffice to establish the superiority of the method of Jacobi and Mayer (compared to the general method) for the integration of a canonical system, at least when $k$ exceeds 2. However, one should not forget that the method supposes essentially that $H$ exists, i.e., that the forces derive from a potential.

## Remarks on the methods of Jacobi and Mayer.

When one employs Jacobi's procedure to integrate an intermediate system ( $r$ ), the integration is even more complicated than is advantageous when one applies Poisson's theorem, so it would hardly seem that one might profit from some simplifications that might result from it in the integration. When one employs Mayer's procedure, the integration is not complicated when the Poisson parentheses give results that are not illusory, but one does not employ those results. As for the first Jacobi method (see pp. 304), there is a gap in it that was filled by Sophus Lie. On that point, we shall refer to the works that were cited before.

## Examples:

I. - The motion of a free material point $x, y, z$ of mass 1 that is subject to a force that derives from a potential $U=y z+z x+x y$.

## 1. - Applying Poisson's theorem.

Since $H$ is equal to $\frac{1}{2}\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right)-(y z+z x+x y)$ here, the canonical equations of motion will be:
(A)

$$
\left\{\begin{array} { l } 
{ \frac { d x } { d t } = p _ { 1 } , } \\
{ \frac { d y } { d t } = p _ { 2 } , } \\
{ \frac { d z } { d t } = p _ { 3 } , }
\end{array} \quad ( B ) \quad \left\{\begin{array}{l}
\frac{d p_{1}}{d t}=y+z \\
\frac{d p_{2}}{d t}=z+x \\
\frac{d p_{3}}{d t}=x+y
\end{array}\right.\right.
$$

It will suffice to know five first integrals of motion that are independent of time. One already has:

$$
f_{0} \equiv\left(p_{1}^{2}+p_{2}^{2}+p_{3}^{2}\right)-2(y z+z x+x y)=h .
$$

If one adds corresponding sides of equations $(B)$ then one will find that:

$$
\frac{d}{d t}\left(p_{1}+p_{2}+p_{3}\right)=2(x+y+z)
$$

so upon multiplying the two sides of that by $p_{1}+p_{2}+p_{3} \equiv d(x+y+z) / d t$, one will get the integral:

$$
f \equiv\left(p_{1}+p_{2}+p_{3}\right)^{2}-2(x+y+z)^{2}=C .
$$

Upon subtracting the first two of equations ( $B$ ), one will then find:

$$
f_{1} \equiv\left(p_{1}-p_{2}\right)^{2}+(x-y)^{2}=c_{1}
$$

which is an integral that will imply two others by permutation, namely, $f_{2}$ and $f_{3}$. Those three integrals $f_{1}, f_{2}, f_{3}$ are distinct, but the integral $f$ is a consequence of them, since one has:

$$
\begin{equation*}
f \equiv 3 f_{1}-f_{1}-f_{2}-f_{3} . \tag{?}
\end{equation*}
$$

Let us see whether Poisson's theorem will give us any new integrals. Let us first associate $f$ and $f_{1}$. If we note that $p_{3}$ and $z$ do not enter into $f_{1}$ and that the derivatives $\frac{\partial f_{1}}{\partial p_{1}}, \frac{\partial f_{1}}{\partial p_{2}}$, on the one hand, and $\frac{\partial f_{1}}{\partial x}, \frac{\partial f_{1}}{\partial y}$, on the other, are equal and of opposite sign, moreover, while $p_{1}, p_{2}$, and $x, y$ enter into $f$ symmetrically, then we will see that $\left(f, f_{1}\right)$ is identically zero. Let us then associate ( $f_{1}$, $f_{2}$ ):

$$
\varphi \equiv\left(f_{1}, f_{2}\right) \equiv p_{1}(z-y)+p_{2}(x-z)+p_{3}(y-x)
$$

Is that integral (which can be obtained directly by summing the three area equalities) distinct from the preceding ones? It is easy to see that it is not. Indeed, write the integrals $f_{1}, f_{2}, f_{3}$ thus:

$$
p_{2}-p_{3}=\sqrt{c_{1}-(y-z)^{2}}, \quad p_{3}-p_{1}=\sqrt{c_{2}-(z-x)^{2}}, \quad p_{1}-p_{2}=\sqrt{c_{3}-(x-y)^{2}} .
$$

If one adds their corresponding sides then one will find that:

$$
0=\sqrt{c_{1}-(y-z)^{2}}+\sqrt{c_{2}-(z-x)^{2}}+\sqrt{c_{3}-(x-y)^{2}},
$$

which is a relation that is independent of the $p_{i}$. On the other hand, if one replaces $p_{1}$ and $p_{2}$ with $p_{3}-\sqrt{c_{2}-(z-x)^{2}}$ and $p_{3}+\sqrt{c_{1}-(y-z)^{2}}$, respectively, in $\varphi$ then that will give:

$$
C=-(z-y) \sqrt{c_{2}-(z-x)^{2}}+(x-z) \sqrt{c_{1}-(y-z)^{2}},
$$

which is a relation that must coincide with the preceding one, since otherwise $y$ and $z$ would be defined as functions of $x$ and only four arbitrary constants $c_{1}, c_{2}, c_{3}, C$. Furthermore, one will soon see that those two relations depend upon only the differences of the variables, namely, $x-y, y-z$. Upon eliminating one of those differences, the other one will disappear, and one will effortlessly arrive at the condition:

$$
4 C^{2}=2 c_{2} c_{3}+2 c_{3} c_{1}+2 c_{1} c_{2}-c_{1}^{2}-c_{2}^{2}-c_{3}^{2},
$$

or rather:

$$
4 \varphi^{2}=2 f_{2} f_{3}+2 f_{3} f_{1}+2 f_{1} f_{2}-f_{1}^{2}-f_{2}^{2}-f_{3}^{2} .
$$

The integral $\varphi$ is symmetric with respect to the three variables $p$ and the three variables $x, y, z$, so one will obtain the same integral by forming the combinations $\left(f_{2}, f_{3}\right)$ and $\left(f_{3}, f_{1}\right)$. Poisson's theorem will not provide any new integral then: The three integrals $f_{1}, f_{2}, f_{3}$ then form a group.

Nonetheless, observe that if one regards only the two integrals $f_{1}, f_{2}$ (combined with $f_{0}$ ) as known then Poisson's theorem will provide a new integral $\varphi$, but only one.

We have thus obtained, by definition, four distinct integrals $f_{0}, f_{1}, f_{2}, f_{3}$ (or $\varphi$ ). The theory of the last multiplier then teaches us that the problem can be solved by quadratures.

## 2. - Applying Jacobi's method.

The problem is to find a complete integral of the equation:

$$
\begin{equation*}
f_{0} \equiv p_{1}^{2}+p_{2}^{2}+p_{3}^{2}-2(y z+z x+x y)=h . \tag{C}
\end{equation*}
$$

In order to do that, employ the first form of Jacobi's method and first determine an integral of the canonical system $(A),(B)$, namely, the integral:

$$
\begin{equation*}
f \equiv\left(p_{1}+p_{2}+p_{3}\right)^{2}-2(x+y+z)^{2}=C . \tag{D}
\end{equation*}
$$

Now calculate a new integral $f_{1}$ of the system $(A),(B)$ and try to deduce an integral $\Psi$ that satisfies the condition that $(f, \Psi) \equiv 0$. Take:

$$
\begin{equation*}
f_{1} \equiv\left(p_{1}-p_{2}\right)^{2}+(x-y)^{2}=c_{1} . \tag{E}
\end{equation*}
$$

One will find that $\left(f, f_{1}\right)$ is identically zero. Therefore, equations $(C),(D),(E)$ define a complete integral of $(C)$. One infers $p_{2}$ as a function of $p_{1}$ from equation $(E)$, substitutes that in $(D)$, and then infers $p_{3}$ as a function of $p_{1}$. Finally, one substitutes that in $(C)$, and one then obtains $p_{1}, p_{2}, p_{3}$ as a function of $x, y, z$, and $h, c, c_{1}$. The complete integral $V=\int p_{1} d x+p_{2} d y+p_{3} d z$ is given by some quadratures that one can carry out until the conclusion with the aid of the change of variables $x+$ $y+z=u, y-z=v, z-x=w$.

Instead of the integral $f_{1}$, one can combine $f_{0}$ and $f$ with the integral $\varphi$ :

$$
\varphi \equiv p_{1}(z-y)+p_{2}(x-z)+p_{3}(y-x)=C .
$$

Indeed, an immediate calculation shows that $(f, \varphi)$ is identically zero. The three equations $(C),(D)$, $\left(E^{\prime}\right)$, are then symmetric with respect to the variables.

Finally, suppose that one starts from the integral $f_{1}$, and not $f$, and one then attempts to deduce an integral $\Psi$ from the second integral $f_{2}$ such that $\left(f_{1}, \Psi\right)$ is identically zero. In order to do that, one forms $\left(f_{1}, f_{2}\right) \equiv \varphi$, and then $\left(f_{1}, \varphi\right)$, which must be expressed as a function of $f_{1}, f_{2}, \varphi$, from the foregoing, and indeed one will get:

$$
\left(f_{1}, \varphi\right) \equiv f_{2}-f_{3} \equiv-f_{1}-2 \sqrt{f_{1} f_{2}-\varphi^{2}} .
$$

One then sets $\Psi=F\left(f_{1}, f_{2}, \varphi\right)$, and $F$ must satisfy the equation:

$$
\begin{equation*}
\frac{\partial F}{\partial f_{2}} \varphi-\frac{\partial F}{\partial \varphi}\left(f_{1}+2 \sqrt{f_{1} f_{2}-\varphi^{2}}\right), \tag{?}
\end{equation*}
$$

which is an equation for which one will easily find an integral $F_{1}$. Upon replacing $f_{1}$ with $c_{1}$ in $F_{1}$ and $f_{2}$ and $\varphi$ with their expressions in terms of $p_{1}, x, y, z$, one will obtain a relation $\Psi=c^{\prime}$ that will determine a complete integral when it is combined with $f_{0}=h, f_{1}=c_{1}$.

Remark. - If one had studied the system of Lagrange equations directly:

$$
x^{\prime \prime}=y+z, \quad y^{\prime \prime}=z+x, \quad z^{\prime \prime}=x+y,
$$

which define the motion of the point, then one would have had to integrate some homogeneous linear equations with constant coefficients that would have shown that in the present case, all of the integrations can be pushed to the limit, which is a result that the preceding method does not exhibit in terms of the variables $x, y, z$. More generally, before one employs Jacobi's method, it is
important to change the variables if one is to perceive the means by which one can separate at least some of the variables in the new equation in the same way that one did in the examples that were treated up to now. For example, an orthogonal change of variables $x, y, z$ will preserve $T$ here and reduce $U$ to the form $\lambda x^{2}+\mu y^{2}+v z^{2}$. Whenever $U$ is a polynomial of degree two (whether homogeneous or not) with respect to $x, y, z$, an orthogonal change of variables will reduce $U$ to the form:

$$
\lambda x^{2}+\mu y^{2}+v z^{2}+l x+m y+n z+d
$$

and the integration will then be immediate.
II. - Motion of a material point that moves in a plane and is subject to a force that derives from the potential $U=\lambda(x+y+t(x-y))$.

In order to minimize the numerical coefficients, we suppose that the point has a mass of $1 / 2$ and that $\lambda$ is equal to $1 / 4$.

We will then have:

$$
H=\frac{1}{4}\left(x^{\prime 2}+y^{\prime 2}\right)-\frac{1}{4}[x+y+t(x-y)]=p_{1}^{2}+p_{2}^{2}-\frac{1}{4}[x+y+t(x-y)],
$$

and we can find a complete integral of the equation:

$$
\begin{equation*}
p+p_{1}^{2}+p_{2}^{2}=\frac{1}{4}[x+y+t(x-y)], \tag{1}
\end{equation*}
$$

in which $p, p_{1}, p_{2}$ are the derivatives $\frac{\partial V}{\partial t}, \frac{\partial V}{\partial x}, \frac{\partial V}{\partial y}$ of a function $V(t, x, y)$.
We shall use the method of Jacobi and Mayer and first attempt to find a first integral $f\left(p_{1}, p_{2}\right.$, $t, x, y)$ of the system:

$$
\frac{4 d p_{1}}{1+t} \equiv \frac{d x}{2 p_{1}}=\frac{4 d p_{2}}{1+t}=\frac{d y}{2 p_{2}} .
$$

The combination:

$$
\frac{2\left(d p_{1}+d p_{2}\right)}{1}=\frac{d x+d y}{2\left(p_{1}+p_{2}\right)}
$$

will immediately give the integral:

$$
\begin{equation*}
\left(p_{1}+p_{2}\right)^{2}-\frac{1}{2}(x+y)=c_{1} . \tag{2}
\end{equation*}
$$

Solve equations (1) and (2) with respect to $p$ and $p_{1}$. That will give:

$$
\left\{\begin{array}{l}
p=-2 p_{2}^{2}+2 p_{2} \sqrt{\frac{1}{2}(x+y)+c_{1}}+\frac{1}{4}[t(x-y)-(x+y)]-c_{1},  \tag{3}\\
p_{1}=-p_{2}+\sqrt{\frac{1}{2}(x+y)+c_{1}} .
\end{array}\right.
$$

Now look for an integral $f\left(p_{2}, t, x, y\right)$ that is common to the two equations:

$$
\left\{\begin{array}{l}
\frac{\partial f}{\partial t}+\frac{1}{2} \frac{\partial f}{\partial p_{2}}\left[\frac{p_{2}}{\sqrt{\frac{1}{2}(x+y)+c_{1}}}-\frac{1}{2}(1+t)\right]+\frac{\partial f}{\partial y}\left[4 p_{2}-2 \sqrt{\frac{1}{2}(x+y)+c_{1}}\right]=0  \tag{4}\\
\frac{\partial f}{\partial x}+\frac{1}{4} \frac{\partial f}{\partial p_{2}} \frac{1}{\sqrt{\frac{1}{2}(x+y)+c_{1}}}+\frac{\partial f}{\partial y}=0
\end{array}\right.
$$

Let us apply Mayer's procedure. Since the coefficients in equations (4) are holomorphic for $t$ $=0$ (in which $p_{2}, x, y, z$ are arbitrary), we set $t=t, x=t \xi$. If $F$ denotes what $f$ will become under a change of variables then one will have:

$$
\frac{\partial F}{\partial t}=\frac{\partial f}{\partial t}+\xi \frac{\partial f}{\partial x}
$$

and one must determine an integral $F$ to the single equation:

$$
\begin{equation*}
\frac{\partial F}{\partial t}+\frac{1}{4} \frac{\partial F}{\partial p_{2}}\left[\frac{2 p_{2}-\xi}{\sqrt{\frac{1}{2}(t \xi+y)+c_{1}}}-(1+t)\right]+\frac{\partial F}{\partial y}\left[4 p_{2}+\xi-2 \sqrt{\frac{1}{2}(t \xi+y)+c_{1}}\right]=0 \tag{5}
\end{equation*}
$$

or rather a first integral of the system:

$$
d t=\frac{4 d p_{2}}{\frac{2 p_{2}-\xi}{\sqrt{\frac{1}{2}(t \xi+y)+c_{1}}}-(1+t)}=\frac{d y}{4 p_{2}+\xi-2 \sqrt{\frac{1}{2}(t \xi+y)+c_{2}}} .
$$

The change of variables $F=-2 p_{2}+\sqrt{\frac{1}{2}(t \xi+y)+c_{1}}$ will immediately exhibit the integral of (5):

$$
F=-\frac{1}{4} t^{2}-2 p_{2}+\sqrt{\frac{1}{2}(t \xi+y)+c_{1}} .
$$

That integral will reduce to $2 p_{2}+\sqrt{\frac{1}{2} y+c_{1}}$ for $t=0$, which is an expression that does not depend upon the variable $\xi$. As a result, it will be an integral of the system (4) when one replaces $\xi$ with $x / t$ in it. Thus, write down the equation:

$$
\begin{equation*}
2 p_{2}+\frac{1}{4} t^{2}-\sqrt{\frac{1}{2}(x+y)+c_{1}}=c_{2} \tag{6}
\end{equation*}
$$

and solve equations (3), (6) for $p, p_{1}, p_{2}$. That will give:

$$
\left\{\begin{array}{l}
p=-\frac{1}{2} c_{1}+\frac{1}{4} t(x+y)-\frac{1}{2}\left(\frac{1}{4} t^{2}-c_{2}\right)^{2}  \tag{7}\\
p_{1}=-\frac{1}{2} c_{2}+\frac{1}{8} t^{2}+\frac{1}{2} \sqrt{\frac{1}{2}(x+y)+c_{1}} \\
p_{2}=\frac{1}{2} c_{2}-\frac{1}{8} t^{2}+\frac{1}{2} \sqrt{\frac{1}{2}(x+y)+c_{1}}
\end{array}\right.
$$

and those equations must define the partial derivatives of a function $V(t, x, y)$. Indeed, one immediately sees that the function:

$$
V=\frac{2}{3}\left[\frac{1}{2}(x+y)+c_{1}\right]^{3 / 2}+\frac{1}{2}(x-y)\left(\frac{1}{4} t^{2}-c_{2}\right)-\frac{t^{5}}{2^{5} \cdot 5}+\frac{1}{12} c_{2} t^{3}-\frac{1}{2}\left(c_{1}+c_{2}^{2}\right) t+\text { const. }
$$

admits $p, p_{1}, p_{2}$ as derivatives.
The problem is then solved, moreover. The motion is defined by the equations:

$$
\frac{\partial V}{\partial c_{1}}=b_{1}, \quad \frac{\partial V}{\partial c_{2}}=b_{2}
$$

i.e.:

$$
\begin{equation*}
\left[\frac{1}{2}(x+y)+c_{1}\right]^{1 / 2}-\frac{1}{2} t=b_{1}, \quad-\frac{1}{2}(x+y)+\frac{1}{12} t^{3}-\frac{1}{2} c_{1} t=b_{2}, \tag{8}
\end{equation*}
$$

here, which are equalities that can be written:

$$
\begin{equation*}
x+y=\frac{1}{2} t^{2}+\alpha t+\beta, \quad x-y=\frac{1}{6} t^{3}+\alpha^{\prime} t+\beta^{\prime} \tag{9}
\end{equation*}
$$

when one sets:

$$
b_{1}=-\frac{1}{2} \alpha, \quad c_{1}=\frac{1}{4}\left(\alpha^{2}-2 \beta\right), \quad b_{2}=-\beta^{\prime}, \quad c_{2}=-\alpha^{\prime} .
$$

Remark. - If one appeals to the ordinary equations of motion, instead of using the method of Jacobi and Mayer, then one will have to integrate the very simple system:

$$
x^{\prime \prime}=\frac{1}{2}(1+t), \quad y^{\prime \prime}=\frac{1}{2}(1-t),
$$

which will immediately give back the equalities (9) upon adding and subtracting and will further simplify when one makes the change of variables $x+y=x_{1}, x-y=y_{1}$, moreover. If one employs the Jacobi method after the latter change of variables then the partial differential equations can be integrated immediately by separating the variables.

In regard to that, it is fitting to observe that, in theory, the method of Jacobi and Mayer is the most advantageous general method, but it can present some grave inconveniences in the applications. Solving for $p_{i}$ can rapidly imply some complications when the variables are not
distinguished from each other by having some special character. Furthermore, the single equation to which Mayer's procedure will lead can depend upon some parameters, and above all, it results from combining several distinct equations, which is a situation that generally masks the simplification, and in particular, any possible separation of variables.

It is then appropriate to employ the method of Jacobi and Mayer only after one has chosen the variables most judiciously and taken into account all of the simplifications that will reveal the most direct procedure. The very elementary example that we just treated shows how complications will quickly appear when one does not take any prior precautions.

## Remark on the canonical systems that are deduced from a first integral.

Consider an arbitrary canonical system:

$$
\begin{equation*}
\frac{d q_{1}}{\frac{\partial H}{\partial p_{1}}}=\frac{d p_{1}}{-\frac{\partial H}{\partial q_{1}}}=\frac{d q_{2}}{\frac{\partial H}{\partial p_{2}}}=\ldots=\frac{d p_{k}}{-\frac{\partial H}{\partial q_{k}}} \tag{1}
\end{equation*}
$$

and let $f\left(q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}\right)$ be a first integral of that system. We know that we can determine a complete integral $V$ of the equation:

$$
H\left(q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=h
$$

that satisfies the equation:

$$
f\left(q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=\alpha
$$

From that, the integral of the canonical system:

$$
\begin{equation*}
\frac{d q_{1}}{\frac{\partial f}{\partial p_{1}}}=\frac{d p_{1}}{-\frac{\partial f}{\partial q_{1}}}=\frac{d q_{2}}{\frac{\partial f}{\partial p_{2}}}=\ldots=\frac{d p_{k}}{-\frac{\partial f}{\partial q_{k}}} \tag{2}
\end{equation*}
$$

can be deduced from that same function $V\left(q_{1}, \ldots, q_{k}, h, a, a_{1}, \ldots, a_{k-2}\right)$. The integral of the system (1) will be defined by the equalities:

$$
\frac{\partial V}{\partial a}=b, \quad \frac{\partial V}{\partial a_{1}}=b_{1}, \quad \ldots, \quad \frac{\partial V}{\partial a_{k-2}}=b_{k-2}, \text { and } \quad p_{i}=\frac{\partial V}{\partial q_{i}}
$$

while the integral of the system (2) will be defined by the analogous equalities:

$$
\frac{\partial V}{\partial h}=c, \quad \frac{\partial V}{\partial a_{1}}=b_{1}, \quad \ldots, \quad \frac{\partial V}{\partial a_{k-2}}=b_{k-2}, \quad \text { and } \quad p_{i}=\frac{\partial V}{\partial q_{i}}
$$

One sees that the relations between the $q_{i}$ that are determined by the systems (1) and (2) admit the common system of integrals:

$$
\frac{\partial V}{\partial a_{1}}=b_{1}, \quad \frac{\partial V}{\partial a_{2}}=b_{2}, \quad \ldots, \quad \frac{\partial V}{\partial a_{k-2}}=b_{k-2}
$$

from which one can generally eliminate $(k-3)$ of the variables $q_{i}$, which will give $(k-2)$ relations:

$$
\begin{gathered}
\varphi_{1}\left(q_{1}, \ldots, q_{k}, h, a, a_{1}, \ldots, a_{k-2}\right)=0, \quad \varphi_{2}\left(q_{1}, \ldots, q_{k}, h, a, a_{1}, \ldots, a_{k-2}\right)=0, \ldots, \\
\varphi_{k-2}\left(q_{1}, \ldots, q_{k}, h, a, a_{1}, \ldots, a_{k-2}\right)=0
\end{gathered}
$$

that are common to (1) and (2).
In particular, assume that the system (1) corresponds to a problem in dynamics ( $H=T-U$ ), and that $f$ is likewise of the form $\mathcal{T}-V$, in which $\mathcal{T}$ is homogeneous of degree two with respect to the $p_{i}$, which do not enter into $V$. At the same stroke, one will solve the given problem of dynamics and the problem for which the canonical function is $\mathcal{T}-V+C(T-U)$, in which $C$ denotes a constant, by finding a complete integral $V$ that is common to the system $H=h, f=\alpha$.

Having done that, if one considers an arbitrary system (1) and an arbitrary integral $f$ then one can say that in general the integration of the system (1) and the integration of the system (2) are two equivalent problems. However, one cannot conclude that when one has determined a certain integral $f$ of a particular system, it will suffice for one to integrate the system (2) in order to know how to integrate the first one. When one has obtained the general integral of (2) in any manner, one can indeed deduce a complete integral of the equation:

$$
f\left(q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=\alpha
$$

but that integral will not generally satisfy the equation $H=h$, and knowing the general integral of the system (2) cannot be of any use in the determination of an integral $V$ that is common to the equations $H=h, f=\alpha$.

For example, suppose that $H$ does not depend upon $q_{1}$. The system (1) will admit the integral $f=p_{1}$, and the system (2) that is deduced from $f$ will reduce to the following one:

$$
\frac{d p_{1}}{d q_{1}}=0, \quad \frac{d q_{2}}{d q_{1}}=0, \quad \ldots, \quad \frac{d p_{h}}{d q_{1}}=0
$$

whose general integral is:

$$
p_{1}=\text { const. }, \quad p_{2}=\text { const. }, \quad \ldots, \quad p_{k}=\text { const. }
$$

However, one cannot deduce a function $V$ from that integral that will verify both of the equations:

$$
H\left(q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=h, \quad \frac{\partial V}{\partial q_{1}}=a,
$$

so one will be reduced to merely finding a complete integral $W\left(q_{2}, \ldots, q_{k}\right)$ of the equation:

$$
H\left(q_{2}, \ldots, q_{k}, a, \frac{\partial W}{\partial q_{2}}, \ldots, \frac{\partial W}{\partial q_{k}}\right)=h .
$$

## Specialized first integrals.

When a function $f\left(t, q_{1}, \ldots, q_{k}, p_{1}, \ldots, p_{k}\right)$ is a first integral of the canonical system:

$$
\begin{equation*}
d t=\frac{d q_{1}}{\frac{\partial H}{\partial p_{1}}}=\frac{-d p_{1}}{\frac{\partial H}{\partial q_{1}}}=\ldots=\frac{d q_{k}}{\frac{\partial H}{\partial p_{k}}}=\frac{-d p_{k}}{\frac{\partial H}{\partial q_{k}}}, \tag{1}
\end{equation*}
$$

it will satisfy the condition:

$$
\frac{\partial f}{\partial t}+(f, H)=0
$$

identically.
Now assume that the preceding relation is not true identically, but is a consequence of the equation $f=0$ :

$$
\frac{\partial f}{\partial t}+(f, H)=M \cdot f
$$

Any motion $q_{i}(t), p_{i}(t)$ that corresponds to the initial conditions $t_{0}, q_{i}^{0}, p_{i}^{0}$ and satisfies the condition that $f\left(t^{0}, \ldots, q_{i}^{0}, \ldots, p_{i}^{0}\right)=0$ will satisfy that condition for any $t$. Indeed, replace the $p_{i}$, $q_{i}$ in $f$ with functions of $t$. That will give:

$$
\frac{d f}{d t}=\frac{\partial f}{\partial t}+(f, H)=M_{1}(t) f,
$$

since the differential equation $d f / d t=M_{1} f$ admits no other integral that is annulled for $t=t_{0}$ than $f \equiv 0$, so the function $f\left(t, \ldots, q_{i}(t), \ldots, p_{i}(t), \ldots\right)$ will be identically zero if it is zero for $t=t_{0}$.

Conversely, if the relation $f=0$ cannot be verified at the instant $t_{0}$ without being verified during all of its motion then one will necessarily have:

$$
\frac{\partial f}{\partial t}+(f, H)=M f
$$

Indeed, any motion that satisfies the condition:

$$
f\left(t^{0}, \ldots, q_{i}^{0}, \ldots, p_{i}^{0}, \ldots\right)=0
$$

must satisfy the condition that:

$$
\frac{d f}{d t}=\frac{\partial f}{\partial t}+(f, H)=0
$$

for the same values $t^{0}, \ldots, q_{i}^{0}, \ldots, p_{i}^{0}, \ldots$, i.e., that relation must be a consequence of the relation $f=0$.

Such a relation $f=0$ is called a specialized first integral. Indeed, one can regard it as something that provides a first integral in which one has given a particular value to the constant. In order to see that, suppose that one has formed the $2 k$ first integrals:

$$
\begin{aligned}
& p_{i}^{0}=\varphi_{i}\left(t_{0}, t, \ldots, p_{i}, \ldots, q_{i}, \ldots\right), \\
& q_{i}^{0}=\psi_{i}\left(t_{0}, t, \ldots, p_{i}, \ldots, q_{i}, \ldots\right),
\end{aligned}
$$

and consider the first integral:

$$
f\left(t^{0}, \ldots, q_{i}^{0}, \ldots, p_{i}^{0}, \ldots\right)=f\left(t_{0}, \ldots, \varphi_{i}, \ldots, \psi_{i}, \ldots\right)=F\left(t, \ldots, p_{i}, \ldots, q_{i}, \ldots\right),
$$

in which $t_{0}$ is a number that generally enters into $F$. An arbitrary motion will verify the condition $F=$ const., and since the constant is annulled with $f\left(t^{0}, \ldots, q_{i}^{0}, \ldots, p_{i}^{0}, \ldots\right)$, any motion whose initial conditions annul $f$ will satisfy the condition that:

$$
F\left(t, \ldots, p_{i}, \ldots, q_{i}, \ldots\right)=0
$$

but, by hypothesis, it will also verify the equality:

$$
f\left(t, \ldots, p_{i}, \ldots, q_{i}, \ldots\right)=0
$$

It follows from this that those two equalities cannot be distinct, since the motions are subject to the condition that $f\left(t^{0}, \ldots, q_{i}^{0}, \ldots, p_{i}^{0}, \ldots\right)=0$, which once more depends upon $2 k-1$ constants. One can then replace the equality $f=0$ with the equality $F=0$. If one so desires, $F$ can again contain $f$ as a factor.

Having made those remarks, suppose that one has determined a specialized first integral $f=0$. From the theorem on page 293, the two equations:

$$
\left\{\begin{array}{c}
\frac{\partial V}{\partial t}+H\left(t, q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0  \tag{2}\\
f\left(t, q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0
\end{array}\right.
$$

will admit an infinitude of common integrals, and in particular, an infinitude of complete integrals (viz., integrals that depend upon $k-1$ arbitrary constants that permit one to give arbitrary values to $k-1$ of the derivatives $\partial V / \partial q_{i}$ ). Knowing one such integral can serve to determine the motions that satisfy the condition that $f=0$.

Indeed, let $F=c$ be a first integral that gives the specialized integral $f=0$ when $c=0$. Write down the equations:

$$
\left\{\begin{array}{c}
\frac{\partial V}{\partial t}+H\left(t, q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0,  \tag{3}\\
F\left(t, q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=c
\end{array}\right.
$$

and let $V_{1}\left(t, q_{1}, \ldots, q_{k}, c, c_{1}, \ldots, c_{k-1}\right)$ be a complete integral of that system. All of the complete integrals $V$ of (2) are obtained by setting $c=0$ in the complete integrals $V_{1}$ of (3). On the other hand, the general motion is defined by the equalities:

$$
\frac{\partial V_{1}}{\partial a}=b, \quad \frac{\partial V_{1}}{\partial a_{1}}=b_{1}, \quad \ldots, \quad \frac{\partial V_{1}}{\partial a_{k-1}}=b_{k-1}, \quad \text { and } \quad p_{i}=\frac{\partial V_{1}}{\partial q_{i}}
$$

Therefore, when one has determined a complete integral of the system (2), say, $V_{1}\left(t, q_{1}, \ldots\right.$, $q_{k}, a_{1}, \ldots, a_{k-1}$ ), one can write down the equalities (4), which will determine $q_{2}, \ldots, q_{k}$, for example, and the $p_{i}$ as functions of $t$ and $q_{1}$ [and $2(k-1)$ arbitrary constants]. It then remains for one to determine $q_{1}$ as a function of $t$, which is a problem that depends upon a second-order equation. When one knows a complete integral of (2), the determination of the motions that satisfy the condition $f=0$ will then come down to the integration of a second-order equation.

## Applying the Legendre transformation to the Jacobi equation.

I will conclude this study of Jacobi's method with a remark concerning the manner by which the parameters $q_{i}, p_{i}$ enter into an arbitrary canonical system and into the corresponding partial differential equation.

It is clear that the variables $p_{i}, q_{i}$ play a symmetric role in the canonical system (1):

$$
\begin{equation*}
\frac{d q_{i}}{d t}=\frac{\partial H}{\partial p_{i}}, \quad \frac{d p_{i}}{d t}=-\frac{\partial H}{\partial q_{i}}, \quad(i=1,2, \ldots, k) \tag{1}
\end{equation*}
$$

since it will suffice to change $H$ into $-H$ if one wishes to replace the $p_{i}$ with the $q_{i}$, and vice versa. From that, it is legitimate to replace the Jacobi equation:

$$
\begin{equation*}
\frac{\partial V}{\partial t}+H\left(t, q_{1}, \ldots, q_{k}, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)=0 \tag{2}
\end{equation*}
$$

with the following one:

$$
\begin{equation*}
\frac{\partial W}{\partial t}-H\left(t, \frac{\partial W}{\partial p_{1}}, \ldots, \frac{\partial W}{\partial p_{k}}, p_{1}, \ldots, p_{k}\right)=0 \tag{3}
\end{equation*}
$$

If one knows a complete integral $W\left(t, p_{1}, \ldots, p_{k}, \alpha_{1}, \ldots, \alpha_{k}\right)$ of the latter equation then one can set:

$$
\frac{\partial W}{\partial \alpha_{1}}=\beta_{1}, \ldots, \frac{\partial W}{\partial \alpha_{k}}=\beta_{k}, \quad q_{i}=\frac{\partial W}{\partial p_{i}} \quad(i=1,2, \ldots, k),
$$

and those equalities will define the general integral of (1).
Moreover, it is quite easy to pass from equation (2) to equation (3) by replacing the variables $q_{i}$ and the function $V$ with the new variables $p_{i}$ and the new function $W$ that are coupled with the first ones by the equalities:

$$
\left\{\begin{align*}
p_{i} & =\frac{\partial V}{\partial q_{i}},  \tag{4}\\
W & =\sum_{i=1}^{k} q_{i} \frac{\partial V}{\partial q_{i}}-V
\end{align*}\right.
$$

Upon repeating the entirely elementary calculation with the aid of which we converted the Lagrange equations to the canonical form (see pp. 120), we will see that those formulas imply the following ones:

$$
\left\{\begin{align*}
q_{i} & =\frac{\partial W}{\partial p_{i}} \\
V & =\sum_{i=1}^{k} p_{i} \frac{\partial W}{\partial p_{i}}-W  \tag{5}\\
\frac{\partial V}{\partial t} & =-\frac{\partial W}{\partial t}
\end{align*}\right.
$$

The change of variables (4) then transforms equation (2) into equation (3). A complete integral $V\left(t, q_{1}, \ldots, q_{k}, \alpha_{1}, \ldots, \alpha_{k}\right)$ of (2) will then correspond to a complete integral $W\left(t, p_{1}, \ldots, p_{k}, a_{1}, \ldots\right.$, $a_{k}$ ) of (3), and one will have the relations:

$$
\frac{\partial V}{\partial a_{i}}=-\frac{\partial W}{\partial \alpha_{i}}
$$

The equations $\frac{\partial V}{\partial a_{i}}=b_{i}, p_{i}=\frac{\partial V}{\partial q_{i}}(i=1,2, \ldots, k)$ imply the equations:

$$
\frac{\partial W}{\partial \alpha_{i}}=-b_{i}, \quad q_{i}=\frac{\partial W}{\partial p_{i}}
$$

The transformation (4) is called the Legendre transformation. It is the simplest of the contact transformations. One can summarize the foregoing by saying that the change of variables $q_{i}=\varpi_{i}$, $p_{i}=\chi_{i}$ will preserve the canonical form of the system (1).

More generally, one can study all of the changes of variables:

$$
\varpi_{i}=\varphi_{i}\left(t, p_{1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right), \quad \chi_{i}=\psi_{i}\left(t, p_{1}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right) \quad(i=1,2, \ldots, k)
$$

that preserve the canonical form of an arbitrary canonical system, or (in another form) all of the transformations:

$$
\begin{aligned}
\varpi_{i} & =\varphi_{i}\left(t, q_{1}, \ldots, q_{k}, V, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right) \quad(i=1,2, \ldots, k), \\
W & =F\left(t, q_{1}, \ldots, q_{k}, V, \frac{\partial V}{\partial q_{1}}, \ldots, \frac{\partial V}{\partial q_{k}}\right)
\end{aligned}
$$

that transform an arbitrary equation (2) (in which $V$ does not occur) into another first-order equation in which $W$ does not occur. Those transformations, which constitute the contact transformations precisely, have been the subject of considerable work by Jacobi and Sophus Lie. One can find the presentation of the latter in the second volume of Lie's Theorie der Transformationsgruppen and a presentation of the former in the previously-cited book by Goursat (Chapter XI).

Let us now address the particular case in which $H$ has the form $T-U$, in which $T$ denotes a homogeneous quadratic form in the $p_{i}$ whose coefficients, as well as $U$, are second-order polynomials in the $q_{i}$ that are independent of $t$ and have no first-degree terms. The function $H$ can be regarded as the canonical function of two distinct problems in mechanics, according to whether one takes the $q_{i}$ or the $p_{i}$ to be the parameters that define the position of the system.

Those two problems can be solved simultaneously because the canonical equations of the first one will coincide with those of the second one when one changes $t$ into $-t$. For example, let:

$$
H=L p_{1}^{2}+2 M p_{1} p_{2}+N p_{2}^{2}
$$

with

$$
L=A q_{1}^{2}+2 B q_{1} q_{2}+C q_{2}^{2}, \quad M=A^{\prime} q_{1}^{2}+2 B^{\prime} q_{1} q_{2}+C^{\prime} q_{2}^{2}, N=A^{\prime \prime} q_{1}^{2}+2 B^{\prime \prime} q_{1} q_{2}+C^{\prime \prime} q_{2}^{2} .
$$

That function $H$ has the canonical form that corresponds to the $d s^{2}$ :
( $\alpha$ )

$$
d s^{2}=\frac{N d q_{1}^{2}-2 M d q_{1} d q_{2}+L d q_{2}^{2}}{L N-M^{2}} .
$$

Now permute the variables $p$ and $q$. $H$ will become:

$$
H^{\prime}=L_{1} p_{1}^{2}+2 M_{1} p_{1} p_{2}+N_{1} p_{2}^{2}
$$

with

$$
L_{1}=A q_{1}^{2}+2 A^{\prime} q_{1} q_{2}+A^{\prime \prime} q_{2}^{2}, \quad M_{1}=B q_{1}^{2}+2 B^{\prime} q_{1} q_{2}+B^{\prime \prime} q_{2}^{2}, N_{1}=C q_{1}^{2}+2 C^{\prime} q_{1} q_{2}+C^{\prime \prime} q_{2}^{2},
$$

and will be the canonical function that corresponds to the $d s^{2}$ :

$$
d s_{1}^{2}=\frac{N_{1} d q_{1}^{2}-2 M_{1} d q_{1} d q_{2}+L_{1} d q_{2}^{2}}{L_{1} N_{1}-M_{1}^{2}} .
$$

Therefore, whenever one has determined the geodesics of the first $d s^{2}$, one will have determined those of the second (and conversely). Indeed, once the geodesics of $d s^{2}$ have been calculated, the canonical system that corresponds to $H$ will be found to have been integrated. One will then know $p_{1}$ and $p_{2}$ as functions of $t$. When one eliminates $t$ from $p_{1}$ and $p_{2}$, the relations will depend upon only two arbitrary constants $p_{2}=F\left(p_{1}, a, b\right)$, and if one replaces $p_{1}$ with $q_{1}$ and $p_{2}$ with $q_{2}$ then the new relation $q_{2}=F_{1}\left(q_{1}, a, b\right)$ will define the geodesics of $d s_{1}^{2}$.

## END OF LECTURE 17


[^0]:    ${ }^{\left({ }^{1}\right)} \Delta^{\prime}$ cannot be zero for any $x, x_{1}, \alpha_{2}, \ldots, \alpha_{n}$, because otherwise $\Delta$ would be identically zero.

[^1]:    $\left(^{2}\right)$ In what follows, I will say that any system of such functions $q_{2}, q_{3}, \ldots, q_{k}$ of $q_{1}$ define a trajectories of $S$.

[^2]:    $\left({ }^{3}\right)$ Equations (5), in which one annuls the $\beta_{i}$, have no other solution than $Q_{1} \equiv Q_{2} \equiv \ldots \equiv Q_{k} \equiv 0$, because their determinant is a power of $\Delta$, and as a result, it will not be zero.

[^3]:    $\left({ }^{4}\right)$ In regard to this, we add that if the $q_{i}$ have been chosen in that way then the discriminant $\Delta$ of $T$ will not be annulled for any real position of the system and for real values of the $q_{i}$. The expression $\sum m_{j}\left(x_{j}^{\prime 2}+y_{j}^{\prime 2}+z_{j}^{\prime 2}\right)$ will then be annulled for some real values of $x_{j}^{\prime}, y_{j}^{\prime}, z_{j}^{\prime}$ that are not all zero.

[^4]:    $\left({ }^{5}\right)$ For the given initial conditions of the system $S$, the vis viva $T$ and the forces are well-defined. However, if the same system of values of the $q_{i}$ corresponds to several positions of $S$ then the $A_{i j}$ and $Q_{i}$ will have several determinations. For a position of $S$ that is taken at random, there will be no ambiguity about which of those determinations that one must choose.

[^5]:    $\left({ }^{6}\right)$ See Picard, Traité d'analyse, Tome II, Chapter XI, page 308.

[^6]:    $\left({ }^{7}\right)$ The trajectory $q_{i}=\varphi_{i}\left(q_{1}\right)$ is a geodesic.

[^7]:    $\left({ }^{8}\right)$ The equilibrium conditions of the system $S$ are obviously $\beta_{i}=0$, and the equalities $q_{i}=a_{i}$ define that equilibrium.

[^8]:    $\left({ }^{9}\right)$ We saw (page 242) that the remarkable trajectories depend upon at most $(k-1)$ parameters.

[^9]:    $\left({ }^{10}\right)$ It is appropriate to observe that the arc of the trajectory $M_{1} M_{2}$ will correspond to a real periodic motion (either true or conjugate).

[^10]:    $\left({ }^{12}\right)$ I intend that to mean that the coordinates $q_{1}, q_{2}, \ldots, q_{k}$ of $M$ differ as little as one desires from the coordinates $a_{1}, a_{2}, \ldots, a_{k}$ of $N^{\prime}$.

[^11]:    $\left({ }^{13}\right)$ One does not regard two motions to be distinct when one of them is deduced from the other by $t$ in $t+$ const.

[^12]:    $\left({ }^{14}\right)$ I intend that to mean that $x, y, z$ remain holomorphic functions of the arc-length $\sigma$.

[^13]:    $\left({ }^{15}\right)$ We gave (page 255) an example in which a singular branch decomposes into an infinitude of segments that tend to $N^{\prime}$ and correspond to just as many periodic motions (which are alternately true or conjugate).

[^14]:    $\left({ }^{16}\right)$ On the subject of higher geometry, see Darboux, (Tome II, Chapters VI, VII, and VIII).

[^15]:    $\left({ }^{17}\right)$ Conversely, if $K$ is independent of $x$, moreover, then [from (a)] the same thing will be true of $f=\sum p_{i} x_{i}^{\prime}-K$.

[^16]:    $\left({ }^{18}\right)$ It is clear that it hardly matters in all of this whether $t$ does or does not figure in the expressions considered.

[^17]:    $\left({ }^{19}\right)$ That is what happens when one associates two of the integrals that are provided by the theorem of the motion of the center of gravity (see pp. 288), say, $f_{1}=\sum p_{i}$ and $f_{2}=\sum q_{i}$, or also the integral $f_{1}=\sum p_{i}$ and the area integral $f_{2}=\sum\left(y_{i} r_{i}-z_{i} q_{i}\right)$. One will then find that $\left(f_{1}, f_{2}\right) \equiv 0$.
    $\left({ }^{20}\right)$ This case presents itself when one associates (see pp. 288) the two integrals $f_{1}=\sum p_{i}$ and $f_{2}=\sum\left(z_{i} p_{i}-x_{i} r_{i}\right)$. One finds that $\left(f_{1}, f_{2}\right)=-\sum r_{i}$, so one has a new integral $f_{3}=r_{i}$. However, $\left(f_{1}, f_{3}\right)$ and $\left(f_{2}, f_{3}\right)$ are zero. Similarly, if one associates the two integrals $f_{1}=\sum\left(y_{i} r_{i}-z_{i} q_{i}\right), f_{2}=\sum\left(z_{i} p_{i}-x_{i} r_{i}\right)$ then one will find that $\left(f_{1}, f_{2}\right)=\sum\left(q_{i} x_{i}-r_{i} y_{i}\right)$, which defines a third integral $f_{3}$, but $\left(f_{1}, f_{3}\right) \equiv-f_{2},\left(f_{2}, f_{3}\right) \equiv f_{1}$.

[^18]:    ${ }^{(21)}$ That case presents itself when one associates (see pp. 288) the two integrals $f_{1}=\sum p_{i}, f_{2}=\sum\left(z_{i} p_{i}-x_{i} r_{i}\right)$. One finds that $\left(f_{1}, f_{2}\right)=-\sum_{i}$, which is a new integral $f_{3}=r_{i}$. However, $\left(f_{1}, f_{3}\right)$ and $\left(f_{2}, f_{3}\right)$ are zero. Similarly, if one associates the two integrals $f_{1}=\sum\left(y_{i} r_{i}-z_{i} q_{i}\right), f_{2}=\sum\left(z_{i} p_{i}-x_{i} r_{i}\right)$ then one will find that $\left(f_{1}, f_{2}\right)=\sum\left(q_{i} x_{i}-r_{i} y_{i}\right)$, which will define a third integral $f_{3}$, but $\left(f_{1}, f_{3}\right) \equiv-f_{2},\left(f_{2}, f_{3}\right) \equiv f_{1}$.

[^19]:    $\left.{ }^{(22}\right)$ That expression coincides with $T+U$ in the case of mechanics.
    $\left({ }^{23}\right)$ That is Jacobi's first method for integrating an arbitrary equation $(A)$. Indeed, one knows that once one knows a complete integral of $(A)$, it is easy to deduce a general integral from that. On the other hand, if one is given a firstorder equation $(A)$ into which $V$ enters explicitly then one can always reduce it to an analogous equation into which $V$ no longer enters, but in which the number of variables is increased by one unit. It suffices to replace the search for functions $V$ with the search for functions $F\left(t, q_{1}, q_{2}, \ldots, q_{k}, V\right)$ that are equal to a constant that defines an integral $V$ of the first equation.

[^20]:    $\left({ }^{24}\right)$ These relations can depend upon constants.

[^21]:    $\left({ }^{25}\right)$ The term complete system applies to even more general systems, moreover. (See Goursat, loc. cit., Chapter II.)

[^22]:    $\left({ }^{26}\right)$ More generally (the conditions of the theorem having been fulfilled), if $f_{1}, f_{2}, \ldots, f_{m}$ depend upon $h$ constants $b_{1}, b_{2}, \ldots, b_{h}$, and if one solves the equations $f_{i}=0$ for the $j$ variables $p_{i}, r$ variables $q_{i}, s$ constants $b_{i}(j+r+s=m)$ then the equalities $\left(F_{i}, F_{j}\right) \equiv 0$ will be verified identically, so $F_{i}=0$ will represent any of the new equations. That will result immediately from the previous argument.

[^23]:    $\left({ }^{27}\right)$ We will see that this condition is not necessary.

[^24]:    $\left({ }^{28}\right)$ Observe that in order to form the system ( $1^{\prime}$ ), it will suffice to find an integral $f_{1}\left(p_{1}, p_{2}, \ldots, p_{k}, q_{1}, \ldots, q_{k}\right)$ of the equation $\left(f_{0}, f_{1}\right)=0$ without having solved $f_{0}$ with respect to one of the $p_{i}$. Upon then solving the system $f_{0}=0, f_{1}=c_{1}$ with respect to $p_{1}, p_{2}$, one will form a system ( $1^{\prime}$ ) for which the condition $\left(f_{0}^{\prime}, f_{1}^{\prime}\right) \equiv 0$ is once more verified: Nevertheless, it is necessary that the solution should be possible.

[^25]:    $\left({ }^{29}\right)$ To simplify the notation, we shall suppress the upper indices (?) on the symbols $f_{l}, F_{l}$.

[^26]:    $\left({ }^{31}\right)$ In order for this solution to be possible, it is necessary and sufficient that the functional determinant $\frac{D\left(\chi_{1}, \ldots, \chi_{2(k-r)}\right)}{D\left(q_{r+1}, \ldots, q_{k}, p_{r+1}, \ldots, p_{k}\right)}$ is not annulled identically for $u_{1}=0$. If that determinant is zero then one of the integrals, say $\chi_{1}$, will be equal to $F\left(\chi_{2}, \ldots, \chi_{2(k-r)}, u_{2}, \ldots, u_{r}\right)$ for $u_{1}=0$, and as a result, it will coincide with the integral $F$. The integrals $\chi_{i}$ will not be distinct then.
    ${ }^{(32)} g_{1}$ never reduces to an absolute constant $A$. Otherwise, the identity $(l)$, in which one sets $\psi_{1}=A$ and gives arbitrary numerical values to $\psi_{2}, \ldots, \psi_{2(k-r)}$, would demand that $\chi$ must be simply a function of $u_{2}, \ldots, u_{r}$.

